# Riemann's Existence Theorem: <br> An elementary approach to moduli 

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## CHAPTER 1

## SCOPE OF THE EXISTENCE THEOREM

This chapter is an overview of this book's three main topics.

- How Riemann's Existence Theorem describes moduli spaces of Riemann surface covers of the Riemann sphere.
- How finite group theory puts practical - for applications - structures into collections of such covers.
- How each finite group generates its own nilpotent theory of fundamental groups, forming systems of moduli spaces with $G_{\mathbb{Q}}$ actions.


## 1. Context for the book

We note L. Ahlfors' satisfying entwining of the algebra and geometry in 1st year graduate complex variables [Ahl79]. This author could do no better than use it as his underpinning. Still, that book leaves the full scope of monodromy a mystery, it prepares little on coordinates describing Riemann surfaces, and none on families of Riemann surfaces.
E. Hille wrote a function theoretic encyclopedia [Hil62]. As a graduate student, I enjoyed how relevant were its historical comments to sophisticed mathematics in the 1800s. For example, mathematicians seeking immortality (in private, of course) might ponder its many serious references to H. Schwarz's work [Sc1890]. Few present day complex variable enthusiasts know the coherence or context of that work. Two authors, G. Springer [Spr57] and R.C. Gunning ([Gun66] and [Gun67]), did great service bringing Riemann surfaces to graduate students by the 1960s. For the former, that was H. Weyl's uniformization approach (as in his projection lemma). For the latter it was the Cartan-Serre vector bundle view of the algebro-differential geometry that works on Riemann surfaces.
E. Neuenschwanden's perspective answers many questions on what took so long for Riemann surfaces to make their mark [Ne81]. He documents contention between Weierstrass's algebraic and Riemann's harmonic function approaches. This is relevant to the relation between Riemann and Abel and Galois. For Weierstrass admits the influence of Abel on his work. Still, one can't see it directly on Riemann. This is despite serious documentation of his intellectual activities, including the direct influence on him of Gauss. Further, [Ne81] leaves unanswered other questions about the assimilation of mathematics.

These modern works have little group theory; not even including the original approaches of Abel, Galois and Riemann. Few presented group theory so dramatically as did H. Weyl. Yet, even Weyl (on quantum mechanics) met resilient resistance to group theory. My convictions are here; I advocate using the power of group theory. Showing how finite and profinite group theory can handle intricate monodromy and moduli, and apply practically to algebra and complex variables,
is my goal. Still, there's a fence to walk. We can't afford to let group theory overwhelm us. Galois was first to note group theory's power. Also, he wrote on its potential to dominate the subject techically.

The introductions of two books, [MM95] and [V̈̈96], show they closely connect through group theory with this book. [Fri94] and [Fri95c] specifically discuss connnections of our topics to $[\mathbf{S e 9 2}]$. These three books concentrate on how Riemann's Existence Theorem applies to the Inverse Galois Problem. By contrast, classical topics appear here more often than in the first two. Also, this author uses standard formulations of the Inverse Galois Problemmuch less. Yet, the reader can find here a leisurely track through Riemann surface theory guided by problems requiring little preparation for their statements, a virtue of the Inverse Galois Problem. My choices often have a long literature before the connection to Riemann surfaces appeared. By occasional referring to topics from these three books, starting in Chap. 4, I have added efficiency to this liesurely pace.

By being leisurely, we (I and the reader) may also consider the struggle of many generations with whether punctured Riemann surfaces and their moduli variation belong to function theory or to algebra. Since it is leisurely, using a style less sophisticated than my papers in the middle 1970's, it might from its opening chapter be mistaken as curiously old-fashioned. Further, my evident hero worship of Riemann can further confuse those who don't know me. What I have tried is an historical model. I attempt to synthesize in two early chapters what might have been the insights of those famous researchers from the 1800s for whom analytic continuation and its applications to algebraic equations was an open extravaganza of intensely studied equations. The complication of mathematics, that Galois remarked on as often as one can do when one is going to disappear long before maturity, overwhelmed all except the technical giants of the time. Yet, from this came synthesis: Abstract approaches that simplified everything for those who could follow them. The people I admire today tend to admire - by aspiration in their own research these very same people. If we aim to please and appeal to Abel, Galois and Riemann on this score, we realize - in rational moments - that is an impossibility. Further, since that is a triumvirate of geniuses, such an appeal detracts from showing why even they struggled, and despite the time that has passed we too, with the whole topic. There is a serious question for mathematics. When does mathematics (versus Riemann) have a firm grasp on a significant subject?

Is it when an elite institution husbands a handful of caretakers of an industry of supporting research? Is it when myriad papers allude to consequent deep theories, even if they don't directly involve the roiling concepts? Is it when some text has nailed the subject completely to a prestigious group's satisfaction? Is it when a blithely confident prestigious group claim the subject's foundations are firm and available to any sincere seeker? Is it when the subject successfully supports several independent and competing schools derived from its basic problems?

We don't know what would convince most research mathematicians of the security of a subject. The author has a point in writing this book; though he cannot easily pick one affirmative viewpoint for the maturity of this book's subject. Its techniques quickly worked to reveal the nature of long standing problems in his hands. On that basis a fair observer might support that the techniques work. Still, there are geniuses beyond Abel, Galois and Riemann who have their viewpoints. Examplars of thinking with great scope and imagination certainly include [An02],
[De89] [Moc96]. We end the book at the wealth of analytic questions and applications raised by Modular Towers, a little before the influence of these writers on the author. So, only a shadow of their influence is here.

All, however, support connecting profinite groups to function theory. That leads to final, painful consideration. Will we, and the world outside mathematics, ever be able to tolerate the inundation that often overwhelms us from the connections bridged by mathematical language?

## 2. A quick summary

A fuller overview follows this sections brief summary.
2.1. A concise description, chapter by chapter. Compact Riemann surfaces as branched covers of a sphere appear in 1st year graduate courses as elementary discussions of multi-valued functions. We expand the usual brief treatment in Chap. 2. This carefully treats analytic continuation to motivate the geometry behind it. It introduces the Existence Theorem sufficiently to get lessons from the theory of abelian algebraic covers of the punctured Riemann sphere (§3.2). It starts with two different definitions of algebraic functions, one from algebraic equations another phrasing from analytic continuation. An imprecise version of Riemann's Existence Theorem is that these describe the same functions. This is an elementary investigation, based on the first half of graduate complex variables.

In this book Riemann's Existence Theorem means the precise statement from Chap. 4. That really organizes all algebraic functions (of $z$ ). Chap. 4 fully develops Riemann's Existence Theorem. It emphasizes data determining a branched cover of the sphere up to equivalence. Abel and Galois started a tradition. Our version: Translate complex analytic and arithmetic geometry problems into group theory through application of forms of Riemann's Existence Theorem.

Advanced texts often append another statement. It is that any compact Riemann surface (Chap. 3) has an analytic (nonconstant) map to $\mathbb{P}_{z}^{1} \stackrel{\text { def }}{=} \mathbb{C} \cup\{\infty\}$, the Riemann sphere. Springer's book [Spr57] dedicates much space to proving this last statement. We rarely use it; our basic data already includes such a function and (given the Riemann-Roch Theorem) includes Springer's goals (see below).

Suppose, however, $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is such an analytic map. Let $z_{0}$ be a particular $z$ value, and consider $X_{z_{0}}$, the fiber of $\varphi$ over $z_{0}$. Then, finding algebraic equations for $X$, necessary for most applications, depends on producing another function $\varphi^{\prime}: X \rightarrow \mathbb{P}_{w}^{1}$ that separates points of $X_{z_{0}}$. The explicit production of such a $\varphi^{\prime}$ is a consequence of uniformization of $X$ by the appropriate simply-connected domain (disk, plane or sphere). As uniformization plays an important role in advanced applications, say, related to $\theta$ functions, we often raise elementary aspects of it.

The globally defined functions, $\varphi$ and $\varphi^{\prime}$ have an algebraic relation $F\left(\varphi, \varphi^{\prime}\right) \equiv 0$ between them: $F \in \mathbb{C}[z, w]$. Let $L_{\varphi^{\prime}} \subset \mathbb{C}$ be the field generated by the ratios of all coefficients of $F$. Let $K$ be a field containing $L_{\varphi^{\prime}}$. A frequent application of this relation $F$ is to give meaning to the expression a $K$ point on $X$. From $F$ and $z_{0} \in K$, there is an equivalence class of permutation representations of the absolute Galois $G_{K}$ of $K$. This comes from its action on points of $X$ over $z_{0}$. Refined applications of covers analyze the dependence of this statement on the choice of $\varphi^{\prime}$.

Chap. 4 shows the following. Let $L_{z}$ be the field generated by the symmetric functions in $\boldsymbol{z}$ (with $\infty$ removed if it appears).
(2.1a) There is a choice of $\varphi^{\prime}$ giving $K$ algebraic over $L_{\boldsymbol{z}}$.
(2.1b) The complete set of minimal fields $\mathcal{L}_{\varphi}$ appearing as $L_{\varphi^{\prime}}$ in the algebraic closure of $L_{\boldsymbol{z}}$ is an intrinsic (moduli) invariant of $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$.
(2.1c) Sometimes (the Existence Theorem shows) $\mathcal{L}_{\varphi}$ consists of a unique field.

When $\boldsymbol{z}$ consists of algebraically independent values, the analysis of $\mathcal{L}_{\varphi}$ includes the moduli (deformation) theory of a cover. That is Part II of the book. Comparing this case with the case $L_{\boldsymbol{z}}=\mathbb{Q}$ (or some other explicit algebraic number field) is tantamount to approaches to the Inverse Galois Problem.

We assume students with one semester each of a graduate algebra course and a graduate complex variables course. Few students master Galois theory from their algebra courses. Thus, we give an analytic continuation approach to showing the field of convergent Puiseux expansions around a point is algebraically closed. This supports many elementary subtopics that could otherwise be baffling. For example, Riemann's Existence Theorem uses an infinite number of incompatible algebraically closed fields containing the field $\mathbb{C}(z)$. Let $\boldsymbol{z}=\left(z_{1}, \ldots, z_{r}\right)$ be a fixed set of points on the sphere. Denote the complement of $\boldsymbol{z}$ on the sphere by $U_{z}$.

Riemann's Existence Theorem is about algebraic functions extensible on $U_{\boldsymbol{z}}$. These are functions with analytic continuations along any path (from an explicit base point) avoiding $\boldsymbol{z}$. At each point $z_{0}$, not in $\boldsymbol{z}$, these algebraic extensible functions embed in the algebraically closed field of Puiseux expansions in $z_{0}$. Isomorphisms between their different embeddings is coded in the fundamental groupoid.

Chap. 2 describes abelian functions of $z$ through analytic continuing branches of the $\log$ function. It demonstrates many basic definitions and some advanced concepts. Among these is that of a group attached to monodromy action. For books motivated by $\theta$ functions and their applications, this book is unusually persistent in emphasizing finite group theory.

Chap. 3 has basics on fundamental groups and permutation representations. Though our definitions and first examples of manifolds are traditional, our aim is to illustrate practical use of deformations of Riemann surfaces. We concentrate on very explicit manifolds. Chap. 5 produces highly structured moduli spaces parametrizing equivalence classes of Riemann surfaces.

Consider the notation around (2.1). For $\boldsymbol{z}$ fixed, and $K=L_{\boldsymbol{z}}$, if $z_{0} \in K$, there is an action of $G_{K}$ on the profinite completion of the fundamental group $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ (Chap. 4). Moduli parameters appear with the following question.

Problem 2.1. What happens with covers of $U_{\boldsymbol{z}}$ as $z$ varies?
First appearances give the following impression.
(2.2a) The fundamental group of $U_{\boldsymbol{z}}$ doesn't change with $\boldsymbol{z}$.
$(2.2 \mathrm{~b}) G_{\mathbb{Q}}$ action changes drastically if you can even consider it varying with $\boldsymbol{z}$.
Both (2.2a) and (2.2b) are wrong.
Suppose we try to write equations (with coefficients in $\boldsymbol{z}$ ) for the deformations of an algebraic function $f=f_{\boldsymbol{z}}$ extensible on $U_{\boldsymbol{z}}$ (Chap. 2). Locally in $\boldsymbol{z}$ this is possible. Going, however, around various closed paths in the space for $\boldsymbol{z}, f_{z}$ might return to a different extension field of $\mathbb{C}(z)$. Riemann's Existence Theorem tells precisely how to calculate which paths return to the original function field (§5.4.1). Hurwitz monodromy action is the phrase for our most important calculations. This produces coordinates for coefficients relating $f_{z}$ algebraically to ( $z, \boldsymbol{z}$ ) (Chap. 5).

Choosing generators and a base point are what allow covering applications of the fundamental group. A response to (2.2a) is that this extra data produces a refined moduli space setup. This motivates a Lie algebra approach to (2.2b) putting the two parts of (2.2) under a common framework. We use ideas from renown papers of Y. Ihara and J.P-Serre and moduli space that give the proper context for the Inverse Galois Problem.

Abelian covers of $U_{\boldsymbol{z}}$ for any $\boldsymbol{z}$ comes from branches of $\log$ (Chap. 2). Ihara studied (parts of the) arithmetic of nilpotent covers of $U_{\boldsymbol{z}}$ when $r=3$ [Iha86]. Nilpotent theory appears in applications to the Inverse Galois Problem. Here it starts from nonsplit nilpotent extensions extending data about covers with any given finite (often simple) group $G$. For $p$ a prime dividing the order of $G$, a univeral totally nonsplit extension ${ }_{p} \tilde{G}$ of $G$ produces sequences of refined moduli spaces (§8.3).
[Fri78] and [Tha86] had common elements: use of the theory of complex multiplication, and an arithmetic philosophy using the braid group. The former used analytic geometry and finite group theory. There is now a natural way to join this to the profinite and function theory approach of the latter. This means joining Modular Towers to the Grothendieck-Teichmüller technology. The tools include extension of Deligne's tangential base points [De89] with insight from Riemann's $\theta$ functions.
2.2. Meaning of the word, elementary in the title. The first two chapters are elementary by most perspectives. Still, understanding Chap. 5 on moduli requires mastery of the first two chapters. The approach is elementary because it allows a newcomer into the area through examples and techniques using finite group theory. Traditionally, for example, with modular curves, one must have serious training in complex analysis. The action happens with automorphic functions on the upper half plane.

Here we often use uniformization from below, replacing the upper half plane and representations of $\mathrm{SL}_{2}(\mathbb{R})$ with the Riemann sphere $\mathbb{P}_{z}^{1}$ and finite group theory. Then, modular curves and their associated towers are an example of the moduli of dihedral group covers. The same technique works by replacing the dihedral group by any finite group. This opens up applications beyond the traditional modular curve approach.

This modular curve generalization uses a construction attached to each prime $p$ dividing the order of a finite group $G$ : The universal $p$-Frattini cover ${ }_{p} \tilde{G}$ of $G$. This especially considers those primes $p$ for which $G$ is $p$-perfect (it has no cyclic quotient of order $p$ ).

Add to this a collection $\mathbf{C}$ of conjugacy classes from $G$ whose elements have order prime to $p$. Then, $(G, p, \mathbf{C})$ produces a sequence of moduli spaces of curves. Example: $G$ is the dihedral $D_{p}$ of order $2 p$ ( $p$ an odd prime) and $\mathbf{C}$ consists of four repetitions of the conjugacy class of involutions. Then, the sequence of moduli spaces is the classical modular curve series $\left\{Y_{1}\left(p^{k+1}\right)\right\}_{k=0}^{\infty}$ : Quotients of the upper half-plane by well-known subgroups denoted $\Gamma_{1}\left(p^{k+1}\right)$ of $\mathrm{PSL}_{2}(\mathbb{Z})$. The $k$ th level of the sequence in this case is $Y_{1}\left(p^{k+1}\right.$. Introducing the generalizing sequences of spaces, Modular Towers, is the book's main advanced topic.

When $G$ is an alternating group $A_{n}(n \geq 4)$, and $p=2$, Modular Tower properties generalize applications of $\theta$ functions. Specifically, in this alternating group case several components may appear in a Modular Tower level. This is unlike
the dihedral case where all levels are connected. We use modular representations of characteristic quotients of ${ }_{p} \tilde{G}(\S 10.2)$. This extension of Schur's theory of universal central extensions connects these components to the famous mod 2 (half-canonical class) invariant from $\theta$ functions.

Function theory, as in cusp forms and Eisenstein series from modular curves also appear here. Since the levels are moduli spaces of curves, we know most about those functions by relating them to $\theta$ functions of curves representing points in the moduli spaces. Such varying $\theta$ functions produce $\theta$-null automorphic forms. Our main examples illustrate this when the moduli spaces are quotients of the upper half plane, giving covers of the classical $j$-line. This exactly corresponds (for any $(G, p))$ to the case $\mathbf{C}$ consists of four conjugacy classes in $G$.

Modular curves, though a guide, are a small portion of the noncongruence quotients of the upper half plane with a tower structure related to a prime $p$. New applications reveal the value of a Riemann's Existence Theorem approach. Function wise it generalizes both the braid group approach to the Inverse Galois Problem and the Tate module.

Early chapters develop detailed motivation for using classical functions. The deeper function theory, however, appears in outline (with exposition on applications related to the literature). Developing this completely is a topic for a later book.

## 3. Early historical motivation

A renown problem from the early 19th century was to express in radicals solutions $x$ of the general $n$th degree polynomial equation

$$
\begin{equation*}
f(x)=x^{n}+a_{1} x^{n-1}+\cdots+a_{n-1} x+a_{n}=0 \tag{3.1}
\end{equation*}
$$

with $f$ of degree $n$ in $x$. The goal specifically asks for solutions $x$ using known functions of the coefficients $a_{1}, \ldots, a_{n}$. The explicitly known functions of the time were what we call radicals.

Traditional books tackle this using Galois theory with pure algebra. They reproduce Galois' Theorem characterizing when a field extension $L / K$ is a subfield of a chain of radical field extensions of $K$. This happens if the Galois closure of $L / K$ has solvable group.

It is a pretty story. Still, Galois' Theorem is not a common object of mathematical pilgrimage (even if Galois is). This treatment hides ingredients that still seize the imagination of modern mathematicians, as it possessed Abel, Galois and Riemann. Abel and Galois recognized group theory for showing, with $a_{i} \mathrm{~s}$ and $n \geq 4$, the field of radical sequences in the $a_{i}$ s do not contain the solutions. Still, these books lack problems motivating present research. Further, the subject's character falls outside the neatly compartmental introduction of rings, groups, modules and elementary classification results of the rest of 1st year graduate algebra. These historically come long after it, leaving the impression Galois theory is both mildly exotic and slightly moribund.
3.1. Consider functions of one variable. To be more explicit turn to complex variables, as did Abel. Instead of $a_{1}, \ldots, a_{n}$ being general, specialize to functions $a_{1}(z), \ldots, a_{n}(z)$ of one complex variable $z$. Assume $a_{1}(z), \ldots, a_{n}(z)$ are in the field $\mathbb{C}(z)$ : rational functions of $z$ with complex coefficients. It is convenient to replace $x$ by a variable $w$ taking complex values. Refer in this specialized form to the equation $f\left(a_{1}(z), \ldots, a_{n}(z), w\right)=m(z, w)=0$.

The left side of (3.1) does not factor into lower degree polynomials over the field $a_{1}, \ldots, a_{n}$ generate. The specialized expression $m(z, w)=0$ may factor over $\mathbb{C}(z)$. To simplify, assume $m$ is an irreducible polynomial in $w$ over $\mathbb{C}(z)$. Analytic continuation displays the $n$ solutions in $w$ as $n$ manifestations of one solution. The manifestations cohere through a group. Here is how it arises.
3.2. Motivating integrals. Critical values $\boldsymbol{z}=z_{1}, \ldots, z_{r}$ of $m$ are places $z^{\prime}$ where $m\left(z^{\prime}, w\right)$ has repeated roots. Fix $z_{0}=z$ not equal to a critical value of $m$. Then the zeros $w$ of $m(z, w)$ have expressions $w_{1}\left(z ; z_{0}\right), \ldots, w_{n}\left(z ; z_{0}\right)$, meromorphic functions in $z$ around $z_{0}$. This holds for any $z_{0}$ outside $z$. So these algebraic functions are extensible on $\mathbb{C} \cup\{\infty\} \backslash\{\boldsymbol{z}\}=U_{\boldsymbol{z}}$ (Chap. 2). The group of $m(z, w)$ (relative to $z$ ) is all permutations of the $w_{i}$ s from continuation around closed paths in $U_{\boldsymbol{z}}$ based at $z_{0}$. Call the $w_{i} \mathrm{~s}$ abelian if this group is abelian.

This study of zeros jibed nicely with another problem of Abel's day: Analyze elementary antiderivatives, like the watershed example $\int \frac{d x}{\sqrt{x^{3}+a x+b}}$. Specifically, what is the dependence of these antiderivatives on the parameters $a$ and $b$ ?

Here $m(z, w)=w^{2}-\left(z^{3}+a z+b\right)$. Write $G(z)=\frac{1}{\sqrt{z^{3}+a z+b}}$ acknowledging (Chap. 2) that plugging in values of $z$ near $z_{0}$ requires choosing one of two functions $G(z)$ analytic in a disc about $z_{0}$ with $G(z)^{-2}=z^{3}+a z+b$. Consider $F(z)$, an antiderivative of $G(z)$, locally. An integral gives $F(z)$. So, it has analytic continuations around $U_{\boldsymbol{z}}$. These continuations produce an abelian group of periods (Chap. 2). Chap. 4 shows the group is $\mathbb{Z} \times \mathbb{Z}$. Further, its fit with the analytic continuations of $G(z)$ appears in the semidirect product $\mathbb{Z} \times \mathbb{Z} \times^{s}\{ \pm 1\}$ (§8). Let $D_{n}$ be the dihedral group of order $2 n$.

Classical modular curves parametrize four branch point $D_{n}$ extensions of $U_{\boldsymbol{z}}$. Galois checked with his theorem for which $n$ these modular curve parameters were solvable functions of the classical $j$ parameter [Rig96, p. 133]. Properties of $F(z)$ entwine integration and the appearance of abelian extensions:
(3.2) $F(z)$ is a versal abelian extensible function on $U_{\boldsymbol{z}}$ with monodromy around $z$ bounded by $G(z)$ (Chap. 4).
Restricting to $U_{\boldsymbol{z}}$ still shows the full scope of Riemann's version of (3.2). The next three sections base a story of his program on analytic continuation.

## 4. Algebraic functions among extensible functions

Denote Laurent series expansions about $z_{0}$ by $\mathcal{L}_{z_{0}}$. Let $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ be extensible (meromorphic) elements of $\mathcal{L}_{z_{0}}$ on $U_{\boldsymbol{z}}$. Call $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ algebraic if it satisfies $m(z, f(z)) \equiv 0$ with $m \in \mathbb{C}[z, w]$ a nonzero polynomial. Characterizing such $f$ through analytic continuation, the main topic of Chap. 2, is the first step to classifying algebraic functions. Any analytic continuation of $f$ around a closed path in $U_{\boldsymbol{z}}$ also gives a zero of $m$. So, there are only finitely many analytic continuations of $f$. Analytic continuations of $f$ along paths whose end points have limits in $\boldsymbol{z}$ take values nowhere dense (a finite set) in the Riemann sphere. This qualitative statement characterizes algebraic $f$. The full force of Riemann's Existence Theorem is in phrasing this through fundamental group representations (Chap. 4). Denote the algebraic elements of $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ by $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$.
4.1. One element of $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ is versal for $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$. There are so many algebraic functions in $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)=\mathcal{E}\left(U_{\boldsymbol{z}}\right)$ (if the cardinality, $r=|\boldsymbol{z}|$ exceeds two).

We can explain little about them by listing their polynomial equations. Yet, there is much structure in this collection.
4.1.1. Setup for uniformizaton. Riemann provided such by finding one function $\tilde{f}_{\boldsymbol{z}}$ giving all of $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$ through a type of Galois correspondence. An outline for this appears in Chap. 3.
(4.1a) Recognize each algebraic function $f \in \mathcal{E}\left(U_{\boldsymbol{z}}\right)$ has an attached topological cover $\varphi_{f}: X_{f} \rightarrow U_{z}$.
(4.1b) Produce a (uni)versal cover $\varphi_{\boldsymbol{z}}: \tilde{U}_{\boldsymbol{z}} \rightarrow U_{\boldsymbol{z}}$ with a discrete group $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ acting on $\tilde{U}_{z}$.
(4.1c) Show $X_{f}$ is a topological quotient of $\tilde{U}_{\boldsymbol{z}}$ by a subgroup of $\pi_{1}\left(U_{\boldsymbol{z}}\right)$.
(4.1d) Show $\tilde{U}_{\boldsymbol{z}}$ has a complex analytic embedding in $\mathbb{C}: h: \tilde{U}_{\boldsymbol{z}} \rightarrow \mathbb{C}$.

As in Chap. 3, (4.1d) produces $\tilde{f}_{z}$ as follows. Let $U_{z_{0}}$ be any disk around $z_{0}$ (on $U_{z}$ ). Cauchy's Theorem (we return soon to that) shows this:
(4.2) Each $g \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ extends to a unique meromorphic function on $U_{z_{0}}$.
4.1.2. $\tilde{U}_{z}$ and Hurwitz equivalence. Riemann's Existence Theorem shows why $\tilde{U}_{\boldsymbol{z}}$ identifies with the upper half plane $\mathbb{H}=\{z \in \mathbb{C} \mid \Im(z)>0\}$. Apply the Existence Theorem (see Chap. 4 or $\S 5.1$ ) with a branch cycle description of form

$$
g_{1}=\left(1 \ldots s_{1}\right) \cdots\left(s_{t-1}+1 \ldots s_{t}\right)
$$

$s_{1}+s_{2}+\cdots+s_{t}=n ; g_{2}=\left(s_{1} s_{2} \ldots s_{t}\right)$ and $g_{3}=\left(g_{1} g_{2}\right)^{-1}$. Count points over branch points: $t+(n-t+1)+1=n+2$.

Uniformize $U_{\{0,1, \infty\}}$ with the classical $\lambda$ function (§7.1.1). Choose $n=r-2$. This produces a genus 0 cover of $\varphi_{\boldsymbol{g}}: X_{\boldsymbol{g}} \rightarrow \mathbb{P}_{z}^{1}$ unramified over $U_{\{0,1, \infty\}}$ with exactly $r$ points over $\{0,1, \infty\}$. Further, $\lambda$ factors through this cover:

$$
\mathbb{H} \rightarrow X_{\boldsymbol{g}} \backslash \varphi_{\boldsymbol{g}}^{-1}(0,1, \infty) \rightarrow U_{\{0,1, \infty\}}
$$

This uniformizes one copy of $\mathbb{P}_{z}^{1}$ minus $r$ points. Deform (differentiably) $X_{\boldsymbol{g}} \backslash$ $\varphi_{\boldsymbol{g}}^{-1}(0,1, \infty)$ to any other copy of $\mathbb{P}_{z}^{1}$ minus $r$ points (Chap. 5).

Regard algebraic functions $f=y$ (of $z$ ) as giving a relation between two variables $x$ and $y$. Classical literature often chooses the isomorphism class of the function field $\mathbb{C}(z, y)$ as the unique goal of an algebraic relation. If $\mathbb{C}(z, y)$ is isomorphic to $\mathbb{C}\left(z^{*}, y^{*}\right)$, this views the algebraic relation between $\left(z^{*}, y^{*}\right)$ (take the minimal polynomial of $y^{*}$ over $\left.\mathbb{C}\left(z^{*}\right)\right)$ as elementary equivalent to the relation between $z$ and $y$. The history of considering algebraic relations had its motivation in integrals. There the most telling invariant of a function field $\mathbb{C}(z, y)$ is the genus $g$ (maximal number of linearly independent holomorphic differentials $\S 6.2$ ) on the function field.

A connected algebraic space parametrizes all algebraic relations of genus $g$ (Chap. 5). Investigating this and subtler problems about algebraic relations suggest a more delicate equivalence between function fields. In addition to the isomorphism of $\mathbb{C}\left(z^{*}, y^{*}\right)$ with $\mathbb{C}(z, y)$, this isomorphism includes that $\mathbb{C}\left(z^{*}\right)=\mathbb{C}(z)$. Call this Hurwitz equivalence. Even in restricting to genus $g$ function fields there are many components to the parameter spaces of Hurwitz equivalences of algebraic relations. Hurwitz (equivalence) spaces all derive from the elementary notion of deforming points as in the construction above for $\tilde{U}_{z}$.
4.1.3. The value of $\tilde{f}_{\boldsymbol{z}}$. Since $\varphi_{\boldsymbol{z}}: \tilde{U}_{\boldsymbol{z}} \rightarrow U_{\boldsymbol{z}}$ is a covering space, $\varphi_{\boldsymbol{z}}^{-1}\left(U_{z_{0}}\right)$ has countably many connected components $\left\{U_{i}\right\}_{i=1}^{\infty}$, each homeomorphic to $U_{z_{0}}$ by restriction of $\varphi_{\boldsymbol{z}}$. Let $\varphi_{1}: U_{1} \rightarrow U_{0}$ be this one-one restriction. Then, (4.1d)
produces the function

$$
\begin{equation*}
\tilde{f}_{\boldsymbol{z}}=h \circ \varphi_{1}^{-1}: U_{0} \rightarrow \mathbb{C} \tag{4.3}
\end{equation*}
$$

This one function distinguishes homotopy classes of paths on $U_{\boldsymbol{z}}$ by analytic continuation. It separates homotopy classes of paths (based at $z_{0}$ ) by its values at end points of analytic continuations. Since $\tilde{U}_{\boldsymbol{z}}$ is simply connected and in $\mathbb{C}$, Riemann's mapping theorem says it is analytically isomorphic to a disk (or to $\mathbb{C}$, if $r=1$ or 2 ) for each $\boldsymbol{z}$.
4.2. Uniformizing from above versus below. Thus, $\tilde{U}_{z}$ is a domain for parametrizing $X_{f}$ for all $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$, as $\boldsymbol{z}$ varies. This complements how we use Riemann's Existence Theorem.
4.2.1. Shortcomings of $h \circ \varphi^{-1}$. The universal covering space helps organize functions and differential forms. Still, algebraists find it hides phenomena close to their interests. For example, $h \circ \varphi_{1}^{-1}$ is neither algebraic nor known: Its values at algebraic points of $U_{\boldsymbol{z}}$ are rarely algebraic. Though based on $\lambda(\tau)$ in $\S 4.1 .2$, it changes with $\boldsymbol{z}$. Yet, it provides no explicit equations for algebraic functions.

Even proving a cover from Riemann's Existence Theorem is algebraic still goes through a hard proof that we now separate from other, more algebraic, observations. Suppose $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is a cover. Let $\varphi_{w}: X \rightarrow \mathbb{P}_{w}^{1}$ be any function separating all points on the fiber $X_{z_{0}}$ over $z_{0}$. Then, $X \rightarrow \mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$ by $x \mapsto\left(\varphi(x), \varphi_{w}(x)\right)$ has closed image birational to $X$ in the algebraic variety $\mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$. Apply Chow's Lemma (Chap. 4) to get that $X$ is algebraic.

Classical construction of $\varphi_{w}$ relies on a uniformization $\mathbb{H} \rightarrow U_{z}$ presenting $U_{z}$ as a quotient $\mathbb{H} / H, H$ a subgroup of $\mathrm{PSL}_{2}(\mathbb{R})$. One must find nontrivial $H$ invariant functions on $\mathbb{H}[\mathbf{K 7 2}$, Chap. III]. Variants are in [Spr57, Chap. 6-10] and [Vö96, Chap. 5]. We rely on the treatment from the last of these references - especially well adapted to Riemann's Existence Theorem. How to find $\varphi_{w}$ (or some related differential form) algebraically appears in many of our examples.
4.2.2. Virtues of $h \circ \varphi^{-1}$. The phrase "abelian theory" means here Riemann's unified generalization of Abel's results. This includes describing functions, abelian covers and the results of integration of differentials on a Riemann surface. It includes Riemann's extension of Cauchy's integral theorem to open Riemann surfaces. We discuss it, and our reason for including a nilpotent theory below. There is no denying the value of $h \circ \varphi^{-1}$.
(4.4a) It organizes tool the abelian and nilpotent theory.
(4.4b) It coordinates analyzing real points on moduli spaces of curves.
(4.4c) It is suspiciously close to being algebraic, producing an algebraic object (a flat $\mathbb{P}^{1}$-bundle) capturing its uniformizing properties.
4.2.3. The Existence Theorem and classical uniformization meet. Each item in (4.4) has Existence Theorem and $\tilde{U}_{\boldsymbol{z}}$ aspects: Uniformization from below versus above. The literature neglects the former, though it is constructive and practical. The latter has had elegant developments.

Both work best as tools for analyzing properties of families (moduli spaces) of curves. They give enhancements when the moduli spaces themselves fit in natural sequences. The abelian theory gave the first such natural sequences. This shows in modular curve sequences ( $\S 8.3$, Chap. 5).
4.2.4. Illustrating with modular curves. When the parameter $r$ (cardinality of $\boldsymbol{z}$ ) is 4 , the comparison between Modular Towers and modular curves is direct. For example, these properties hold for Modular Towers when $r=4$.

- Their levels are curves.
- They include modular curve towers and come with an essential prime $p$ : Its powers correspond to Modular Tower levels.
- They lie over the classical $j$-line and have useful cusps over $j=\infty$.
- All levels are moduli spaces, with variants corresponding to structures going with modular curve notation $X_{0}\left(p^{k+1}\right), X_{1}\left(p^{k+1}\right)$ and $X\left(p^{k+1}\right)$.
Any finite group $G$ and prime $p$ dividing $|G|$ produces many Modular Towers; many more than there are modular curve towers. The name Modular Tower comes from this comparison and the group (modular representation ) theory that appears in their analysis.

An elementary comparison occurs in analyzing real points on a Modular Tower. Through Riemann's Existence Theorem this gives the essential data about cusps. From that come their geometric properties (Chap. 5), including genuses of their components. This is especially interesting when the finite group $G$ producing the Modular Tower is simple and the prime $p$ is 2 . We now discuss the Existence Theorem, then the abelian theory.

## 5. $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$ and data from groups

Riemann's Existence Theorem (Chap. 4) compactifies $\varphi_{f}: X_{f} \rightarrow U_{\boldsymbol{z}}$ to a ramified cover of Riemann surfaces $\bar{\varphi}_{f}: \bar{X}_{f} \rightarrow \mathbb{P}_{z}^{1}$. It then turns the process around by using special generators of the fundamental group $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ of $U_{z}$. From these it produces all elements of $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$.
5.1. Identifying a fundamental group requires generators. Suppose $G$ is a finite transitive subgroup of $S_{n}$. A surjective homomorphism $\psi: \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right) \rightarrow G$ canonically produces a cover $X_{\psi} \rightarrow U_{\boldsymbol{z}}$ from homotopy classes of paths. We don't need generators of $\pi_{1}\left(U_{\boldsymbol{z}}\right)$ to define these covers Chap. 3. They, however, handily list all such homomorphisms and therefore all such covers. Convenient listing of covers allows explicitly computing properties of Hurwitz spaces (Chap. 5).

The collections of $r$ paths (based at $z_{0}$ ) we call classical generators of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ appear in Chap. 3. Points in $z$ produce conjugacy classes $\mathbf{C}_{\boldsymbol{z}}$ in $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Classical generators are homotopy classes of paths respectively representing these conjugacy classes. Choose representing paths that pair wise meet only at their beginning and end point $z_{0}$. Label one as $\bar{g}_{1}$. Label the others from their having a clockwise order in leaving the point $z_{0}$. These $r$ paths $\bar{g}_{1}, \ldots, \bar{g}_{r}$ now satisfy

## (5.1) $\bar{g}_{1} \bar{g}_{2} \cdots \bar{g}_{r}=1$ : The product-one condition.

Invariants of Hurwitz space components appear from (5.1) (§10.1 illustrates). Classical generators - satisfying these conditions -automatically generate the fundamental group (Chap. 3). Solving for $\bar{g}_{r}$ presents the fundamental group as a free group on $r-1$ generators. Yet, that violates the product-one symmetry. So, that free group presentation appears only in stray computations.

This part of Riemann's theory works very well. It successfully applies to many problems. These require some finite group theory. It is the center of the first third of the book. Polynomial equations describe algebraic curves. This is what gives structure allowing fields of definitions and interpreting rational points. The

Riemann's Existence Theorem approach, however, emphasizes effective group theory over manipulating explicit equations. Exercises and examples illustrate this (Chap. 3, Chap. 4, Chap. 9).
5.2. Changing classical generators. There is no canonical set of classical generators for $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. The necessary variation of this choice produces the braid and mapping class groups (Chap. 5). This complication enriches mathematics. Still, it requires explanation.

The second third of the book organizes collections of elements from $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$. This allows $\boldsymbol{z}$ to vary. Sets of $r$ (unordered) distinct points on $\mathbb{P}_{z}^{1}$ have a topology and analytic structure extending that of $\mathbb{P}_{z}^{1}$. This set is $\mathbb{P}^{r} \backslash D_{r}=U_{r}$ : Projective $r$ space minus the discriminant locus (Chap. 5). Think of $U_{r}$ as monic polynomials of degree either $r$ or $r-1$ with distinct roots. Or, consider it the quotient of $\left(\mathbb{P}_{z}^{1}\right)^{r} \backslash \Delta_{r}$ by permutation action of $S_{r}$, the symmetric group of degree $r$, on ordered $r$-tuples of points. Here $\Delta_{r}$ is $r$-tuples with distinct coordinates.

The fundamental group of $U_{r}$ is the degree $r$ Hurwitz monodromy group $H_{r}$ (Chap. 5), an Artin braid group quotient. A permutation representation of $H_{r}$ produces the space of deformations of $X_{f}$. These are unreduced Hurwitz spaces.

A given function $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$ deforms in many ways as $\boldsymbol{z}$ varies. Local deformation, however, of $\bar{\varphi}_{f}: \bar{X}_{f} \rightarrow \mathbb{P}_{z}^{1}$ is unique along any path. This allows analyzing parameters for these moduli spaces. Yet, it leads further from explicit equations. To paraphrase Joni Mitchell's "Both Sides Now" (from the 60's): Something's lost and something's gained in putting equations away. Explicit functions, however, return with the abelian and nilpotent theory.
5.3. Moving $z$, even with $z_{0}$ fixed, forces changing generators. Picture: $z_{1}$ and $z_{2}$ follow semicircles, producing

$$
\begin{equation*}
Q_{1}:\left(\bar{g}_{1}, \ldots, \bar{g}_{r}\right) \mapsto\left(\bar{g}_{1} \bar{g}_{2} \bar{g}_{1}^{-1}, \bar{g}_{1}, \ldots, \bar{g}_{r}\right) \tag{5.2}
\end{equation*}
$$

Replacing 1 by $i \leq r-1$ gives the full generating collection $Q_{1}, \ldots, Q_{r-1}$ of the Hurwitz monodromy group $H_{r}$ (Chap. 5). The $H_{r}$ action from (5.2) on classical generators is the technical tool for describing families of covers.

Let $G$ be a fixed finite group. Assume these further ingredients.
(5.3a) $\boldsymbol{z}^{\prime}$ is a specific point of $U_{r}$.
(5.3b) $\psi_{\boldsymbol{z}^{\prime}}: \pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right) \rightarrow G$ is a specific surjective homomorphism to $G$ using classical generators $\bar{g}_{1}, \ldots, \bar{g}_{r}(\S 5.1)$.
(5.3c) $T: G \rightarrow S_{n}, n$ an integer, is a faithful permutation representation.

Then, $\psi_{\boldsymbol{z}^{\prime}}$ gives a finite (ramified) cover $\varphi_{G, T, \boldsymbol{z}^{\prime}}=\varphi_{\boldsymbol{z}^{\prime}}: X_{\boldsymbol{z}^{\prime}} \rightarrow \mathbb{P}_{z}^{1}$ of Riemann surfaces of degree $n$. The images of $\bar{g}_{1}, \ldots, \bar{g}_{r}$ give generators $g_{1}, \ldots, g_{r}$ of $G \leq S_{n}$, with an associated set of $r$ conjugacy classes $\mathbf{C}$ in $G$. Riemann's Existence Theorem labels covers by $g_{1}, \ldots, g_{r}$ (branch cycles). It gives $\varphi_{G, T, \boldsymbol{z}^{\prime}}$ as an equivalence relation on homotopy classes of paths based at $z_{0}$. Suppose $\boldsymbol{z}^{\prime}$ moves to nearby $\boldsymbol{z}^{\prime \prime}$, with $z_{0} \in \mathbb{P}_{z}^{1}$ and paths representing $\bar{g}_{1}, \ldots, \bar{g}_{r}$ fixed. Then, there is a unique isomorphism of $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right)$ and $\pi_{1}\left(U_{\boldsymbol{z}^{\prime \prime}}, z_{0}\right)$ commuting with their maps to $G$.

An automorphism $\alpha$ of $\pi_{1}\left(U_{z}, z_{0}\right)$ sends generators to new generators, changing $\psi_{\boldsymbol{z}}$ to $\psi_{\boldsymbol{z}} \circ \alpha$. Inner automorphisms of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$, however, produce covers equivalent to the old cover. It is moduli of covers we use; equivalence two homomorphisms if they differ by an inner automorphism. Further, only automorphisms
from the Hurwitz monodromy group $H_{r}$ send classical generators to classical generators (possibly changing the intrinsic order of the paths). Such automorphisms arise from deforming the pair $\left(\boldsymbol{z}, z_{0}\right)$ along closed paths in $U_{r}$. They preserve the conjugacy classes of classical generators. So, $\mathbf{C}$, the conjugacy class set in $G$, is an $H_{r}$ invariant of any given homomorphism $\psi$.
5.4. The moduli spaces appear. The Nielsen class of ( $G, \mathbf{C}$ ) (Chap. 5) consists of $r$-tuples $\left(g_{1}, \ldots, g_{r}\right)$ satisfying the product-one condition attached to $(G, \mathbf{C})$. The Existence Theorem uses classical generators of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ to produce equivalence class of covers.
5.4.1. Writing equations in $\boldsymbol{z}$. The Nielsen class $\operatorname{Ni}(G, \mathbf{C})$ has entries in a set of conjugacy classes $\mathbf{C}$ in $G$, independent of the braid action. Thus, $H_{r}$ acts on elements of $\mathrm{Ni}(G, \mathbf{C})$ (similar to (5.2)). An aside: We need to quotient by conjugation from $G$. Here is how to think of this action.

Let $\varphi_{0}: X_{0} \rightarrow \mathbb{P}_{z}^{1}$ be a cover from the Existence Theorem using $\boldsymbol{g} \in \mathrm{Ni}(G, \mathbf{C})$. Take the branch points to be $\boldsymbol{z}_{0}$. What if someone asks for explicit equations for this cover? That could mean either:
(5.4a) equations just for $\varphi_{0}$; or
(5.4b) equations for $\varphi_{\boldsymbol{z}}: X_{\boldsymbol{z}} \rightarrow \mathbb{P}_{z}^{1}$, with branch points $\boldsymbol{z}$, valid for $\boldsymbol{z}$ near $\boldsymbol{z}_{0}$ (where it specializes to $\varphi_{0}$ ).
Don't those seem like asking too little? Why concentrate on one set of branch points $\boldsymbol{z}_{0}$, or even on a neighborhood of $\boldsymbol{z}_{0}$ ? You'd want $\varphi_{\boldsymbol{z}}$ valid for all $\boldsymbol{z} \in U_{r}$. If, however, this were possible, then analytically continuing $\varphi_{z}$ around any closed path $\mathcal{P}$ in $U_{r}$ would return you to $\varphi_{0}$.

The homotopy class of $\mathcal{P}$ is an element $Q_{\mathcal{P}}$ of $H_{r}$. Further, Chap. 5 shows the cover at the end of $\mathcal{P}$ has a branch cycle description $(\boldsymbol{g}) Q_{\mathcal{P}}$. (Compute that with the starting classical generators of $\pi_{1}\left(U_{z_{0}}\right)$.) So, finding equations for $\varphi_{\boldsymbol{z}}$ valid for all $\boldsymbol{z}$ requires $(\boldsymbol{g}) Q_{\mathcal{P}}$ be $\boldsymbol{g}$ (modulo conjugation by $G$ or closely related). This you can check: Is $(\boldsymbol{g}) Q$ essentially $\boldsymbol{g}$ for all $Q \in H_{r}$. Example: Consider

$$
\boldsymbol{g}=((123),(321),(145),(154)) \in \operatorname{Ni}\left(A_{5}, \mathbf{C}_{3^{4}}\right)
$$

(§10.1). Then $(\boldsymbol{g}) Q_{2}=((1,23),(245),(321),(154))$. This is not conjugate to $\boldsymbol{g}$ even under $S_{5}$. So, as typical when $r \geq 4$, there are no such equations for $\varphi_{\boldsymbol{z}}$.
5.4.2. Analytic continuations of $\varphi_{\boldsymbol{z}_{0}}$. Nontrivial $H_{4}$ action means coefficients of equations for $\varphi_{\boldsymbol{z}}$ act as coordinates for a nontrivial cover of $U_{r}$. What cover?

It comes from the action of $H_{r}$, the fundamental group of $U_{r}$, on $\mathrm{Ni}(G, \mathbf{C})$ produced by covering space theory. Notation for this cover depends on the equivalence used for elements of the Nielsen class (as in (5.5)). Typical notation is $\mathcal{H}(G, \mathbf{C}, T)$. Each point of $\mathcal{H}(G, \mathbf{C}, T)$ corresponds to an equivalence class of covers: A point over $\boldsymbol{z} \in U_{r}$ is an element from $\operatorname{Ni}(G, \mathbf{C})$ attached to $\boldsymbol{z}$. Then, $\mathcal{H}(G, \mathbf{C}, T)$ itself covers the space $U_{r}$ of distinct unordered $r$-tuples of points from $\mathbb{P}_{z}^{1}$ (Chap. 5).

Various equivalences among covers produce different versions of this space. Two predominate in early applications. Denote the subgroup of $S_{n}$ normalizing $G$ and permuting the conjugacy classes in $\mathbf{C}$ by $N_{S_{n}}(G, \mathbf{C})$.
(5.5a) $\mathcal{H}(G, \mathbf{C})^{\text {in }}: T$ is the regular representation and the Galois cover comes with a fixed isomorphism between its Galois group and $G$ (inner spaces).
(5.5b) $\mathcal{H}(G, \mathbf{C}, T)^{\text {abs }}: T$ any faithful representation, with $r$-tuples equivalenced by $N_{S_{n}}(G, \mathbf{C})$ conjugation (absolute spaces).

See the main example of this chapter at $\S 10.1$.
The combinatorial groups in Chap. 5 have long histories: the Artin braid group, the Hurwitz monodromy group and the mapping class group. As in Chap. 4, we give formal proofs. Pictures appear only to convey conceptual symbolic data. Absolute spaces are the work horses in applications (Chap. 9). Inner spaces, however, directly connect the Inverse Galois Problem to generalizations of modular curves (§7.4).
5.4.3. Statics and dynamics of a cover. In the game of mentally writing equations for a cover, why would one cover be more significant than another? Many historical applications, such as the Inverse Galois Problem, consider a cover with equations over $\mathbb{Q}$ as most significant. For example, many arithmetic problems gain solutions if one can produce a cover with a particular monodromy group over $\mathbb{Q}(z)$ or over $\mathbb{Q}$. Such a cover provides solutions to related problems over another field by extending its equations to that field.

We picture such a cover $\varphi_{0}: X_{0} \rightarrow \mathbb{P}_{z}^{1}$ as being at the crossroads of a network of roads. The real points on $\mathcal{H}(G, \mathbf{C})^{\text {in }}$ would go through the point corresponding to $\varphi_{0}$, as would all $p$-adic points for every prime $p$. Concentrate on a real point, $\boldsymbol{p}_{0} \in \mathcal{H}(G, \mathbf{C})^{\text {in }}$ corresponding to a cover $\varphi_{0}$ over $\mathbb{R}$. To get a measure of the potential energy of this point we measure its distance from boundary points on $\mathcal{H}(G, \mathbf{C})^{\text {in }}$. Developing such a measure, depends on measuring something that goes to 0 as we deform $\varphi_{0}$ along a real component going to a boundary point, and the measuring coordinates must be canonical functions of the coordinates of the point $\boldsymbol{p}$ as it moves from $\boldsymbol{p}_{0}$ to the chosen boundary point.

The theory of abelian covers on $\bar{X}_{0}$ gives classical functions that we can use for making such measurements. As easily this could be on $\bar{X}_{0}$ minus a finite number of points, as with $U_{\boldsymbol{z}}$. Still, in the compact case, functions in $\mathcal{E}\left(\bar{X}_{0}\right)$ with finitely many analytic continuations are algebraic.

## 6. Abelian theory on $\bar{X}_{f}$ and integration

Let $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Suppose analytic continuations $f_{\gamma}$ of $f(z)$ have this property. (6.1) $f_{\gamma}(z)=f(z)$ for each closed path $\gamma$ based at $z_{0}$.

Rather than extensible, Chap. 2 calls $f$ extendible. Denote extendible elements of $\mathcal{E}\left(U_{\boldsymbol{z}}\right)$ by $\mathcal{E}\left(U_{\boldsymbol{z}}\right)^{\text {ext }}$.

Consider $f \in \mathcal{E}\left(U_{\boldsymbol{z}}\right)^{\text {ext }}$. Cauchy's Theorem in $U_{\boldsymbol{z}}$ shows precisely the nature of integrals $f(z) d z$ around certain closed paths. Since these are integrals, assume without further mention the paths miss any poles of $f d z$. Let $\boldsymbol{z}_{f, \infty}$ be the set of these poles. Assume for simplicity it is a finite set (appropriate for algebraic functions) which may include $\infty: z^{n} d z$ has a pole of order $n+2$ at $\infty$.

The definition of integral makes sense. Let $F(z)$ be an antiderivative of $f$ in a neighborhood of $z_{0}$. For any (simplicial) path $\gamma:[0,1] \rightarrow U_{\boldsymbol{z}}$, take the indefinite integral to the end point of $\gamma$ to be $F_{\gamma}$ (Chap. 2).

Cauchy's Residue Theorem: Let $\gamma$ be a closed path homologous to 0 in $U_{z}$. Compute $\int_{\gamma} f(z) d z$ from the winding number of $\gamma$ and residue of $f$ at each $z^{\prime} \in \boldsymbol{z}_{f, \infty}$ (Chap. 2). Winding numbers are values of integrals $\int_{\gamma} \omega$ where $\omega$ is a differential form - logarithmic, or of 3rd kind - taking the shape $\frac{1}{2 \pi i} \frac{d z}{z-z^{\prime}}=\omega_{z^{\prime}}$ with $z^{\prime} \in$ $\boldsymbol{z}_{f, \infty}$. Also, winding numbers appear in the definition of being homologous to 0 : The path has winding number 0 about each point in $\boldsymbol{z}$.
6.1. Changes of significance for algebraic $f$. Here is a paraphrase of Cauchy. Suppose (6.1) holds. Then, poles of $f$ and the map $\gamma \rightarrow\left(\int_{\gamma} \omega_{z_{1}}, \ldots, \int_{\gamma} \omega_{z_{r}}\right)$ determine $\int_{\gamma} f(z) d z$ when $\gamma$ is closed and homologous to 0 .

Suppose, however, $f$ is both algebraic and extendible. That means it is a rational function on $\mathbb{P}_{z}^{1}$. Then, there is no significant difference between the points in $\boldsymbol{z}$ and those in $\boldsymbol{z}_{f, \infty}$. By combining them both in the set $\boldsymbol{z}_{f, \infty}$ this allows dropping the homologous to 0 condition. We may consider integrals around any closed path.

Riemann made a more abstract change. Antiderivatives of $\omega_{z_{i}}$ are (up to an additive constant) branches of $\log \left(z-z_{i}\right), i=1, \ldots, r$. Recognizing Cauchy's Theorem as a statement entirely about integrals of meromorphic differentials (not of functions) immediately allowed generalizations. Here is what the abelian theory does for $U_{\boldsymbol{z}}$ (Chap. 2).
(6.2a) It gives explicit differentials providing details on integrals of any meromorphic differentials around any closed paths.
(6.2b) It describes elements of $\mathcal{E}\left(U_{\boldsymbol{z}}\right)^{\text {alg }}$ with associated group abelian.

Chap. 2 does $(6.2 b)$ by corresponding such functions to an $r$-tuple in $(\mathbb{Q} / \mathbb{Z})^{r}$ with entries summing to 0 .
6.2. Extending Cauchy's Theorem to $\bar{X}_{f}$. Riemann extended Cauchy's Theorem to $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$ not satisfying (6.1). Compatible with (6.2), he extended it to meromorphic differentials on $\bar{X}_{f}$. This emphasis on differentials over functions didn't throw functions out. They were still there through the definition of $d f$, the differential and $d f / f$, the logarithmic differential of $f$ (Chap. 3).

The serious step was analyzing the space of holomorphic (or first kind) differentials on $\bar{X}_{f}$ (Chap. 3, Chap. 4): differentials with no poles anywhere. Standard notation for this $g=g\left(\bar{X}_{f}\right)$ dimensional space over $\mathbb{C}$ is $\Gamma\left(\bar{X}_{f}, \Omega^{1}\right)$ : Global sections of the sheaf of holomorphic differentials on $\bar{X}_{f}$. The genus $g$ of $\bar{X}_{f}$ now attaches to $f=f(z)$, toward pinning its place among algebraic functions of $z$.

Guidance came from the Abel-Jacobi-Legendre differentials like $\frac{d z}{\sqrt{z^{3}+a z+b}}$ from $\S 3.2$. Just giving the dimension of $\Gamma\left(\bar{X}_{f}, \Omega^{1}\right)$ called for a more abstract approach. Riemann needed a full basis to solve the Jacobi-Inversion problem. Relying on coordinates from $\mathbb{P}_{z}^{1}$ was a confining kludge.

With points removed from $\bar{X}_{f}$, add further logarithmic (or 3rd kind) differentials. In $U_{\boldsymbol{z}}$, the (vector-)space of logarithmic differentials has a preferred basis by reference to classical generators of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)(\S 5.1)$.

Extending this to $\bar{X}_{f}$ still leaves an infinite set of choices for a $\Gamma\left(\bar{X}_{f}, \Omega^{1}\right)$ basis, with all choices related by the action of a group: The symplectic group $\mathrm{Sp}_{2 g}(\mathbb{Z})$. Different basis choices correspond to different choices of $2 g$ closed paths whose homology classes determine integration of any meromorphic differential around closed paths. This is an imprecise statement of Cauchy's Theorem on $\bar{X}_{f}$.

As with $U_{\boldsymbol{z}}$, there is a notion of classical generators. With $U_{\boldsymbol{z}}$ the paths were nonintersecting, except at the base point. On $\bar{X}_{f}$ classical generators signifies normalizing information about the intersection of these $2 g$ paths. Given classical generators for $U_{z}$ there is a process for producing classical generators on $X_{f}$. This provides explicit actions of appropriate subgroups of $H_{r}$ on the homology of $\bar{X}_{f}$. Suppose $\bar{X}_{f}$ appears in the moduli space of curves of genus $g$. Then, the whole action may well give $\mathrm{Sp}_{2 g}(\mathbb{Z})$. On Hurwitz spaces, however, the data is more refined. The significant group action may be much smaller.
6.3. Jacobians and generalizing Abel's Theorem. Suppose $\omega_{1}, \ldots, \omega_{g}$ is a specific $\Gamma\left(\bar{X}_{f}, \Omega^{1}\right)$ basis. As with $U_{z}$, Cauchy's Theorem on $\bar{X}_{f}$ builds from this data an abelian group. In this case it is a compact complex torus $J\left(\bar{X}_{f}\right)$, the Jacobian of $\bar{X}_{f}$. Follow Mumford's view [Mum76, p. 58-67]. Consider the space of locally defined holomorphic tangent vectors for $\bar{X}_{f}$ as dual to locally defined holomorphic differential forms (Chap. 3). Then, paths are dual to holomorphic differentials (by integration). The problem is to interpret a dual to global holomorphic differentials. This generalizes the Abel-Jacobi approach to Cauchy's Theorem and it produces an abelian covering theory.

Let $h$ be any meromorphic function on $\bar{X}_{f}$ (Chap. 3) of degree $u$. Then, $h$ : $\bar{X}_{f} \rightarrow \mathbb{P}_{z}^{1}$ has zeros $x_{1}^{0}, \ldots, x_{u}^{0}$ and poles $x_{1}^{\infty}, \ldots, x_{u}^{\infty}$. A mysterious identification then occurs: $\bar{X}_{f}$ appears in $J\left(\bar{X}_{f}\right)$. So, each zero $x_{i}^{0}$ and pole of $h$ produces a point in $J\left(\bar{X}_{f}\right)$. List these as $\boldsymbol{p}_{x_{i}^{0}}, \boldsymbol{p}_{x_{i}^{\infty}}, i=1, \ldots, u$.
6.3.1. Logarithmic differentials. Yet, finding the $\boldsymbol{p}$ s doesn't require giving $h$. It only needs points $x_{1}^{0}, \ldots, x_{u}^{0}$ and $x_{1}^{\infty}, \ldots, x_{u}^{\infty}$ on $\bar{X}_{f}$ viewed as inside $J\left(\bar{X}_{f}\right)$. Define $\left[D_{\boldsymbol{x}}\right]=\left[D\left(\boldsymbol{p}_{x_{i}^{0}}, \boldsymbol{p}_{x_{i}^{\infty}}, i=1, \ldots, u\right)\right]$ as the sum of all the $\boldsymbol{p}_{x_{i}^{0}} \mathrm{~S}$ minus all the $\boldsymbol{p}_{x_{i}^{\infty} \mathrm{S}}$ on $J\left(\bar{X}_{f}\right)$. To say $[D]$ is zero means it is the origin of $J\left(\bar{X}_{f}\right)$. Abel's Theorem (generalized) says existence of $h$ with these zeros and poles characterizes exactly when $[D]$ is zero.

If $h$ exists, consider the logarithmic derivative $d h / h$. This is a meromorphic differential of 3rd kind with pure imaginary periods. Even if $h$ doesn't exist, given the divisor $D_{\boldsymbol{x}}$ above, the following holds.
(6.3) There is a unique differential $\omega_{\boldsymbol{x}}$ with residue divisor $D_{\boldsymbol{x}}$ having pure imaginary periods (Chap. 4).
6.3.2. Coordinates from holomorphic differentials. Suppose $[D]$ is not zero, but $m[D]$ is zero on $J\left(\bar{X}_{f}\right)$ for some integer $m$. Then, repeating all the zeros and poles $m$ times produces a function $h$ on $\bar{X}_{f}$. The $m$ th root of $h$ defines an abelian unramified cover $Y \rightarrow \bar{X}_{f}$. So, the abelian theory of $\bar{X}_{f}$ appears from this version of Cauchy's Theorem. Riemann produced $\theta=\theta_{\bar{X}_{f}}$ functions to provide global coordinates (uniformization) for this construction. They are functions on $\mathbb{C}^{g}(\S 6.5)$.

Many mathematical items on $\bar{X}_{f}$ appear constructively from this. This includes functions and meromorphic differentials (with particular zeros and poles). This was a central goal in generalizing Abel's Theorem: To provide Abel(-Jacobian) constructions for a general Riemann surface. For the function $h$ it has this look:

$$
\begin{equation*}
h(x)=\prod_{i=1}^{u} \theta\left(\int_{x_{i}^{0}}^{x} \boldsymbol{\omega}\right) / \prod_{i=1}^{u} \theta\left(\int_{x_{i}^{\infty}}^{x} \boldsymbol{\omega}\right) . \tag{6.4}
\end{equation*}
$$

In $\theta$ you see $g$ coordinates; the $i t h$ entry is $\int_{x_{i}^{0}}^{x} \omega_{i}$. Each holds an integral over one basis element from $\boldsymbol{\omega}$. Integration paths join respective points on $\bar{X}_{f}$ 's universal covering space. The integrals make sense up to integration around closed paths. So, they define a point in $J\left(\bar{X}_{f}\right)$.

Even if $h$ doesn't exist, the logarithmic differential of (6.4) does. It gives the third kind differential from (6.3). Here you see the differential equation defining $\theta$ functions. In the expression for $h$, replace $\int_{x_{i}^{\infty}}^{x} \boldsymbol{\omega}$ by a vector $\boldsymbol{w}$ in the universal covering space of the Jacobian. Form the logarithmic differential of it: $d \theta(\boldsymbol{w}) / \theta(\boldsymbol{w})$. Translations by periods will change it by addition of a constant. With $\nabla$ the gradient in $\boldsymbol{w}, \nabla(\nabla \theta(\boldsymbol{w}) / \theta(\boldsymbol{w}))$ is invariant under the lattice of periods.

Thus, $J\left(\bar{X}_{f}\right)$ provides transparent coordinates for differentials, and their periods, through a mysterious embedding of $\bar{X}_{f}$ in it. Then, objects from the abelian structure on $J\left(\bar{X}_{f}\right)$ restrict to $\bar{X}_{f}(\S 10.6)$. To use, however, Riemann's theory an algebraist faces two major complications.
6.4. Complication 1: The role of $f$. Suppose $\bar{X}_{f}$ varies in the Hurwitz space $\mathcal{H}(G, \mathbf{C})$ attached to $(G, \mathbf{C})$. It moves along a path in $U_{r}$ with the coordinates for $\boldsymbol{z}$. Is Riemann's theory sufficiently algebraic to express the changes using equations with coefficients in the the point of $\mathcal{H}(G, \mathbf{C})$ corresponding to $\bar{X}_{f}$. Answer: It is algebraic in many ways, though rarely will coordinates from $\mathcal{H}(G, \mathbf{C})$ support all the identifications. Here is why.
6.4.1. The Picard components. There are three geometric ingredients in Riemann's theory: $J\left(\bar{X}_{f}\right), \bar{X}_{f}$ and the zero $(\Theta)$ divisor of the function $\theta=\theta_{\bar{X}_{f}}$ (§6.5). The first identifies with divisor classes $\operatorname{Pic}^{0}\left(\bar{X}_{f}\right)=\mathrm{Pic}_{f}^{0}$ of degree 0 on $\bar{X}_{f}$ (Chap. 4). The second embeds naturally (algebraically) in $\mathrm{Pic}_{f}^{1}$, divisor classes of degree 1 on $\bar{X}_{f}$. Then, $\Theta_{f}$ is the dimension $g-1$ variety of positive divisor classes in $\mathrm{Pic}_{f}^{g-1}$.

Further, $\mathrm{Pic}_{f}^{g}$ interprets the Riemann-Roch Theorem and the Jacobi Inversion Problem geometrically (Chap. 4). It takes its group structure from adding two positive divisors of degree $g$ together modulo linear equivalence. Weil used this for an algebraic construction of $\mathrm{Pic}_{f}^{0}$ years after his thesis. His principle: The nearly well defined addition on positive divisors produced a unique complete algebraic group on the homogeneous space of divisor classes. Therefore $\operatorname{Pic}_{f}^{0}$ is almost the symmetric product of $\bar{X}_{f}$ taken $g$ times. Riemann's theory was an inspiration to Weil's 1928 thesis (§10.6). Still, Weil was not certain until later that $\operatorname{Pic}_{f}^{0}$ and $\bar{X}_{f}$ have the same field of definition. This reminds that what now looks obvious is the result of many mathematical stories.
6.4.2. Half-canonical classes. All Picard components $\operatorname{Pic}_{f}^{k}$ are pair wise analytically isomorphic. Yet, finding an isomorphism analytic in the Hurwitz space coordinates may require moving to a cover of the Hurwitz space (§10.6).

Applying Riemann's theory directly requires having $\bar{X}_{f}$ and the $\Theta_{f}$ divisor on $\operatorname{Pic}_{f}^{0}$. For example, suppose there is an analytic assignment of a divisor class of degree $g-1$ on each curve $\bar{X}_{f}$ in the Hurwitz family. Then, translation of $\Theta_{f}$ by this divisor class puts it in $\mathrm{Pic}_{f}^{0}$. Here it would be available to construct the $\theta$ function. Convenient for this might be a half-canonical class: two times gives divisors for meromorphic differentials (Chap. 4).

Places marked by $\oplus$ in the Constellation Table of $\S 10.1$ signify inner Hurwitz spaces components that support such an assignment of half-canonical classes. This example shows how the Schur multiplier of a finite group appears in describing connected components of Hurwitz spaces (§10.2). It is a taste of the nilpotent theory arising in Modular Towers (§8.3). One last subtlety, however, occurs. Only some half-canonical translates work to give a formula like (6.4). They must be odd; the linear system has odd dimension (Chap. 4). This includes that $\theta(\mathbf{0})=0$ : When you plug in $x=x_{i}^{0}$ you expect $h\left(x_{i}^{0}\right)=0$. For the correct multiplicity of a zero on the right of (6.4), the gradient of the $\theta$ at $\mathbf{0}$ also must be nonzero. Such half-canonical classes always exist (Chap. 4).

Half-canonical classes, however, attached to $\oplus$ components in $\S 10.1$ are even. Sometimes they provide nontrivial $\theta$-nulls along the moduli space.

Riemann was even less algebraic in relating $\bar{X}_{f}$ and its Jacobian. He used coordinates from $\tilde{X}_{f}$, its universal covering space, to uniformize $\bar{X}_{f}$.
6.5. Complication 2: $\tilde{X}_{f}$ and nilpotent covers. The analytic isomorphism class of $\tilde{X}_{f}$ depends on the genus $g$ of $\bar{X}_{f}$. If $g=0$ it is the sphere, if $g=1$ it is $\mathbb{C}$ and it is the upper half plane $\mathbb{H}$ (or disk) if $g \geq 2$. As with $U_{\boldsymbol{z}}$ (§4.1), suppose we accept that $\tilde{X}_{f}$ is an analytic subspace of the Riemann sphere. Then, this comes from the Riemann mapping theorem. Still, it is not the uniformizing space we would expect. That would be $\tilde{X}_{f}^{\text {ab }}$, the quotient of $\tilde{X}_{f}$ by the subgroup of $\pi_{1}\left(\bar{X}_{f}\right)$ generated by commutators. This is the maximal quotient of $\tilde{X}_{f}$ that is an abelian cover of $\bar{X}_{f}$.
6.5.1. Abelian Frattini covers. Mathematics rarely looks directly at $\tilde{X}_{f}^{\text {ab }}$. It embeds in the universal covering space $\mathbb{C}^{g}$ of $J\left(\bar{X}_{f}\right)$. It is on $\mathbb{C}_{g}$ that $\theta_{\bar{X}_{f}}$ lives with its zeros, the $\Theta$ divisor, meeting $\tilde{X}_{f}^{\text {ab }}$ transversally. Periods of differentials on $\bar{X}_{f}$ translate $\tilde{X}_{f}^{\text {ab }}$ into itself. Yet, it is sufficiently complicated there seems to be no device for picturing it.

There are two models for picturing this. A standard picture shows the complex structure on a complex torus (like the Jacobian). It is of a fundamental domain (parallelpiped) in $\mathbb{C}^{g}$. Then, $2 g$ vectors representing generators of the lattice defining the complex torus (Chap. 3) give the sides of the parallelpiped. Inside this sits the pullback of $\bar{X}_{f}$. The geometry for this picture uses geodesics (straight lines) from the flat (Euclidean) metric defining distances on the complex torus.

Assume the genus of $\bar{X}_{f}$ is at least 2 . Then, the universal covering $\tilde{X}_{f}$ of $\bar{X}_{f}$ is the upper half plane $\tilde{X}_{f}$. A standard picture for $\bar{X}_{f}$ appears by grace of this having the structure of a negatively curved space. Geodesics here provide a polygonal outline of a set representing points of $\bar{X}_{f}$ (Chap. 4). Since $\tilde{X}_{f} \rightarrow \tilde{X}_{f}^{\text {ab }}$ is unramified, $\tilde{X}_{f}^{\mathrm{ab}}$ inherits a metric tensor with constant negative curvature. Yet, it sits snuggly in a flat space. Every finite abelian (unramified) cover $Y$ of $\bar{X}_{f}$ is a quotient of $\tilde{X}_{f}^{\text {ab }}$; it is a minimal cover of $\bar{X}_{f}$ with that property. Recall: We started with $\varphi_{f}: \bar{X}_{f} \rightarrow \mathbb{P}_{z}^{1}$. Assume it is a Galois cover, with group $G$.

Let $\mathcal{G}_{f}$ denote the abelian covers $\psi: Y \rightarrow \bar{X}_{f}$ with $\psi_{f}=\psi \circ \varphi_{f}: Y \rightarrow \mathbb{P}_{z}^{1}$ also Galois. Call $\psi_{f}$ a (relatively abelian) Frattini cover if the following holds. For any sequence $Y \rightarrow W \rightarrow \mathbb{P}_{z}^{1}$, of covers with $W \neq \mathbb{P}_{z}^{1}$, there is always a proper cover of $\mathbb{P}_{z}^{1}$ that $W \rightarrow \mathbb{P}_{z}^{1}$ and $\bar{X}_{f} \rightarrow \mathbb{P}_{z}^{1}$ factor through. A Frattini cover has no differentials and functions that pull back from covers disjoint from $\varphi_{f}$, so its function theory isn't accessible by knowing smaller degree covers. The most mysterious quotients of $\tilde{X}_{f}^{\text {ab }}$ are these relatively abelian Frattini covers.

This Frattini cover notion does not require an abelian cover $\psi$. Still, a Frattini cover arises always from $\psi$ being a Galois cover with nilpotent (a product of its $p$-Sylows) group.
6.5.2. No universal nilpotent cover. Relatively nilpotent Frattini covers produce natural sequences of moduli spaces generalizing sequences of modular curve covers (§8.3). Further, these moduli space sequences interpret many expectations about the regular version of the Inverse Galois Problem (§8). Relatively nilpotent covers and especially relatively Frattini covers bring up a combination of group theory and function theory. This includes many problems around new aspects of
the abelian theory using the Frattini property. This book explores aspects of it through these sequences of moduli spaces.

A complete understanding of all nilpotent (versus abelian) covers of $\bar{X}_{f}$ requires new, recent, ideas. An immediate difficulty is that there is no $\tilde{X}_{f}^{\text {nil }}$ similar to $\tilde{X}_{f}^{\text {ab }}$. Equivalently, no nontrivial subgroup of $\Gamma_{0}=\pi_{1}\left(\bar{X}_{f}, x_{0}\right)$ is in the intersection of all iterates of commutators in this group.

That is, let $\Gamma_{k}<\pi_{1}\left(\bar{X}_{f}, x_{0}\right)$ be elements of form $\left(g_{1}\left(g_{2}\left(\ldots g_{k-1}, g_{k}\right) \ldots\right)\right.$ with $g_{1}, \ldots, g_{k} \in \Gamma_{0}$. Only $\{1\}$ is in all the $\Gamma_{k}$ s. So, putting structure on the complete collection of algebraic nilpotent covers of $\bar{X}_{f}$ requires profinite limits. First consider how profinite limits appear in $G_{\mathbb{Q}}$ acting on points of the moduli spaces.

## 7. Acting with $G_{\mathbb{Q}}$

What changes in replacing $\boldsymbol{z}$ by $\boldsymbol{z}^{\prime}$, another $r$-tuple of elements? You might expect the fundamental group of $U_{\boldsymbol{z}}$ to tell nothing about changes. As a group it remains the same. We don't, however, use it as an abstract group. Its generators appear directly in applications. Changing $z$ forces changing generators. Yet, we understand the braiding changes from $H_{r}(\S 5.2)$. From elementary principles they give a profinite guide for action of $G_{\mathbb{Q}}$.
7.1. Acting on Laurent series. Suppose $\sigma \in G_{\mathbb{Q}}$ and $z_{0} \in \mathbb{Q}$. Assume $f(z)=\sum_{n=N}^{\infty} a_{n}\left(z-z_{0}\right)^{n}$ has coefficients in $\overline{\mathbb{Q}}$. Then, $\sigma$ acts on the $a_{n}$ s, producing $f_{\sigma}$. The hypothesis, however, of algebraic coefficients won't hold for $\tilde{f}_{z}$ from (4.3).
7.1.1. Setup for a test Case: $r=3$. Suppose $z_{1}, z_{2}, z_{3}$ are in $\mathbb{Q}$. Change the variable $z$ by an element of $\mathrm{SL}_{2}(\mathbb{Z})$ to map $\left\{z_{1}, z_{2}, z_{3}\right\}$ in some order to $\{0,1, \infty\}$. Six different permutations $\alpha \in \mathrm{SL}_{2}(\mathbb{Z})$ do this, depending on the order we choose. Composing $\tilde{U}_{z} \rightarrow U_{z}$ with one of these produces $\lambda: \tilde{U}_{z} \rightarrow U_{0,1, \infty}=\mathbb{P}_{\lambda}^{1} \backslash\{0,1, \infty\}$. Riemann's uniformization appears from a classical function, $\lambda: \mathbb{H} \rightarrow U_{0,1, \infty}$ (Chap. 4).
7.1.2. Uses for $\lambda(\tau)$. Periods of an antiderivative of $F(z)$ form an additive subgroup of $\mathbb{C}$ isomorphic to $\mathbb{Z} \times \mathbb{Z}(\S 3.2)$. In that notation, consider

$$
m(z, w)=w^{2}-z(z-1)(z-\lambda)
$$

with $\lambda \in \mathbb{P}_{z}^{1} \backslash\{0,1, \infty\}$. Choose $\tau \in \mathbb{H}$ so the function $\lambda$ takes $\tau$ to the value $\lambda$ (appearing in $m(z, w)$ ). Identify $\mathbb{Z} \times \mathbb{Z}$ with the subgroup $H_{\tau}$ of $\mathbb{C}$ that 1 and $\tau$ generate. Other choices of $\tau$ give the same lattice $H_{\tau}$. It only depends on $\lambda$. Let $\Gamma(2)$ be the group of integral matrices congruent to the identity matrix modulo 2. Suppose $\lambda\left(\tau_{0}\right)=\lambda_{0}$. Then, $\tau \mapsto \lambda(\tau)$ has as preimage of $\lambda_{0}$ the set $\Gamma(2)\left(\tau_{0}\right)=\left\{\alpha\left(\tau_{0}\right) \mid \alpha \in \Gamma(2)\right\}: \lambda$ uniformizes $\mathbb{H} / \Gamma(2)$.

Picard used $\lambda$ to show any nonconstant function $f(z)$ meromorphic on $\mathbb{C}$ excludes at most three values. Assume otherwise, and $f(\mathbb{C})$ excludes $0,1, \infty$. Then the monodromy theorem (Chap. 3) analytically continues $\lambda^{-1} \circ f$ to a function $\mathbb{C} \rightarrow \mathbb{H}$. The maximum modulus principle prevents existence of nonconstant holomorphic function maps $\mathbb{C}$ into the upper half plane. This contradiction shows $f$ must be constant [Ahl79, p. 307].
7.1.3. Another valuable function. Ordering the coordinates of $z$ violates some of our goals. The origins of the subject kept that in mind. Use the notation $U_{\lambda: 0,1, \infty}$ when the variable for $U_{0,1, \infty}$ is $\lambda$. Six elements of $\mathrm{PSL}_{2}(\mathbb{Z})$, forming a subgroup $S$, leave stable the set $\{0,1, \infty\}$. Then, $S$ acts on $U_{\lambda: 0,1, \infty}$. The quotient is $\mathbb{P}_{j}^{1} \backslash\{\infty\}=U_{j: \infty}$. The composite from $\mathbb{H} \rightarrow U_{j: \infty}$ is a Galois cover with group $\mathrm{PSL}_{2}(\mathbb{Z})$ (Chap. 4). It is ramified (not a topological cover) over fixed points of
elements in $\mathrm{SL}_{2}(\mathbb{Z})$ with eigenvalues 4 th or 6 th roots of 1 . We use $j(\tau)$ to display how Modular Towers of reduced Hurwitz spaces when $r=4$ (four elements in $\mathbf{C}$ ) generalize classical modular curves.
7.1.4. $G_{\mathbb{Q}}$ won't directly act on $\lambda$ and $j$. A theorem of Schneider-Siegel says $\tau\left(z_{0}\right)$ and $z_{0}$ are simultaneously algebraic only if $\tau$ is the ration of periods for an elliptic curve with complex multiplication. Therefore, even the constant term in the expansion of $\lambda^{-1}(z)$ around $z_{0}$ won't often by algebraic. That illustrates the extent previous generations sought to prove properties of $\lambda(\tau)$. Here, however, it shows using $\tilde{f}_{z}$ directly for the action of $G_{\mathbb{Q}}$ won't work.
7.2. Profinite fundamental groups. Suppose $X \rightarrow U_{\boldsymbol{z}}$ is a finite (unramified) cover, and $\boldsymbol{z}$ consists of algebraic points. Then, $X=X_{f}$ where $f$ has the following properties (Chap. 4).
(7.1a) It is defined by a nontrivial polynomial equation $m(z, f(z)) \equiv 0$.
(7.1b) $m=m(z, w)$ has algebraic coefficients.
(7.1c) $\frac{\partial m}{\partial w}\left(z_{0}\right)$ and $m\left(z_{0}, w\right)$ have no simultaneous zeros.

Apply the implicit function theorem (Chap. 2). It says $m(z, w)$ has $\operatorname{deg}_{w}(m)$ distinct zeros in $\mathcal{L}_{z_{0}}$. Conclude: Coefficients of $f(z)$ around $z_{0}$ are algebraic.
7.2.1. Grothendieck's Alternative. Define $\sigma \in G_{\mathbb{Q}}$ acting on a path $\gamma$ through what the result does to algebraic functions $f$ :

$$
f \mapsto f_{\sigma^{-1} \circ \gamma \circ \sigma}=f_{\gamma^{\sigma}}
$$

In words: Apply $\sigma^{-1}$ to the coefficients of $f$, analytically continue $f$ around $\gamma$ and then apply $\sigma$ to the coefficients of the result. The effect of $\gamma$ on algebraic functions determines it. So this determines $\gamma^{\sigma}$.

Problem 7.1. What does $\gamma^{\sigma}$ look like?
Only if $\sigma$ is complex conjugation $\epsilon$ will there be a path $\gamma^{\prime}$ (independent of $f$ ) so that represent $f_{\gamma^{\sigma}}=f_{\gamma^{\prime}}$. To see this, apply the theorem of Artin-Schreier: $\sigma$, if not complex conjugation $\epsilon$, either has infinite order or it is $\mu \epsilon \mu^{-1}$ where all powers of $\mu$ give distinct conjugates of $\epsilon$. Further, $\sigma$ and $\mu$ generate an uncountable subgroup of $G_{\mathbb{Q}}$. If all the $\gamma^{\sigma}$ s were paths, $\left\{\gamma^{\sigma^{\prime}}\right\}_{\sigma^{\prime} \in\langle\sigma\rangle}$ would have to be a countable, therefore finite, set. Simple considerations show this is impossible.
7.2.2. Where can we put $\gamma^{\sigma}$ ? Let $\overline{\mathbb{Q}}$ be the algebraic closure of $\mathbb{Q}$ in $\mathbb{C}$. The collection $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$ is in the Laurent series about $z_{0}$. With no loss we're allowed to assume the coefficients are in $\overline{\mathbb{Q}}$.

This gives an ordering: $f \leq g$ if $\overline{\mathbb{Q}}(z, g) \supset \overline{\mathbb{Q}}(z, f)$. Action of a path on $\overline{\mathbb{Q}}(z, g)$ determines its action on $\overline{\mathbb{Q}}(z, f)$. So, paths act on the equivalence classes and respect this ordering. Each equivalence class defines a specific function field inside $\mathcal{L}_{z_{0}}$. It is the exact data you get from a cover and a point on the cover over $z_{0}$. The ordering allows considering $\mathscr{P}_{z_{0}}$, projective systems of (algebraic) points over $z_{0}$. Thus, paths act on $\mathbb{P}_{z_{0}}$ (Chap. 4 or [Ihar91, p. 104]).

Proposition 7.2. This action on $\mathbb{P}_{z_{0}}$ determines paths in $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. The collection $\left\{\gamma^{\sigma}\right\}_{\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right), \sigma \in G_{Q}}$ also acts on $\mathbb{P}_{z_{0}}$. Define $\pi_{1}^{\text {alg }}$ to be the projective completion of this action. Then, $\pi_{1}^{\text {alg }}$ is the completion of $\pi_{1}$ by all normal subgroups of finite index. Further, $G_{\mathbb{Q}}$ acts on this.
7.3. Extending $G_{\mathbb{Q}}$ action. Extend the homomorphism $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right) \rightarrow G$ to $\psi_{\boldsymbol{z}^{\prime}, z_{0}}: \pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right)^{\text {alg }} \rightarrow G$. As a profinite group, $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right)^{\text {alg }}$ is also a free group on $r$ (topological) generators modulo one relation. Here, however, there are many more sets of classical generators.

For $G_{\mathbb{Q}}$ to act requires $\boldsymbol{z}$ is stable under $G_{\mathbb{Q}}$. Then, $G_{\mathbb{Q}}$ acts on $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right)^{\text {alg }}$ through Ihara's pro-braid group if $z_{0} \in \mathbb{Q}$. Again, recognize this action through its effect on classical generators of $\pi_{1}\left(U_{\boldsymbol{z}}^{\prime}, z_{0}\right)^{\text {alg }}$. Dependence on $z_{0}$ is so subtle, that any two distinct choices of $z_{0}$ give different actions. One remedy is to consider only the induced action of $G_{\mathbb{Q}}$ modulo inner automorphisms by $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}, z_{0}}\right)^{\text {alg }}$. Two further points guide investigations.
(7.2a) Unless the cover $X_{\boldsymbol{z}^{\prime}, z_{0}} \rightarrow \mathbb{P}_{z}^{1}$ coming from $\psi_{\boldsymbol{z}^{\prime}, z_{0}}$ is Galois and defined (with its automorphisms) over $\mathbb{Q}$, the action of $G_{\mathbb{Q}}$ won't respect $\psi_{z^{\prime}, z_{0}}$.
(7.2b) The action is so big, interesting properties of $G_{\mathbb{Q}}$ are hard to detect at the level of finite covers.
7.4. Motivation from the Inverse Galois Problem. Consider a finite group $G$ and the regular version of the Inverse Problem. It says for some $\boldsymbol{z}, G$ should be the group of a cover of $U_{z}$ with it and its automorphisms over $\mathbb{Q}$. That is, $G$ should be an $r$-branch point realization over $\mathbb{Q}$. To find $r$ and this cover needs structure.

You won't want to do one group at a time. So, we look at various quotients of $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}, z_{0}}\right)^{\text {alg }}$ with classical generators up to an action by $H_{r}$. Then, use $G_{\mathbb{Q}}$ action to investigate when there might be a value of $r$ and a corresponding $\boldsymbol{z}^{\prime}$ to realize such a quotient over $\mathbb{Q}$. Rather, however, than taking finite group quotients of $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}, z_{0}}\right)^{\text {alg }}$, take them maximally Frattini. Then dependence of $G_{\mathbb{Q}}$ action on $z^{\prime}$ has some uniformity. This gives the application generalizing modular curves Chap. 5 calls Modular Towers.

Start with a finite group $G$. Call a surjective homomorphism $\mu: H \rightarrow G$ Frattini if for any subgroup $H^{*} \leq H, \mu\left(H^{*}\right)=G$ implies $H^{*}=H$. This is the exact group translation of the cover property from $\S 6.5$.1. Suppose $\mu$ corresponds to a sequence of covers $\mu^{*}: X \rightarrow X / \operatorname{ker}(\mu) \rightarrow X / H$. Then, any proper cover $W$ appearing in the factorization $X \rightarrow X / H$ must factor properly through the cover $X / \operatorname{ker}(\mu) \rightarrow X / H$. A profinite group $\tilde{G}$ gives the maximal Frattini cover of $G$. All other group covers of $\mu: H \rightarrow G$ are targets for the map $\tilde{G} \rightarrow G$. Given $\psi: \pi_{1}\left(U_{\boldsymbol{z}}\right)^{\text {alg }} \rightarrow G$, a significant geometric invariant of $\psi$ is the set of maximal Frattini quotients of $\pi_{1}\left(U_{\boldsymbol{z}}\right)^{\text {alg }}$ (quotients of $\tilde{G}$ ) appearing as factors of $\psi$. These Frattini invariants interpret properties of the levels of Modular Towers. Their simplest instances refine Riemann's theory of $\theta$ characteristics (§10.1). They give many implications for the Inverse Galois Problem.

Conjugacy classes $\mathbf{C}$ hit by classical generators separate these homomorphisms discretely. This data gives structure to the problem. A preliminary investigation with $(G, \mathbf{C})$ from the Branch Cycle Lemma (Chap. 9, see $\S 8.2$ ) produces a necessary condition for a $(G, \mathbf{C})$ realization (over $\mathbb{Q}$ ). It is that $\mathbf{C}$ be a rational union of conjugacy classes.

## 8. Extensible nilpotent functions and the group $\tilde{G}$

We explain the universal Frattini cover $\tilde{G}$ of $G$ following the guide of Abel. He solved an inverse problem to part of the expression by radicals problem. This
produced dihedral group extensions, labeled by parameters still appearing in treatments of modular curves. For a prime $p, \mathbb{Z}_{p}$ denotes the $p$-adic numbers. Suppose $A$ and $B$ are two abelian groups. Assume elements of $A$ act as automorphisms of $B: a \in A$ acts on $b \in B$ giving $a(b)$. Then, form a group on $A \times B\left(\right.$ called $\left.A \times{ }^{s} B\right)$ using multiplication of $2 \times 2$ matrices:

$$
\left(\begin{array}{cc}
a & b \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
a^{\prime} & b^{\prime} \\
0 & 1
\end{array}\right)=\left(\begin{array}{cc}
a+a^{\prime} & a b^{\prime}+b \\
0 & 1
\end{array}\right)
$$

8.1. A guide from dihedral groups. Case: $G=D_{p}=\mathbb{Z} / p \times^{s}\{ \pm 1\}$ has $\mathbb{Z}_{p} \times^{s}\{ \pm 1\}$ and $\mathbb{Z}_{p} \times^{s} \mathbb{Z}_{2}$ as the pieces of its universal Frattini cover. Patch these together as a fiber product over $D_{p}$. This generalizes: For each prime $p$ dividing $|G|$, there is a universal $p$-Frattini cover ${ }_{p} \tilde{G}$ (Chap. 5). You can deal with one prime at a time. So, for investigating the arithmetic properties of quotients of $\pi_{1}\left(U_{z}\right)^{\text {alg }}$, consider the biggest quotients compatible with $r$ and $\mathbf{C}$ satisfying the Branch Cycle Lemma. Let $p$ be a prime. Recall: A conjugacy class in a finite group is called $p^{\prime}$ if its elements have order prime to $p$.

Certain properties of ${ }_{p} \tilde{G}$ suggest levels of a tower of moduli spaces.
(8.1a) ${ }_{p} \tilde{G} \rightarrow G$ has a pro-free pro-p group $\operatorname{ker}_{0}$ as kernel.
(8.1b) It has a characteristic sequence of quotients $G_{k}, k=0,1, \ldots$
(8.1c) Each $p^{\prime}$-conjugacy class of $G$ lifts uniquely to a $p^{\prime}$-conjugacy class of ${ }_{p} \tilde{G}$.
(8.1d) Elements of $G_{k}$ whose images in $G$ generate, already generate $G_{k}$.

Form $\operatorname{ker}_{1}$ as the closed subgroup of $\operatorname{ker}_{0}$ generated by $\operatorname{ker}_{0}^{p}$ and the commutators $\left(\operatorname{ker}_{0}, \operatorname{ker}_{0}\right)$. This gives $G_{1}$ in (8.1b) as the quotient ${ }_{p} \tilde{G} / \operatorname{ker}_{1}$. Continue inductively to form the other $G_{k} \mathrm{~s}$.
8.2. Applying the Branch Cycle Lemma. When there is profinite data, or over $\mathbb{R}$ or $\mathbb{Q}_{p}$, the explicit formula from the Branch Cycle Lemma is valuable.

Suppose $\sigma \in G_{\mathbb{Q}}$ maps to $n_{\sigma} \in \hat{\mathbb{Z}}^{*}=G\left(\mathbb{Q}^{\text {cyc }} / \mathbb{Q}\right)$. Find $\pi \in S_{r}$ to satisfy $z_{i}^{\sigma}=z_{(i) \pi}$. Then, a $(G, \mathbf{C})$ realization (over $\mathbb{Q}$ at $\left.\boldsymbol{z}\right)$ implies

$$
\begin{equation*}
C_{(i) \pi}^{n_{\sigma}}=C_{i}, i=1, \ldots, r \tag{8.2}
\end{equation*}
$$

Suppose the following:
(8.3) $\mathbf{C}$ consists of $r$ conjugacy classes whose elements have orders prime to $p$.

Note: Classes $\mathbf{C}$ from $G$ uniquely extend to $p^{\prime}$ classes in all $G_{k}$ s. Also, suppose $(G, \mathbf{C})$ passes Branch Cycle test (8.1). Then, so does $\left(G_{k}, \mathbf{C}\right)$ for all values of $k$. This illustrates a phenomenon: The groups $G_{k}$ are similar. So, they produce a guiding question.

Question 8.1. Are the $G_{k} \mathrm{~s}$ so similar their realizations fall to the Inverse Galois Problem with a $k$-free bound on the number of branch points?

The answer is conjecturally "No!" If you bound the number of branch points, there should be a bound on the values of $k$ for which $G_{k}$ has a $K$ regular realization where $K$ is a number field. Making this bound explicit, however, is another matter. The Mazur-Merel Theorem is well-known. It says, for any number field $K$, there is an explicit bound $C_{K}$ on $p^{k+1}$ so that for $p^{k+1}>C_{K}$, there are no non-cusp rational points on the modular curve $X_{1}\left(p^{k+1}\right)$. Below we see this interprets as the easiest special case of this conjecture: There are but finitely many four branch point, dihedral group involution realizations. The first step in the process forces us into
investigating the structure of some Modular Tower. An H-M (Harbater-Mumford) representative of $(G, \mathbf{C})$ is an $r$-tuple $\boldsymbol{g} \in \mathbf{C}$ with this property:
(8.4) $\langle\boldsymbol{g}\rangle=G$ and $g_{2 i-1}=g_{2 i}^{-1}, i=1, \ldots, s$ with $r=2 s$.

Approach the following statement by considering $r^{\prime}$ to be very large (say, two trillion). Then, consider if you can see a difference between the following cases.
(8.5a) $G$ is the monster (or use your favorite simple group) and $p=2$.
(8.5b) $G$ is $D_{5}$ and $p=5$.

ThEOREM 8.2. Fix $r^{\prime}$. Suppose there are $\left(G_{k}, \mathbf{C}_{k}\right)$ realizations over $\mathbb{Q}$ with $r_{k} \leq r^{\prime}$ conjugacy classes in $\mathbf{C}_{k}$, for each $k \geq 0$. Then, there exists $r \leq r^{\prime}$ and $p^{\prime}$-conjugacy classes $\mathbf{C}$ with $\left(G_{k}, \mathbf{C}\right)$ realizations over $\mathbb{Q}$ for all $k$.

If $p=2$, each $\left(G_{k}, \mathbf{C}\right)$ realization falls on a Hurwitz space component corresponding to an $H_{r}$ orbit containing $H-M$ representatives.
8.3. Thm. 8.2 and Modular Towers. Thm. 8.2 (Chap. 5) produces $p^{\prime}$ conjugacy classes $\mathbf{C}$ in ${ }_{p} \tilde{G}$ and a sequence $\left\{\boldsymbol{z}_{k}\right\}_{k=0}^{\infty}$ of $\mathbb{Q}$-stable unordered $r$-tuples of distinct points from $\mathbb{P}_{z}^{1}$. This sequence has the property that $z_{k}$ lies under a $\left(G_{k}, \mathbf{C}\right)$ realization. Further, suppose $p=2$. Then, the attached homomorphisms $\pi_{1}\left(U_{\boldsymbol{z}_{k}}\right)^{\text {alg }} \rightarrow{ }_{p} \tilde{G}$ send classical generators of $\pi_{1}\left(U_{\boldsymbol{z}_{k}}\right)^{\text {alg }}$ to H-M representatives in ${ }_{p} \tilde{G}$ so the induced quotient to $G_{k}$ has $G_{\mathbb{Q}}$-stable kernel.

Chap. 5 shows how this system of realizations fits into a system of moduli spaces generalizing classical modular curves. Consider all maps $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right) \rightarrow{ }_{p} \tilde{G}$ with generators $\boldsymbol{g}$ mapping to $\mathbf{C}$ as $\boldsymbol{z}$ runs over $U_{r}\left(z_{0} \notin \boldsymbol{z}\right)$. For each $k$ this produces an affine algebraic variety $\mathcal{H}_{k}$. Its $\mathbb{C}$ points correspond to equivalence classes of maps $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right) \rightarrow G_{k}$ (with $\boldsymbol{z}$ variable). The group $\mathrm{GL}_{2}(\mathbb{C})$ acts on these spaces. The quotient is another affine variety $\mathcal{H}_{k}^{\mathrm{rd}}$, level $k$ of the Modular Tower for $(G, \mathbf{C}, p)$.

A significant case: $G=D_{p}$ ( $p$ odd), $p$ the prime and $\mathbf{C}$ is $r=4$ repetitions of the conjugacy class of involutions (elements of order 2) in $D_{p}$. Then, $\mathcal{H}_{k}^{\text {rd }}$ is the modular curve $X_{1}\left(p^{k+1}\right)$ minus its cusps. Each case with $r=4$ produces a tower of curves, respective quotients of the upper half plane by finite index subgroups of $\mathrm{PSL}_{2}(\mathbb{Z})$. Usually the Modular Tower levels are noncongruence covers. They always have a useful moduli space structure.
8.4. A diophantine view of a nilpotent theory. Generalizations of theorems of Mazur and Serre now have formulations through the action of $G_{\mathbb{Q}}$ on projective systems of points on the spaces

$$
\begin{equation*}
\cdots \rightarrow \mathcal{H}_{k+1}^{\mathrm{rd}} \rightarrow \mathcal{H}_{k}^{\mathrm{rd}} \rightarrow \cdots \rightarrow \mathcal{H}_{0}^{\mathrm{rd}} \rightarrow U_{r}^{\mathrm{rd}}=J_{r} . \tag{8.6}
\end{equation*}
$$

Conjecture 8.3 (Main Conjecture). Suppose ( $G, \mathbf{C}, p$ ) is data for a Modular Tower. Assume $G$ is centerless and does not have $\mathbb{Z} / p$ as a quotient. For $k$ large, $\mathcal{H}_{k}^{\text {rd }}$ has no $\mathbb{Q}$ points.
8.4.1. Interpreting the Main Conjecture. Thus, $\mathbb{Q}$ realizations of $G_{k}$ require increasing large sets of conjugacy classes for $k$ large. This is more refined information than from any known versions of the Branch Cycle Lemma. If $p=2$, Thm. 8.2 says rational points will appear only on $\mathrm{H}-\mathrm{M}$ components of the sequence, and this refines the problem immensely. Changing $\mathbb{Q}$ to another number field $K$ requires significant generalization (Chap. 5).

Here is a response to the setup of cases from (8.5). Both require information on the geometry of Modular Tower levels we don't know yet. The dihedral group case
(with $r$ equal two trillion) looks easier because it translates to statements about classical moduli spaces: The moduli of cyclic $5^{k+1}$ degree covers of hyperelliptic curves (of genus $1,000,000,000,000-1$ ). No one knows if this space is without $\mathbb{Q}$ points for large $k$. Suppose the curves in the family have genus 1 . Then we know much since the Modular Tower levels are modular curves.

Yet, with the monster, there could be surprises. For example, for $\left(A_{6}, \mathbf{C}_{3^{5}}\right)$ with $p=2$, there are no $\mathbb{Q}$ points at level 1 of the Modular Tower. Reason: There are no points at level 1 at all, the result of the $\otimes$ symbol at $(6,5)$ in the Constellation Table of $\S 10.1$. The case $r=4$ gets much attention for problems that immediately generalize those for modular curves (§10.5).
8.4.2. Nilpotency from projective systems of points. Let $X$ be a compact Riemann surface. Denote the pro-p quotient of the fundamental group of $X$ by $\pi_{1}(X)^{(p)}$. When this group appears only up to inner automorphism, we drop the notation for the base point. Thm. 8.2 includes a nilpotent theory. Consider one of the homomorphisms $\psi_{\boldsymbol{z}}: \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }} \rightarrow{ }_{p} \tilde{G}$ mapping a fixed set of classical generators of into the $p^{\prime}$-conjugacy classes $\mathbf{C}$.

Let $X_{0} \rightarrow \mathbb{P}_{z}^{1}$ be the $G$ quotient cover from this homomorphism. For investigating all possible such maps $\psi_{\boldsymbol{z}}$, note it factors through a smaller quotient group of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$. This is an extension $M_{\boldsymbol{z}}$ (independent of $\psi_{\boldsymbol{z}}$ as a group extension) of $G=G_{0}$ by $\pi_{1}\left(X_{0}\right)^{(p)}$.

Call two such homomorphisms $M_{z} \rightarrow{ }_{p} \tilde{G} \rightarrow G_{0}$ inner equivalent if they differ by inner automorphisms from $\operatorname{ker}_{0}$ in (8.1a). Suppose $X_{0}=X_{p} \rightarrow \mathbb{P}_{z}^{1}$ corresponds to $\boldsymbol{p} \in \mathcal{H}_{0}^{\text {rd }}$. Projective systems of points on the Modular Tower over $\boldsymbol{p}$ correspond to inner homomorphism classes of $M_{z} \rightarrow{ }_{p} \tilde{G} \rightarrow G_{0}$. Shorten this phrase to a point on the Modular Tower. In this case refer to $M_{z}$ as $M_{\boldsymbol{p}}$. Let the set of inner homomorphism classes be $\mathcal{T}_{\boldsymbol{p}}$.

Homomorphisms factoring through ${ }_{p} \tilde{G}$, surjective to $G_{0}$, map surjectively to ${ }_{p} \tilde{G}($ from $(8.1 \mathrm{~d}))$. Let $g=g\left(X_{\boldsymbol{p}}\right)$ be the genus of $X_{\boldsymbol{p}}$ - transparent from $\mathbf{C}$ by the Riemann-Hurwitz formula. So, $\pi_{1}\left(X_{\boldsymbol{p}}\right)^{(p)}$ is a free pro- $p$ group on $2 g$ generators modulo one commutator relation.
8.4.3. $G_{\mathbb{Q}}$ action on $\operatorname{Ni}\left({ }_{p} \tilde{G}, \mathbf{C}\right)^{\text {in }}$. The notion of Nielsen class (§5.4) applies uniformly to $\left({ }_{p} \tilde{G}, \mathbf{C}\right)$. Its absolute and inner versions inherit an $H_{r}$ action. Orbits for this action correspond to projective systems of components at the levels of the Modular Tower. Reducing this action modulo ker $_{0}$ maps each orbit to an $H_{r}$ orbit at level 0 . Components of $\mathcal{H}_{k}^{\text {rd }}$ (over $\overline{\mathbb{Q}}$ ) map among each other by $G_{\mathbb{Q}}$ acting on the coefficients of their equations.

We don't often see equations for these moduli spaces. So, figuring this action from the data is one of our main problems. From this, regard $G_{\mathbb{Q}}$ as acting on the $H_{r}$ orbits in $\operatorname{Ni}\left({ }_{p} \tilde{G}, \mathbf{C}\right)^{\text {in }}$. In the $A_{n}$ examples of $\S 10.1$, there are components at a finite level $k$ that have no projective system of components above them. This could happen with any $(G, \mathbf{C})$. The invariant in $\S 9.1$ catches these obstructed components precisely, when you can compute it (Chap. 5).

Problem 8.4. Compute the $G_{\mathbb{Q}}$ action on $H_{r}$ orbits of $\operatorname{Ni}\left({ }_{p} \tilde{G}, \mathbf{C}\right)^{\text {in }}$. Also, compute the pattern of chains of obstructed components.
8.4.4. A nilpotent Tate Grassmanian. For $G$ any finite group this theory has a large pro-nilpotent part. Thus, it generalizes the abelian theory setup.

Suppose $\boldsymbol{z} \in U_{r}$ lies below a $\mathbb{Q}$ point $\boldsymbol{p} \in \mathcal{H}_{0}^{\text {rd }}$. Then, $G_{\mathbb{Q}}$ acts on $\pi_{1}\left(X_{\boldsymbol{p}}\right)^{(p)}$ (modulo inner automorphisms) as a quotient of the action on $\pi_{1}\left(X_{\boldsymbol{p}}\right)^{\text {alg }}$. Act by $G_{\mathbb{Q}}$ on the quotient of $\pi_{1}\left(X_{\boldsymbol{p}}\right)^{(p)}$ by the closed subgroup of commutators. Denote this quotient by $T_{\boldsymbol{p}}$, the Tate module for $p$. This gives the theory of abelian covers of $X_{\boldsymbol{p}}$ with group order a power of $p$. Its relation to the Jacobian of $X_{\boldsymbol{p}}$ is clear. It is the projective system of points of $p$-power order on the Jacobian.

Continue the actions of $G_{\mathbb{Q}}$. Suppose $\alpha$ is in $\mathcal{T}_{\boldsymbol{p}}$. Then, $\sigma \in G_{\mathbb{Q}}$ acts on $\alpha$ (on the right) through the composition $\alpha \circ \sigma$ (Chap. 9). There is a Lie algebra structure on $\pi_{1}\left(X_{\boldsymbol{p}}\right)^{(p)}$. Using it and the Weil pairing allows dualizing these maps. The result is $\mathcal{T}_{\boldsymbol{p}}^{*}$, a nilpotent version of $G_{\mathbb{Q}}$ acting on a Grassmannian of a Tate module of the Jacobian for $X_{p}$ (Chap. 5).

One goal of Modular Towers is to provide small actions for $G_{\mathbb{Q}}$. Modular Towers retains the feel of finite groups. Though a generalization of modular curves, the group theory reminds of situations yielding groups as Galois groups. Chap. 9 reviews achievements of that program, appearing in detail in [Se92], [MM95] and [Vö96] (see [Fri94]). In particular, the Dettweiler-Völklein generalization of Katz's rigid tuples [DVo98] pushes realization of Chevalley simple groups to a new place. It produces many cases with $G_{0}$ simple where $\mathbb{Q}$ points are dense in $\mathcal{H}_{0}^{\text {rd }}$.

These give a setting for $\widehat{\mathcal{G T}}$ relations close to the Inverse Galois Problem territory. Yet, the pro-finite elements of Modular Towers are like those of modular curve towers, suitable for checking the effect of these constraints. One goal is to see if $\widehat{\mathcal{G} T}$ relations force significant quotients of ${ }_{p} \tilde{G}$ to have $\mathbb{Q}$ realizations.

## 9. The Grothendieck-Teichmüller group

When $G_{\mathbb{Q}}$ acts on fundamental groups related to moduli spaces, that action preserves underlying geometry. Often that geometry is not obvious to us. So, asking what to expect from a $G_{\mathbb{Q}}$ action has us delving more deeply to where the geometry appears. The principle everyone uses occurs in divining components of a moduli space. The expectation is $G_{\mathbb{Q}}$ should map these components among each other, unless a geometric reason prevents it.
9.1. Moduli spaces with several components. The Constellation Table of $\S 10.1$ illustrates this. Superficially the two components appearing at the locus $(n, r)(r \geq n)$ have much in common. Action of $G_{\mathbb{Q}}$, however, on their equations leaves them fixed. Setup: The only alternative is it maps one of them to the other, because their union is a moduli space. Finish: The Schur multiplier invariant gives a geometric condition separating the components (§10.2.2).

Does $G_{\mathbb{Q}}$ have relations appearing everywhere in moduli space actions? These would induce relations for $G_{\mathbb{Q}}$ acting on all related moduli spaces (Chap. 5). The Grothendieck-Teichmüller group offers such relations. We discuss now the implication of these for the Inverse Galois Problem. Recall the space $J_{r}=U_{r} / \mathrm{PGL}_{2}(\mathbb{C})$ and its relative $\Lambda_{r}=\left(\mathbb{P}_{z}^{1}\right)^{r} \backslash \Delta_{r}$ when $r=4: \Lambda_{4}=U_{\lambda: 0,1, \infty}(\S 7.1 .1)$.
$\S 4.1$ has a description of the extensible algebraic functions $\mathcal{E}\left(\Lambda_{4}, \lambda_{0}\right)^{\text {alg }}$. Each starts from a Laurent series in $\lambda_{0}$ that analytically continues along any path in $\Lambda_{4}$.
9.2. Deligne's tangential base points. Deligne suggested an extra structure to $\mathcal{E}\left(\Lambda_{4}, \lambda_{0}\right)^{\text {alg }}$ by expanding the choices of base point [De89]. The elements of $\mathcal{L}_{\lambda_{0}}$ sit inside an algebraically closed field $\mathcal{P}_{z_{0}}$, convergent Puiseux expansions
around $\lambda_{0}$ (Chap. 2). They look like Laurent series in $\left(\lambda-\lambda_{0}\right)^{1 / e}$ for some integer $e$. They don't, however, work as functions in a neighborhood of $\lambda_{0}$ (Chap. 2).

Give the special case $\lambda^{1 / e}$ meaning by making it take positive values along the real axis pointing from 0 to 1 . This produces an analytic expression convergent in a neighborhood of any point on the positive real axis between 0 and 1. An alternative would ask $\lambda^{1 / e}$ to take positive values along the real axis in the negative direction from 0 to $-\infty$.

Distinguish between those two choices. Extend the meaning of the first to all Puiseux expansions about 0 using the notation $\mathcal{P}_{\overline{01}}$. Each produces a meromorphic function defined near 0 to the right of 0 . Similarly, for the second choice use the notation $\mathcal{P}_{\overline{0 \infty}}$. Each element in this defines a meromorphic function near 0 to the left of 0 . To be explicit, choose an open disk (on $\mathbb{P}_{\lambda}^{1}$ ). It should be symmetric about the real axis, tangent to the imaginary axis and contain part of the real axis from 0 to 1 (Chap. 2). Denote this disk $D_{\overline{01}}$.

For any $i$ and $j$, distinct elements from $\{0,1, \infty\}$ form the similar set of functions $\mathcal{P}_{\overline{i j}}$. The ordering from $\S 7.2 .2$ on algebraic functions in $\mathcal{L}_{0}$ extends to algebraic elements of $\mathcal{P}_{\overline{01}}$. So does the action of $G_{\mathbb{Q}}$ extend (Chap. 4).

Denote the set of ordered arrows by $\mathbb{B}$. Label the linear fractional transformations that permute $\{0,1, \infty\}: t_{\overline{i j}}$ takes $i$ to $0, j$ to 1 and $k$ to $\infty$. Apply $t_{\overline{i j}}{ }^{-1}$ to $D_{\overline{01}}$ to get similar disks $D_{\overline{i j}}$ attached to $\mathcal{P}_{\overline{i j}}$.

Principle 9.1 (Branch Extensibility). Consider $f \in \mathcal{E}\left(\Lambda_{4}, \lambda_{0}\right)^{\text {alg }}$ and $i, j$ distinct elements from $\{0,1, \infty\}$. Suppose $\gamma:[0,1] \rightarrow \Lambda_{4}$ is a path with $\gamma(0)=\lambda_{0}$ and $\gamma(1)$ in $D_{\overline{i j}}$. Then, there exists a unique $F_{f_{\gamma}} \in \mathcal{P}_{\overline{i j}}$ restricting to $f_{\gamma}$. The collection of order preserving maps on the equivalence classes of fields $\mathbb{C}\left(\lambda, F_{f_{\gamma}}\right)$ is $\pi_{\overline{01}}=\pi_{1}\left(\Lambda_{4}, \overline{01}\right)^{\text {alg }}$. It has a natural $G_{\mathbb{Q}}$ action (Chap. 4).

Let $x$ be a clockwise circle ([Ihar91] takes counterclockwise; see comments of $\S 11)$ around 0 meeting $D_{\overline{01}}$. It represents an element of $\pi_{\overline{01}}$ from Princ. 9.1. For example, suppose in the definition of $F_{f_{\gamma}}$ that $\gamma(1)$ is on $x$. Take $F=F_{f_{\gamma}}$ equal to $h\left(\lambda^{1 / e}\right)$ with $h$ meromorphic around 0 . Let $\zeta_{e}=e^{\frac{2 \pi i}{e}}$.

The effect of $x$ on $F$ is the substitution $\lambda^{1 / e} \mapsto \zeta_{e}^{-1} \lambda^{1 / e}$. So, $\sigma^{-1} \circ x \circ \sigma$ (following $\S 7.2 .1)$ gives this sequence of operations on a power series. Act on coefficients with $\sigma^{-1}$, then substitute $\zeta_{e}^{-1} \lambda^{1 / e}$, then act by $\sigma$ on the resulting coefficients. Use the notation of $\S 8.2$ : $n_{\sigma}$ is restriction of $\sigma$ to cyclotomic numbers. The total effect is the substitution $\lambda^{1 / e} \mapsto \zeta_{e}^{-n_{\sigma}} z^{1 / e}$. So, $x^{\sigma}=x^{n_{\sigma}}$.
9.3. The first two relations. Following [AnIh88], the $t_{\overline{i j}} \mathrm{~s}$ act on Puiseux expansions. So, they give maps among the fundamental groups $\pi_{\overline{i j}}$.
9.3.1. Continuations from $\overline{01}$ to $\overline{10}$. Extend this to the fundamental groupoid (Chap. 3), to give $\pi_{\overline{01 \overline{10}}}=\pi_{1}\left(\Lambda_{4} ; \overline{01}, \overline{10}\right)$. Let $\gamma_{p}:[0,1] \rightarrow \Lambda_{4}$ be a path running along $\mathbb{R} \cup\{\infty\}$ from 0 toward 1 , with $\gamma_{p}(0) \in D_{\overline{01}}$ and $\gamma(1) \in D_{\overline{10}}$. As with $x$ it defines an element of $\pi_{\overline{01 \overline{10}}}$.

Let $x^{\prime}$ be the transform of $x$ by $t_{\overline{\overline{10}}}(\lambda \mapsto 1-\lambda)$. Take $y=\gamma_{p} \circ x^{\prime} \circ \gamma_{p}^{-1}$. Then, $y$ represents an element of $\pi_{\overline{01}}$. Even easier than $x, \sigma^{-1} \circ y \circ \sigma$ has the effect of $\gamma_{p}^{\sigma}\left(x^{\prime}\right)^{n_{\sigma}}\left(\gamma^{-1}\right)^{\sigma}$. Let $m_{\sigma}=\gamma_{p}^{\sigma} \gamma_{p}^{-1}$. Then $y^{\sigma}$ equals $m_{\sigma} y^{n_{\sigma}} m_{\sigma}^{-1}$.

Since $x$ and $y$ are topological generators of $\pi_{\overline{01}}$, the effect of $\sigma$ on them determines the action of $\sigma$. It makes sense to write $m_{\sigma}(x, y)$. If $P_{1}$ and $P_{2}$ are two both homotopy classes of paths with the same end points, then they are conjugate.

Even though this is a profinite group, apply this to $\gamma_{p}$ and $\gamma_{p}^{-1}$. Therefore, $m_{\sigma}$ is a commutator in the pro-free group $x$ and $y$ generate.
9.3.2. The product-one relation. Most significant is what $\sigma$ does to $x y$. Equivalently: What is $z^{\sigma}$, with $z=(x y)^{-1}$ the 3rd element in a product-one relation (as in (5.1)). The formula for this comes from the first two Drinfeld-Ihara relations:
(9.1a) $m_{\sigma}(x, y) m_{\sigma}(y, x)=1$; and with $u_{\sigma}=\frac{n_{\sigma}-1}{2}$,
(9.1b) $m_{\sigma}(z, x) z^{u_{\sigma}} m_{\sigma}(y, z) y^{u_{\sigma}} m_{\sigma}(x, y) x^{u_{\sigma}}=1$.

Apply $t_{\overline{10}}$ to $m_{\sigma}(x, y)$ to see (9.1a). Let $r$ be the half-circle from the center of $D_{\overline{10}}$ to the center of $D_{\overline{1 \infty}}$ going clockwise. Then, $r$ defines an element of $\pi_{\overline{10}, \overline{1 \infty}}$. Expression (9.1b) comes from applying $\sigma$ to the geometric relation

$$
t_{\overline{1 \infty}}{ }^{2}\left(r \circ \gamma_{p}\right) \circ t_{\overline{1 \infty}}\left(r \circ \gamma_{p}\right) \circ\left(r \circ \gamma_{p}\right)=1
$$

We left out the famous 5-cycle relation [Ihar91, p. 107]. It forcefully appears soon.
9.3.3. Return of the $j$-line. There is a conspicuous quotient of the fundamental group of $\pi_{1}\left(\mathbb{P}_{j: 0,1, \infty}^{1}\right)(\S 7.1 .3)$. It has generators $\gamma$ :

$$
\gamma_{0}=q_{1} q_{2}, \quad \gamma_{1}=q_{1} q_{2} q_{1} \text { and } \gamma_{\infty}=q_{2}
$$

from a quotient of $H_{4}$ (Chap. 5 ; see $\S 5.3$ ). These satisfy
(9.2) $\gamma_{0}^{3}=1, \gamma_{1}^{2}=1, \gamma_{0} \gamma_{1} \gamma_{\infty}=1$; the group $\left\langle\gamma_{0}, \gamma_{1}\right\rangle$ is $\operatorname{PSL}_{2}(\mathbb{Z})$.

When $r=4$ a reduced Hurwitz space has a Riemann's Existence Theorem description coming from these generators acting on a reduced Nielsen class (Chap. 5). The geometry of the reduced Hurwitz spaces $\left\{\mathcal{H}_{k}^{\text {rd }}\right\}_{k=0}^{\infty}$ shows from analyzing $\gamma$. Most crucial are disjoint cycles of $\gamma_{\infty, k}$, the result of $\gamma_{\infty}$ in its action on $\mathrm{Ni}_{k}^{\mathrm{rd}}$.

Principle 9.2 (Cusp Principle). Each disjoint cycle of $\gamma_{\infty, k}$ corresponds to a cusp point for $\overline{\mathcal{H}}_{k}^{\text {rd }}$ over $j=\infty$. Further, each cusp has its own geometry.
9.4. Detecting $\widehat{\mathcal{G T}}$ at the level of a Modular Tower. Relations (9.1) have versions for action of $G_{\mathbb{Q}}$ on $\boldsymbol{\gamma}$. Yet, we must generalize them beyond their present shape to have them suit the geometry of a Modular Tower. Here is why.
9.4.1. Viewing tangential base points from $\mathbb{P}^{4}$. Deligne's tangential base points come from components of real points on $\left(\mathbb{P}_{z}^{1}\right)^{4} \backslash \Delta_{4}=U^{4}$. An example is $R_{z_{1}, z_{2}, z_{3}, z_{4}}$ : 4-tuples of distinct points on $\mathbb{R} \cup\{\infty\}=\mathbb{R}_{\infty}$ where the four points are in the same order as $(0,1, \infty,-1)$ around the circle. Rearrangements from permuting the elements $\left\{z_{1}, z_{2}, z_{3}\right\}$ produce new connected components. To get to $\mathbb{B}, \bmod$ out by the subgroup of $\mathrm{PSL}_{2}(\mathbb{R})$ stabilizing each component.

Formulas similar to (9.1) allow working directly with $H_{4}$. Replace elements of $\mathbb{B}$ with the image of $R_{z_{1}, z_{2}, z_{3}, z_{4}}$ in $\mathbb{P}^{4} \backslash D_{4}=U_{4}$. This is in [IM95] which also treats higher values of $r$.
9.4.2. Other real component configurations. The sets $R_{z_{1}, z_{2}, z_{3}, z_{4}}$ often fail to capture the cusp geometry on a Modular Tower. Here is an example. Real points on level 1 of the $\left(A_{5}, \mathbf{C}_{3^{4}}\right)$ Modular Tower $(\S 10.1, \S 10.3)$ lie on the genus 12 component of $\overline{\mathcal{H}}_{1}^{\text {rd }}$. Denote that $\overline{\mathcal{H}}_{1}^{+}$.

Real points on $\overline{\mathcal{H}}_{1}^{+}$collect in eight disjoint components, each associated to a cusp (of width 20). Four attach to $\mathrm{H}-\mathrm{M}$ representatives in this Nielsen class (§8.2). Let $C P_{z_{1}, z_{2}, z_{3}, z_{4}}=\left\{z_{1}, z_{2} \in \mathbb{H} \mid z_{3}=\bar{z}_{1}, z_{4}=\bar{z}_{2}\right\}$ : Two sets of complex conjugate pairs of points, with the first two in the upper half plane. The preimage in the inner Hurwitz space of these eight components is 32 real components. Each lies over a locus of real points in $U_{4}$ with preimage in $U^{4}$ of type $C P_{z_{1}, z_{2}, z_{3}, z_{4}}$.

There are $\left(\frac{4}{2}\right)=6$ choices for which two coordinates to put in the upper half plane. Then, counting possible lower half plane pairings with $z_{1}$ gives a total of 12 such real components of $U_{4}$. Action of $\mathrm{PSL}_{2}(\mathbb{R})$ on $C P_{z_{1}, z_{2}, z_{3}, z_{4}}$ has a hyperbolic description (Chap. 2). Put $z_{1}$ at $i$ under this action, so orbits of $z_{2}$ are points at a fixed (hyperbolic) distance from $i$.

To choose an explicit representative from each orbit, take a (hyperbolic) circle from $i$ to $i+1$ : A half circle perpendicular to the real axis through $i$ and $i+1$. Then, the interval description for Deligne's tangential base points has as analog the portion of the circle from $i$ to the right of $i$ going through $i+1$.
9.5. Variants of the Drinfeld-Ihara relations Chap. 9. There was a first definition of $\widehat{\mathcal{G T}}$. It was a subgroup of the automorphisms $\operatorname{Aut}\left(\hat{F}_{2}\right)$ of the profree group on generators $x$ and $y$.
9.5.1. $\widehat{\mathcal{G T}}$ : A moving target. $\widehat{\mathcal{G T}}$ 's elements are automorphisms of the form $(x, y) \mapsto\left(x^{n}, m y^{n} m^{-1}\right): n \in \hat{\mathbb{Z}}^{*}$ and $m \in\left(\hat{F}_{2}, \hat{F}_{2}\right)$ satisfy relations 9.1 (with the 5 -cycle relation). The composition of two is another automorphism. That this composition also satisfies the relations is more serious. This gives a group structure to such pairs $(n, m)$.

After the first definition, there was speculation $\widehat{\mathcal{G} T}$ might contain $G_{\mathbb{Q}}$ as an open subgroup. These days, however, $\widehat{\mathcal{G T}}$ presents a moving target. Recent joint work of Nakamura and Schneps reveals new relations satisfied by the image of $G_{\mathbb{Q}}$ in $\widehat{\mathcal{G T}}$. Its unclear whether to relabel $\widehat{\mathcal{G T}}$ appropriately for these relations or to start indexing a sequence of $\widehat{\mathcal{G T}}$-like groups. Yet, there are still only few of them and each is precious.
9.5.2. Cusps producing other base points. Consider $G_{\mathbb{Q}}$ acting as permutations on $H_{r}$ orbits of a reduced Nielsen class (§8.4.2).

Several steps are necessary to include $\widehat{\mathcal{G T}}$ type relations (Chap. 9). First: Develop corresponding relations from tangential base points using components like $C P_{z_{1}, z_{2}, z_{3}, z_{4}}$ (as suggested in [Fri95a, App. C-D]).

Second: Complete comparing with $\widehat{\mathcal{G T}}$ by extending this action to $\operatorname{Ni}\left({ }_{p} \tilde{G}, \mathbf{C}\right)^{\text {rd }}$. This works because $H_{4}$ acting on generators of the 4-punctured sphere identifies with a subgroup $H_{5}^{\prime}$ of $H_{5}$. We explain.

As in $\S 5.2$, consider $\left(\mathbb{P}_{z}^{1}\right)^{5} \backslash \Delta_{5}=U^{5}$. There is a fibration, $U^{5} \rightarrow U^{4}$ by projection on the first four coordinates. Embed $S_{4}$ in $S_{5}$ as the permutations leaving 5 fixed. Then, $S_{4}$ acts to give a new fibration, $U_{4} \times \mathbb{P}_{z_{5}}^{1} \backslash D_{5}^{\prime} \rightarrow U_{4}$ with $D_{5}^{\prime}$ the image of $\Delta_{5}$ in $\mathbb{P}^{4} \times \mathbb{P}^{1}$ (Chap. $5,[\mathbf{B F 8 2}],[\mathbf{D F r} 99]$ for the $\mathbb{R}$ analysis). Even without this quotient, analogs of all $\widehat{\mathcal{G T}}$ relations would appear. The fiber is a copy of $\mathbb{P}^{1}$ minus four points, with classical generators identified with words in $Q_{1}, \ldots, Q_{5}$. So, even analogs of the 5 -cycle relation (Chap. 9) show in identifying the $G_{\mathbb{Q}}$ action on $\operatorname{Ni}\left({ }_{p} \tilde{G}, \mathbf{C}\right)^{\text {rd }}$ when $r=4$.

Comparison between $\widehat{\mathcal{G T}}$ and Modular Towers then has these practical goals. Use all cusps on a Modular Tower to define the $\widehat{\mathcal{G T}}$ attached to that Modular Tower.

Problem 9.3. What do $\widehat{\mathcal{G T}}$ relations applied to Modular Towers detect about $\mathbb{Q}$ orbits on $\operatorname{Ni}\left({ }_{p} \tilde{G}, \mathbf{C}\right)^{\text {rd }}$. Compare with the Branch Cycle Lemma and $\omega$ (§10.2.2) invariant combination?

We conclude by tying together four advanced goals of the research motivating this book. It is convenient to do this by joining classical $\theta$-functions to Modular

Towers. Each diophantine element of this section gives specific detailed results on the Modular Towers of this example (Chap. 5).

## 10. Combining the Existence Theorem and $\theta$ functions

The first Hurwitz spaces were moduli spaces of simple branched covers. In this case the Hurwitz spaces are connected. An easy application of the Riemann-Roch theorem then shows connectedness of the moduli space of curves of genus $g$.
10.1. Theta functions and Hurwitz spaces. An example with many applications comes from covers with alternating groups $A_{n}$ as monodromy groups. Take $A_{n}, n \geq 4$, with the prime $p=2$ and 3 -cycles ( $r$ of them) as data for a Modular Tower. The usual representation $T_{n}$ gives an absolute space of degree $n$ covers with group $A_{n}$. There is a corresponding inner space of Galois covers (as in (5.5)). The following diagram displays the complete set of inner Hurwitz space components at level 0 of their Modular Tower.

Locations in this diagram have an attached integer pair $(n, r)$. The location shows components of the inner Hurwitz space for $\left(A_{n}, \mathbf{C}_{3^{r}}\right)$. Abbreviate this to $\mathcal{H}_{n, r}^{\text {in }}$. The corresponding absolute spaces would be for the data $\left(A_{n}, \mathbf{C}_{3^{r}}, T_{n}\right)$, or $\mathcal{H}_{n, r}^{\text {abs }}$. The group is the alternating group $A_{n}$. Conjugacy classes are $r$ repetitions of 3 -cycles. There is a famous spin group cover of $A_{n}, \tilde{A}_{n}$ where $\tilde{A}_{n} \rightarrow A_{n}$ is a central nonsplit extension with kernel $\mathbb{Z} / 2$. The universal 2-Frattini cover of $A_{n}$ (as in (8.1) automatically factors through $\tilde{A}_{n}$. This is a special case of a general phenomenon. The universal $p$-Frattini cover ${ }_{p} \tilde{G}$ of any perfect group $G$ factors through the universal central extension of $G$.

Labels for rows are by the genuses of the degree $n$ covers. The relation between the spaces $\mathcal{H}_{n, r}^{\text {abs }}$ and $\mathcal{H}_{n, r}^{\text {in }}$ comes from a corollary in [Fri96].

Proposition 10.1 (Absolute-Inner). The natural map $\mathcal{H}_{n, r}^{\mathrm{in}} \rightarrow \mathcal{H}_{n, r}^{\text {abs }}$ has degree 2. Each component of the former maps to a corresponding component of the latter.
10.1.1. Explanation of the symbols. Two primitive icons appear. The symbol $\otimes$ corresponds to a(n irreducible) component whose points represent covers $\hat{X} \rightarrow \mathbb{P}^{1}$ with this property. A special degree two unramified cover $\hat{Y} \rightarrow \hat{X}$ satisfies
(10.1) $\hat{Y} \rightarrow \mathbb{P}^{1}$ composed from $\hat{Y} \rightarrow \hat{X}$ and $\hat{X} \rightarrow \mathbb{P}^{1}$ is Galois with group $\tilde{A}_{n}$.

Then, $\oplus$ denotes a component of covers in $\mathcal{H}_{n, r}^{\text {in }}$ having no such $\hat{Y}$ cover. Excluding the genus 0 row, all rows have exactly two components. One is of $\otimes$ type,

Table 1. Constellation of spaces $\mathcal{H}\left(A_{n}, \mathbf{C}_{3^{r}}\right)^{\text {in }}$
Components correspond to lifting invariant values.
Genus at $(n, r)$ of degree $n$ cover: $g=r-n+1$
of the Galois cover: $g^{*}=\frac{(r-3) n!}{6}$

| $\xrightarrow{g \geq 1}$ | $\otimes \oplus$ | $\otimes \oplus$ | $\ldots$ | $\otimes \oplus$ | $\otimes \oplus$ | $\stackrel{1 \leq g}{\longleftrightarrow}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\xrightarrow{g=0}$ | $\otimes$ | $\oplus$ | $\ldots$ | $\otimes$ | $\oplus$ | $\stackrel{0=g}{\longleftrightarrow}$ |
| $n \geq 4$ | $n=4$ | $n=5$ | $\ldots$ | $n$ even | $n$ odd | $4 \leq n$ |

the other of $\oplus$ type. The spin group cover of alternating groups reveals its presence in components of Hurwitz spaces.
10.1.2. Subtleties about Schur multipliers. This phenomenon holds in general. Schur multipliers of finite groups produce distinct components of the Hurwitz space. For each conjugacy class C in $\mathbf{C}$, let $b_{\mathrm{C}}$ be its multiplicity of appearance in $\mathbf{C}$. A generalization of a Conway-Parker result has as hypothesis that $b_{\mathrm{C}}$ is suitably large for all C in C. Conclusion: Distinct components in level $k$ of a Modular Tower correspond exactly to elements in a subgroup of the Schur multiplier.

Yet, whether $b_{\mathrm{C}}$ is suitably large depends on $G_{k}$ (or on $k$ ) with $G=G_{0}$ fixed. This is the issue of $\S 10.2$. The Constellation Table shows level 0 of Modular Towers for all alternating groups with $p=2$ and 3 -cycle conjugacy classes.

Further, covers in one component differ from those in another in a simple striking way. Suppose $\hat{\varphi}_{\boldsymbol{p}}: \hat{X}_{\boldsymbol{p}} \rightarrow \mathbb{P}_{z}^{1}$ is a cover attached to $\boldsymbol{p} \in \mathcal{H}_{n, r}^{\mathrm{in}}$. Then the differential $d \hat{\varphi}_{\boldsymbol{p}}$ has a divisor of form $2 \hat{D}_{\boldsymbol{p}}$. (This happens whenever all elements in the conjugacy classes $\mathbf{C}$ have odd order.) The divisor $\hat{D}_{\boldsymbol{p}}$ is canonically defined over $\boldsymbol{p}$. Let $\operatorname{dim}\left(\hat{D}_{\boldsymbol{p}}\right)$ be the dimension of the space of meromorphic functions $h$ on $\hat{X}_{\boldsymbol{p}}$ for which $(h)+D_{\boldsymbol{p}} \geq 0$ (Chap. 3, Chap. 4).

So, $\hat{D}_{\boldsymbol{p}}$ defines a half-canonical divisor at each point on $\mathcal{H}_{n, r}^{\mathrm{in}}$, and a halfcanonical divisor class on $\mathcal{H}_{n, r}^{\text {in,rd }}$. A formula of Fried-Serre ([Fri96], [Ser90b]) says the components of $\mathcal{H}_{n, r}^{\text {in,rd }}$ separate according to $\operatorname{dim}\left(\hat{D}_{\boldsymbol{p}}\right)$ modulo 2 . For $r \geq n$, there is an $\oplus$ component of even half-canonical classes, the other of odd. For the components of $\mathcal{H}_{n, r}^{\text {abs }}$, define a similar divisor $D_{\boldsymbol{p}}$. Then, the formula for even or odd half-canonical classes is $\operatorname{dim}\left(D_{\boldsymbol{p}}\right)+r\left[\operatorname{Ser} \mathbf{9 0 b}\right.$, Thm. 3]. Note: When $X_{\boldsymbol{p}}$ has genus 0 , $\operatorname{dim}\left(D_{\boldsymbol{p}}\right)$ is 0 . Alternating $\otimes$ and $\oplus$ signs in the first row of the Constellation Table correspond exactly to $r$. $\S 10.2$ shows this is a small piece of an invariant applying to every Modular Tower.
10.2. Conjugacy class products. Examples show the Branch Cycle Lemma and $\omega$ invariant ( $\S 10.2 .2$ ) combination work well in this profinite context. Still, computing $\omega$ is not yet easy.
10.2.1. How modular representations appear. Computing the $\omega$ invariant for a Modular Tower relies on modular representation theory. The $\omega$ invariant is trivial for the usual modular curve tower. Here the kernel of ${ }_{p} \tilde{G} \rightarrow G$ is one dimensional $\left({ }_{p} \tilde{G}=\mathbb{Z}_{p} \times^{s}\{ \pm 1\}\right.$ and $\left.\operatorname{ker}_{0}=p \mathbb{Z}_{p}\right)$. It is, however, more interesting for Modular Towers in the Constellation Table of $\S 10.1$.

Consider the location $(5,4)$. Four repetitions of the conjugacy class $\mathrm{C}_{3}$ of 3cycles appear there. Here consider it a conjugacy class in ${ }_{2} \tilde{A}_{5}\left({ }_{p} \tilde{G}\right.$ for $A_{5}$ and $\left.p=2\right)$. As above, let $\mathrm{C}_{3}^{4}$ be all products of four elements from $\mathrm{C}_{3}$. Let $M_{k}$ be $G_{k} / G_{k+1}$, the $G_{k}$ module associated to level $k$. For any $G_{k}$ submodule $M_{k}^{\prime}$ of $M_{k}$ there is a quotient $G_{k+1} / M_{k}^{\prime}=G_{k}^{\prime}$. A special case is when $M_{k} / M_{k}^{\prime}=W_{k}$ is maximal for $G_{k}$ acting trivially on it.

Suppose $G_{0}$ is a perfect group (includes any simple group). Then, $W_{k}$ is the maximal exponent $p$ Schur multiplier of $G_{k}$ and $G_{k}^{\prime}=R_{k}$, the representation cover of $G_{k}$ (Chap. 5). This $A_{5}$ case has $p=2$ and $R_{k} \rightarrow G_{k}$ has kernel $\mathbb{Z} / 2$ for each $k$.

Let $O$ be an $H_{r}$ orbit of $\operatorname{Ni}\left(G_{k}, \mathbf{C}\right)$ with $\boldsymbol{g}$ a representative. Since $R_{k} \rightarrow G_{k}$ is a central extension, $\boldsymbol{g}$ has a unique lift to $\tilde{\boldsymbol{g}} \in R_{k}^{r} \cap \mathbf{C}$. If the product-one condition holds for $\tilde{\boldsymbol{g}}$, then it is in $\mathrm{Ni}\left(R_{k}, \mathbf{C}\right)$. Otherwise let $s(\boldsymbol{g}) \in W_{k}$ be the product of the
$\tilde{\boldsymbol{g}}$ entries. Running over all such orbits $O$ creates a subset $\mathbf{O b s}_{1, k}=\mathbf{O b s}_{1}\left(G_{k}, \mathbf{C}\right)$ of $W_{k}$ not containing the identity.

Suppose $O$ is an $H_{r}$ orbit with $s(\boldsymbol{g})=1$. Consider $M_{k+1} \lesseqgtr M_{k}^{\prime} \leq W_{k}$, with $M_{k}^{\prime}$ a $G_{k}$ submodule. Call $O$ obstructed at $M_{k}^{\prime}$ if these two properties hold.
(10.2a) $\boldsymbol{g}$ lifts to $\mathrm{Ni}\left(G_{k}^{\prime}, \mathbf{C}\right)$, but not to $\mathrm{Ni}\left(G_{k+1}, \mathbf{C}\right)$.
(10.2b) $M_{k}^{\prime}$ is minimal with this property.

From [FrK97, $\S 2]$ (or Chap. 5), $G_{k+1} / M_{k}^{\prime}$ has a nontrivial $p$ part to its Schur multiplier. Also, $M_{k}^{\prime}$ contains a proper submodule distinct from $\mathbf{1}_{G_{k}}$. Under these assumptions (running over allowable orbits $O$ ) put $M_{k}^{\prime}$ in the set $\mathbf{O b s}_{2, k}$. We state a problem only $\mathbf{C}=\mathbf{C}_{3^{r}}$ (general formulation in Chap. 5). The answer is not known even if $r=4$.

Problem 10.2 (Commutator Problem). Fix $r \geq 4$ even. What are the elements of $\mathbf{O b s}_{1, k}=\mathrm{C}_{3}^{r} \cap W_{k} \backslash\{1\}$ ? Suppose $k$ is large. Is this set just the identity? Then, the same question for $\mathbf{O b s}_{2, k}$ where we ask if it is empty for $k$ large.
10.2.2. Interpreting Problem 10.2. Notice the problem is about commutators. Suppose $r$ is even and C is any conjugacy class with $\mathrm{C}=\mathrm{C}^{-1}$. Then, elements of $\mathrm{C}^{r}$ are products of $r / 2$ commutators of form $\left(g, g^{\prime}\right)$ with $g, g^{\prime} \in \mathrm{C}$. Now assume $G_{0}$ is a perfect group. Then, so are the $G_{k} \mathrm{~s}$ for all $k$. Therefore, for $r$ large, all elements of $G_{k}$ are in $\mathrm{C}^{r}$. The crucial elements are in $W_{k}$ ? For example, make a graph on the group $G_{k}$. Elements of $G_{k}$ are the vertices, and edges are pairs $g_{1}, g_{2} \in R_{k}$ with $g_{1} g_{2}^{-1} \in \mathrm{C}$. As a function of $k$, what is the minimal distance between 1 and $W_{k} \backslash\{1\}$ ?

The sets $\mathbf{O b s} s_{1, k}$ and $\mathbf{O b s}{ }_{2, k}$ give a version of the $\omega$ invariant ( $\S 10.2 .2$, Chap. 5 , [Fri95a, Part III], [Ser90a]). This big invariant $\omega(O)$ is a collection of conjugacy classes in the kernel of ${ }_{p} \tilde{G} \rightarrow G_{0}$. An $H_{4}$ orbit that contributes to the sets $\mathbf{O b s}$ is obstructed; $O$ has nothing above it at level $k+1$. Suppose we know these sets and they determine the $\overline{\mathbb{Q}}$ components of $\mathcal{H}_{k}^{\text {rd }}$. Then, it is easy to compute definition fields of obstructed components contributing to $\mathbf{O b s}_{i, k} i=1,2$.
10.3. The diophantine effect of few components. Take $r=4$. Chap. 5 shows the genus of components in the sequence (8.6) goes up with $k$. That suffices to prove Conj. 8.3 when $r=4$. For example, level 0 of the $\left(A_{5}, \mathbf{C}_{3^{4}}\right)$ Modular Tower (§10.1) has one genus 0 component. Yet, level 1 has two components of respective genuses 12 and 9. The latter is obstructed [BFr02].

This one example illustrates the influence of Schur multipliers (equivalent to distinguishing $\theta$ characteristics). Why no obstructed component at level 0 , and then such appears at level 1? The Schur multiplier presence at level 1 comes from two same length (1152) $H_{4}$ orbits on $\mathrm{Ni}_{1}^{\mathrm{in}}$. So, the inner Hurwitz space has two absolutely irreducible components of the same degree as covers of $U_{4}$. Yet, they aren't conjugate under $G_{\mathbb{Q}}$. The $H_{4}$ orbits gave distinct permutation representations that show profoundly in the cusps of the reduced spaces cover $\mathbb{P}_{j}^{1}$.

Suppose $r=4$ and all components at some level of a Modular Tower have genus least 2 . This assures only finitely many points (no matter what is the number field $K$ ) at some level $k$. That does come from Falting's Theorem (the former Mordell Conjecture [Fal83]), though there are other older techniques that are more explicit about computing the exceptional values of $k[\mathbf{F r} \mathbf{0 2}, \S 5]$.

What, however, will help analyze levels of a Modular Tower when $r \geq 5$; they are no longer curves? We don't know. It would be valuable to show level $k$ components are varieties of general type for large $k$. Then, according to a conjecture of Lang, rational points on that level would lie in a lower dimensional subset. That would be progress, though not the quality of Conjecture 8.3.

More to the point would be a canonical height on a Modular Tower. Having in print background for developing this is an important goal of this book.
10.4. Height functions. Let $K$ be a number field. Let $\mathcal{H}_{k}^{\dagger}$ be the unobstructed components of $\mathcal{H}_{k}^{\text {rd }}(\S 10.2 .1)$. The goal is a function $H_{G, \mathbf{C}}: \mathcal{H}_{0}^{\text {rd }} \rightarrow \mathbb{R}$ whose properties prove Main Conjecture 8.3. That's simple enough and too much to expect. So, following [Fal83], aim for a finiteness result. Consider finding functions $H_{k}: \mathcal{H}_{0}^{\text {rd }} \rightarrow \mathbb{R}, k=0, \ldots$, with these properties.
(10.3a) $H_{k}(\boldsymbol{p})$ is nondecreasing in $k$ for each fixed $\boldsymbol{p}$.
(10.3b) For $k$ large it is positive on a nontrivial Zariski open subset $V_{k}$ of $\mathcal{H}_{0}^{\text {rd }}$.
(10.3c) $H_{k}$ is a sum of local height functions, one for each prime dividing $|G|$.
(10.3d) There are no $K$ points on $\mathcal{H}_{k}^{\dagger}$ over $V_{k}$.
(10.3e) When $r=4, \mathcal{H}_{k}^{\dagger}$ consists of finitely many curves. For $k$ large, $H_{k}$ should detect that the genus of all components of $\mathcal{H}_{k}^{\dagger}$ has gone beyond one.

Should such a function be effective? Bounding $k$ with $H_{k}$ not positive on an open set is only one critical problem. As important is to describe the nonordinary (see $\S 10.5$ ) locus that is the intersection of $\cap_{k}\left(\mathcal{H}_{0}^{\dagger} \backslash V_{k}\right)$. There also must be an overall measure using the branch points. The primes dividing $|G|$ contribute heavily to a measure of how branch points behave.
10.5. Introducing nonordinary points. We prefer to think of Conj. 8.3 as the Main Operating Conjecture. It's value is to find failures in nonobvious places. These would provide astounding realizations for Inverse Galois Problem. [FKVo98] has an example of a Chevalley group $G_{0}=\operatorname{PGL}_{n}(p)$ (with certain special $p$ and $n$ and conjugacy classes $\mathbf{C}$ (with $r=5$ ). The $p$-adic version $G^{\dagger}$ is a $p$-Frattini cover of $G_{0}$ (a common phenomenon, attested to in [Ser86]). It has characteristic quotients $G_{k}^{\dagger}$ formed as in (8.1). Then, there is a projective system of $\left(G_{k}^{\dagger}, \mathbf{C}\right)$ realizations (over some number field $K$ ).

Since $G^{\dagger}$ is a $p$-Frattini cover of $G_{0}$, it is the image of a map ${ }_{p} \tilde{G} \rightarrow G^{\dagger}$. Let ker* be the kernel of this map. So, this gives a $K$ point on a significant Modular Tower quotient. There is exactly one point in $\mathcal{H}^{\text {in,rd }}\left(G_{0}, \mathbf{C}\right)$ under a $K$ point in the tower. It would be proper to call such a point extraordinary. The literature, however, uses the name nonordinary point. Justifying that name, and locating nonordinary points and there corresponding Modular Tower quotients is a topic motivated by classical problems.
10.5.1. $\mathbb{R}$ contribution to height. Cusps of $\mathcal{H}_{k}^{\mathrm{rd}}$ guide us to the behavior at the real prime. Cusps attached to $\mathrm{H}-\mathrm{M}$ representatives give a degeneracy that goes with $\mathbb{R}$ contribution to the height function. This is what happens at level 1 of the $(4,5)$ location. Elementary techniques of $[\mathbf{B F r} \mathbf{0 2}]$ and $[\mathbf{D F r} 90]$ use uniformization of $\mathbb{R}$ points on Hurwitz spaces.

The less elementary part comes from interpreting them with group theory. Combining this with tangential base points as in $\S 9.4 .2$ allows analyzing new functions on a Modular Tower. This includes the even $\theta$-nulls from $\S 6.4 .2$, which relate to other functions:
(10.4a) half-canonical differentials on the space $\mathcal{H}^{\text {in,rd }}$; and (10.4b) Scholl's Eisenstein series associated with cusps [Scho86].

The cusp tangential base point geometry allows quantifying the amount of degeneracy as points of the moduli space approach the cusp. Cusps attached to $\mathrm{H}-\mathrm{M}$ representatives (as in (8.4)) support a total degeneracy. Including contributions for all cusps is still an open problem.
10.5.2. Combining geometry and function theory. Here is one development with modular curve precedents. Consider a Modular Tower (with $r=4$ ) and a degree 0 divisor $D$ supported in cusps of a component at some level of the tower. Sometimes such a $D$ generates a torsion group on the Jacobian of the tower component. Cases include when the support of $D$ consists of cusps associated to H-M representatives (as in $\oplus$ components of $\S 10.1$ ). We give a brief outline.

Under the hypotheses, consider the automorphic function $\theta_{0}$ on the reduced Hurwitz space coming from the $\theta$-nulls along the fibers of the family. Scholl associates to $D$ a sum $E_{D}$ of Eisenstein series. Since $D$ is a divisor on the curve giving the Modular Tower component, it corresponds to a logarithmic differential on this curve ( $\S 6.3 .1$ ). This is $E_{D}$.

So, following (6.4), our goal is to evaluate $E_{D}$ using $\theta$ functions. An example place would be the level 0 component $\mathcal{H}_{0:(5,4)}^{\text {rd }}=\mathcal{H}_{0}^{\text {rd }}$ of the Modular Tower at locus $(5,4)$ of the $\S 10.1$ Constellation Table. This component has genus 0 . Its Jacobian is trivial. So we don't mean a $\theta$ function on $\mathcal{H}_{0}^{\text {rd }}$ (or on $\mathcal{H}_{0}^{\text {odd }}$ where this computation really happens, see $\S 10.6$ ). Yet, it is much more than a genus 0 curve. It is a moduli space from whose points we gather data.

Evaluate a significant 3 rd kind differential such as $E_{D}$ from $\theta_{0}$ at each cusp tangential base point (as in $\S 9.4 .2$ ) in the support of $D$. As $\theta_{0}$ is canonically defined for a family over $\mathbb{Q}$, its expansion at the cusps has algebraic coefficients. A Theorem of Waldschmidt [Wa79] interprets this algebraic coefficient statement. It is equivalent to $D$ generating a torsion group in the Jacobian.

Since these components are moduli spaces, this has interpretations for the Inverse Galois Problem. Here is a low-brow corollary of the geometry in this story. There are exactly three regular $\mathbf{C}_{3^{4}}$ realizations (up to $\mathrm{SL}_{2}(\mathbb{Q})$ action) of the spin group cover of $A_{5}$. These realizations correspond to three points on the genus 1 pullback of $\mathcal{H}_{0}^{\text {rd }}$ to the $\lambda$-line. The cusps there generate a group of order 12 over $\mathbb{Q}$. Nine of those points are cusps, but three are not.

A bigger story, however, requires considering a curve $\hat{X}_{\boldsymbol{p}}$ (of genus 21) corresponding to a point $\boldsymbol{p}$ in the real locus of a tangential base point of $\S 9.4 .2$ type. Calculation of $E_{D}$ gives a measure of how $\hat{X}_{\boldsymbol{p}}$ degenerates (into unions of copies of $\mathbb{P}_{z}^{1}$ ) as it pushs along toward evaluation at the tangential base point. It is a bigger story because function theory informs about cusps on all projective systems of components in the Modular Tower. Height data involves all levels of a Modular Tower. Chap. 9 tells that story, related to [R77], [CTT98] and [GR78].

This focused example brings together function theory, geometry and arithmetic on a Modular Tower. It illustrates many potential applications of Modular Towers.
10.5.3. $p$ contribution to the height. This investigation comes from restricting the action of $G_{\mathbb{Q}}$ to $G_{\mathbb{Q}_{p}}, p$ the prime of the Modular Tower. After Hasse's invariant, the idea of nonordinary points for $p$ started with Serre-Tate theory ([Se72], [Se68]). Ihara used Hasse's invariant in examples that still inform us [Ihar00]. Mochizuki's use of canonical Frobenius elements defines the meaning of ordinary (and nonordinary?) directly [Moc96]. His theory, however, must descend from the moduli space of curves of genus $g$ to the moduli spaces in a Modular Tower. Defining and identifying nonordinary points on a Modular Tower is at the top of the problems this text aims at (Chap. 9).

In Ihara's approach the theory is entirely nilpotent. He has $p$-adic versions of classical functions. Especially, such have appeared from the action of $G_{\mathbb{Q}}$ on the second commutator quotient of $\pi_{\overline{01}}=\pi_{1}\left(\Lambda_{4}, \overline{01}\right)^{(p)}(\S 9.1)$. Coordinates arise from going to the induced Lie algebra actions. The rubric comes from Gassner-Magnus matrices. These give coordinates for the Lie algebra of an automorphism group acting on the second graded term of the Lie algebra of $\pi_{\overline{01}}^{(p)}$. Abelian covers of $\Lambda_{4}$ are Fermat curves. Similar to the discussion of $\S 8.4 .4$, this is a $p$-adic Lie algebra acting on the $p$-Tate module of Fermat curves. [Ihar91] describes the appearance of Jacobi sum grössencharacters.

These use partials (in the Lie algebra) of $m_{\sigma}(x, y)$ from (9.1). The IharaDrinfeld relations are vital here. Nakamura connects Ihara's example and another case: Replace $\Lambda_{4}$ by an elliptic curve minus one point. When it is an elliptic curve with bad degeneration at $p,[\mathbf{N a} 98]$ produces a Tate Eisenstein series. This is a logarithmic partial of Ihara's series. For some examples from the Constellation Table, the real Eisenstein series of $\S 10.5 .1$ have $p$-adic parallels to Nakamura's examples. This is what we mean by function theory on the nilpotent part of Modular Towers.

The nonlinear part, from $G_{0}$ still has a classical function relation as with $\theta$ invariants in $\S 10.1$. The nilpotent part, in examples, produces global functions on the moduli space. Specifically we expect these functions, at least those from $\mathrm{H}-\mathrm{M}$ representative cusps, to tell us about nonordinary points.
10.6. Weil's distributions. Look at (6.4) again. Weil's thesis constructed an analog of it: $(h(x)) \equiv \prod_{i=1}^{u} \theta_{x_{i}^{0}}^{w}(x) / \prod_{i=1}^{u} \theta_{x_{i}^{\infty}}^{w}(x)$. Here is its meaning. Both sides are fractional ideals in the ring of integers $\mathcal{O}_{K}$ of a number field $K$. The $\equiv$ sign means the left and right are equal up to a bounded fractional ideal. The left side is the principal fractional ideal that $h(x)$ generates. Most important, of course, are the functions $\theta_{x^{\prime}}^{w}: x \mapsto \theta_{x^{\prime}}(x)$ maps $K$ points $x$ into integral ideals. This function is defined only up to $\equiv$. Weil's distribution theorem allowed he (and Siegel [Si29]) to perform diophantine magic.

This works to define part of the height data for the commutative quotient of a Modular Tower. We explain. Denote the commutator subgroup of a profinite group $H$ by $(H, H)$. Replace inner homomorphism classes of $M_{\boldsymbol{p}} \rightarrow{ }_{p} \tilde{G}$ in $\S 8.4 .2$ by the sequence $M_{\boldsymbol{p}} /\left(\pi_{1}\left(X_{\boldsymbol{p}}\right), \pi_{1}\left(X_{\boldsymbol{p}}\right)\right) \rightarrow{ }_{p} \tilde{G} /\left(\operatorname{ker}_{0}, \operatorname{ker}_{0}\right)$. The question is now a refined question about subspaces of the Tate module of $J\left(X_{\boldsymbol{p}}\right)$.
[Si29] starts with a crude set of reductions by going to a finite extension of $K$. Doing this point-by-point along a Hurwitz space would be a disaster. Canonical heights avoid this. Here is a related allusion to the odd half-canonical classes.

Following a comment from $\S 6.4 .2$, we should replace $\mathcal{H}_{0}^{\text {rd }}$ by its pullback to a space $\mathcal{H}_{0}^{\text {odd }}$. Points of $\mathcal{H}_{0}^{\text {odd }}$ are pairs, $\boldsymbol{p} \in \mathcal{H}_{0}^{\text {rd }}$ with an odd half-canonical class on
$X_{\boldsymbol{p}}$. When the general point of $\mathcal{H}_{0}^{\text {odd }}$ carries a non-degenerate half-canonical class (§6.4.2) this starts an effective analysis. We still don't know what to do in the general case.

Here is a final word on even half-canonical classes. The locations in the Constellation Table with $\oplus$ support even half-canonical classes varying analytically with the coordinates of the Hurwitz space. Suppose the attached $\theta_{\boldsymbol{p}}$ is not zero at the origin of $J\left(\hat{X}_{\boldsymbol{p}}\right)$. Then, taking its value at the origin provides an automorphic form (the meaning is precise and conventional when $r=4$ ) on $\mathcal{H}_{0}^{\mathrm{in}, \mathrm{rd}}$ whose value appears in inspecting properties of the cusps.
10.7. Prelude to the general case? Level 1 of Constellation Table Modular Towers has further surprises related to the Schur multipliers of the level 1 groups. These illustrate practical applications of the nilpotent extension theory of covers (Chap. 9). There are lessons for group theory and geometry.

One is that nilpotent extensions (of any given group, simple or solvable) occur in many constructions with underlying geometric meaning. Such events don't naturally extend to solvable extensions much less to general (pro-)finite group theory. Consider lessons from the dihedral group $D_{p}$ and its association with the modular curve case of Modular Towers. It has a natural series of groups by changing the prime $p$ to any other prime: vary $p$ among primes. That isn't, however, so special.
10.7.1. Hecke operators. Consider the notation arising from $\S 8.1$ for the dihedral group $D_{p}=\mathbb{Z} / p \times^{s}\{ \pm 1\}$. Let $p^{\prime}$ be a prime distinct from 2 or $p$. The famous Hecke Operators of modular curve theory come from there being several values of $j\left(\tau_{1}\right), \ldots, j\left(\tau_{p^{\prime}+1}\right)$ for $\tau \in \mathbb{H}$ where $j\left(p^{\prime} \tau_{j}\right)$ is a particular value. This produces an algebraic correspondence represented by a curve $T_{p^{\prime}}$ on $X_{0}(N) \times X_{0}(N)$. A natural correspondence automatically induces an action on holomorphic differentials and cohomology, etc. Significantly, this correspondence produces a lift of the Frobenius correspondence from characteristic $p^{\prime}$ : The Eichler-Shimura congruence formula.

Here is how to interpret this from a Modular Tower viewpoint. Consider the Nielsen class $\mathrm{Ni}\left(D_{p}, \mathbf{C}_{2^{4}}\right)^{\text {abs }}=\mathrm{Ni}_{0}$. Suppose $\boldsymbol{g} \in \mathrm{Ni}_{0}$ is the branch cycle description of a cover $f_{\boldsymbol{g}}: \mathbb{P}_{w}^{1} \rightarrow \mathbb{P}_{z}^{1}$ with $D_{p}$ as monodromy group and involutions as branch cycles. This description comes from a choice of classical generators of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Then, the Galois closure of $f_{\boldsymbol{g}}$ is an elliptic curve $E$ which has a canonical degree $p$ isogeny to another elliptic curve $E^{\prime}$. Let $A_{p^{\prime}}$ be any cyclic subgroup of $p^{\prime}$ order on $E$ and let $A_{p^{\prime}}^{\prime}$ be its image in $E^{\prime}$. The morphism $E / A_{p^{\prime}} \rightarrow E^{\prime} / A_{p^{\prime}}^{\prime}$ modulo multiplication by -1 produces a new rational function $f_{\boldsymbol{g}, p^{\prime}}$. This is the genesis of the Hecke theory. It won't extend easily to a general Modular Tower. Yet, there are other candidates for constructions like the above.

Let $H$ be any finite group acting irreducibly on a $\mathbb{Z}$ module $V$ of rank $m$. Consider conjugacy classes $\mathbf{C}$ of $H$. (Take $H=\{ \pm 1\}$ and $V=\mathbb{Z}$ to get the dihedral group situation.) Consider the semi-direct product $V \times{ }^{s} H$ and then for each prime $p$, take $V / p V \times{ }^{s} H=H_{p}$. Suppose $(p,|H|)=1$. Extend the conjugacy classes to $H_{p}$. Then, apply the Modular Tower construction to ( $\left.H_{p}, \mathbf{C}, p\right)$.

Let $p^{\prime}$ be a prime distinct from those dividing $|H|$ and $p$. Add in $V / p^{\prime} V$ with an $H$ action to get $V / p G V \times V / p^{\prime} V \times^{s} H$ with an extension of the branch cycles $\mathbf{C}$ to this. This produces situations analogous to that for Hecke operators. This remains unexplored territory. A few examples will encourage further exploration. Examples of this type should give Modular Towers uniformizing natural collections of varieties defined over $\mathbb{Q}$, when the Branch Cycle Lemma conditions imply $\mathbb{Q}$
structures (§8.2). Given $H$ what varieties have such a natural uniformization? We haven't developed the expertise to consider this in detail. The value of making such a formulation is that all the arithmetic (including rational point statements) will fall under a uniform rubric. This would include using the Main Conjecture 8.3 on Modular Towers.
10.7.2. Separating the nilpotent tail and the nonnilpotent quotient. Group extensions of a given $G_{0}$ by a solvable group behave no better than general extensions of $G_{0}$. Roughly, the only distinction between solvable (excluding nilpotent) and general groups is that only cyclic groups appear as simple composition factors in solvable groups. That is the author's belief. With it goes the feeling that each finite nonnilpotent group $G_{0}$ generates its own intrinsic geometry. The discrete invariants of $\S 10.2$ capture much of this.

Then, there is a rich function theory appearing in the geometry from the nilpotent tail of a Modular Tower (as in §10.5). Together they separate the nilpotent tail from the nonnilpotent quotient. We believe this separation is natural and inevitable, and will never be breached. Further, our diophantine experience with problems involving solvable groups is that they belong more with the nonnilpotent quotient than with the nilpotent tail. We intend these comments to raise questions about modern understanding of Galois' famous theorem.

## 11. Aids to the reader and choice of actions

Expression numbers go from the left margin and most running lists use latin letters. For example, item 3 of expression 2 of section 5 of chapter 4 is (5.2c). Reference to it in another chapter would use the variant Chap. 4 (5.2c). Lemmas, corollaries, theorems, remarks, definitions, and examples fall under one collection of numbers: Definition 3 of section 9, written Def. 9.3, might follow Ex. 9.2. Figures have their own numbering system. Exercises appear as the last section in each chapter. References to these follow a special notation: Exercise 3, part c) appearing in section 9 appears as [9.3c]. Again, the chapter is given if it is in another chapter than that being read. Bibliographical items have notational shorthand for the author's(s') name(s), followed by a pinpoint reference, the usual $\mathrm{IA}_{\mathrm{E}} \mathrm{X}$ scheme. Like [Ahl79, p. 31].

There is sufficient material for a year course around two themes: fundamental groups in complex analytic geometry and families of Riemann surfaces. A third semester of complex analysis might cover just Chap. 2, Chap. 3 and Chap. 4. One year of complex analysis and one semester of graduate algebra are sufficient background. We assume undergraduate topology, as in a junior-senior analysis course, for proper background for the treatment of fundamental groups (Chap. 3).

The author spent much time considering on which side permutation groups would act. He chose the right side as the primary action side. That is, when $g \in S_{n}$ is an element of the symmetric group acting on integers $\{1, \ldots, n\}$, usually we write $g$ applied to $i$ as $(i) g$. It is not possible to be universally consistent. It is so typical to act with matrices on the left that with matrix groups we follow the usual convention. In making this decision there were these problems:

- Eventually, no matter the starting side, situations force simultaneous action on the other side.
- Group products in fundamental groups work with permutation representations only if you act on the right.
- Finite group theorists in the United States act on the right.

Many students trained by such books as [Lan71] and [Jac85] put group actions on the left. Neither book, however, does enough group theory training to dissuade from the need to spend considerable further time. Of course, there are always notational ways around the difficulties in any one situation.

The exposition on Riemann-Roch and the Picard groups in Chap. 4 quotes such sources as [FaKr90], [Mum76] and [Se59]. In addition, later examples quote finite group theory results outside the scope of this book. This goes with the book's aspiration to teach group theory interpretation, rather than detail. It simplifies exposition on examples to use references to [Vö96], in place of lengthy computations. The differences between the two books are large, ours geometric, while [Vö96] is more group theoretic. We, however, spend as much time on group theory. Our intention is to teach its use through examples to a generation of students interested in using Riemann surfaces who have little training in group theory. Still, the reader will recognize the two authors had more than a passing acquaintance.

## 12. Poetry and Mathematics

In the solipsistic world of mathematics, there are still many who find the subject matter of moduli of covers - that this book tackles - beautiful. The author agrees, with reservations.

Mathematics isn't poetry though Keats gave us hope it might be!
A thing of beauty is a joy for ever: its loveliness increases; it will never Pass into nothingness; but still will keep a bower quiet for us ...: From Endymion
12.1. The grandest virtues. The grandest virtue of mathematics is its modularity; That it builds from pieces. Second: That it lasts so long. An ingredient here is its independence of the framing secular language used. One easily sees the appearance of pythagorean triples in the Rind Papyrus. Yet, few would appreciate that the pyramid architect Imhotep was a god to the Egyptian Middle Kingdom.

Still, the converse of Keats' rhyme may not hold. The persistence of mathematics does not imply its beauty. The Durants suggest:

Poetry makes of language and feeling a music that cannot be heard across the frontiers of speech. [Du54, p. 77]
Independent of my abilities with written and spoken German, I can thrill to the simplifying structure Riemann brought to algebraic functions. Though I never think to tack a new verse onto Endymion, adding consequentially to Abel, Galois and Riemann is an ever present goal.
12.2. The eye of the beholder. Mathematical colleagues often don't appreciate the goals of other areas. One $\theta$ function adherant can't imagine the value of preoccupation with diophantine properties of large fields, and vice-versa. I'm speaking of co-writers I've known for over 30 years. It is one example of many.

If mathematicians sincerely fail to see the beauty of each others' grand enterprises, how could the world at large have the language and intellectual base to agree with what we think beautiful? In practice it is extremely difficult to explain the beauty of mathematics, even on occasion to a Nobel Prize winning Chemist; or to nonmathematical graduates of our elite institutions. Our perceptions can fail from
not appreciating the depth of what we already know before we address our papers. Failure to recognize the absorbed contribution of previous generations has much to do with the present hubris of today's mathematical community.

In particular, we (collectively) learned much from Abel, Galois and Riemann, though the first two produced very few theorems, and the third influenced mathematics through something strikingly beyond theorems. Abel and Galois used group theoretic interpretation to bring simplicity to an area littered with facts labeled as theorems. Riemann created coordinates for analyzing the details of a world of baffling geometries. All inherited and enhanced the goal of synthesizing algebra and geometry that Lagrange first articulated. In the age of specialization, we still recognize the coherency of mathematics in large part because of these people.

Mathematics is the only language supporting rich neologisms that bears its topics unadulterated to other areas and other generations. It overwhelms us locally in our seminars and colloquiums. Our students rail against what they think its incoherence, though its free inundating associations cause far more problems. The world, however, slowly accustoms to it, long forgetting - especially in related sciences - what a miracle of persistence is wrought by the foundation of clear definition. Definition that more than spotlights a resonant example; fluid definition that takes on new shape in each generation. In its fluidity it lasts and lasts and lasts. So we are certain, will the ideas of Abel, Galois and Riemann.
12.3. Two afterthoughts. The following found its motivation from $\theta$ functions and diophantine properties of large fields. There is an exact sequence [FrV92]:

$$
1 \rightarrow \tilde{F}_{\omega} \rightarrow G_{\mathbb{Q}} \rightarrow \prod_{n=2}^{\infty} S_{n} \rightarrow 1
$$

The group on the left is the profree group on a countable number of generators. The group on the right is the direct product of the symmetric groups, one copy for each integer. The absolute Galois group is caught between two known groups.

Here is a paraphrase from [Fri99, Acknowledgements]
The 20th century of mathematics belongs to group theory applications; I don't mean just Lie groups or classifications.

## CHAPTER 2

## ANALYTIC CONTINUATION

## 1. Why Riemann's Existence Theorem?

We start with two different definitions of algebraic functions. An imprecise version of Riemann's Existence Theorem is that these describe the same set of functions. Chap. 2 has two goals. First: To define and show the relevance of analytic continuation in defining algebraic functions. Second: To illustrate points about Riemann's Existence Theorem in elementary situations supporting the main ideas. Our examples are abelian algebraic functions. They come from analytic continuation of a branch of the log function. This also shows how integration relates algebraic functions to crucial functions that aren't algebraic. These examples depend only on homology classes, rather than homotopy classes, of paths. The slow treatment here quickens in Chap. 4 to show how Riemann's approach organized algebraic functions without intellectual inundation.
1.1. Introduction to algebraic functions. The complex numbers are $\mathbb{C}$, the nonzero complex numbers $\mathbb{C}^{*}$ and the reals $\mathbb{R}$. We start with analytic (more generally, meromorphic) functions defined on an open connected set $D$, a domain on $\mathbb{P}_{z}^{1}=\mathbb{C} \cup\{\infty\}$, the Riemann sphere: $\S 4.6$ defines analytic and meromorphic. The standard complex variable is $z$. When $D$ is a disk, a function $f(z)$ analytic on $D$ has a presentation as a convergent power series about the center $z_{0}$ of $D$. The first part of the book describes algebraic functions (of $z$ ). Let $D$ be any domain in $\mathbb{P}_{z}^{1}$ and $z_{0}, z^{\prime} \in D$. Denote (continuous) paths beginning at $z_{0}$ and ending at $z^{\prime}$ by $\Pi_{1}\left(D, z_{0}, z^{\prime}\right)(\S 2.2 .2)$. Use $\Pi_{1}\left(D, z_{0}\right)$ for closed paths in $D$ based at $z_{0}$. For any finite set $\boldsymbol{z}=\left\{z_{1}, \ldots, z_{r}\right\} \subset \mathbb{P}_{z}^{1}$ denote $\mathbb{P}_{z}^{1} \backslash\{\boldsymbol{z}\}$ by $U_{\boldsymbol{z}}$.
1.1.1. Riemann's definition of algebraic functions. Suppose $f(z)$ is analytic in a neighborhood of $z_{0}$. Call $f$ algebraic if some finite set $\boldsymbol{z} \subset \mathbb{P}_{z}^{1}$ has these properties.
(1.1a) An analytic continuation (Def. 4.1) of $f(z)$ along each $\lambda \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ exists. Call this $f_{\lambda}(z)$. Let $\mathcal{A}_{f}\left(U_{\boldsymbol{z}}\right)$ be the collection $\left\{f_{\lambda}\right\}_{\lambda \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)}$.
(1.1b) The set $\mathcal{A}_{f}\left(U_{z}\right)$ is finite.
(1.1c) For $z^{\prime} \in \boldsymbol{z}$, limit values of $f_{\lambda}$ along $\lambda \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}, z^{\prime}\right)$ is a finite set.
1.1.2. Standard definition of algebraic functions. There is another definition of algebraic function (of $z$ ). Suppose $f(z)$ is analytic on a disk $D$. It is algebraic if some polynomial $m(z, w) \in \mathbb{C}[z, w]$ (nonconstant in $w$ ) satisfies
(1.2) $m(z, f(z)) \equiv 0$ for all $z \in D$.

This chapter explains (1.1) and its equivalence with (1.2) (Prop. 7.3).
Simple examples illustrate (1.1) and (1.2). These often appear briefly in a first course in complex variables. Though they give only algebraic functions with abelian monodromy group, they hint how Chap. 4 lists all algebraic functions.

We review elementary field theory as it applies to $f(z)$ satisfying (1.2). With no loss assume $m(z, w)$ in (1.2) is irreducible in the ring $\mathbb{C}[z, w]$ ([9.8]) and $f(z)$ satisfies (1.2). Any graduate algebra book is proper for this review, including [Lan71], [Jac85] and [Isa94]. The latter, with the best treatment of permutation representations and group theory, will be our basic reference. [Isa94, Chap. 17] contains material supporting the comments of $\S 1.2$.
1.2. Equivalence of algebraic functions of $z$. Let $\mathbb{C}(z)$ be the field of the rational functions in $z$. Its elements $u(z)$ consist of ratios $P_{1}(z) / P_{2}(z)$ with $P_{1}, P_{2} \in \mathbb{C}[z]$. Standard notation denotes the greatest common divisor of $P_{1}$ and $P_{2}$ as $\left(P_{1}, P_{2}\right)$. Suppose $P_{1}$ and $P_{2}$ have no common nonconstant factor: Write this as $\left(P_{1}, P_{2}\right)=1$. Then the integer degree of $u(z), \operatorname{deg}(u)$, is $\max \left(\operatorname{deg}\left(P_{1}\right), \operatorname{deg}\left(P_{2}\right)\right)$. The Euclidean algorithm finds the greatest common divisor of $P_{1}$ and $P_{2}$. Factor this out to compute $\operatorname{deg}(u)$. This degree is also the degree of the field extension $\mathbb{C}(z)$ over $\mathbb{C}(u(z)):[\mathbb{C}(z): \mathbb{C}(u(z))][9.3]$.

Suppose $L$ and $K$ are fields with $K \subset L$. The degree $[L: K]$ of $L / K$ is the dimension of $L$ as a vector space over $K$. Assume $L=K(\alpha)$ for some $\alpha \in L$. Then, [ $L: K$ ] is the maximal number of linearly independent powers of $\alpha$ over $K$ : the degree of $\alpha$ over $K$. This degree is also the minimal positive degree of an irreducible polynomial $f_{\alpha}(w) \in K[w]$ having $\alpha$ as a zero. Up to multiplication by a nonzero element of $K, f_{\alpha}(w)$ is unique. If $L / K$ is a field extension, $\alpha \in L$ is algebraic over $K$ if $[K(\alpha): K]<\infty$.
1.2.1. The degree of $\mathbb{C}(z) / \mathbb{C}(u(z))$. Introduce variables $z^{\prime}$ and $w^{\prime}$. Write $u\left(w^{\prime}\right)$ as $P_{1}\left(w^{\prime}\right) / P_{2}\left(w^{\prime}\right)$ with $\left(P_{1}, P_{2}\right)=1$, and

$$
m\left(z^{\prime}, w^{\prime}\right)=P_{1}\left(w^{\prime}\right)-z^{\prime} P_{2}\left(w^{\prime}\right) \in \mathbb{C}\left[z^{\prime}, w^{\prime}\right]
$$

Then, $m\left(z^{\prime}, w^{\prime}\right)$ is irreducible of degree $n=\max \left(\operatorname{deg}\left(P_{1}\right), \operatorname{deg}\left(P_{2}\right)\right)$ [9.3]. Consider $m\left(z^{\prime}, w^{\prime}\right)$ as a polynomial in $w^{\prime}$ with coefficients in the field $\mathbb{C}\left(z^{\prime}\right)$. Let $w^{\prime \prime}$ be a zero of this polynomial in some algebraic closure of $\mathbb{C}\left(z^{\prime}\right)=K$. Then, $L=\mathbb{C}\left(z^{\prime}\right)\left(w^{\prime \prime}\right)=$ $\mathbb{C}\left(w^{\prime \prime}\right)$ is the quotient field of the integral domain $R=K\left[w^{\prime}\right] /\left(m\left(z^{\prime}, w^{\prime}\right)\right)$. It is a degree $n$ extension of $\mathbb{C}\left(z^{\prime}\right)$. Now $\mathbb{C}\left(z^{\prime}\right)$ is isomorphic to $\mathbb{C}(u(z))$ : map $z^{\prime}$ to $u(z)$. Map $w^{\prime \prime}$ to $z$ to extend this to an isomorphism of $L$ with $\mathbb{C}(z)$.
1.2.2. Degree of function fields over $\mathbb{C}(z)$. $\S 1.2 .1$ uses Cauchy's abstract production of $\mathbb{C}\left(z^{\prime}\right)\left(w^{\prime \prime}\right)$ with $w^{\prime \prime}$ a zero of $m\left(z^{\prime}, w^{\prime}\right)$ [Isa94, Lem. 17.18]. It, however, explicitly identifies $w^{\prime \prime}$ with $z$ and $z^{\prime}$ with $u(z)$. Putting $L$ in $\mathbb{C}(z)$, a genus 0 or pure transcendental field over $\mathbb{C}$, is convenient for seeing the algebraic relation between functions - like $z^{\prime}$ and $w^{\prime \prime}$.

Now assume $f(z)$ is any algebraic function according to (1.2). Similarly construct $L=\mathbb{C}(z, f(z))$, a degree $\operatorname{deg}_{w}(m(z, w))$ field extension of the rational functions $\mathbb{C}(z)$. This is the algebraic function field of $m$ (or of $f$ ). Call any $f^{*} \in L$ with $L=\mathbb{C}\left(z, f^{*}\right)$ a primitive generator of $L / \mathbb{C}(z)$. (Or, $f$ is just a primitive generator when reference to $z$ is clear.)
1.2.3. Equivalence of presentations of $L / \mathbb{C}(z)$. Infinitely many algebraic functions $f$ gives the same field $L$ up to isomorphism as an extension of $\mathbb{C}(z)$. Within a fixed algebraic closure of $\mathbb{C}(z)$ it is abstractly easy to list all primitive generators of $L$. They have the form $f^{*}=g\left(z, f_{k}\right)$ with $f_{k}$ any other zero of $m(z, w)$ and $g(z, u) \in \mathbb{C}(z)[u]$. To assure $\mathbb{C}\left(z, f^{*}\right)=L$ add that $\left[\mathbb{C}\left(z, f^{*}\right): \mathbb{C}(z)\right]=[L: \mathbb{C}(z)]$. Riemann's Existence Theorem lists algebraic extensions of $\mathbb{C}(z)$ efficiently by listing the isomorphism class of extensions $L / \mathbb{C}(z)$ and not specific algebraic functions.

Suppose $\mathbb{C}(f(z))$ contains $z$. Then, $L=\mathbb{C}(f(z))$ is pure transcendental. So, it is easy to list (without repetition) generating algebraic functions. Even, however, when the total degree of $m$ is as small as $3, L$ usually is not pure transcendental field $[9.10 \mathrm{~g}]$. While listing generating functions of $L$ is then harder, it isn't our main problem. To identify when two function field extensions $L_{1} / \mathbb{C}(z)$ and $L_{2} / \mathbb{C}(z)$ are (or are not) isomorphic is more important. Two questions arise: Is $L_{1}$ isomorphic to $L_{2}$ ? If so, does the isomorphism leave $\mathbb{C}(z)$ fixed?

Abel handled these questions for cubic equations. His results would have been easy if $L$ was pure transcendental. This book includes applying Riemann's extension of Abel's Theorems. Riemann's Existence Theorem is the start of this extension.

Riemann's Existence Theorem foregoes having all algebraic functions within one convenient algebraic closure. There may be no unique algebraic closure of $\mathbb{C}(z)$ so useful as $\mathbb{C}$. $\S 1.3$ introduces the infinite collection of incompatible algebraically closed fields appearing in Riemann's Existence Theorem. Every algebraic function $f(z)$ appears in each of them.
1.3. Puiseux expansions. Consider the Laurent field $\mathcal{L}_{z^{\prime}}$ consisting of series $f(z)=\sum_{n=N}^{\infty} a_{n}\left(z-z^{\prime}\right)^{n}$, with $N$ any integer, possibly negative, where $f(z)\left(z-z^{\prime}\right)^{-N}$ is convergent in some disk about $z^{\prime}$. Elements of $\mathcal{L}_{z^{\prime}}$ define functions meromorphic at $z^{\prime}$. Then, $\mathcal{L}_{z^{\prime}}$ is a field, containing $\mathbb{C}(z)$ and we are familiar with it. It isn't, however, algebraically closed. To remedy that, for any positive integer $e$ form $\mathcal{P}_{z^{\prime}, e}$, convergent series in a variable $u_{e}$. Think of $u_{e}$ as $\left(z-z^{\prime}\right)^{1 / e}: u_{e}^{e}=z-z^{\prime}$.

For $e \mid e^{*}$ let $t=e^{*} / e$. Map $\mathcal{P}_{z^{\prime}, e}$ into $\mathcal{P}_{z^{\prime}, e^{*}}$ by substituting $u_{e^{*}}^{t}$ for $u_{e}$. Regard the union $\cup_{e=1}^{\infty} \mathcal{P}_{z^{\prime}, e}=\mathcal{P}_{z^{\prime}}$ as a field, the direct limit of the fields $\cup_{e=1}^{\infty} \mathcal{P}_{z^{\prime}, e}$ with its set of compatible generators $\left\{u_{e}\right\}_{e=1}^{\infty}$. Details on the following are in [9.9].

Lemma 1.1. Suppose $\mathcal{P}^{*} / \mathcal{L}_{z^{\prime}}$ is any field extension generated by a sequence of elements $\left\{u_{e}^{*}\right\}_{e=1}^{\infty}$ with these properties.
(1.3a) $u_{e}^{*}$ is a solution of the equation $u^{e}=z-z^{\prime}$.
(1.3b) $\left(u_{e e^{\prime}}^{*}\right)^{e^{\prime}}=u_{e}^{*}$ for all positive integers $e, e^{\prime}$ : compatibility condition.

Then, $u_{e} \mapsto u_{e}^{*}$ gives a canonical isomorphism between $\mathcal{P}^{*}$ and $\mathcal{P}_{z^{\prime}}$ that is the identity on $\mathcal{L}_{z^{\prime}}$. In particular, automorphisms of the Galois extension $\mathcal{P}_{z^{\prime}} / \mathcal{L}_{z^{\prime}}$ correspond one-one with compatible systems of roots of 1 .

The field $\mathcal{P}_{z^{\prime}}$ of Puiseux expansions around $z^{\prime}$ provides an explicit algebraically closed field extension of $\mathbb{C}(z)$. It is clear fractional exponents are necessary for an algebraic closure. It is harder to see they give an algebraically closed field (Cor. 7.5). The fields $\mathcal{P}_{z^{\prime}}$ and $\mathcal{P}_{z^{\prime \prime}}$ are isomorphic. Such an isomorphism, however, restricts to mapping $\mathbb{C}(z) \rightarrow \mathbb{C}(z)$ by $z \mapsto z-\left(z^{\prime \prime}-z^{\prime}\right)$. For comparing all algebraic functions of $z$ we usually must regard these algebraically closed fields as distinct. Each, in its own way contains the field of algebraic functions (using either (1.1) or (1.2)).

Comparing expressions for a given algebraic function embedded in different Puiseux fields leads to our precise version of Riemann's Existence Theorem.
1.4. Monodromy groups and the genus. Both definitions (1.1) and (1.2) readily attach a group $G_{f}$ to any algebraic function $f(z)$. Using an irreducible $m(z, w)$ from (1.2) (with $m(z, f(z)) \equiv 0$ ) it is the group of the splitting field of $m(z, w)$ over $\mathbb{C}(z)$ ([Isa94, p. 267] and [9.5]). The order of this group is the degree of the splitting field extension over $\mathbb{C}(z)$. Efficient use of group theory gives more structured information than describing field extensions. Knowing something
about the Galois group is usually better information than comes from looking at polynomial coefficients.
$\S 4.4 .1$ gives a geometric construction for $G_{f}$. Chap. 4 has this group as its main theme. This group reveals $\mathcal{A}_{f}(D)$ from (1.1) as the complete set of zeros $w$ of $m(z, w)$ (Prop. 6.4). Then, $G_{f}$ acts through analytic continuation. This representation of $G_{f}$ on $\mathcal{A}_{f}(D)$ (of degree $\operatorname{deg}_{w}(m(z, w)$ ) is discrete data from $f$. Discrete here means the group $G_{f}$ does not change with continuous changes in $\boldsymbol{z}$.

Every algebraic function $f$ has another integer attached to it, the genus of its function field (Chap. 4). If $L=\mathbb{C}(z, f(z))$ is isomorphic to $\mathbb{C}(t)$ for some $t \in L$, it has genus 0 as above. This means all genus 0 function fields are abstractly isomorphic. Note: The integer $[L: \mathbb{C}(z)]$ is rarely a good clue for computing the genus [9.3]. Abel's results allow viewing genus 1 function fields as similar to genus 0 function fields, though that similarity has limits. Crucial: Unlike genus 0 fields, there are many isomorphism classes of genus 1 function fields (over $\mathbb{C}$ ).

Abel's results allow listing isomorphism classes of genus 1 function fields, exactly as we list points of $\mathbb{P}_{z}^{1}$. That is, with a classical parameter $j$ replacing $z$, finite values of $j$ correspond one-one to isomorphism classes of genus 1 function fields. As with $\mathbb{P}_{z}^{1}$ the value $j=\infty$ requires special consideration. Even if $L$ has genus 1, we don't easily find where its corresponding $j$ value is in this list. Still, for many problems this is a satisfactory theory.

Riemann generalized much of Abel's Theorem to function fields of all genuses. Most difficult was his analog, for genus greater than 1 , of a parameter space for isomorphism classes of fields. Variants on its study continue today, and this book is an example.
1.5. Advantages of Riemann's definition. Defining branches of $z^{\frac{1}{e}}$ (§8.3) on any disk $D$ in $\mathbb{C} \backslash\{0\}$ gives a practical introduction to analytic continuation. This gives the simplest algebraic functions. Still, how would we have located $w=f(z)$ satisfying $f(z)^{5}-2 z f(z)+1=0$ by a similar definition? The field $\mathbb{C}(z, f(z))$, like $\mathbb{C}\left(z^{\frac{1}{e}}\right)$, is pure transcendental [9.3]. Yet, this is not obvious from a Puiseux expansion of $f(z)$ around some point.

Suppose $f(z)$ is a convergent power series satisfying (1.1). Can we expect to find data appropriate to its description?: The set $\boldsymbol{z}$ of exceptional values, and the finite group expressing there are but finitely many analytic continuations around closed paths. Excluding elementary examples, the Riemann's Existence Theorem approach suggests it doesn't pay to give functions by their power series. Elliptic functions (Chap. 4 §7.1) are a good example where the functions are explicit, though power series don't give their definition. Riemann's Existence Theorem uses group data to replace power series information about $f(z)$.

This is practical, computable information about algebraic equations making Riemann's approach useful to the rest of mathematics. Especially it gives a way to track complete collections of related algebraic functions. This is the story of moduli of families of covers. Abel used the modular function classical texts call $j(\tau)$ where $\tau$ is a complex number in the upper half plane. We refine and generalize this theme.

## 2. Paths

We assume elementary properties of the complete fields, the real numbers $\mathbb{R}$ and the complex numbers $\mathbb{C}$ as in [Rud76, Chap. 1], [Ahl79, §1.1-1.3].
2.1. Notation from calculus. For each positive integer $n$, let $\mathbb{R}^{n}$ (resp. $\mathbb{C}^{n}$ ) be the set of ordered $n$-tuples $\boldsymbol{x}=\left(x_{1}, \ldots, x_{n}\right)$ (resp. $\left.\boldsymbol{z}=\left(z_{1}, \ldots, z_{n}\right)\right)$ of real (resp. complex) numbers. The set $\mathbb{R}^{n}$ is a vector space over $\mathbb{R}$ : addition of $\boldsymbol{x}=\left(x_{1}, \ldots, x_{n}\right)$ and $\boldsymbol{y}=\left(y_{1}, \ldots, y_{n}\right)$ gives $\left(x_{1}+y_{1}, \ldots, x_{n}+y_{n}\right)$; and scalar multiplication of $\boldsymbol{x}$ by $r \in \mathbb{R}$ gives $r \boldsymbol{x}=\left(r x_{1}, \ldots, r x_{n}\right)$. The zero element (origin) of $\mathbb{R}^{n}$ is $\mathbf{0}=(0, \ldots, 0)$. The inner product of $\boldsymbol{x}$ and $\boldsymbol{y}$ is $\boldsymbol{x} \cdot \boldsymbol{y}=\sum_{i=1}^{n} x_{i} y_{i}$. The law of cosines (from high school trigonometry) interprets the dot product $\cdot$ to give the expression $|\boldsymbol{x}||\boldsymbol{y}| \cos (\theta)$ where $\theta$ is the (counter clockwise) angle from the side from $\mathbf{0}$ to $\boldsymbol{x}$ to the side from $\mathbf{0}$ to $\boldsymbol{y}$ in (a/the) plane containing $\mathbf{0}, \boldsymbol{x}, \boldsymbol{y}$. Define the distance between points $\boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^{n}$ to be

$$
|x-y|=\sqrt{(x-y) \cdot(x-y)}
$$

Here are simple properties of the distance function.
(2.1a) $|\boldsymbol{x}| \geq 0$ and $|\boldsymbol{x}|=0$ if and only if $\boldsymbol{x}=\mathbf{0}$.
(2.1b) $|\boldsymbol{x}-\boldsymbol{z}| \leq|\boldsymbol{x}-\boldsymbol{y}|+|\boldsymbol{y}-\boldsymbol{z}|$ for $\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{z} \in \mathbb{R}^{n}$ : the triangle inequality.

Thus, the distance function gives a metric on $\mathbb{R}^{n}$.
2.2. Elementary properties and paths. Multiplication of complex numbers is crucial, especially that each nonzero complex number has a multiplicative inverse. Still, vector calculus often appears in the study of analytic functions using the topological identification of $\mathbb{R}^{2}$ with $\mathbb{C}$. In standard coordinates: $(x, y) \in \mathbb{R}^{2} \mapsto x+i y=z \in \mathbb{C}$. Rephrase multiplication of complex numbers on elements of $\mathbb{R}^{2}: z_{1} \leftrightarrow\left(x_{1}, y_{1}\right)$ and $z_{2} \leftrightarrow\left(x_{2}, y_{2}\right)$ gives the association $z_{1} z_{2} \leftrightarrow\left(x_{1} x_{2}-y_{1} y_{2}, x_{1} y_{2}+x_{2} y_{1}\right)$. Beyond these properties we gradually introduce statements from a one semester graduate course in complex variables. Paths and integration, however, are so important, we pause for notation around integration of 1 -forms and Riemannian metrics.

For $a, b \in \mathbb{R}, a<b,[a, b]$ denotes the closed interval of $\mathbb{R}$ with $a$ and $b$ as end points. A path in $\mathbb{R}^{n}$ consists of a continuous map $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$ for some choice of $a$ and $b$ with $a<b$. That is, for each $t \in[a, b]$, there is a range value $\gamma(t)$, the point on the path at time $t$.

Integration around paths turns computations into first year calculus integrals or derivatives. Such integration extends to manifolds (Chap. 3) because they are pieces of $\mathbb{R}^{n}$ tied together. Since $\gamma(t)$ is a point of $\mathbb{R}^{n}$, it has coordinates. One standard notation for these coordinates is $\left(f_{1}(t), \ldots, f_{n}(t)\right)$ ( $f$ is for function). Another possible notation is $\left(x_{1}(t), \ldots, x_{n}(t)\right)$. We prefer $\left(\gamma_{1}(t), \ldots, \gamma_{n}(t)\right)$. The points $\gamma(a)$ and $\gamma(b)$ are, respectively, the initial and end points of the path. The path $\gamma$ is closed if $\gamma(a)=\gamma(b)$.
2.2.1. Derivatives of a path. Call $\gamma$ differentiable if

$$
\frac{d \gamma(t)}{d t}=\left(\frac{d \gamma_{1}(t)}{d t}, \ldots, \frac{d \gamma_{n}(t)}{d t}\right)
$$

the tangent vector to $\gamma$ at $t$, exists and is continuous for each $t \in[a, b]$. (Use onesided limits at the end points.) Reminder: $\frac{d \gamma(t)}{d t}$ is a point in $\mathbb{R}^{n}$. Interpret it as giving a direction and speed (length of the vector $\frac{d \gamma(t)}{d t}$ ) of travel along the path $\gamma$ at time $t$. We always insist $\gamma$ is continuous (to be a path).

Definition 2.1. Let $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$ be a path. For $a \leq a^{\prime}<b^{\prime} \leq b$ denote the restriction of $\gamma$ to $\left[a^{\prime}, b^{\prime}\right]$ by $\gamma_{\left[a^{\prime}, b^{\prime}\right]}$. Call $\gamma$ simplicial if for some integer $m$ there exist
$t_{0}=a<t_{1}<\cdots<t_{m-1}<t_{m}=b$ with $\gamma_{\left[t_{i}, t_{i+1}\right]}$ differentiable, $i=0, \ldots, m-1$. This includes $\gamma_{\left[t_{i}, t_{i+1}\right]}$ having a one-sided derivative at the end points.
2.2.2. Paths and connectedness. The notation $\Pi_{1}\left(X, x_{0}, x_{1}\right)$ denotes the collection of (continuous) paths in a topological space $X$, starting at $x_{0}$ and ending $x_{1}$. Write $\Pi_{1}\left(X, x_{0}\right)$ when $x_{0}=x_{1}$. We often need paths in integrals to be simplicial. When necessary, the text assumes this implicitly for $\gamma$, though we may merely write $\gamma \in \Pi_{1}\left(X, x_{0}, x_{1}\right)$. For analytic continuation, or integrating meromorphic differentials, simplicialness is necessary only for paths satisfying explicit conditions as in (Rem. 4.4). One subtle use of simplicial paths is to give classical generators of the fundamental group of $U_{\boldsymbol{z}}$ (Chap. 4).

If $\Pi_{1}\left(X, x_{0}, x_{1}\right)$ is nonempty, then $x_{1}$ is path-connected to $x_{0}$. This is an equivalence relation, and the equivalence classes are the path-connected components of $X$. For subspaces of manifolds (Chap. 3; in particular, subspaces of $\mathbb{R}^{n}$ ), the pathconnected components are the same as the connected components. Further, for our examples, using simplicial paths would define the same components. [Ahl79, p. 54-58] discusses connectedness at greater length.
2.3. Integrals along a simplicial path. Using simplicial paths guarantees existence of various integrals, including arc length and line integrals along $\gamma$. We explain this. Let $\gamma$ be a simplicial path in $\mathbb{R}^{n}$. Consider $T_{\gamma}:[a, b] \rightarrow \mathbb{R}^{2 n}$ defined by $t \mapsto\left(\gamma(t), \frac{d \gamma(t)}{d t}\right)$. Suppose $F=F(\boldsymbol{x}, \boldsymbol{y})=F\left(x_{1}, \ldots, x_{n}, y_{1}, \ldots, y_{n}\right)$ is defined and continuous on an open set containing the range of $T$. The integral

$$
\begin{equation*}
\int_{\gamma} F \stackrel{\text { def }}{=} \int_{a}^{b} F \circ T_{\gamma} d t \tag{2.2}
\end{equation*}
$$

exists, though $\frac{d}{d t}\left(\gamma_{i}(t)\right)$ may be undefined for finitely many $t[\mathbf{R u d 7 6}$, p. 126]. Here are two traditional cases.
(2.3a) $F=\sqrt{\boldsymbol{y} \cdot Q(\boldsymbol{x})(\boldsymbol{y})}$ with $Q(\boldsymbol{x}): \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ by $\boldsymbol{y} \mapsto Q(\boldsymbol{x})(\boldsymbol{y})$ linear in $\boldsymbol{y}$, where $Q(\boldsymbol{x})$ is a symmetric and positive definite matrix for each $\boldsymbol{x}$.
(2.3b) $F=G(\boldsymbol{x}) \cdot \boldsymbol{y}$ with $G=\left(G_{1}(\boldsymbol{x}), \ldots, G_{n}(\boldsymbol{x})\right): \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ a continuous function (vector field) defined on the range of $\gamma$.
Definition 2.2. Suppose $\gamma$ is a one-one function onto its range. Case (2.3a) of (2.2) is the arc length of $\gamma$ relative to the infinitesimal metric $Q(\boldsymbol{x})$ at $\boldsymbol{x}$. [9.19] explains the value of tensor form for metrics. In case (2.3b), (2.2) is the line integral of the differential one form $G \cdot d \boldsymbol{x}=\sum_{i=1}^{n} G_{i}(\boldsymbol{x}) d x_{i}$ along $\gamma$.

Here is the crucial point of these examples. Suppose we change $\gamma$ to another parameterization $\gamma^{*}$ of the same set. Then, (2.2) doesn't change modulo these conditions: $\gamma^{*}$ is one-one in case (2.3a); and $\gamma^{*}$ has the same beginning and end points as $\gamma$ in case (2.3b). Proving this uses Lemma 2.3 [9.19b].

Recall from vector calculus, the physical meaning of (2.3b). It is the work done in moving a particle along the path parametrized by $\gamma$ against the force field $G$. Here is the formula for computing integrals of such differential expressions along $\gamma$ :

$$
\begin{equation*}
\int_{\gamma} \sum_{i=1}^{n} G_{i}(\boldsymbol{x}) d x_{i} \stackrel{\text { def }}{=} \sum_{i=1}^{n} \int_{a}^{b} G_{i}\left(\gamma_{1}, \ldots, \gamma_{n}\right) \frac{d \gamma_{i}}{d t} d t \tag{2.4}
\end{equation*}
$$

Tensor form of a metric defines distance along $\gamma$ from an integral of positive functions [9.19]. The triangle inequality is automatic: $\int_{a}^{b} f(t) d t+\int_{b}^{c} f(t) d t \geq \int_{a}^{c} f(t) d t$ if $f(t) \geq 0$ for $t \in[a, c]$.

Lemma 2.3 (Change of Variable Formula). Let $\gamma:[c, d] \rightarrow \mathbb{R}$ be a simplicial path. Assume $f: \mathbb{R} \rightarrow \mathbb{R}$ is continuous, defined on the range of $\gamma$ and $a=\gamma(c)$, $b=\gamma(d)$. Then,

$$
\int_{a}^{b} f(x) d x=\int_{c}^{d} f(\gamma(t)) \frac{d}{d t}(\gamma(t)) d t
$$

Proof. This is a variant on [Apo57, p. 216]. Let $F(x)=\int_{a}^{x} f(t) d t$ for $x$ in the range of $\gamma$, and $H(x)=\int_{c}^{x} f(\gamma(t)) \frac{d}{d t}(\gamma(t)) d t$. The functions $F(\gamma(x))$ and $H(x)$ are both continuous. Excluding finitely many $x$, the chain rule shows they have the same derivatives. Thus $H(x)-F(\gamma(x))$ is a constant evaluated by taking $x=c$ :

$$
H(c)-F(\gamma(c))=H(c)-F(a)=0-0=0
$$

The formula follows by taking $x=d$.
Apostol notes: "Many texts prove the preceding theorem under the added hypothesis that $\frac{d \gamma(t)}{d t}$ is never zero on $[c, d]$. The interval joining $a$ to $b$ need not be the image of $[c, d]$ under $\gamma$."
2.4. Relation between integrals and analytic functions. Integration theory is the heart of complex variables. Equations, algebraic or differential, with coefficients analytic on a domain $D$, define the classical functions of complex variables. By a domain we mean an open connected topological subspace of a given topological space. The first examples of the subject are domains in $\mathbb{C}$, the complex plane. As we use them, we will remind of most basics from a first semester graduate complex variables course.

This chapter refers to basic material of [Ahl79]. The notation $\mathcal{H}(D)$ denotes the ring (integral domain [9.8a]) of functions analytic (equivalently, holomorphic) on $D$. With $R$ any ring, let $R[w]$ be polynomials in $w$ with coefficients in $R$.
2.4.1. Analytic Functions. The definition of analytic function reflects how the chain rule works for a composition of an analytic function and a path. Assume $\lambda:[a, b] \rightarrow D$ is any differentiable path: $t \mapsto \lambda_{1}(t)+i \lambda_{2}(t)$ has $\lambda_{1}=\Re(\lambda)$ and $\lambda_{2}=\Im(\lambda)$, differentiable on the interval $[a, b]$.

Definition 2.4. Suppose $z_{0} \in D, t_{0} \in[a, b]$ and $\lambda:[a, b] \rightarrow D$ is any path, differentiable at $t_{0}$, for which $\lambda\left(t_{0}\right)=z_{0}$. Then, $f(z)$ defined on $D$ is analytic at $z_{0}$ if there exists a complex number $M+i N$ dependent only on $f$ and $z_{0}$ with

$$
\begin{equation*}
\frac{d}{d t}(f \circ \lambda)\left(t_{0}\right)=(M+i N) \frac{d \lambda}{d t}\left(t_{0}\right) \tag{2.5}
\end{equation*}
$$

To compute the derivative on the left, assume $f(z)=u(x, y)+i v(x, y)$ has partial derivatives (not necessarily continuous) and use the chain rule.

Apply (2.5) to $t \mapsto z_{0}+\left(t-t_{0}\right) \boldsymbol{v}$ in two cases: $\boldsymbol{v}=1$ and $\boldsymbol{v}=i$. This produces $t_{\text {wo }}$ expressions for each of $M$ and $N$. That $M$ and $N$ could satisfy both expressions is equivalent to the Cauchy-Riemann equations:

$$
\begin{equation*}
M=\frac{\partial u}{\partial x}=\frac{\partial v}{\partial y} \text { and } N=-\frac{\partial u}{\partial y}=\frac{\partial v}{\partial x} \tag{2.6}
\end{equation*}
$$

with each expression evaluated at $\lambda\left(t_{0}\right)$.
2.4.2. The notation $f^{\prime}(z)$. To accentuate that the expression $M+i N$ comes from $f$ alone, denote it by $f^{\prime}(z)$ or $\frac{d f}{d z}$. It only, however, exists for functions satisfying the Cauchy-Riemann equations. Here are ways it is like a derivative.
(2.7a) It fits in the chain rule for $\frac{d}{d t}$ of $f(\lambda(t))$ like a derivative.
(2.7b) Directional derivative $D_{\boldsymbol{v}}$ of $f(z)$ in the direction $\boldsymbol{v}$ works as does the gradient for a general function $F: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}: D_{\boldsymbol{v}}(f)\left(z_{0}\right)=f^{\prime}\left(z_{0}\right) v$ is $\frac{d}{d t}\left(f\left(z_{0}+t v\right)\right)(0)$. Check equivalence of this with being analytic!
(2.7c) Analytic composites $\mathbb{C} \xrightarrow{h} \mathbb{C} \xrightarrow{g} \mathbb{C}$ have a simple chain rule [Con78, p. 35]:

$$
\frac{d}{d z}(g \circ h)(z)=\frac{d g}{d w}(w)_{\mid w=h(z)} \frac{d h}{d z}(z)
$$

$$
\begin{equation*}
f^{\prime}(z) d z \text { acts like the differential 1-form } h^{\prime}(x) d x \text { in first year calculus. } \tag{2.7~d}
\end{equation*}
$$

2.5. More explanation of differential forms. First, consider (2.7d) in more detail. The fundamental theorem of calculus says $\int_{a}^{b} h^{\prime}(x) d x=h(b)-h(a)$. A partial analog for integration on $\mathbb{C}$ considers $f^{\prime}(z) d z$, with $f$ analytic. We say $f$ is a primitive (or antiderivative) of $f^{\prime}$. The outcome is the same. Let $z_{a}$ and $z_{b}$ be two points in $D$. Then, let $\lambda:[a, b] \rightarrow D$ be a piecewise differentiable path from $z_{a}$ to $z_{b}$. [Con78, Ch. IV, Th. 1.18]:

$$
\begin{equation*}
\int_{\lambda} f^{\prime}(z) d z=\int_{a}^{b} f^{\prime}(\lambda(t)) \frac{d}{d t}(\lambda(t)) d t=f\left(z_{b}\right)-f\left(z_{a}\right) . \tag{2.8}
\end{equation*}
$$

DEfinition 2.5 (Differential forms). Suppose $m, n: \mathbb{C} \rightarrow \mathbb{C}$ are continuous on $D$, though maybe not analytic. The symbol $m(z) d x+n(z) d y$ is a differential (complex 1-form) on $D$. Closed, locally exact and exact differentials appear later.

A differential 1-form is analytic (or holomorphic) if on each disc in $D$ it has the form $f(z) d z$ with $f(z)$ analytic. We also use meromorphic differentials: $f$ is meromorphic on $D$. [Con78, p. 63] introduces only the differential 1-forms $m(z) d z$, $\left(m(z)\right.$ may not be analytic). It often uses $\int_{\lambda} f$ to substitute for $\int_{\lambda} f d z$. These have the form above: Write $d z$ as $d x+i d y$. They don't, however, include all differential 1-forms $m(z) d x+n(z) d y$.

It is convenient to change variables from $(x, y)$ to $(z, \bar{z})$ to write differentials in the form $u(z) d z+n(z) d \bar{z}$ with $\bar{z}=x-i y$ (and $d \bar{z}=d x-i d y$ ). Chap. 3 Lem. 5.6 formulates the several complex variable version of the next lemma. Call a function anti-holomorphic if about each point it has a power series expression in $\bar{z}$.

Lemma 2.6. The operator $\frac{1}{2}\left(\frac{\partial}{\partial x}-i \frac{\partial}{\partial y}\right)$ maps $z$ to 1 and $\bar{z}$ to 0 . So, it extends the action of $\frac{\partial}{\partial z}$ on holomorphic functions, and it kills anti-holomorphic functions. Similarly, $\frac{1}{2}\left(\frac{\partial}{\partial x}+i \frac{\partial}{\partial y}\right)$ extends the action of $\frac{\partial}{\partial \bar{z}}$ from anti-holomorphic functions to all differentiable functions.

If $f$ is a differentiable function, the expression for the total differential $d f=$ $\frac{\partial f}{\partial x} d x+\frac{\partial f}{\partial y} d y$ is the same as $\frac{\partial f}{\partial z} d z+\frac{\partial f}{\partial \bar{z}} d z$.

Proof. Everything is from the definitions. The sums defining $\frac{\partial}{\partial z}$ and $\frac{\partial}{\partial \bar{z}}$ act on differentiable functions. For the last equality in differentials, check that $\frac{\partial f}{\partial z} d z+\frac{\partial f}{\partial \bar{z}} d z$, written in $x$ and $y$, gives $\frac{\partial f}{\partial x} d x+\frac{\partial f}{\partial y} d y$.

## 3. Branch of $\log (z)$ along a path

Let $D$ be a domain in $\mathbb{C}^{*}$. Denote a path $\gamma:[a, b] \rightarrow D$ by just $\gamma$. A power series $\sum_{n=0}^{\infty} \frac{z^{n}}{n!}$ defines the exponential function $e^{z}$.
3.1. How $e^{z}$ defines branches of $\log (z)$. The exponential has properties so valuable for explicit computation that many parts of mathematics find functions generalizing it. This chapter practices with the exponential function how that works. Here are basic properties of $e^{z}$.
(3.1a) $e^{0}=1$ and $e^{z_{1}+z_{2}}=e^{z_{1}} e^{z_{2}}: e^{z}$ gives a homomorphism $\mathbb{C} \rightarrow \mathbb{C}^{*}$.
(3.1b) $e^{x+i y}=e^{x}(\cos (y)+i \sin (y))$.

In particular, the exact values $w \in \mathbb{C}$ with $e^{w}=1$ are in the set $\{n 2 \pi i \mid n \in \mathbb{Z}\}$. Variants of the following definition appear throughout this chapter.

Definition 3.1. Suppose $h(t)$ is a continuous function defined on $[a, b]$ satisfying $e^{h(t)}=\gamma(t)$. Call $h$ a branch of $\log (z)$ (or, of log) along $\gamma$.

For $z_{0} \in D$, let $\gamma:[a, b] \rightarrow z_{0}$ be the constant path. Suppose $w=w_{0}$ is one solution of $e^{w}=z_{0}$. Then, all solutions are $\left\{w_{0}+n 2 \pi i\right\}$ : possible values of a branch of $\log h(z)$ at $z_{0}$. An easier definition is of a branch of $\log$ on the domain $D$. This is a continuous function $H: D \rightarrow \mathbb{C}$ satisfying $e^{H(z)}=z$ for all $z \in D$ : a right inverse to the exponential function. It is necessary to assure $0 \notin D ; e^{H(0)}=0$ has no solution $H(0)$ because $e^{z}$ never equals 0 .
3.2. Questions about branches of log. The two definitions raise the following questions. Variants apply to the general topic of analytic continuation.
(3.2a) What is the relation between Def. 3.1 and the definition of $H$ ?
(3.2b) When does a branch of $\log$ exist along $\gamma$, and if it exists how many such branches are there?
(3.2c) How does Def. 3.1 give a simple criterion for the existence of $H$ (on $D$ )?
(3.2d) What integrals naturally associate with interpreting existence of $H(z)$ ?
(3.2e) What natural geometric relation between $\mathbb{C}^{*}$ and $\mathbb{C}$ codifies the answers to the previous questions?
Prop. 3.2 answers questions (3.2a), (3.2b) and (3.2c). Then, Prop. 3.5 answers those remaining. These arguments motivate the theory of Riemann surface covers and their moduli. We never use classical language referring to branch cuts (except in a simple example for its historical utility). In the proposition, unless otherwise said, assume $[a, b]$ is the domain of any path.

Proposition 3.2. Suppose $H(z)$ is a branch of $\log$ on D. Fix $z_{0} \in D$. Then, $h^{\dagger}(t)=H(\gamma(t))$ is a branch of log along $\gamma$. Further, suppose $h(t)$ is a branch of $\log$ along $\gamma$. Then, for $t_{0} \in[a, b]$ there is a branch $H$ of $\log$ on a neighborhood of $\gamma\left(t_{0}\right)$ with $H(\gamma(t))=h(t)$ for $t$ close to $t_{0}$.

Even if there is no branch of log on D, the following hold.
(3.3a) There is always a branch $h(t)$ of $\log$ along $\gamma$.
(3.3b) For $h^{*}(t)$ any branch of $\log$ along $\gamma, h(t)-h^{*}(t)$ is constant on $[a, b]$.
(3.3c) $h(t)+2 \pi i m, m \in \mathbb{Z}$, gives the complete set of branches of log along $\gamma$.
(3.3d) There is a branch $H(z)$ of $\log$ on $D$ precisely if for each $\gamma \in \Pi_{1}\left(D, z_{0}\right)$, $h(b)=h(a)$ for $h$ some branch of $\log$ along $\gamma$.
3.3. Proof of Prop. 3.2. If $e^{H(z)} \equiv z$ for $z \in D$, then $e^{H(\gamma(t))} \equiv \gamma(t)$ for $t \in[a, b]$ as in the proposition statement. Thus, $h^{\dagger}(\gamma(t))$ is a branch of log along $\gamma$.

Now suppose $h^{*}(t)$ is any branch of log along $\gamma$. Then,

$$
e^{h(t)} / e^{h^{*}(t)}=e^{h(t)-h^{*}(t)}=\gamma(t) / \gamma(t) \equiv 1
$$

for $t \in[a, b]$. So, the continuous function $F(t)=h(t)-h^{*}(t)$ maps the connected set $[a, b]$ into the topological subspace $2 \pi i \mathbb{Z}$ of $i \mathbb{R}$. The range of a connected set under a continuous function is connected. This shows the range of $F(t)$ is a single point; $F(t)$ is constant on $[a, b]$.

Suppose $z_{0} \in \mathbb{C}$ satisfies $e^{z_{0}}=\gamma(a)$. The rest of the proof has three parts, corresponding to patching pieces of branches of log along $\gamma$.
3.3.1. Extending a branch of log on a subpath. Suppose $\left[a^{\prime}, b^{\prime}\right] \subset[a, b]$. Then, restriction of $\gamma$ to $\left[a^{\prime}, b^{\prime}\right]$ produces a new path, $\gamma_{\left[a^{\prime}, b^{\prime}\right]}$. Let $h_{t_{0}}(t)$ be a branch of log along $\gamma_{\left[a, t_{0}\right]}$ for $t_{0} \in[a, b]$ with $t_{0}<b$.

A classical construction produces a branch $H(z)$ of $\log$ in any sector

$$
S_{\theta_{1}, \theta_{2}}=\left\{r e^{i \theta} \mid \theta_{1}<\theta<\theta_{2}\right\} \text { with } \theta_{2}-\theta_{1} \leq 2 \pi[9.7 a]
$$

Any disk in $\mathbb{C}^{*}$ is in some sector. Restrict $H$ to a disk around $\gamma\left(t_{0}\right)=z_{0}$ and translate it by an integer multiple of $2 \pi i$ to assume $H\left(z_{0}\right)=h\left(t_{0}\right)$. From above, $H(\gamma(t))$ is a branch of log along $\gamma$ restricted to $\left[t_{0}-\epsilon, t_{0}+\epsilon\right]$ for $\epsilon>0$ small. Since $H\left(z_{0}\right)=h\left(t_{0}\right)$, these two branches of $\log$ are equal on $\gamma_{\left[t_{0}-\epsilon, t_{0}\right]}$. If $t_{0}+\epsilon \leq b$, this defines a branch of log along $\gamma_{\left[a, t_{0}+\epsilon\right]}$ :

$$
h_{t_{0}+\epsilon}(t)= \begin{cases}h_{t_{0}}(t) & \text { for } t \in\left[a, t_{0}\right]  \tag{3.4}\\ H(\gamma(t)) & \text { for } t \in\left[t_{0}, t_{0}+\epsilon\right]\end{cases}
$$

We say $h_{t_{0}+\epsilon}$ extends $h_{t_{0}}$.
3.3.2. Sequences of extensions of branch of log. Suppose $t_{0}<t_{1}<\cdots<b$ and $h_{i}(t)$ is a branch of $\log$ along $\gamma_{\left[a, t_{i}\right]}$, with $h_{i}(a)=z_{0}$ for each $i$. Then, from the first part of the proof, $h_{i+1}$ extends $h_{i}$. As the $t_{i}$ s are increasing and bounded, they have a limit point, $t^{*}$. Define $h_{t^{*}}$ by this formula: for $t<t^{*}, h_{t^{*}}(t)=h_{i}(t)$ where $t<t_{i}$; and $h_{t^{*}}\left(t^{*}\right)=\lim _{i} h_{i}\left(t_{i}\right)$. Note: The left side is independent of $i$. The right side has a limit because it is a Cauchy sequence.
3.3.3. Completing existence of branch of log. §3.3.2 shows there is a maximal $t^{\prime}$ having a branch of $\log h_{t^{\prime}}$ along $\gamma_{\left[a, t^{\prime}\right]}$. Then, if $t^{\prime}<b, \S 3.3 .1$ gives an extension to $\gamma_{\left[a, t^{\prime}+\epsilon\right]}$ for some $\epsilon>0$. Thus, $t^{\prime}=b$. That completes proving existence of the extension. Criterion (3.3d) for a branch of $\log$ on a domain is a special case of Lemma 4.12. This depends only on the notion of multiplying paths.

Suppose, as in Prop. 3.2, $h$ is a branch of $\log$ along $\gamma$. For $t \in[a, b]$ there is a neighborhood $D_{t}$ of $\gamma(t)$ and a branch $H_{t}(z)$ of $\log$ on $D_{t}$ satisfying this property.
(3.5) $H\left(\gamma\left(t^{\prime}\right)\right)=h\left(t^{\prime}\right)$ for $t^{\prime}$ close to $t$.

This matches Def. 4.1: There is an analytic continuation of $H_{a}(z)$ along $\gamma$.
Example 3.3 (Branch of $\log$ along a circle). The function $t \mapsto e^{2 \pi i t}=\gamma(t)$, $t \in[0,1]$, parametrizes the counterclockwise unit circle. Let $\epsilon>0$ be small. As in $[9.7 \mathrm{a}], H_{\epsilon}\left(r e^{2 \pi i t}\right)=\ln (|r|)+2 \pi t$ is a branch of $\log$ for all $z$ of form $r e^{2 \pi i t}$, $0 \leq t \leq 1-\epsilon$. So, $h_{\epsilon}(t)=2 \pi i t$ is a branch of $\log$ along $\gamma_{[0,1-\epsilon]}$. Like the proof of Prop. 3.2, $h(t)=2 \pi i t$ extends $h_{\epsilon}$ to be a branch of log along $\gamma$.
3.4. Branch of $\log$ as a primitive. Let $g: D \rightarrow D^{\prime}$ by $w \mapsto g(w)$ be continuous. Assume $g\left(w_{0}\right)=z_{0}$ with $w_{0} \in D$ and $\gamma:[a, b] \rightarrow D^{\prime}$ has $\gamma(a)=z_{0}$.

Definition 3.4. Consider $\gamma^{*}:[a, b] \rightarrow D$ with $\gamma^{*}(a)=w_{0}$. Call it a lift (relative to $g$ ) of $\gamma\left(\right.$ based at $\left.w_{0}\right)$ if $g\left(\gamma^{*}(t)\right)=\gamma(t)$ for all $t \in[a, b]$.
$\S 4.4$ has explicit notation for multiplying paths, as in $\gamma \cdot \gamma^{\prime}$. Let $D$ be a domain in $\mathbb{C}^{*} ; f(z)=1 / z$ is analytic in $D$. Suppose $\gamma \in \Pi_{1}\left(D, z_{0}, z^{\prime}\right)$ and $\Delta_{z^{\prime}}$ is a disc in $D$ about $z^{\prime}$. For $z \in \Delta_{z^{\prime}}$ define $F_{1}(z)$ as $\int_{\gamma \cdot \gamma^{\prime}} \frac{d z}{z}$ where $\gamma^{\prime}$ is any path from $z^{\prime}$ to $z$ in $\Delta_{z^{\prime}}$. The discussion before Def. 5.1 has the precise definition of winding number.

Proposition 3.5. Given $\gamma, F_{1}(z)=F_{1, \gamma}(z)$ depends only on the end point of $\gamma^{\prime}$. Also, $\frac{d F_{1}}{d z}=\frac{1}{z}$ for all $z \in \Delta_{z^{\prime}}$. In particular, $F_{1}(z)$ differs by a constant from a branch of $\log$ along $\gamma \cdot \gamma^{\prime}$. Suppose $\gamma_{1}$ and $\gamma_{2}$ have the same end points. Then, $F_{1, \gamma_{1}}-F_{1, \gamma_{2}}=2 \pi$ im with $m$ the winding number of $\gamma_{1} \cdot \gamma_{2}^{-1}$ about the origin.

Consider $\psi: \mathbb{C} \rightarrow \mathbb{C}^{*}$ by $w \mapsto e^{w}$. Suppose $\gamma:[a, b] \rightarrow \mathbb{C}^{*}$ has beginning point $z_{0}$ with $e^{w_{0}}=z_{0}$. Then, a branch of $\log$ along $\gamma$ (with initial value $w_{0}$ ) is a lift of $\gamma$ (starting at $w_{0}$; relative to $\psi$ ). Let $D^{*}$ be the connected component of $\psi^{-1}(D)$ through $w_{0}$. Then, there is a branch of $\log$ on $D$ with value $w_{0}$ at $z_{0}$ exactly when $\psi$ is one-one to $D$ on $D^{*}$.

The first part requires Cauchy's Theorem ([Ahl79, p. 141, Cor. 1], [Con78, p. 84]). This typifies how integration of analytic functions arises. Abel and Riemann based information on differentials; in Riemann's Existence Theorem they are a substantial subplot.

Proposition 3.6 (Cauchy's Theorem on a disk). Suppose $D$ is a domain in $\mathbb{P}_{z}^{1}$ and $f(z)$ is analytic on $D$. Further, assume $D$ is either analytically isomorphic to $\mathbb{C}$ or to a disk. Then, $\int_{\gamma} f(z) d z=0$ for each closed path in $D$.

Proof of Prop. 3.5. Integration of $f(z)=1 / z$ along paths in $\mathbb{C}^{*}$ analytically continues a primitive for $f$ at the initial point. Thus, to prove $F_{1}(z)$ is independent of $\gamma^{\prime}$ only requires showing the integral is 0 for any closed path $\gamma^{\prime}$ in $\Delta_{z^{\prime}}$. This, follows from Prop. 3.6. The remainder follows by plugging in a lift $\gamma^{*}$ of $\gamma: e^{\gamma^{*}(t)}=\gamma(t)$ for $t \in[a, b]$. By definition $\gamma^{*}$ gives a branch of $\log$ along $\gamma$.

## 4. Analytic continuation along a path

Suppose $f(z)$ is a branch of $\log$ on a domain $D \subset \mathbb{C}^{*}$. Since $e^{z}$ is analytic on $\mathbb{C}$, Def. 3.1 provides analytic continuation of $f(z)$ along any path in $\mathbb{C}^{*}$. It does so using an equation $e^{w}=z$ to force the desired extension. The following generalizes Def. 3.1 (see $\S 6.1$ ). It requires no equation for extending an analytic function.
4.1. Definition of analytic continuation. Suppose $f$ is analytic in a neighborhood $U_{z_{0}} \subset D$ of $z_{0}$ and $\gamma:[a, b] \rightarrow D$ is a path in $D$ based at $z_{0}$.

DEFINITION 4.1 (Analytic continuation of $f$ along $\gamma$ ). Let $f^{*}:[a, b] \rightarrow \mathbb{C}$ be a continuous function with the following properties.
(4.1a) $f^{*}(t)=f(\gamma(t))$ for $t$ close to $a$ (in $[a, b]$ ).
(4.1b) For each $t^{\prime} \in[a, b]$, there is a function $h_{t^{\prime}}(z)$ analytic on a disk $D_{t^{\prime}}$ about $\gamma\left(t^{\prime}\right)$ with $h_{t^{\prime}}(\gamma(t))=f^{*}(t)$ for $t$ near $t^{\prime}$ (in $[a, b]$ ).
If such an $f^{*}$ exists, this definition produces $h_{t^{\prime}}(z)$. This is the analytic continuation of $f$ to $t^{\prime}$. It is an analytic function in some neighborhood of $\gamma\left(t^{\prime}\right)$. Usually, however, the important reference is to the end function $h_{b}(z)$, analytic in a neighborhood of $\gamma(b)$. This we call $f_{\gamma}(z)=f_{\gamma}$, analytic continuation of $f$ (along $\gamma$ ).

Note: $f^{*}(t)$ determines all data for an analytic continuation. It is unique: its difference from another function suiting (4.1) must be constant (hint of [9.8a]). Again, there is a related definition.

Suppose $\hat{f}: D \rightarrow \mathbb{C}$ satisfies $\hat{f}(z)=f(z)$ for all $z \in U_{z_{0}}$. We call $\hat{f}$ an analytic continuation or extension of $f$ to $D$.

REMARK 4.2. Let $\gamma:[a, b] \rightarrow \mathbb{P}_{z}^{1}$ be a nonconstant path. Here is an example of a function analytic at $\gamma(a)$ with no analytic continuation along $\gamma$. Assume $\gamma\left(t^{\prime}\right) \neq \gamma(a)$ for $t^{\prime}$ close to $a$ and let $f$ be a branch of $\log \left(z-\gamma\left(t^{\prime}\right)\right)$ about $\gamma(a)$. Algebraic functions, and others, like branches of log, analytically continue along any path missing some finite set $\boldsymbol{z}$ of points on $\mathbb{P}_{z}^{1}$. Def. 4.5 introduces $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$, analytic functions around $z_{0}$ that are extensible if we avoid $\boldsymbol{z}$.
4.2. Practical analytic continuation. Analytic functions have a power series expression around each point of their domain. This converges in any disc not containing a singularity of the analytic function [Ahl79, p. 179, Thm. 3].
4.2.1. Using disks of convergence. In Def. 4.1, for example, consider $\gamma$ with range a segment of the real axis. Assume also $f^{*}$ is real-valued along $\gamma$ with continuous derivatives of all order. Then, an analytic function restricts to $f^{*}$ along $\gamma$ if and only if $f^{*}$ has a Taylor series around each point. This gives a practical alternative definition of analytic continuation using polygonal paths like $\gamma^{*}$ in the next lemma. Notation is from Def. 4.1.

Lemma 4.3. The following is equivalent to $f$ having an analytic continuation along $\gamma$. There exists a partition $a=t_{0}<t_{0}^{*}<t_{1}<t_{1}^{*}<\cdots<t_{n-1}^{*}<t_{n}=b$ of $[a, b]$, disks $D_{i}$ centered about $\gamma\left(t_{i}\right)$ and $f_{i} \in \mathcal{H}\left(D_{i}\right)$ with these properties.
(4.2a) $D_{i} \cap D_{i+1} \neq \emptyset$ and $f_{i}(z)=f_{i+1}(z)$ for $z \in D_{i} \cap D_{i+1}$.
(4.2b) $\gamma(t) \in D_{i}$ for $t \in\left[t_{i}, t_{i}^{*}\right], \gamma(t) \in D_{i+1}$ for $t \in\left[t_{i}^{*}, t_{i+1}\right], i=0, \ldots, n-1$.
(4.2c) $f_{0}(z)=f(z)$ for $z \in D_{0}$.

Further, let $\gamma^{*}$ be the path following consecutive line segments $\gamma\left(t_{i}\right)$ to $\gamma\left(t_{i}^{*}\right)$, then $\gamma\left(t_{i}^{*}\right)$ to $\gamma\left(t_{i+1}\right), i=0, \ldots, n-1$. Then, $f_{\gamma^{*}}=f_{\gamma}$.

Proof. Suppose we have the pairs $\left(D_{i}, f_{i}\right), i=1, \ldots, n$, and the partition of $[a, b]$. This gives an analytic continuation of $f$ along $\gamma$ by the following formula:

$$
f^{*}(t)= \begin{cases}f_{i}(\gamma(t)) & \text { for } t \in\left[t_{i}, t_{i}^{*}\right] \\ f_{i+1}(\gamma(t)) & \text { for } t \in\left[t_{i}^{*}, t_{i+1}\right]\end{cases}
$$

Then, $f^{*}(t)$ provides an analytic continuation from Def. 4.1.
Follow notation of $\S 3.3 .1$. Inductively consider analytic continuation of $f$ to the end point of $\gamma_{\left[a, t_{i}\right]}$ (and $\gamma_{\left[a, t_{i}^{*}\right]}$ ). Set up the induction by showing this is analytic continuation of $f$ to the end point of $\gamma_{\left[a, t_{i}\right]}^{*}\left(\right.$ and $\left.\gamma_{\left[a, t_{i}^{*}\right]}^{*}\right)$. The essential point is $f_{i}$ exists on a disk containing the range of $\gamma$ on $\left[t_{i}, t_{i}^{*}\right]$. So, $f_{i}$ in a neighborhood of $\gamma\left(t_{i}^{*}\right)$ analytically continues $f_{i}$ (from a neighborhood of $\gamma\left(t_{i}\right)$ ) along any path entirely within $D_{i}$. Then, at the end points of $\gamma$ and $\gamma^{*}, f_{\gamma^{*}}=f_{\gamma}$.

Now assume we have an analytic continuation of $f$ along $\gamma$. Completing the lemma requires creating $\left(D_{i}, f_{i}\right)$ for a corresponding partition of $[a, b]$. Since the range of $\gamma$ is compact, the distance between $\gamma(t)$ and $\gamma\left(t^{\prime}\right)$ is a uniformly continuous function of $\left(t, t^{\prime}\right)$. So, for $d^{\prime}>0$ there exists $d>0$ with $\left|\gamma(t)-\gamma\left(t^{\prime}\right)\right|<d^{\prime}$ if $\left|t-t^{\prime}\right|<d$. Choose $d^{\prime}$ with the following property.
(4.3) For each $t^{\prime} \in[a, b]$, there is a disk of radius no more than $d^{\prime}$ around $\gamma\left(t^{\prime}\right)$ supporting analytic $h_{t^{\prime}}(z)$ as in Def. 4.1.

Compactness of the range of $\gamma$ produces such a $d^{\prime}$. Use $d$ from the above comment. Partition $[a, b]$ so $\left|t_{i}-t_{i}^{*}\right|$ and $\left|t_{i}^{*}-t_{i+1}\right|$ are at most $d$. Then, inductively show this partition has the desired properties.

REmARK 4.4 (Nonsimplicial paths). $\S 4.6$ extends Lemma 4.3 to $D \subset \mathbb{P}_{z}^{1}$. There geodesic paths on $\mathbb{P}_{z}^{1}$ might replace polygonal paths: its pieces are arcs on longitudinal circles. The proof extends with no change.

Lem. 4.3 makes no assumption paths are simplicial. Chap. 3 applies the lemma to general continuous paths. A simplicial assumption allows integrating general differential 1-forms or for computing arc length. Still, suppose $\omega=f(z) d z$ is an analytic 1-form in a neighborhood of $z_{0}$ and $\gamma:[a, b] \rightarrow \mathbb{C}$ is a (continuous, not necessarily simplicial) path with beginning point $z_{0}$.

Let $D$ be any domain containing the range of $\gamma$ in which $f$ extends analytically along each path. Lemma 4.3 produces a simplicial (or polygonal) path $\gamma^{*}$ in $D$ (notice $D$ contains no potential poles of $f$ ) along which integration of $f$ is defined. Let $F(z)$ be an antiderivative of $f(z)$. Analytic continuation of $F(z)$ along $\gamma^{*}$ allows defining $\int_{\gamma} \omega$ to be $F(\gamma(b))-F(\gamma(a))$.
4.2.2. The word monodromy. Monodromy isn't in Webster's dictionary. It is in [Ahl79, p. 295] and [Con78, p. 219] in the statement of the Monodromy Theorem (§8.2 and Chap. 3 Prop. 6.11). The Oxford English Dictionary references exactly the same theorem. It gives it the following meaning:

The characteristic property: If the argument returns by any path to its original value, the function also returns to its original value.
We extend that to include regions where a function may not return to its original value. For this we add group data that accounts for the nonreturn. The loose name for that structure is monodromy action, though we often drop the last word.

The simplest setup for discussing monodromy starts with these elements:
(4.4a) a domain $D$ and $z_{0} \in D$
(4.4b) a closed path $\lambda$ based at $z_{0}$
(4.4c) $f(z)$ analytic in a neighborhood of $z_{0}$
(4.4d) $f$ has an analytic continuation around $\lambda$

Then, analytic continuation around $\lambda$ produces a (possibly) new function, $f_{\lambda}$ analytic in a neighborhood of $z_{0}$.

Definition 4.5 (Extensibility). Assume the setup of (4.4) for every closed path in $D$. Call such an $f$ extensible in $D:(f, D)=\left(f, D, z_{0}\right)$ is extensible. This is a neologism, differing from the notion $f$ has an extension (is extendible) to $D$. Denote the complete set of extensible functions in $D$ (based at $z_{0}$ ) by $\mathcal{E}\left(D, z_{0}\right)$.

By assumption $\mathcal{E}\left(D, z_{0}\right) \subset \mathcal{L}_{z_{0}}$. So, field operations like multiplication and taking ratios make sense. Suppose $f, g \in \mathcal{E}\left(D, z_{0}\right)$. Recall the notation $\mathbb{C}[z, u, v]$ for polynomials in $z, u, v$. Define $\mathbb{C}[z, f, g]=R$ to be $\{\alpha(z, f, g)$ with $\alpha \in \mathbb{C}[z, u, v]\}$.

Lemma 4.6. With the above assumptions, the ring $R$ consists of extensible functions. For any $\lambda \in \Pi_{1}\left(D, z_{0}\right), \alpha_{\lambda}=\alpha\left(z, f_{\lambda}, g_{\lambda}\right)$.

Assume $f \in \mathcal{E}\left(D, z_{0}\right)$ and $D$ is analytically isomorphic to a disk (or to $\mathbb{C}$ ). Then, $f$ is extendible (restriction of an analytic function) on $D$.

Proof. For the first part, show the last result for $f+g$ and $f g$. Every element in $R$ is built from such algebraic operations. Now consider the case $D$ is a disk. Cauchy's Integral formula for an analytic function says a power series for an analytic
function converges up to a singularity on its boundary of convergence. Consider $f \in \mathcal{E}\left(D, z_{0}\right)$ with $z_{0}$ the center of the disk $D$.

Suppose the power series for $f$ converges only on a disk of radius smaller than $D$. Then, analytic continuation of $f$ to some singular boundary point fails. This is contrary to $f \in \mathcal{E}\left(D, z_{0}\right)$.

More generally, let $\beta: D \rightarrow \Delta$ be an analytic isomorphism of $D$ with a disk. Then, $\left(f \circ \beta^{-1}, \Delta, \beta\left(z_{0}\right)\right)$ extends to $F(z)$, and $F(\beta(z))$ extends $f$.

Remark 4.7. Webster's dictionary defines extensible to mean capable of being extended, whether in length or breadth; susceptible of enlargement. That agrees with our definition. Still, it has extendible as a synonym of extensible, whereas we distinguish between the two words.
4.2.3. Meromorphic extensibility. It simplifies many discussions to allow meromorphic functions in $\mathcal{E}\left(D, z_{0}\right)$. Even on $U_{\boldsymbol{z}}$, in considering $f \in \mathcal{E}\left(D, z_{0}\right)$, we eventually remove $z^{\prime}$ from $z$ if it is only a pole of $f$. The simplest way is to allow in $\mathcal{E}\left(D, z_{0}\right)$ functions $f$ having for each path $\gamma$ some $g \in \mathbb{C}(z)$ with $g(z) f(z)$ extensible along $\gamma$ as in Def. 4.5. Technical proofs would use extensibility of $g(z) f(z)$ and analytic continuation to the end point of $\gamma$ would be $(g(z) f(z))_{\gamma} / g(z)$. The result, of course, could have a pole at the end of the path.

In Def. 4.1 there is an auxiliary function $f^{*}:[a, b] \rightarrow \mathbb{C}: f^{*}(t)=f(\gamma(t))$, the values of $f$ along $\gamma$. Extending $f^{*}$ to allow poles requires allowing maps into $\mathbb{P}_{z}^{1}$.

For example: If $g(z)$ is a branch of $\log$ at $z_{0}=1$, we allow $g(z) /(z-1)$ in $\mathcal{E}\left(\mathbb{C}^{*}, 1\right)$. Unless there is a reason to be careful about poles, most discussions will proceed as with extensibility of analytic functions. Integrals and primitives of a function require such care (§4.3). Occasions may need extending this definition to include infinitely many poles.
4.2.4. Conjugates of $f$. Assume $f \in \mathcal{E}\left(D, z_{0}\right)$. Even if $\lambda$ isn't closed, $f_{\lambda}$ has meaning for any path $\lambda$ in $D$ based at $z_{0}$. This produces conjugates of $f$ (in $D$ ) or the monodromy range of $\left(f, D, z_{0}\right)$ :

$$
\begin{equation*}
\mathcal{A}_{f}\left(D, z_{0}\right)=\mathcal{A}_{f}(D)=\left\{f_{\lambda}(z)\right\}_{\lambda \in \Pi_{1}\left(D, z_{0}\right)} \tag{4.5}
\end{equation*}
$$

Regard $f_{\lambda_{1}}, f_{\lambda_{2}} \in \mathcal{A}_{f}(D)$ as equal if are the same function near $z_{0}$. As in [9.8a], $f_{\lambda_{1}}$ and $f_{\lambda_{2}}$ are then equal in any neighborhood of $z_{0}$ where they are meromorphic. Prop. 7.3 implies conjugate here is exactly as in basic Galois Theory. Suppose $h \in K[x]$ an irreducible polynomial over a field $K$ and $h(\alpha)=0$. Then, the full collection of zeros of $h$ are the conjugates of $\alpha$.

Recall the Laurent series field $\mathcal{L}_{z_{0}}$ (about $z_{0}$ ). This consists of ratios of power series convergent around $z_{0}$. The ring $\mathcal{A}_{f}\left(D, z_{0}\right)$ is in $\mathcal{L}_{z_{0}}$. So we may form the composite field $\mathbb{C}\left(\mathcal{A}_{f}\left(D, z_{0}\right)\right)$ these functions generate. Still, not all elements of $\mathbb{C}\left(\mathcal{A}_{f}\left(D, z_{0}\right)\right)$ are in $\mathcal{E}\left(D, z_{0}\right)$ unless $f$ is algebraic.

Lemma 4.8. If $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ is algebraic (as in (1.2)), then $1 / f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$. So, the field $\mathbb{C}(z, f)$ that $z$ and $f$ generate in $\mathcal{L}_{z_{0}}$ is in $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$.

Proof. This requires showing extensions of $f$ have only finitely many zeros. Suppose $f$ satisfies an equation $m(z, f(z))$ with $m \in \mathbb{C}[z, w]$. Then, $\operatorname{deg}_{z}(m)$ bounds the number of solutions of $m(z, 0)=0$. That shows $f(z)$ has only finitely many zeros among its analytic continuations, so $1 / f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$.

Prop. 7.3, showing equivalence of (1.1) and (1.2), lets Lem. 4.8 apply without reservation to algebraic functions.
4.3. A branch of a primitive. Continue notation from $\S 4.1$. Suppose $F(z)$ is a primitive of $f(z)$ in $U_{z_{0}}: \frac{d F}{d z}=f(z)$. This discussion does require care on extensibility of meromorphic functions as in $\S 4.6$. If $f$ is meromorphic in $D$, and $z^{\prime} \in D$, write $f$ as $h_{1}(z)+f_{1}(z)$ with these properties.
(4.6a) $f_{1}$ is analytic in a neighborhood of $z^{\prime}$.
(4.6b) $h_{1}(z)=\frac{1}{z-z^{\prime}} m_{z^{\prime}}\left(\frac{1}{z-z^{\prime}}\right)$ with $m_{z^{\prime}}(z) \in \mathbb{C}[z]$ ( $\equiv 0$ for $f$ analytic at $\left.z^{\prime}\right)$.

Then, the residue of $f$ at $z^{\prime} \in D$ is $m_{z^{\prime}}(0)$.
Definition 4.9. Consider $f \in \mathcal{E}\left(D, z_{0}\right), z^{\prime} \in D$ and a path $\gamma:[a, b] \rightarrow D$ based at $z_{0}$. Denote the restriction $\gamma_{[a, t]}$ to $[a, t]$ by $\gamma_{t}$. We say $f$ has no residue along $\gamma$ if $f_{\gamma_{t}}$ has no residue for each $t \in[a, b]$.

A (branch of) primitive of $f(z)$ along $\lambda:[a, b] \rightarrow D$ is an analytic continuation $\hat{F}_{\lambda}$ of $F(z)$ along $\lambda$. We also label it by $\hat{F}:[a, b] \rightarrow D$.

Lemma 4.10. Assume $f \in \mathcal{E}\left(D, z_{0}\right)$. Then, $f$ has a primitive in a neighborhood of $z_{0}$ when it has no residue at $z_{0}$. Let $\gamma:[a, b] \rightarrow D$ be a path in $D$ along which $f$ has no residue. Then there exists a primitive $\hat{F}:[a, b] \rightarrow \mathbb{C}$ of $f$ along $\gamma$. Further, for $c \in \mathbb{C}$, there is a unique such $\hat{F}$ with $\hat{F}(a)=c$.

Proof. Get a primitive for $f$ in a neighborhood of $z_{0}$ from a primitive for each term in the Laurent series for $f$ around $z_{0}$. The function $z^{k}$ has a primitive $\frac{1}{k+1} z^{k+1}$ if $k \neq-1$. The discussion from $\S 3.4$ has done overkill on showing $z^{-1}$ has no primitive. That is, $f$ must have 0 as residue at $z_{0}$ to have a primitive. Further, by assumption every analytic continuation of $f$ (in $D$ ) has this property.

Let $D_{0}$ be a disk centered at $z_{0}$ and contained in $D$. By assumption $f(z)$ has no residue along any path in $D$. So, it has a primitive $F(z)=F_{0}(z)$ in this disk; integrate the power series for $f(z)$ term by term. The primitive is unique up to addition of a constant.

Now apply the notation of Lemma 4.3. Similarly, there exists $F_{i}(z)$, a primitive of $f_{i}(z)$ in $D_{i}, i=1, \ldots, n$. Since $f_{i}=f_{i+1}$ in $D_{i} \cap D_{i+1}, F_{i}(z)$ and $F_{i+1}$ have equal derivatives on this intersection. Thus, $F_{i}-F_{i+1}$ is a constant on $D_{i} \cap D_{i+1}$. This sets up for an induction. Assume $k$ is an integer for which $F_{0}(z), \ldots, F_{k}(z)$ give an analytic continuation of $F(z)$ along $\gamma_{\left[a, t_{k}\right]}$. Let $F_{k+1}$ be the function we just produced, where $F_{k}-F_{k+1}=b$ for $z \in D_{k} \cap D_{k+1}$. Now replace $F_{k+1}$ by $F_{k+1}+b$. Continue inductively on $k$ to conclude the result.
4.4. Continuation along products of paths. Let $\lambda_{1}:[a, b] \rightarrow D$ be a path where $\lambda_{1}(a)=z_{0}$ and $\lambda_{1}(b)=z_{1}$. Assume $\lambda_{2}:\left[a^{*}, b^{*}\right] \rightarrow D$ is another path and $\lambda_{1}(b)=\lambda_{2}\left(a^{*}\right)$. Create a new path $\lambda_{1} \cdot \lambda_{2} \stackrel{\text { def }}{=} \lambda^{\dagger}:\left[a, b+b^{*}-a^{*}\right] \rightarrow D:$

$$
\lambda^{\dagger}= \begin{cases}\lambda_{1}(t) & \text { for } t \in[a, b]  \tag{4.7}\\ \lambda_{2}\left(t+a^{*}-b\right) & \text { for } t \in\left[b, b+b^{*}-a^{*}\right]\end{cases}
$$

The proof of Lemma 4.12 includes detailed notation for a sequence of analytic continuations. Use that notation for details of the following lemma. Given a path $\lambda$, denote the path $t \mapsto \lambda(b-t+a), t \in[a, b]$, by $\lambda^{-1}$, the inverse of $\lambda$. If $\lambda$ is simplicial so is $\lambda^{-1}$. Continue notation for the function $f$ and let $f_{1}=f_{\lambda_{1}}$ be analytic continuation of $f$ along a path $\lambda_{1}$.

Lemma 4.11. For paths $\lambda_{1}, \lambda_{2}$ and $\lambda_{3}$, assume the end point of $\lambda_{i}$ equals the beginning point of $\lambda_{i+1}, i=1,2$. Analytic continuation of $f_{1}$ along $\lambda_{2}, f_{2}=\left(f_{1}\right)_{\lambda_{2}}$,
is the analytic continuation $f_{\lambda_{1} \cdot \lambda_{2}}$ of $f$ along $\lambda_{1} \cdot \lambda_{2}$. Then, $f_{\left(\lambda_{1} \cdot \lambda_{2}\right) \cdot \lambda_{3}}=f_{\lambda_{1} \cdot\left(\lambda_{2} \cdot \lambda_{3}\right)}$ giving unambiguous meaning to $f_{\lambda_{1} \cdot \lambda_{2} \cdot \lambda_{3}}$. Also, $f_{\lambda \cdot \lambda-1}=f$.

As in §2.3, $\int_{\lambda_{1} \cdot \lambda_{2}} F d z=\int_{\lambda_{1}} F d z+\int_{\lambda_{2}} F d z$, Further, $\int_{\lambda \cdot \lambda^{-1}} F d z=0$.
While $\lambda \cdot \lambda^{-1}$ isn't the constant path (at $\lambda(a)$ ), Lemma 4.11 lists situations where it acts as if it is.

Lemma 4.12. Suppose $f \in \mathcal{E}\left(D, z_{0}\right)$. Let $\lambda^{*}$ be any path with beginning point $z_{0}$ and end point $z_{1}$. Let $f_{1}=f_{\lambda^{*}}$. There is a one-one map between $\mathcal{A}_{f}\left(D, z_{0}\right)$ and $\mathcal{A}_{f_{1}}\left(D, z_{1}\right)$. Also, $f$ is extendible to $D$ if and only if $f_{\lambda}=f$ for each $\lambda \in \Pi_{1}\left(D, z_{0}\right)$.
$\S 4.5$ has the proof of Lemma 4.12. It says there is an analytic function $\hat{f}$ on $D$ restricting to $f$ around $z_{0}$ exactly when $\mathcal{A}_{f}\left(D, z_{0}\right)$ has a single element. Then, monodromy action on $(f, D)$, or (if $D$ is clear, on $f$ ) is trivial.
4.4.1. A permutation representation. For $f \in \mathcal{E}\left(D, z_{0}\right)$ and $\lambda \in \Pi_{1}\left(D, z_{0}\right)$, Lemma 4.11 gives a permutation of $\mathcal{A}_{f}\left(D, z_{0}\right)$ by $h \mapsto h_{\lambda}$ for $h \in \mathcal{A}_{f}\left(D, z_{0}\right)$. Denote $h_{\lambda}$ by $(h) T(\lambda)$ to distinguish $T(\lambda)$ as a permutation of the set $\mathcal{A}_{f}\left(D, z_{0}\right)$. According to Lemma 4.11,

$$
\begin{equation*}
\left((h) T\left(\lambda_{1}\right)\right) T\left(\lambda_{2}\right)=(h) T\left(\lambda_{1}\right) \circ T\left(\lambda_{2}\right)=(h) T\left(\lambda_{1} \cdot \lambda_{2}\right), \tag{4.8}
\end{equation*}
$$

for $\lambda_{1}, \lambda_{2} \in \Pi_{1}\left(D, z_{0}\right)$.
That is, analytic continuation gives a homomorphism from the semi-group (set with multiplication) $\Pi_{1}\left(D, z_{0}\right)$ to permutations on $\mathcal{A}_{f}\left(D, z_{0}\right)$. From Lem. 4.11, the permutation $T(\lambda)$ has $T\left(\lambda^{-1}\right)$ as its inverse permutation. So, the image set of permutations is a group. Call it the monodromy group $G_{f, D}$ of $(f, D)$.

Chap. 3 puts an equivalence relation, homotopy, on $\Pi_{1}\left(D, z_{0}\right)$ to produce the fundamental group $\pi_{1}\left(D, z_{0}\right)$. In particular, from those results $T$ produces a permutation representation of $\pi_{1}\left(D, z_{0}\right)$. This chapter's elementary examples depend only on homology classes of $\Pi_{1}\left(D, z_{0}\right)$ ( $\S 5$ and [9.12]; Chap. $3 \S 6.2$ has the comparison).
4.5. Proof of Lemma 4.12. We show unique analytic continuation to the end points of each closed path implies $f$ extends analytically to $D$. First, we construct the map between $\mathcal{A}_{f}\left(D, z_{0}\right)$ and $\mathcal{A}_{f_{1}}\left(D, z_{1}\right)$ based on $\lambda^{*}$ as in the lemma. Then, $\mathcal{A}_{f}\left(D, z_{0}\right)$ consists of a single element if and only if $\mathcal{A}_{f_{1}}\left(D, z_{1}\right)$ does. Then, we construct $F$, the extension of $f$.
4.5.1. Identifying $\mathcal{A}_{f}\left(D, z_{0}\right)$ and $\mathcal{A}_{f_{1}}\left(D, z_{1}\right)$. Given $h=f_{\lambda} \in \mathcal{A}_{f}\left(D, z_{0}\right)$, apply Lemma 4.11 several times to produce this chain:

$$
\begin{align*}
h_{\lambda^{*}} & =f_{\lambda \cdot \lambda^{*}}=  \tag{4.9}\\
f_{\lambda^{*} \cdot\left(\lambda^{*}\right)^{-1} \cdot \lambda^{*} \cdot \lambda^{*}} & =\left(f_{1}\right)_{\left(\lambda^{*}\right)^{-1} \cdot \lambda \cdot \lambda^{*}},
\end{align*}
$$

since $\left(\lambda^{*}\right)^{-1} \cdot \lambda \cdot \lambda^{*} \in \Pi_{1}\left(D, z_{1}\right)$. This gives a map from $\mathcal{A}_{f}\left(D, z_{0}\right)$ to $\mathcal{A}_{f_{1}}\left(D, z_{1}\right)$ : Conjugating paths based at $z_{0}$ by $\lambda^{*}$.

Map in the other direction by conjugating by $\left(\lambda^{*}\right)^{-1}$. These maps between $\mathcal{A}_{f}\left(D, z_{0}\right)$ and $\mathcal{A}_{f_{1}}\left(D, z_{1}\right)$ are inverse to each other. That is, conjugating $\mathcal{A}_{f}\left(D, z_{0}\right)$ by $\lambda^{*} \cdot\left(\lambda^{*}\right)^{-1}$ acts trivially on $\mathcal{A}_{f}\left(D, z_{0}\right)$ (from in Lemma 4.11). Conclude: Monodromy action on $f$ (in $D$ ) is trivial if and only the same holds for $f_{\lambda^{*}}$.
4.5.2. Extending $f$ to be analytic on $D$. We prove the last statement of the lemma. Suppose $f$ extends to $\hat{f}$ analytic on $D$. Then uniqueness of analytic continuation shows $f_{\lambda}(\lambda(t))=\hat{f}(\lambda(t))$ for each $t$ near $b\left(\lambda \in \Pi_{1}\left(D, z_{0}\right)\right)$.

Now suppose $f_{\lambda}=f$ for each $\lambda \in \Pi_{1}\left(D, z_{0}\right)$. For $z^{\prime} \in D$, assume $z$ is in a disk neighborhood about $z^{\prime}$ entirely contained in $D$. Set $\hat{f}(z)$ equal to $f_{\lambda}(z)$
with $\lambda:[a, b] \rightarrow D$ a path where $\lambda(a)=z_{0}$ and $\lambda(b)=z^{\prime}$. Lem. 4.6 says $f_{\lambda}$ extends to be analytic in the whole disk neighborhood. So this defines $f_{\lambda}(z)$. Let $\lambda^{*}:\left[a^{*}, b^{*}\right] \rightarrow D$ be another such path with the same end points. We have only to show $f_{\lambda^{*}}(z)=f_{\lambda}(z)$.

Then, $\lambda^{\dagger}=\lambda^{-1} \cdot \lambda^{*}$ is a closed path based at $\lambda(b)$. From §4.5.1, analytic continuation of $f_{\lambda}$ around $\lambda^{\dagger}$ equals $f_{\lambda}(z)$. It also equals analytic continuation of $f_{\lambda}$ along $\lambda^{-1}$ followed by analytic continuation of $f$ along $\lambda^{*}$. The result of these analytic continuations is $f_{\lambda^{*}}$. This proves the desired equalities.
4.6. Extending analytic continuation to $\mathbb{P}_{z}^{1}$. Similar definitions work for meromorphic functions in a domain, including analytically continuing meromorphic functions. It simplifies results of Chap. 3 to systematically extend paths into $\mathbb{P}_{z}^{1}$. Recall: A neighborhood basis of open sets around each point gives the topology on a space. Around $\infty$ the neighborhood basis consists of sets of form $N \cup\{\infty\}$ where $N$ is the complement of any closed set in $\mathbb{C}$.

EXAMPle 4.13 (Meromorphic functions). Suppose for some disc $\Delta_{z_{0}}$ about $z_{0}$, $D \cap \Delta_{z_{0}}=\Delta_{z_{0}} \backslash\left\{z_{0}\right\}$. That is, $z_{0}$ is an isolated boundary point of a domain $D$. Further, assume $f$ is analytic on $D$ and it extends to a meromorphic function at $z_{0}$. That means $\lim _{z \mapsto z_{0}}\left(z-z_{0}\right)^{n+1} f(z)=0$ for some $n \in \mathbb{Z}$ [Con78, p. 109]. The minimal such $n$ allows expressing $f(z)$ as $\left(z-z_{0}\right)^{n} h(z)$ with $h$ holomorphic and nonzero in a neighborhood of $z_{0}$. If the minimal $n$ is negative, then $f$ has a pole of order $n$. Define $F: D \cup\left\{z_{0}\right\} \rightarrow \mathbb{P}_{z}^{1}$ by this formula:

$$
F(z)= \begin{cases}f(z) & \text { for } z \in D  \tag{4.10}\\ \infty & \text { for } z=z_{0}\end{cases}
$$

Continuity of $F$ is equivalent to continuity of $z \mapsto 1 / F(z)$ around $z_{0}$. This function is continuous at $z_{0}$ (taking the value 0 ). So it is continuous around $z_{0}$.

DEFINITION 4.14 (Analytic maps to $\mathbb{P}_{z}^{1}$ ). Suppose $f: D \rightarrow \mathbb{P}_{z}^{1}$ is analytic. Assume $z_{0}$ is an isolated boundary point of $D$ and $f$ extends to be meromorphic in a neighborhood of $z_{0}$. Then, we say the extension $F: D \rightarrow \mathbb{P}_{z}^{1}$ is analytic. If $f\left(z_{0}\right)=\infty$, this means $z \mapsto 1 / f(z)$ (with $z_{0} \mapsto 0$ ) is analytic in a neighborhood of $z_{0}$. Also, suppose $\infty$ is an isolated boundary point of $D$ on $\mathbb{P}_{z}^{1}$. Let $D^{\prime}$ be the image of $D$ under $z \mapsto 1 / z$. Then, $f$ extends analytically to $F: D \cup\{\infty\} \rightarrow \mathbb{P}_{z}^{1}$ if $g(z)=f(1 / z)$ extends analytically to $D^{\prime} \cup\{0\}$ in a neighborhood of 0 .

Those functions $f: \mathbb{P}_{z}^{1} \rightarrow \mathbb{P}_{z}^{1}$ analytic everywhere are the rational functions $\mathbb{C}(z)$ in $z[9.3 \mathrm{f}]$. Extending Lem. 4.10 to allow any $D$ in $\mathbb{P}_{z}^{1}$ only requires clarifying what will be the residue at $\infty$. This allows integrations of analytic functions $f: D \rightarrow \mathbb{P}_{z}^{1}$ along paths for any domain $D$ in $\mathbb{P}_{z}^{1}$.

Definition 4.15. By definition a function $f(z)$ meromorphic in a neighborhood of $\infty$ is in $\mathcal{L}_{\infty}$, Laurent series in $1 / z: f(z)=g(1 / z)$ with $g \in \mathcal{L}_{0}$. The residue at $\infty$ is the coefficient of $z$ in $\frac{-g(z)}{z^{2}}$.

For example, $f(z)=1 / z$ has residue -1 at $\infty$. So, it has no primitive at $\infty$.
This chapter's examples explicitly compute conjugates of special functions $f$. Riemann's Existence Theorem turns this around when $D$ is $U_{\boldsymbol{z}}=\mathbb{P}_{z}^{1} \backslash\{\boldsymbol{z}\}$. Running over all algebraic $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$, Chap. 4 describes all possible permutations of the sets $\mathcal{A}_{f}\left(U_{\boldsymbol{z}}, z_{0}\right)$. The goal will be to recognize $f$ by the permutations that come from applying $\Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Then Riemann's Existence Theorem produces (algebraic) $f$ realizing a given labeling. It doesn't, however, give $f$ explicitly; it only exists.

Given such an $f$, suppose $g \in \mathbb{C}(z, f)$ and $\mathbb{C}(z, g)=\mathbb{C}(z, f)$ : $f$ and $g$ are primitive generators of this field (over $z ; \S 1.2 .2$ ). §1.2 gives $u(w), v(w) \in \mathbb{C}(z)[w]$ with $g=u(f)$ and $f=v(g)$. Here is a particular case of Lem. 4.6.

Lemma 4.16. For $\lambda \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$, $g_{\lambda}=u\left(f_{\lambda}\right)$ and $f_{\lambda}=v\left(g_{\lambda}\right)$.

## 5. Winding numbers and homology

Winding numbers appear in $\S 3.4$. Here is the formal definition for the winding number of the closed path $\gamma$ (in $\mathbb{C}$, not passing through $z^{\prime}$ ) about $z^{\prime}$ :

$$
n_{z^{\prime}}(\gamma)=\frac{1}{2 \pi i} \int_{\gamma} \frac{d z}{z-z^{\prime}}
$$

This definition alone would justify complex variables; it defines this winding for any path avoiding $z^{\prime}$.

Definition 5.1. Suppose $D$ is a domain in $\mathbb{C}, z_{0} \in D$ and $\gamma_{1}, \gamma_{2} \in \Pi_{1}\left(D, z_{0}\right)$ have the same winding numbers about each point in $\mathbb{C} \backslash D$. We say they are homologous (in $D$ ). A path is homologous to 0 if all winding numbers for points in $\mathbb{C} \backslash D$ are 0 . It is obvious this forms an equivalence relation on $\Pi_{1}\left(D, z_{0}\right)$. Denote the equivalence classes by $H_{1}(D)$ : the (first) homology group of $D$.
5.1. Extending Def. 5.1. Suppose $\gamma_{1}, \gamma_{2} \in \Pi_{1}\left(D, z_{0}, z_{1}\right)$. Extend the definition of homologous paths: $\gamma_{1}$ and $\gamma_{2}$ are homologous if the closed path $\gamma=\gamma_{1} \cdot \gamma_{2}^{-1}$ is homologous to 0 . Suppose $\gamma$ is a closed path in $\mathbb{C}$. Use the notation $\mathbb{P}_{z}^{1} \backslash \gamma$ for the complement of the range of $\gamma$ in $\mathbb{P}_{z}^{1}$. If $z^{\prime} \in \mathbb{C} \backslash \gamma$, we have a winding number $n_{z^{\prime}}(\gamma)$ of $\gamma$ about $z^{\prime}$. If $\gamma_{1}, \gamma_{2} \in \Pi_{1}\left(D, z_{0}\right)$, then $\gamma_{1} \cdot \gamma_{2}$ is homologous to $\gamma_{2} \cdot \gamma_{1}$. This is because all winding numbers are from computations of integrals in Lem. 4.11. For $\gamma$ a closed path in $\mathbb{P}_{z}^{1}$ denote the complement of the range of $\gamma$ by $\mathbb{P}_{z}^{1} \backslash \gamma$.

Lemma 5.2. In the previous notation, let $U_{1}, \ldots, U_{r^{\prime}}$ be the connected components of $\mathbb{P}_{z}^{1} \backslash \gamma$. One of these, say $U_{r^{\prime}}$ includes $\infty$. Then $n_{z^{\prime}}(\gamma)$ is a constant function of $z^{\prime}$ (with $\gamma$ fixed) as $z^{\prime}$ runs over a connected component of $\mathbb{P}_{z}^{1} \backslash \gamma$. So, if $z^{\prime} \in U_{r^{\prime}} \backslash\{\infty\}$, then $n_{z^{\prime}}(\gamma)=0$.

Let $n_{i}(\gamma)$ be the winding number of $\gamma$ around any point in $U_{i}, i=1, \ldots, r^{\prime}$. Suppose $D \subset \mathbb{C}$ is any domain containing the range of $\gamma$. Any connected component of $\mathbb{C} \backslash D$ is in one of the $U_{i}$ s. Denote the set of integers $i$ with $U_{i}$ containing a component of $\mathbb{C} \backslash D$ by $I_{D}$. Then, the function $i \in I_{D} \mapsto n_{i}(\gamma)$ determines the homology class of $\gamma$ in $D$.

Proof. This follows immediately by noticing $g\left(z^{\prime}\right)=\int_{\gamma} \frac{d z}{z-z^{\prime}}$ is an analytic (and therefore continuous) function on $\mathbb{P}_{z}^{1} \backslash \gamma$. Its values are in $2 \pi i \mathbb{Z}$, a discrete set. So, it is constant on each connected component of $\mathbb{P}_{z}^{1} \backslash \gamma$ (proof of Prop. 3.2).

Now suppose $z^{\prime} \in U_{r} \backslash\{\infty\}$. Then, some big disc $\Delta^{\prime}$ contains all of (the range of) $\gamma$. Let $z^{\prime \prime}$ be any other point in $U_{r} \backslash\{\infty\}$ outside $\Delta^{\prime}$. A previous observation shows $n_{z^{\prime}}(\gamma)=n_{z^{\prime \prime}}(\gamma)$. Further, $g(z)=1 /\left(z-z^{\prime \prime}\right)$ is analytic in $\Delta^{\prime}$. Apply Cauchy's Theorem 3.6 to conclude $n_{z^{\prime \prime}}(\gamma)=0$.

Finally, consider the function $i \in I_{D} \mapsto n_{i}(\gamma)$. This determines the winding numbers of $\gamma$ on each connected component of $\mathbb{C} \backslash D$. This, in turn determines the homology class of $\gamma$.

Denote the image of $\gamma$ in $H_{1}(D)$ by $[\gamma]_{h}$. We understand that a tuple of integers from Lemma 5.2 may be our best interpretation. Further, additivity of winding numbers gives $\left[\gamma_{1} \cdot \gamma_{2}\right]_{h}=\left[\gamma_{1}\right]_{h}+\left[\gamma_{2}\right]_{h}$.
5.2. Homology for domains including $\infty$. Def. 5.1 doesn't include defining homologous paths if a domain in $\mathbb{P}_{z}^{1}$ includes $\infty$. (This includes allowing the paths to go through $\infty$.) Several adjustments allow extending the definition. Chap. 3 has a general approach, one that will not put $\infty$ in a special place. Here we follow implications from a standard complex variables course.
5.2.1. Use linear transformations. If $z^{\prime} \in \mathbb{P}_{z}^{1} \backslash D$ and $\infty \in D$, choose a linear (fractional) transformation $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$ mapping $z^{\prime}$ to $\infty$ [9.14]. Since $\gamma_{1}, \gamma_{2}$ are paths in $D, \alpha \circ \gamma_{1}$ and $\alpha \circ \gamma_{2}$ don't go through $\infty$. Now, apply Def. 5.1 to $\alpha \circ \gamma_{1}$ and $\alpha \circ \gamma_{2}$ relative to $\alpha(D)$. To justify this, check that $\alpha \circ \gamma_{1} \cdot\left(\alpha \circ \gamma_{2}\right)^{-1}$ being homologous to 0 doesn't depend on $\alpha$ [9.14e]. If $D=\mathbb{P}_{z}^{1}$, declare all closed paths to be homologous to 0 .

There is one obvious problem. Suppose $\psi_{D_{1}, D_{2}}: D_{1} \subset D_{2}$ is the inclusion map. Yet, you have already chosen points $z_{i}^{\prime} \in \mathbb{C} \backslash D_{i}$ for reverting homology to a winding number computation, with $z_{1}^{\prime} \neq z_{2}^{\prime}$. Then, we lose having an explicit map $\bar{\psi}_{D_{1}, D_{2}}: H_{1}\left(D_{1}\right) \rightarrow H_{1}\left(D_{2}\right)$ induced from paths in $D_{1}$ also being paths in $D_{2}$.
5.2.2. Excising $\infty$. Assume $\infty \in D, z_{0} \in D \backslash\{\infty\}$ and $\Delta_{\infty}$ is some closed disk about $\infty$ lying entirely in $D$. Regard $\mathbb{P}_{z}^{1}$ as an actual sphere (in $\mathbb{R}^{3}$ ). Assume the radius of $\Delta_{\infty}$ is one unit (see $\S 5.4 .1$ ). Let $\Delta_{\infty, s}$ be the closed disk about $\infty$ of radius $s, 0<s \leq 1$. Let $D_{\infty}=D \backslash\{\infty\}$. Now, $H_{1}\left(D_{\infty}\right)$ has meaning from Def. 5.1.

Let $U_{1}, \ldots, U_{r}$ be the connected components of $\mathbb{C} \backslash D$. Each defines a winding number for $\gamma \in \Pi_{1}\left(D_{\infty}, z_{0}\right)$. Use notation from Lemma 5.2:

$$
\gamma \in \Pi_{1}\left(D_{\infty}, z_{0}\right) \mapsto[\gamma]_{h}=\left(n_{1}(\gamma), \ldots, n_{r}(\gamma)\right) \in \mathbb{Z}^{r}
$$

Define $H_{1}(D)$ by extending $[\gamma]_{h}$ to paths in $\Pi_{1}\left(D_{\infty}, z_{0}\right)$ going through $\infty$. For this, consider the submodule $M_{r}$ of $\mathbb{Z}^{r}$ that $\boldsymbol{v}_{r}=(1,1, \ldots, 1) \in \mathbb{Z}^{r}$ generates.

Suppose $\gamma \in \Pi_{1}\left(D, z_{0}\right)$ goes through $\infty$. Apply Lemma 4.3 to replace $\gamma$ by a geodesic path $\gamma^{*}$ in $D$ (Rem. 4.4) with these properties.
(5.1a) $\gamma$ and $\gamma^{*}$ have the same end points.
(5.1b) If $f \in \mathcal{E}\left(D, z_{0}\right)$, then $f_{\gamma}=f_{\gamma^{*}}$.

If $\gamma^{*}$ doesn't go through $\infty$, precede as below. Otherwise, If $\gamma^{*}$ goes through $\infty$ then it does so only finitely many times. It is the product of a finite number of paths $\gamma^{\prime}$ with the property there is a neighborhood of $\infty, \Delta_{s_{0}} \subset D$, which $\gamma^{\prime}$ returns to and leaves just once. With no loss assume there exists $a<t_{1}<t_{2}<b$ with $\gamma(t) \in \Delta_{s_{0}}$ for $t \in\left[t_{1}, t_{2}\right]$ and $\gamma(t) \notin \Delta_{s_{0}}$ for $t$ outside this interval. Therefore, $\gamma\left(t_{1}\right)$ and $\gamma\left(t_{2}\right)$ are on the boundary $\partial \Delta_{s_{0}}$ of $\Delta_{s_{0}}$. There are two paths on $\partial \Delta_{s_{0}}$ going at constant speed from $\gamma\left(t_{1}\right)$ to $\gamma\left(t_{2}\right)$. Let $\tau$ be one of these. Form a new path, $\gamma^{*}$ from $\gamma$ using this formula:

$$
\gamma^{*}(t)= \begin{cases}\gamma(t) & \text { for } t \in\left[a, t_{1}\right]  \tag{5.2}\\ \tau(t) & \text { for } t \in\left[t_{1}, t_{2}\right] \\ \gamma(t) & \text { for } t \in\left[t_{2}, b\right]\end{cases}
$$

Then, $\left[\gamma^{*}\right]_{h} \in H_{1}\left(D_{\infty}\right)$.
Definition 5.3. In the above, when $\infty \in D$, define $H(D)$ to be $H_{1}\left(D_{\infty}\right) / M_{r}$. Denote the canonical map $H_{1}\left(D_{\infty}\right) \rightarrow H(D)$ by $\psi$. Extend to $[\gamma]_{h}$ : Take $\psi\left(\left[\gamma^{*}\right]_{h}\right)$ to be its image in $H(D)$. Prop. 5.4 completes why this is well defined.
5.3. Computing $H_{1}(D)$ for explicit domains. The word explicit has only subjective meaning. It depends on personally interpreting what it means to know data. Still, consider $U_{\boldsymbol{z}}=\mathbb{P}_{z}^{1} \backslash\{\boldsymbol{z}\}$ for some set of $r$ points $\boldsymbol{z}$. Then, giving $\boldsymbol{z}$ explicitly has comfortable interpretation from experience.

This generalizes to when $\mathbb{P}_{z}^{1} \backslash D$ has $r$ connected components, $C_{1}, \ldots, C_{r}$. Our treatment tacitly assumes $r$ is finite. Then, interpret giving $D$ explicitly as knowing simple closed paths bounding each of the $C_{i} \mathrm{~s}$. Such paths might be circles or polygons with explicit beginning and end points. Given these conditions, computing the homology class of an explicit path in $D$ uses calculations within our experience.

The next proposition specializes a statement in Chap. 4 with homotopy classes replacing homology classes. It gives $\infty$ a special status, that Chap. 4 will not. Simple examples, like [9.10], illustrate having $\infty$ play a special role.

Suppose $D$ is a domain in $\mathbb{P}_{z}^{1}$ whose complement $C(D)$ in $\mathbb{P}_{z}^{1}$ has $r>0$ connected components $C_{1}, \ldots, C_{r}=C(D)_{1}, \ldots, C(D)_{r}$. Denote this ordering of the components as $J_{D}$ with the proviso $C_{r}=C_{\infty}$ is the component containing $\infty$ if $D \subset \mathbb{C}$. If $\infty \in D$, add $C_{\infty}$ by including the empty set $\emptyset$ as the last position. Write $D_{\infty}$ for $D \backslash \infty$. As in $\S 5.2 .1$, consider an inclusion map $\psi_{D_{1}, D_{2}}: D_{1} \subset D_{2}$.

Each connected component of $\mathbb{P}_{z}^{1} \backslash D_{2}$ is in some connected component of $\mathbb{P}_{z}^{1} \backslash D_{1}$. (If $\infty \in D_{2}$ regard $\emptyset$ as $C\left(D_{2}\right)_{\infty}$.) This induces a map $\psi_{D_{1}, D_{2}}^{\dagger}: J_{D_{2}} \rightarrow J_{D_{1}}$. The module $M_{r}$ is from §5.2.2. Recall the definition of a residue of a meromorphic function $f$ at a point $z^{\prime} \in D$ from (4.6).

Proposition 5.4. Suppose $D$ is a domain in $\mathbb{P}_{z}^{1}$ where $C(D)$ has $r$ connected components. Then, $H_{1}(D)$ is isomorphic to $\mathbb{Z}^{r-1}$. If $\infty \in C(D)$, then $\gamma \in \Pi_{1}\left(D, z_{0}\right) \mapsto[\gamma]_{h}$ of Prop. 5.2 and Def. 5.3 gives this isomorphism explicitly. If $\infty \in D$, this identifies $H_{1}(D)$ with $\mathbb{Z}^{r} / M_{r}$ (isomorphic to $H_{1}\left(D_{\infty}\right) / M_{r}$ ), also isomorphic to $\mathbb{Z}^{r-1}$.

Suppose $C\left(D^{\prime}\right)$ has $r^{\prime}$ components and $D \subset D^{\prime}$, with $\infty \in C\left(D^{\prime}\right)$. Then, these isomorphisms induce $\mathbb{Z}^{r-1} \rightarrow \mathbb{Z}^{r^{\prime}-1}$ where $n_{1}, \ldots, n_{r-1} \mapsto m_{1}, \ldots, m_{r^{\prime}-1}$ by

$$
m_{j}=\sum_{i \in J_{D^{\prime}}, \psi_{D, D^{\prime}}^{\dagger}(i)=j} n_{i}
$$

Assume $f$ is meromorphic in $D$ and $\gamma \in \Pi_{1}\left(D, z_{0}\right)$ passes through no residue of $f$. Then, $\int_{\gamma} f(z) d z$ depends only on $[\gamma]_{h}$ and the residues of $f$ at points in $D$.
5.4. Proof of Prop. 5.4. Let $z_{0} \in D$. As above, denote the $r$ connected components of $\mathbb{P}_{z}^{1} \backslash D$ by $C_{1}, \ldots, C_{r}$. First assume $\infty \in C_{r}$. For each $i, 1 \leq i \leq r-1$, there is a closed path $\gamma_{i}=\delta_{i} \cdot \bar{\gamma}_{i} \cdot \delta_{i}^{-1} \in \Pi_{1}\left(D, z_{0}\right)$ with the following description.
(5.3a) $\delta_{i}:[0,1] \rightarrow D$ and $\bar{\gamma}_{i}:[0,1] \rightarrow D$ are paths with $\bar{\gamma}_{i}$ closed.
(5.3b) $\delta_{i}(0)=z_{0}$ and $\delta_{i}(1)=\bar{\gamma}_{i}(0)$.
(5.3c) $\bar{\gamma}_{i}$ has winding number 1 around each point in $C_{i}$.
(5.3d) $\bar{\gamma}_{i}$ has winding number 0 around each point in $C_{j}, j \neq i$.
5.4.1. Construction of $\bar{\gamma}_{i}$. Our construction of $\gamma_{i}$ is similar to that of [Ahl79, p. 140]. Again use the metric topology on $\mathbb{P}_{z}^{1}$ identifying it with a sphere in $\mathbb{R}^{3}$ with coordinates $(r, u, v)$. So, $z_{0} \in \mathbb{P}_{z}^{1}$ corresponds to $\left(r_{0}, u_{0}, v_{0}\right) \in \mathbb{R}^{3}$. Each point of the sphere has a vector pointing outward, perpendicular to the tangent plane to the sphere at $\left(r_{0}, u_{0}, v_{0}\right)$. Further, in any disk on the sphere around $\left(r_{0}, u_{0}, v_{0}\right)$, the boundary of this disk has a well-defined orientation around $\left(r_{0}, u_{0}, v_{0}\right)$. We
take it counterclockwise around the outward normal to the disk at its center. This orientation applies to any simple closed path in the disk [9.17].

Components of $C(D)$ are closed, disjoint (and bounded). Let $d\left(z_{i}, z_{j}\right)$ be the distance (along the minor arc) between $z_{i} \in C_{i}$ and $z_{j} \in C_{j}$. The function $1 / d\left(z_{i}, z_{j}\right)$ has a minimum on $C_{i} \times C_{j}$. Running over all $i$ and $j$ let $\delta$ be at most $1 / \sqrt{2}$ times the smallest of these minimums. Form a grid on $\mathbb{P}_{z}^{1}$ of equally spaced longitudes and latitudes, with spacing at most $\delta$. The closed (spherical) squares (and triangles) of this grid each meet at most one component of $C(D)$.

Let $Q$ be one of the closed grid squares. Its boundary orientation is counter clockwise around any outward normal to an interior point of $Q$ [9.17e]. Define $\bar{Q}_{i}$ to be the union of all $Q \mathrm{~s}$ meeting $C_{i}$. Such a $Q$ meets none of the $C_{j} \mathrm{~s}$ with $j \neq i$. Let $\bar{\gamma}_{i}$ be the topological boundary of $\bar{Q}_{i}$. This is the union of bounding sides - oriented counter clockwise from the paths bounding the $Q$ s - to squares of $\bar{Q}_{i}$. Also, $\bar{Q}_{i}$ includes only sides appearing in exactly one $Q$. Such a side has three (or two, if the grid element is by chance a triangle) other sides of grid squares meeting each vertex. Exactly one side is in $D$ and on another square in $\bar{Q}_{i}$. So, each vertex has an adjoining segment of $\bar{\gamma}_{i} ; \bar{\gamma}_{i}$ is a simple closed (oriented) path.
5.4.2. Winding numbers of $\bar{\gamma}_{i}$. Choose any square $Q^{*}$ in $\bar{Q}_{i}$ and any point $z^{\prime} \in Q^{*} \cap C_{i}$. The winding number of $\bar{\gamma}_{i}$ about $z^{\prime}$ is

$$
n_{i}\left(\bar{\gamma}_{i}\right)=n_{z^{\prime}}\left(\bar{\gamma}_{i}\right)=\sum_{Q \in \bar{Q}_{i}} n_{z^{\prime}}(\partial Q)=n_{z^{\prime}}\left(\partial Q^{*}\right)=1 .
$$

Similarly, $n_{j}\left(\bar{\gamma}_{i}\right)=0$ for $j \neq i$. Winding numbers of the path $\gamma_{i}$ with respect to the $C_{j} \mathrm{~s}$ are the same as for $\bar{\gamma}_{i}$. This is from their definition as an integral (5.3); the integral along $\delta_{i}$ cancels with the integral along $\delta_{i}^{-1}$.

Suppose $\infty \in C(D)$. Let $\gamma$ be any closed path in $D$. To $\gamma$ associate the $r$-tuple $\left(n_{1}(\gamma), \ldots, n_{r}(\gamma)\right) \in \mathbb{Z}^{r}$. Then, the path $\prod_{i=1}^{r} \gamma_{i}^{n_{i}}$ is homologous to $\gamma$. Thus, the winding number map is onto $\mathbb{Z}^{r-1}$. This completes Prop. 5.4 for $\infty \in C(D)$.
5.4.3. The case $\infty \in D$. Consider the map $H_{1}\left(D_{\infty}\right) \rightarrow H_{1}\left(D_{\infty}\right) / M_{r}=H_{1}(D)$. The latter is the definition of $H_{1}(D)$. So we comment only on why the image of $\gamma \in \Pi_{1}\left(D, z_{0}\right)$ depends only on the path $\gamma^{*}$ from (5.2). There were two stages to forming $\gamma^{*}$. The first replaced $\gamma$ by a geodesic path where (5.1) gives its relation to $\gamma$. Suppose $\gamma_{1}$ and $\gamma_{2}$ are two such choices. Then, $f_{\gamma_{1}}=f_{\gamma_{2}}$ for any $f$ extensible to all of $D$. In particular, this applies to $f$ a branch of $\log \left(\frac{z-z_{i}}{z-z_{j}}\right)$ with $z_{i} \in C_{i}$. Its analytic continuations around $\gamma_{1}$ and $\gamma_{2}$ are the same. Therefore, if neither $\gamma_{1}$ nor $\gamma_{2}$ go through $\infty$, the winding numbers of $\gamma_{1} \gamma_{2}^{-1}$ with respect to all components of the complement of $D$ are the same.

Then, we adjusted the geodesic path to a new path $\gamma^{*}$ which for certain did not go through $\infty$. There were, however, two such choices for $\gamma^{*}$. Label these $\gamma_{1}^{*}$ and $\gamma_{2}^{*}$. Let $\delta$ be the parametrized boundary $\partial \Delta_{s_{0}}$ of $\Delta_{s_{0}}$. Then $\delta=\tau_{1} \cdot \tau_{2}$ with $\tau_{1}$ going from $\gamma\left(t_{1}\right)$ to $\gamma\left(t_{2}\right)$ and $\tau_{2}$ going (in the same direction) from $\gamma\left(t_{2}\right)$ to $\gamma\left(t_{1}\right)$. For simplicity assume $\delta$ goes clockwise around $\infty$ (as in §5.4.1).

Then, $\gamma_{1}^{*}=\gamma_{\left[a, t_{1}\right]} \cdot \tau_{1} \cdot \gamma_{\left[t_{2}, b\right]}$ and $\gamma_{2}^{*}=\gamma_{\left[a, t_{1}\right]} \cdot \tau_{2}^{-1} \cdot \gamma_{\left[t_{2}, b\right]}$. Integrals determine homology classes in $H_{1}\left(D_{\infty}\right)$. From Lemma 4.11, $\gamma_{2}^{*}$ and

$$
\gamma_{2}^{\prime}=\gamma_{\left[a, t_{1}\right]} \cdot \tau_{2}^{-1} \cdot \tau_{1}^{-1} \cdot \tau_{1} \cdot \gamma_{\left[t_{2}, b\right]}
$$

have the same homology class. So, $\left[\gamma_{2}^{*}\right]_{h}-\left[\gamma_{1}^{*}\right]_{h}$ is $\left[\tau_{1} \cdot \tau_{2}\right]_{h}$. From Cauchy's Theorem $3.6,\left[\tau_{1} \cdot \tau_{2}\right]_{h}$ is independent of $s_{0}$. On the other hand, $\delta$ bounds the disk complement
of $\Delta_{s_{0}}$ in the counter clockwise direction. By assumption that disk contains all components of $C(D)$. So, $n_{z^{\prime}}(\delta)=1$ as $z^{\prime}$ runs over points in all components of $C(D):\left[\gamma_{2}^{*}\right]_{h}-\left[\gamma_{1}^{*}\right]_{h}=(1, \ldots, 1)$. This shows the images of $\left[\gamma_{1}^{*}\right]_{h}$ and $\left[\gamma_{2}^{*}\right]_{h}$ in $H_{1}(D)$ are the same. That is, $M_{r}$ measures exactly the discrepancy in substituting $\gamma^{*}$ for the original path.
5.4.4. Integrals along homologically trivial paths. Now assume $f$ is meromorphic in $D$. It suffices to show the following. If $\gamma_{1}, \gamma_{2} \in \Pi_{1}\left(D, z_{0}\right)$, and $\gamma=\gamma_{1} \cdot \gamma_{2}^{-1}$ is homologous to 0 , then $\int_{\gamma_{1}} f d z-\int_{\gamma_{2}} f d z=\int_{\gamma} f d z$ depends only on the residues of $f$ in $D$. Let $R_{f}$ be the poles of $f$ for which $f$ has nonzero residues. If $\infty \in C(D)$, and $\gamma \in \Pi_{1}\left(D, z_{0}\right)$ is homologically trivial, then Cauchy's Residue Theorem ([Ahl79, p. 149] or [Con78, p. 112]) says $\int_{\gamma} f d z$ is $\sum_{z^{\prime} \in R_{f}} n_{z^{\prime}}(\gamma) \operatorname{Res}_{z}^{\prime}(f)$. This is the result we want, at least if $\infty \in C(D)$. We won't need to consider the possibility of $f$ having infinitely many nonzero residues.

A reduction of the Residue Theorem to the case $f$ is analytic in $D$ is algebraic. Cauchy's Theorem in this case may be the most important result from first year complex variables. We state it and a generalization for use later.

Definition 5.5. Suppose $u, v: D \rightarrow \mathbb{C}$ are continuous (though maybe not analytic). The differential 1-form $\omega=u(z) d x+v(z) d y$ is locally exact if for each $z_{0} \in D$, there exists $F_{z_{0}}(z)=F(z)$ in a neighborhood of $z_{0}$ with these properties.
(5.4a) $F(z)$ has continuous partial derivatives.
(5.4b) $\frac{\partial F}{\partial x}=u(z)$ and $\frac{\partial F}{\partial y}=v(z)$.

THEOREM 5.6. Suppose $f$ is analytic in $D$, and $\gamma \in \Pi_{1}\left(D, z_{0}\right)$ is homologous to 0 in $D$. Then, $\int_{\gamma} f d z=0$. More generally, this holds with any locally exact differential $\omega$ on $D$ replacing $f d z[$ Ahl79, p. 144, Thm. 16].

Thm. 5.6 holds even if $\infty \in D$ [9.13a]. If we only assume $f \in \mathcal{E}\left(D, z_{0}\right)$, then $\int_{\gamma} f d z, \gamma \in \Pi_{1}\left(D, z_{0}\right)$, usually depends on more than the residues of $f$ and $[\gamma]_{h} \in H_{1}(D)[9.13 \mathrm{~d}]$.

## 6. Branch of solutions of $m(z, w)=0$

This section discusses the implicit function theorem. It is the key ingredient for showing a function satisfying (1.2) satisfies (1.1),
6.1. Branch of inverse of $f(z)$. Suppose $f(z)$ is meromorphic on $D$ and has range $D^{\prime}$. A branch of (right) inverse of $f(z)$ on $D^{\prime}$ is a continuous function $g: D^{\prime} \rightarrow D$ with $f \circ g(z)=z$ for $z \in D^{\prime}$.

DEFINITION 6.1 (Branch of inverse of $f$ along a path). Let $\gamma:[a, b] \rightarrow D$ be a path and $f \in \mathcal{E}\left(D, z_{0}\right)$. Let $g(z)$ be a branch of inverse of $f(z)$ in a neighborhood of $z_{0}$. Then a branch of (right) inverse of $f$ along $\gamma$ is an analytic continuation of $g(z)$ along $\gamma$.

We now change the variable $z$ to $w$, and discuss functions analytic in $w$. This sets notation for the full implicit function theorem. Suppose $f(w)$ is analytic in a neighborhood $\Delta_{w_{0}}$ of $w_{0}$, and $f\left(w_{0}\right)=z_{0}$. For a given fixed $z$, assume $\partial \Delta_{w_{0}}$ passes through no zero or pole of $f(w)-z$ (as a function of $w$ ). Then,

$$
\begin{equation*}
n_{z}=\frac{1}{2 \pi i} \int_{\partial \Delta_{w_{0}}} \frac{f^{\prime}(w) d w}{f(w)-z} \text { and } g(z)=\frac{1}{2 \pi i} \int_{\partial \Delta_{w_{0}}} \frac{w f^{\prime}(w) d w}{f(w)-z} \tag{6.1}
\end{equation*}
$$

count the number $n_{z}$ (resp. the sum $g(z)$ ) of zeros of $f(w)-z$ in $\Delta_{w_{0}}$. By Leibniz's theorem, compute the derivative of $g(z)$ by applying $\frac{\partial}{\partial z}$ under the integral sign (see $\S 7.1)$. So, $g(z)$ is analytic in $z$ for $z$ close to $z_{0}$.

Lemma 6.2. Suppose $f(w)-z_{0}$ has exactly one zero (and no poles) in a neighborhood $\Delta_{w_{0}}$ of $w_{0}$. For $z$ sufficiently close to $z_{0}, f(w)-z$ also has only one zero (and no poles). Thus, the second expression of (6.1) defines a branch $g(z)$ of the inverse of $f(z)$ locally.

The proof of the implicit function theorem in $\S 6.2$ includes the proof of Lemma 6.2.
6.1.1. Branch of $f(z)^{\frac{1}{e}}$ along a path. For $e$ a positive integer, we use the inverse of the $e$ th power map in a general form. This returns to branch of log.

Suppose $f$ is meromorphic in a domain $D$. Let $\gamma:[a, b] \rightarrow D$ be any path whose range misses all zeros and poles of $f(z)$. Then, define a branch of $\log (f(z))$ along $\gamma$ to be a continuous function $h(t)$, for which $e^{h(t)}=f(\gamma(t)), t \in[a, b]$. Existence of a branch of $\log (f(z))$ along such $\gamma$ follows from Prop. 3.2. It is the same as a branch of $\log$ along the path $f \circ \gamma:[a, b] \rightarrow f(D)$.

Define a branch of $f(z)^{\frac{1}{e}}$ along $\gamma$ using $h(t)$ a branch of $\log (f(z))$ along $\gamma$ :

$$
\begin{equation*}
e^{h(t) / e} \stackrel{\text { def }}{=} \operatorname{Br}\left((f(z))^{\frac{1}{e}}\right)(\gamma(t)) \tag{6.2}
\end{equation*}
$$

The left side has a clear meaning. Define the right side to be the value of the branch at $\gamma(t)$. Check: The left of (6.2) to the $e$ th power is $f(\gamma(t))$, as expected. As before, there are $e$ such branches.

Applying Prop. 3.2 gives a unique branch $h(t)$ having a specific value $h(a)$ equal to one of the $e$ th roots of $f(\gamma(a))$.
6.1.2. Local inverses of rational functions. Suppose $f=f_{1} / f_{2} \in \mathbb{C}(w)$ with $\left(f_{1}, f_{2}\right)=1$. Consider the set $X_{f}=\left\{(z, w) \in \mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1} \mid f(w)-z=0\right\}$. Each point $\left(z_{0}, w_{0}\right)$ on $\mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$ has a basis of open sets; each set in the basis is the product of an open set around $z_{0}$ and an open set around $w_{0}$. Intersect those open sets with $X_{f}$ to get neighborhoods of points of $X_{f}$. We discuss for which $\left(z_{0}, w_{0}\right)$ there exists $g(z)$ analytic in a neighborhood of $z_{0}$ satisfying
(6.3) $g\left(z_{0}\right)=w_{0}$ and $f(g(z))=z$.

That is, $g$ produces a local parametrization of a neighborhood of $\left(f\left(w_{0}\right), w_{0}\right)$ by $z \mapsto(z, g(z)):(z, g(z))$ is on $X_{f}$ because $f(g(z))-z \equiv 0$.

There is a global parametrization of $X_{f}$ by $w \mapsto(f(w), w): f(w)-f(w) \equiv 0$. This parametrization, however, isn't as a function of $z$. It is insistent reference to $z$ as the parameter that gives coherent information about the algebraic function $g(z)$.

Lemma 6.2 says points $\left(z_{0}, w_{0}\right)$ with a multiplicity one zero $w_{0}$ of $f(w)-z_{0}$ have neighborhoods projecting one-one to the $z$-line: $(z, g(z)) \mapsto z$. Assume $z_{0} \neq \infty$. Then, $w_{0}$ is a multiplicity one zero of $f_{1}(w)-z_{0} f_{2}(w)$. If this doesn't hold, then $w_{0}$ is a zero of $f_{1}(w)-z_{0} f_{2}(w)$ and its derivative $f_{1}^{\prime}(w)-z_{0} f_{2}^{\prime}(w)$ in $w$. Call it a critical value. Eliminate $z_{0}$.
(6.4) Critical values of $w_{0}$ are zeros of $f_{1}(w) f_{2}^{\prime}(w)-f_{2}(w) f_{1}^{\prime}(w)$.

In particular, there are at $\operatorname{most} \operatorname{deg}\left(f_{1}\right)+\operatorname{deg}\left(f_{2}\right)-1$ critical values of $w_{0}$ (or of $z_{0}$ ). [9.4] precisely defines critical values when $w_{0}$ is a pole of $f$.
6.1.3. Abel's application. Apply the chain rule to $f(g(z)) \equiv z$ :

$$
\begin{equation*}
\frac{d f}{d w}_{\left.\right|_{w=g(z)}} \frac{d g}{d z}=1 \tag{6.5}
\end{equation*}
$$

Therefore, $\frac{d g}{d z}=1 /\left.\frac{d f}{d w}\right|_{w=g(z)}$. This is the complex variable variant of how first year calculus computes an antiderivative of inverse trigonometric functions. Abel applied this to a (right) inverse of a branch of primitive from the following integral

$$
\begin{equation*}
\int_{\gamma} \frac{d z}{\left(z^{3}+c z+d\right)^{\frac{1}{2}}} \tag{6.6}
\end{equation*}
$$

with $c, d \in \mathbb{C}$ (Chap. 4 §7.1). Use (6.2) to interpret $h(z) d z=\frac{d z}{\left(z^{3}+c z+d\right)^{\frac{1}{2}}}$ around some base point $z_{0}: h(z)$ is a branch of $\left(z^{3}+c z+d\right)^{-\frac{1}{2}}$. Let $f(z)$ be a primitive for $h(z) d z$. Apply (6.5) to $f(g(z))=z$ (special case of (7.3)):

$$
\begin{equation*}
\frac{d g(z)}{d z}=\left(g(z)^{3}+c g(z)+d\right)^{\frac{1}{2}} \tag{6.7}
\end{equation*}
$$

Let $\boldsymbol{z}\left\{z_{1}, z_{2}, z_{3}, \infty\right\}$, the three zeros of $z^{3}+c z+d$ and $\infty$. Analytic continuation of $\left(z^{3}+c z+d\right)^{-\frac{1}{2}}$ and its primitive $f(z)=f(z ; c, d)$ produce the collection $\mathcal{A}_{f}\left(U_{\boldsymbol{z}}\right)$. First year calculus computes the inverse of a primitive of $h_{1}(z)=\left(z^{2}+c z+d\right)^{-\frac{1}{2}}$, recognizing it from the trigonometric function $\sin (z)$. This has a unique analytic continuation everywhere in $\mathbb{C}$. Abel discovered the same was true for the inverse $g(z)=g(z ; c, d)$ of $f(z ; c, d)$; it extends everywhere in $\mathbb{C}$. Many conclusions follow.

This example will inspire later topics. For example, dependence of $g(z)=$ $g(z ; c, d)$ on $(c, d)$ usefully distinguishes between algebraic curves defined by $w^{2}-$ $z^{3}+c z+d$ as a function of $(c, d)$ (Chap. $\left.4 \S 7.1\right)$. For each $(c, d), g(z ; c, d)$ is to the exponential function as (6.6) is to a branch of $\log (z)$.
6.2. Implicit function theorem. Consider $m(z, w) \in \mathcal{H}(D)[w]$ (a polynomial in $w$ with coefficients in $\mathcal{H}(D))$. Suppose $g(z)$ is analytic on $D$ and $m(z, g(z)) \equiv$ 0 . We discuss paths $\gamma \rightarrow D$ along which there is an analytic continuation of $g(z)$. Such paths should exclude $z^{\prime}$ having a $w^{\prime}$ with
(6.8) $m\left(z^{\prime}, w^{\prime}\right)=0$ and $\frac{\partial m}{\partial w}\left(z^{\prime}, w^{\prime}\right)=0$.

Riemann's Existence Theorem produces the Riemann surface attached to $g(z)$ (Chap. 4). Data for the Riemann surface include information about all embeddings of $\mathbb{C}(z, g(z))$ in Puiseux fields. This important, though lesser data, is available from the proof that Puiseux fields are algebraically closed (§7.3). Given a polynomial $m(z, w)$ it is theoretically possible, though not always practical, to compute exactly the Puiseux embeddings of $\mathbb{C}(z, g(z))$ from $m$.
6.2.1. Branch and critical points. A branch of solutions to $m(z, w)$ along $\gamma$ is an analytic continuation of $g(z)$ along $\gamma$. Such analytic continuations avoid points $z^{\prime}$ having $w^{\prime}$ satisfying (6.8). Prop. 6.4 references $h_{0} \in \mathbb{C}[z]$ in the expression

$$
\begin{equation*}
m(z, w)=h_{0}(z) w^{n}+h_{1}(z) w^{n-1}+\cdots+h_{n}(z) \tag{6.9}
\end{equation*}
$$

If $z^{\prime}$ is a zero of $h_{0}, m\left(z^{\prime}, w\right)$ has degree lower than $n$ in $w$.
Definition 6.3 (Branch point of $(m, w)$ ). A point $\left(z^{\prime}, w^{\prime}\right)$ is critical for $(m, w)$ if it satisfies (6.8). Call $z^{\prime} \in \mathbb{C}$ a branch point of $(m, w)$ if either there exists $w^{\prime}$ with $\left(z^{\prime}, w^{\prime}\right)$ a critical point or $\operatorname{deg}\left(m\left(z^{\prime}, w\right)\right)<\operatorname{deg}_{w}(m(z, w))=n$.

Suppose $z^{\prime}$ is not a branch point of $(m, w)$. Then, there are exactly $n$ distinct values $w^{\prime}$ with $m\left(z^{\prime}, w^{\prime}\right)=0$. The substitutions $z \mapsto 1 / z$ and/or $w \mapsto 1 / w$ allows extending the definition of critical points of $m(z, w)$ to include $z^{\prime}$ and/or $w^{\prime}$ equal to $\infty$ (see [9.4] and [9.11]). Use the notation of (6.9) and $U_{\boldsymbol{z}}=\mathbb{P}_{z}^{1} \backslash \boldsymbol{z}$.
6.2.2. Algebraic according to (1.2) implies (1.1). Now we see that algebraic by the equation definition implies algebraic by the analytic continuation definition.

Proposition 6.4. Suppose $\boldsymbol{z}$ includes $\infty$ and all branch points of $(m, w)$. Assume $\left(z_{0}, w_{0}\right)$ satisfies the first equation of (6.8), but $z_{0} \notin \boldsymbol{z}$. Then, there is a $g(z)$ analytic near $z_{0}$ with $m(z, g(z)) \equiv 0$ and $g\left(z_{0}\right)=w_{0}$. For $\gamma \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right), g(z)$ analytically continues along $\gamma$ and $m\left(z, g_{\gamma}(z)\right) \equiv 0$ (near the end point of $\gamma$ ).

If $m(z, w) \in \mathbb{C}[z, w]$ is irreducible, then $\boldsymbol{z}$ is a finite set. There are exactly $n$ branches of solutions of $m(z, w)$ along any $\gamma \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ (and exactly $n$ elements of $\mathcal{A}_{g}\left(U_{z}\right)$ ). Conclude: $X_{m}=\left\{(z, w) \in \mathbb{C} \times \mathbb{C} \mid m(z, w)=0, z \in U_{z}\right\}$ is connected and $g$ is algebraic according to (1.1).
The proof takes up $\S 7.1$. Then we get complete equivalence between (1.1) and (1.2).

## 7. Equivalence of the two definitions of algebraic

We show ( $m, w$ ) has only finitely many branch points if $m \in \mathbb{C}[z, w]$.
Lemma 7.1. Assume $m \in \mathcal{H}(D)[w]$ and $\operatorname{deg}_{w}(m)=n>0$. Suppose there is no domain $D^{\prime} \subset D$ in which all $z^{\prime} \in D^{\prime}$ are branch points. Then, the branch points of $(m, w)$ have no accumulation point in $D$. Further, if $m \in \mathbb{C}[z, w]$, either $m$ and $\frac{\partial m}{\partial w}$ have a common factor, or $(m, w)$ has only finitely many branch points.

Proof. Suppose the lemma is false, and $z^{\prime}$ is such an accumulation point. Let $\Delta_{z^{\prime}} \subset D$ be a disk around $z^{\prime}$. So, in this disk there is a sequence of pairs $\left(z_{j}, w_{j}\right)$, $j=1,2, \ldots$ with these properties:
(7.1) $\quad w_{j}$ is a multiple zero of $m\left(z_{j}, w\right)$ and $\lim _{j \mapsto \infty} z_{j}=z^{\prime}$.

Let $\mathcal{R}_{z^{\prime}}$ be the ring of power series in $z$ convergent in a neighborhood of $z^{\prime}$. Then, $\mathcal{R}_{z^{\prime}}$ is a principle ideal domain.

Regard $m$ and $\frac{\partial m}{\partial w}$ as polynomials in $w$ with coefficients in $\mathcal{R}_{z^{\prime}}$. Apply the Euclidean algorithm [9.11]. It produces the greatest common divisor $m_{1}(w)$ of $m$ and $\frac{\partial m}{\partial w}$ in the form $a(z, w) m+b(z, w) \frac{\partial m}{\partial w}=m_{1}(z, w)$, a nonzero polynomial. These polynomials in $w$ have coefficients in $\mathcal{H}\left(D^{\prime}\right)$ with $D^{\prime}$ a neighborhood of $z^{\prime}$.

If $\operatorname{deg}_{w}\left(m_{1}\right) \geq 1$ for each $z^{\prime} \in D^{\prime}$, a zero $w^{\prime}$ of $m_{1}\left(z^{\prime}, w\right)$ gives a common zero of $m\left(z^{\prime}, w\right)$ and $\frac{\partial m}{\partial w}\left(z^{\prime}, w\right)$. This is contrary to our assumption. So $\operatorname{deg}_{w}\left(m_{1}\right)=0$ and the $z_{j}$ s are zeros of $m_{1}$, an analytic function of $z$, accumulating at $z^{\prime}$. So, $m_{1}$ is identically zero contrary to a previous observation.

Apply the Euclidean algorithm to the case $m \in \mathbb{C}[z, w]$. Conclude: If $m$ and $\frac{\partial m}{\partial w}$ have no common factor, then $m_{1}$ is a polynomial in $z$, and all branch points are zeros of it. Thus, there are only finitely many such zeros.
7.1. Proof of Prop. 6.4. Assume $\left(z_{0}, w_{0}\right)$ is not a critical point of $(m, w)$. Let $g(z)$ be $\frac{1}{2 \pi i} \int_{C} w \frac{\partial m}{\partial w}(z, w) d w / m(z, w)$ for each $z$ close to $z_{0}$ with $C$ a counter clockwise circle suitably close to $w_{0}$. We show there are neighborhoods, $U_{z_{0}}$ of $z_{0}$ and $U_{w_{0}}$ of $w_{0}$, with $U_{z_{0}} \times U_{w_{0}}$ free of critical points of $(m, w)$.

To do this, extend Lemma 7.1. Simplify notation by taking $z_{0}=0$ and $w_{0}=0$. Then, $m(0, w) \neq 0$ for $0<|w|<r_{1}$. As $z \mapsto 0, m(z, w) \mapsto m(0, w)$ uniformly with respect to $w$. So, there exists $r<r_{2}<r_{1}$ with $|m(z, w)-m(0, w)|<|m(0, w)|$ for $|z|<r_{2}$ and $|w|<r$. By Rouche's Theorem [Con78, p. 125], $m(z, w)$ and $m(0, w)$ have the same number of zeros in $|w|<r$. So, $m(z, w)$ has a single zero in this region and $g(z)$ gives it.

With $C$ fixed and $z$ close to (but not equal) $z_{0}$, apply $\frac{\partial}{\partial z}$ under the integral giving $g(z)$ to compute its derivative. The partial derivative of $w \frac{\partial m}{\partial w}(z, w) / m(z, w)$ exists and is continuous. Thus, Leibniz's rule [Con78, p. 68] says this gives $\frac{d g}{d z}$, showing it is analytic.

Now consider analytic continuation of $g(z)$ along any path in $U_{\boldsymbol{z}}$. This is the same as the proof of Prop. 3.2 starting at $\S 3.3 .1$. The key ingredient was analytically continuing $g(z)$ beyond the end point of any given path. We have the tools now for that. If $\gamma:[a, b] \rightarrow U_{\boldsymbol{z}}$ is any path, there is a neighborhood of $\gamma(b)$ and $g_{1}(z)$ analytic in this neighborhood with $g_{1}(\gamma(b))$ the value of the extension of $g(z)$ to the end point. As in that proof, since $m\left(\gamma(t), g_{1}(\gamma(t))\right) \equiv 0$ for $t$ close to $b, m\left(z, g_{1}(z)\right) \equiv 0$ for all $z$ with $g_{1}(z)$ defined.

This leaves showing that as $\gamma$ runs over $\Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right), g_{\gamma}$ runs over all $n$ branches $g_{1}, \ldots, g_{n}$ of solutions of $m(z, w)$ around $z_{0}$. Suppose, however, it runs over only the subset $g_{1}, \ldots, g_{t}$ with $t<n$. Consider

$$
M(z, w) \stackrel{\text { def }}{=} \quad \begin{align*}
& \prod_{i=1}^{t}\left(w-g_{i}(z)\right)=  \tag{7.2}\\
& \\
& w^{t}-G_{1}(z) w^{t-1}+G_{2}(z) w^{t-2}+\cdots+(-1)^{t} G_{t}(z)
\end{align*}
$$

Each $G_{i}(z)$ is a symmetric polynomial $S_{i}\left(w_{1}, \ldots, w_{t}\right)$ in $w_{1}, \ldots, w_{t}$ evaluated at $\left(g_{1}, \ldots, g_{t}\right)$. So, $G_{i} \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ (Lem. 4.6).

By assumption, for $\gamma \in \Pi_{1}\left(U_{z}, z_{0}\right), g_{1, \gamma}, \ldots, g_{t, \gamma}$ is a permutation of $g_{1}, \ldots, g_{t}$. Thus, $G_{i, \gamma}=S_{i}\left(g_{1, \gamma}, \ldots, g_{t, \gamma}\right)=S_{i}\left(g_{1}, \ldots, g_{t}\right)$ (Lem. 4.6). So, $\mathcal{A}_{G_{i}}\left(U_{z}\right)$ contains a single element, $i=1, \ldots, t$. Apply Riemann's removable singularity theorem [Ahl79, p. 124] exactly as in the proof of Cor. 7.5. Conclude: Singularities of $G_{i}$ in $\mathbb{P}_{z}^{1}$ are at worst poles. So $G_{i}$ is a rational function in $z: M(z, w) \in \mathbb{C}(z)[w]$.

Plug in $g_{1}(z)=g(z), M(z, g(z)) \equiv 0$. Therefore, $M$ is an irreducible polynomial for $g(z)$ over $\mathbb{C}(z)$ of degree $t<n$. This is contrary to the function field being of degree $n$. This contradiction proves the transitivity statement and concludes the proof of Prop. 6.4. The $n$ elements of $\mathcal{A}_{g}\left(U_{\boldsymbol{z}}\right)$ give the $n$ values $w^{\prime}$ satisfying $m\left(z_{0}, w\right)=0$. So, as $\lambda$ runs over closed paths for which $g_{\lambda}\left(z_{0}\right)=w^{\prime}$, this connects all the points of $X_{m}$ lying over $z_{0}$. Therefore, analytic continuation along the connected set $U_{\boldsymbol{z}}$ connects all the points of $X_{m}$. For future use, here is the lemma hidden in this argument.

Lemma 7.2. Suppose $f(z)$ is analytic in a neighborhood of $z_{0} \notin \boldsymbol{z}$ with $\boldsymbol{z}$ the branch points of $m(z, w) \in \mathbb{C}[z, w]$ and $m(z, f(z)) \equiv 0$. Let $g \in \mathbb{C}(z, f(z))$ and assume $g_{\lambda}=g$ for each $\lambda \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Then, $g \in \mathbb{C}(z)$.
7.2. The converse and integrals along paths. Assume $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$. If $f$ satisfies (1.1) we see it satisfies a nontrivial polynomial equation. Let $f_{1}, \ldots, f_{n}$ be the conjugates of $f$. Apply to $f_{1}, \ldots, f_{n}$ the argument in (7.2) for $g_{1}, \ldots, g_{t}$.

Proposition 7.3. The definitions (1.1) and (1.2) are equivalent.
Assume $m(z, g(z)) \equiv 0$, as in Prop. 6.4. Analytic continuation of $g(z)$ along $\gamma:[a, b] \rightarrow U_{z}$ produces $t \mapsto h(t)$, continuous; $h(t)$ is one of the $n$ distinct values $w^{\prime}$ of $m\left(\gamma(t), w^{\prime}\right)=0$. For $n_{1}, n_{2} \in \mathbb{C}[z, w]$, let $n_{1}(z, w) / n_{2}(z, w)=n(z, w)$. Define the integral of $n(z, g(z))$ along $\gamma$ :

$$
\begin{equation*}
\int_{\gamma} n(z, g(z)) d z \stackrel{\text { def }}{=} \int_{a}^{b} n(\gamma(t), h(t)) d t \tag{7.3}
\end{equation*}
$$

Avoid paths through zeros of $n_{2}$ to assure the integral exists.
7.3. $\mathcal{P}_{z^{\prime}}$ is algebraically closed. Let $\Delta_{z^{\prime}}$ be a closed disk in $\mathbb{P}_{z}^{1}$ centered at $z^{\prime}$. Denote $\Delta_{z^{\prime}} \backslash\left\{z^{\prime}\right\}$ by $\Delta_{z^{\prime}}^{0}$. We show analytic continuations of $f(z) \in \mathcal{E}\left(\Delta_{z^{\prime}}^{0}, z_{0}\right)$ depend only on analytic continuation of $f$ on a circle about $z^{\prime}$. This will show $\mathcal{P}_{z^{\prime}}$ is algebraically closed. Let $\delta$ be the counter clockwise circle about $z^{\prime}$ through $z_{0}$.

Proposition 7.4. If $\lambda \in \Pi_{1}\left(\Delta_{z^{\prime}}^{0}, z_{0}\right)$ has winding number $n_{z^{\prime}}(\lambda)=e(\lambda)$, then $f_{\lambda}=f_{\delta^{e}(\lambda)}$.

Prop. 7.4 gives the complete theory of Riemann surface covers of a punctured disk (in Chap. 3). The proof of Prop 7.4 is in $\S 7.4$.

Corollary 7.5. As in Prop. 7.4, assume $f \in \mathcal{E}\left(\Delta_{z^{\prime}}^{0}, z_{0}\right)$ is algebraic over $\mathcal{L}_{z^{\prime}}$. Let $e=e_{f}$ be the minimal positive integer with $f_{\delta e}(z)=f(z)\left(\right.$ near $\left.z_{0}\right)$. Then, $f \in \mathcal{P}_{z^{\prime}, e}$ and $\mathcal{L}_{z^{\prime}}(f) / \mathcal{L}_{z^{\prime}}$ is isomorphic to $\mathcal{P}_{z^{\prime}, e} / \mathcal{L}_{z^{\prime}}$. In particular, the Puiseux expansion field $\mathcal{P}_{z^{\prime}}$ is algebraically closed. Algebraic functions in $\mathcal{P}_{z^{\prime}, e}$ consist of composites $h(\alpha(z))$ with $h$ algebraic in $\mathcal{L}_{z^{\prime}}$ and $\alpha(z)$ in the set $\left\{\left(z-z^{\prime}\right)^{1 / e}\right\}_{e=1}^{\infty}$.

Proof. If $f(z)$ is algebraic over $\mathcal{P}_{z^{\prime}}$, then it satisfies an equation of degree $n$ with coefficients in $\mathcal{P}_{z^{\prime}}$. There are only a finite number of coefficients. With no loss assume these are in $\mathcal{P}_{z^{\prime}, e^{\prime}}$ for some $e^{\prime} ; f$ is algebraic over $\mathcal{P}_{z^{\prime}, e^{\prime}}$. We want to show $f \in \mathcal{P}_{z^{\prime}, e^{\prime} e}$ for some $e$.

Replace $u_{e^{\prime}}=\left(z-z^{\prime}\right)^{1 / e^{\prime}}$ by $z-z^{\prime}$ everywhere in the equation for $f(z)$ to revert this to where $f$ is algebraic over $\mathcal{P}_{z^{\prime}}$. Or, use this usual algebra observation: If $f$ is algebraic over $\mathcal{P}_{z^{\prime}, e^{\prime}}$, since $\mathcal{P}_{z^{\prime}, e^{\prime}}$ is algebraic over $\mathcal{L}_{z^{\prime}}$, the degree of $f$ is finite over $\mathcal{L}_{z^{\prime}}$, equal to $\left[\mathcal{P}_{z^{\prime}, e^{\prime}}(f): \mathcal{P}_{z^{\prime}, e^{\prime}}\right]\left[\mathcal{P}_{z^{\prime}, e^{\prime}}: \mathcal{P}_{z^{\prime}}\right](\S 1.2)$.

Suppose $f \in \mathcal{E}\left(\Delta_{z^{\prime}}, z_{0}\right)$. Also, $m(f(z)) \equiv 0$ for $z \in \Delta_{z^{\prime}}$ with $m(w) \in \mathcal{L}_{z^{\prime}}[w]$ and $\lambda \in \Pi_{1}\left(\Delta_{z^{\prime}}, z_{0}\right)$. Then, $f_{\lambda}$ is another zero of $m(w)[9.8 c]$. Let $\operatorname{deg}_{w}(m(w))=n$. Then $f_{\lambda^{e}}=f$ for some integer $e \leq n$. Choose $e$ minimal. Then, use $\delta$ as in Prop. 7.4. It shows $e$ is the minimal integer with $f_{\delta e}=f$.

For simplicity, assume $z^{\prime}=0\left(\Delta_{z^{\prime}}=\Delta_{0}\right)$ with $w_{0}$ a solution of $w_{0}^{e}=z_{0}$. Let $\Delta_{1}$ be the preimage of $\Delta_{0}$ by the map $\psi: u \rightarrow u^{e}: \Delta_{1}^{0}$ the preimage of $\Delta_{0}^{0}$. Finally, let $\delta_{1}$ be the counter clockwise circle through $w_{0}$ around 0 in $\Delta_{1}^{0}$. Then,

$$
f \circ \psi_{\delta_{1}}(u)=f_{\delta^{e}}(\psi(u))=f(\psi(u))
$$

Apply Prop. 7.4 to $\left(f \circ \psi, \Delta_{1}^{0}, w_{0}\right)$ to conclude $f \circ \psi_{\gamma}=f \circ \psi$ for $\gamma \in \Pi_{1}\left(\Delta_{1}^{0}, w_{0}\right)$. Lemma 4.12 implies $f \circ \psi$ is analytic in $\Delta_{1}^{0}$. Replace $z$ by $u^{e}$ in the coefficients of $m(w)$. Let $\mathcal{L}_{0, u}$ be convergent Laurent series in $u$ around $u=0$. This gives $m_{1}(w) \in \mathcal{L}_{0, u}[w]$ and $m_{1}(f \circ \psi(u)) \equiv 0$. So, as $u \mapsto 0, f \circ \psi(u)$ goes to one of finitely many values on the Riemann sphere.

Apply Riemann's removable singularity theorem [Ahl79, p. 124]: $f \circ \psi$ extends to an analytic function $\Delta_{1} \rightarrow \mathbb{C} \cup\{\infty\}$. That is, $f \circ \psi$ is analytic in $u$ with $u^{e}=z$. As in $[9.9 \mathrm{~g}]$, this embeds the function field $\mathbb{C}(z, f(z))$ into $\mathcal{P}_{z^{\prime}, e}$. As $f(z)$ has $e$ conjugates over $\mathcal{L}_{z^{\prime}},\left[\mathcal{L}_{z^{\prime}}(f(z)): \mathcal{L}_{z^{\prime}}\right]$ is at least $e$. As $\mathcal{L}_{z^{\prime}}(f(z))$ is a subfield of $\mathcal{P}_{z^{\prime}, e}$, with $\left[\mathcal{P}_{z^{\prime}, e}: \mathcal{L}_{z^{\prime}}\right]=e$, the two fields are equal. This concludes the proof.
7.4. Proof of Prop. 7.4. Let $\lambda \in \Pi_{1}\left(\Delta_{z^{\prime}}^{0}, z_{0}\right)$ have winding number $n_{z^{\prime}}(\lambda)$ around $z^{\prime}$. The proof is in parts for later use. They consist of preliminary notation and description; explicit contraction of $\lambda$ to a path having range the points of $\delta$; and an observation on analytic continuation around such a path. Lemma 4.3 assures $f_{\lambda}=f_{\lambda^{*}}$ with $\lambda^{*}$ a polygonal path. So, with no loss assume $\lambda$ is polygonal.
7.4.1. Notational simplifications. The range of $\lambda$ is compact, and it does not include $z^{\prime}$. So, there is a minimal distance $r_{0}$ between $z^{\prime}$ and the range of $\lambda$. Let $A$ be an annulus around $z^{\prime}$ with inner radius $r^{\prime}<r_{0}$ and outer radius $R^{\prime}$ giving the boundary of $\Delta_{z^{\prime}}$. For simplicity assume $z^{\prime} \neq \infty$ and the disk $\Delta_{z^{\prime}}$ is in the complex plane, rather than on the Riemann sphere. Since circles go to circles by stereographic projection, the only adjustment to use the Riemann sphere would be to compose the description of the sets here with stereographic projection. Also, for simplicity, assume $z_{0}-z^{\prime}=r_{0} e^{2 \pi \theta_{0}}$ has $\theta_{0}=0$.
7.4.2. Description of $A$. The point $z_{v}=z^{\prime}+r_{0} e^{2 \pi i v}$ lies on $\delta$. We also use $z_{v}^{-}=z^{\prime}+r^{\prime} e^{2 \pi i v}$ and $z_{v}^{+}=z^{\prime}+R^{\prime} e^{2 \pi i v}$. The points of the line segment cut by a ray from $z^{\prime}$ to $z_{v}^{+}$meet $A$ in the set

$$
L_{v}=\left\{z_{v}-s\left(z_{v}^{-}-z_{v}\right) \mid s \in[-1,0]\right\} \cup\left\{z_{v}+s\left(z_{v}^{+}-z_{v}\right) \mid s \in[0,1]\right\}
$$

Thus the annulus is the union of the points on $L_{v}, v \in[0,1]$. Reference the point on $L_{v}$ corresponding to $s \in[-1,1]$ by $L_{v}(s)$.
7.4.3. Contraction of $A$ to $\delta$. Define $\Gamma: A \times[0,1] \rightarrow A$ by

$$
\Gamma\left(L_{v}(s), u\right)= \begin{cases}z_{v}-(1-u) s\left(z_{v}^{-}-z_{v}\right) & \text { for } s \in[-1,0] \\ z_{v}+(1-u) s\left(z_{v}^{+}-z_{v}\right) & \text { for } s \in[0,1]\end{cases}
$$

Finally, for each $u \in[0,1]$ we have a path $\gamma_{u}:[a, b] \rightarrow A$ :

$$
t \mapsto \gamma_{u}(t)=\Gamma(\gamma(t), u)
$$

Note: $\gamma_{0}(t)=\gamma(t)$ and $\gamma_{1}(t)$ has range in the points of $\delta$. Further, $\gamma_{1}(t)$, being the contraction of a polygonal path to $\delta$ changes direction but finitely many times. Take $f$ as in the statement of Prop. 7.4. Conclude easily: $f_{\gamma_{1}}=f_{\delta^{e_{1}}}$ with $e_{1}$ the winding number of $\gamma_{1}$ around $z^{\prime}$.
7.4.4. $f_{\gamma_{u}}$ constant in $u \in[0,1]$. For $u \in[0,1]$ consider the continuous function $f_{u}^{*}(t)$ giving analytic continuation (according to Def. 4.1) along $\gamma_{u}$. Let $h_{u, t}$ be the analytic function with restriction to $\gamma_{u}\left(t^{\prime}\right)$ giving $f_{u}^{*}\left(t^{\prime}\right)$ for $t^{\prime}$ close to $t$.

Lemma 4.3 says for $\left(u^{\prime}, t^{\prime}\right)$ close to $(u, t), h_{u, t}$ restricts to $\gamma_{u^{\prime}}\left(t^{\prime}\right)$ to give $f_{u^{\prime}}^{*}\left(t^{\prime}\right)$. Since $f_{u^{\prime}}^{*}\left(t^{\prime}\right)$ is a composition of two continuous functions $\gamma_{u^{\prime}}\left(t^{\prime}\right)$ and $h_{u, t}$, it is continuous. Thus, $f_{u}^{*}(b)$ is a continuous function of $u$. As $f_{u}^{*}(b)$ is in the discrete set of end values of the analytic continuations of $f$ in $\Delta_{z^{\prime}}^{0}$, it is constant in $u$.

Since $z_{0}$ is not a branch point of the algebraic function $f$, the end value $f_{u}^{*}(b)$ determines $f_{\gamma_{u}}$. So, $f_{\gamma_{1}}=f_{\gamma}$, to conclude the proof of the proposition.
7.5. Ramification indices, branch cycles and inertia groups. Consider $L / \mathbb{C}(z)$, a finite extension. Let $z^{\prime} \in \mathbb{P}_{z}^{1}$ and let $\mu: L \rightarrow \mathcal{P}_{z^{\prime}}$ be an embedding of $L$ into Puiseux expansions about $z^{\prime}$. As in [9.9], let $\zeta_{e}=e^{2 \pi i / e}$ for $e \geq 1$ an integer.

Definition 7.6. The ramification index of $\left(L, z^{\prime}, \mu\right)$ is the minimal integer $e=e\left(L, z^{\prime}, \mu\right)$ for which $\mathcal{P}_{z^{\prime}, e}$ contains $\mu(L)$.
7.5.1. A crucial automorphism. Let $\hat{L}$ be the Galois closure of $L / \mathbb{C}(z)$. Cor. 7.5 says there is an integer $\hat{e}$ giving an embedding $\psi: \hat{L} \rightarrow \mathcal{P}_{z^{\prime}, e}$ fixed on $\mathbb{C}(z)$. Here is how $\psi$ produces a conjugacy class in $G(\hat{L} / \mathbb{C}(z))$ depending only on $z^{\prime}$. Let $g_{z^{\prime}}$ be the automorphism of $\mathcal{P}_{z^{\prime}, \hat{e}}$ mapping $\left(z-z^{\prime}\right)^{1 / \hat{e}}$ to $\zeta_{\hat{e}}^{-1}\left(z-z^{\prime}\right)^{1 / \hat{e}}$. This is restriction of a topological generator of the group of the whole algebraic closure.

Denote invertible integers modulo $e$ by $(\mathbb{Z} / e)^{*}$. Consider compatible sequences of integers $m_{e} \in \mathbb{Z} / e^{*}, e \geq 1$ : $m_{e e^{\prime}} \bmod e=m_{e}$ for all integers $e, e^{\prime}$. Denote this
collection $\hat{\mathbb{Z}}^{*}$. Similarly, $\hat{\mathbb{Z}}$ is the compatible collection of $m_{e} \in \mathbb{Z} / e$. Then, $\hat{\mathbb{Z}}$ is a topological ring whose (multiplicative) units are $\mathbb{Z}^{*}[\mathbf{F J 8 6}$, Chap. 1].

Remark 7.7 (Use of the $p$-adics). Here is a reminder of the algebra for writing elements of $\hat{\mathbb{Z}}^{*}$. First: Consider only $e$ that are powers of a particular prime $p$. Then, the compatible sequences $\left\{m_{p^{k}}^{\prime}\right\}_{k=1}^{\infty}$ analogous to $\hat{\mathbb{Z}}$ is $\mathbb{Z}_{p}$, the $p$-adic integers. These satisfy $m_{k}^{\prime} \in \mathbb{Z} / p^{k}$, with $m_{k+1}^{\prime}=m_{k}^{\prime} \bmod p^{k}$ with $k=1, \ldots$. The direct product of the $\mathbb{Z}_{p}$ s over primes $p$ is $\hat{\mathbb{Z}}$. The direct product of the units $\hat{\mathbb{Z}}_{p}^{*}$ of $\hat{\mathbb{Z}}$ is $\hat{\mathbb{Z}}^{*}$. Symbolically write elements of $\hat{\mathbb{Z}}_{p}^{*}$ as series $a_{0}+a_{1} p+a_{2} p^{2}+\cdots$. Here $1 \leq a_{0} \leq p-1$ and $0 \leq a_{i} \leq p-1$ are arbitrary. Without this procedure, excluding 1 and -1 , it might be hard to list any elements of $\hat{\mathbb{Z}}^{*}$.

Lemma 7.8. The automorphism $g_{z^{\prime}}$ maps $\mathcal{P}_{z^{\prime}, e}$ into itself for each e. Its effect on $\mathcal{P}_{z^{\prime}, e e^{\prime}}$ extends its effect on $\mathcal{P}_{z^{\prime}, e}$.

Let $\sigma$ be any automorphism of $\mathcal{P}_{z^{\prime}}$ fixed on $\mathcal{L}_{z^{\prime}}$. The effect of $\sigma$ on $\mathcal{P}_{z^{\prime}, e}$ is the same as $g_{z^{\prime}}^{m_{e}}$ for some $m_{e} \in(\mathbb{Z} / e)^{*}$. So, $\sigma$ corresponds to an element of $\hat{\mathbb{Z}}^{*}$.

Proof. This requires checking the effect of $g_{z^{\prime}}$ on generators of the field extensions. By definition, $g_{z^{\prime}}\left(z-z^{\prime}\right)^{1 / e e^{\prime}}=\zeta_{e e^{\prime}}^{-1}\left(z-z^{\prime}\right)^{1 / e e^{\prime}}$. Put both sides to the power $e^{\prime}$ and then apply $g_{z^{\prime}}$. As $g_{z^{\prime}}$ is a field automorphism,

$$
\left(g_{z^{\prime}}\left(\left(z-z^{\prime}\right)^{1 / e e^{\prime}}\right)\right)^{e^{\prime}}=g_{z^{\prime}}\left(\left(\left(z-z^{\prime}\right)^{1 / e e^{\prime}}\right)^{e^{\prime}}\right)=g_{z^{\prime}}\left(\left(z-z^{\prime}\right)^{1 / e}\right)
$$

Yet, $\left(g_{z^{\prime}}\left(\left(z-z^{\prime}\right)^{1 / e e^{\prime}}\right)\right)^{e^{\prime}}=\left(\zeta_{e e^{\prime}}^{-1}\left(z-z^{\prime}\right)^{1 / e e^{\prime}}\right)^{e^{\prime}}$. As $\zeta_{e e^{\prime}}^{e^{\prime}}=\zeta_{e}$ (by definition), this concludes the first part.

Powers of $g_{z^{\prime}}$ give the group of the degree $e$ extension $\mathcal{P}_{z^{\prime}, e} / \mathcal{L}_{z^{\prime}}[9.9 \mathrm{~d}]$. So, $\sigma$ restricted to $\mathcal{P}_{z^{\prime}, e}$ equals $g_{z^{\prime}}^{m_{e}}$ for some $m_{e} \in(\mathbb{Z} / e)^{*}$. Let $\sigma_{e}$ be restriction of $\sigma$ to $\mathcal{P}_{z^{\prime}, e}$. Compatibility of these $m_{e} \mathrm{~s}$ is from $\sigma_{e}$ being restriction of $\sigma_{e e^{\prime}}$ to $\mathcal{P}_{z^{\prime}, e}$.
7.5.2. Embeddings and branch cycles. Continue the discussion starting §7.5.1. Restrict $g_{z^{\prime}}$ to $\hat{L}$. Since $\hat{L} / \mathbb{C}(z)$ is Galois and $g_{z^{\prime}}$ fixes $\mathbb{C}(z)$, this gives an automorphism $g_{z^{\prime}, \psi}$ of $\hat{L}$. Denote this element of $G(\hat{L} / \mathbb{C}(z))=G$ by $g_{z^{\prime}, \psi}$. It depends on $\psi$, the choice of the embedding. Call it the branch cycle attached to the pair $\left(z^{\prime}, \psi\right)$.

Lemma 7.9. For $z^{\prime} \in \mathbb{P}_{z}^{1},[\hat{L}: \mathbb{C}(z)]$ distinct embeddings $\psi: \hat{L} \rightarrow \mathcal{P}_{z^{\prime}, \hat{e}}$ leave $\mathbb{C}(z)$ fixed. As $\psi$ runs over such embeddings, $g_{z^{\prime}, \psi}$ runs over a conjugacy class in $G$. Suppose $f(z)$, meromorphic about a nonbranch point $z_{0}$, satisfies $m(z, f(z)) \equiv 0$, $m \in \mathbb{C}[z, w]$. So, $z^{\prime} \in \mathbb{P}_{z}^{1}$ produces a conjugacy class $\mathrm{C}_{z^{\prime}}$ of $G=G(\hat{L} / \mathbb{C}(z))$. With $\boldsymbol{z}$ the branch points of $(m, w)$, for each $z^{\prime} \notin \boldsymbol{z}, \mathrm{C}_{z^{\prime}}=\{1\}$.

Let $\delta$ be a clockwise (closed) circle around $z^{\prime} \in \boldsymbol{z}$ bounding a closed disk $\Delta_{z^{\prime}}$. Assume $\Delta_{z^{\prime}}\left(\right.$ excluding possibly $\left.z^{\prime}\right)$ contains no other branch point of $(m, z)$ and $z_{0} \in \Delta_{z^{\prime}}$. Let $f_{1}, \ldots, f_{n}$ be a complete list of conjugates of $f$. Denote analytic continuation of $f_{j}$ around $\delta$ by $f_{j, \delta}$. Then, for some choice of $\psi, g_{z^{\prime}, \psi}$ maps this set to $f_{1, \delta}, \ldots, f_{n, \delta}$.

Proof. Cor. 7.5 produces one embedding, $\psi: \hat{L} \rightarrow \mathcal{P}_{z^{\prime}, \hat{e}}$. Let $\alpha$ run over the automorphisms of $\hat{L}$ fixed on $\mathbb{C}(z)$. Then, $\psi \circ \alpha: \hat{L} \rightarrow \mathcal{P}_{z^{\prime}, \hat{e}}$ runs over $[\hat{L}: \mathbb{C}(z)]$ embeddings of $\hat{L}$ into the algebraic closure of $\mathcal{L}_{z^{\prime}}$ fixed on $\mathbb{C}(z)$. Galois theory says this is the exact number of embeddings possible. So we have listed them all.

Consider the effect on $g_{z^{\prime}, \psi}$ of composing $\psi$ with $\alpha$. The new automorphism is

$$
g_{z^{\prime}, \psi \circ \alpha}=(\psi \circ \alpha)^{-1} \circ g_{z^{\prime}} \circ(\psi \circ \alpha)=\alpha^{-1} g_{z^{\prime}, \psi} \alpha
$$

That is, $g_{z^{\prime}, \psi \circ \alpha}$ runs over the conjugacy class of $g_{z^{\prime}, \psi}$ in $G$ as $\alpha$ runs over $G$.
Regard elements $f_{1}, \ldots, f_{n}$ as in $\mathcal{L}_{z_{0}}$. Let $h(z)$ be a branch of $\left(z-z^{\prime}\right)^{1 / \hat{e}}$ defined in this neighborhood of $z_{0}$. Giving an embedding of $\hat{L}$ (fixed on $\mathbb{C}(z)$ ) into $\mathcal{P}_{z^{\prime}, \hat{e}}$ is equivalent to giving an embedding of $\hat{L}$ mapping $f_{1}, \ldots, f_{n}$ into power series $g_{1}(h(z)), \ldots, g_{n}(h(z))$ in $h(z), g_{1}, \ldots, g_{n} \in \mathcal{L}_{0}$. Analytic continuation of $g_{1}(h(z)), \ldots, g_{n}(h(z))$ around $\delta$ maps $g_{i}(h(z))$ to $g_{i}\left(\zeta_{\hat{e}}^{-1} h(z)\right)$. This is the effect of restriction of $g_{z^{\prime}}$ on the embedding of the $f_{i} \mathrm{~s}$ in the Puiseux expansions.
7.5.3. Branch cycles and inertia groups. Choosing $\zeta_{\hat{e}}^{-1}$ (rather than $\zeta_{\hat{e}}$ ) in the definition of $g_{z^{\prime}}$ is convenient (later). This assures $\delta$ in Lem. 7.9 is a clockwise path. The conjugacy class $\mathrm{C}_{z^{\prime}}$ in Lem. 7.9 is crucial to precise formulations of Riemann's Existence Theorem. This is the branch cycle conjugacy class attached to $z^{\prime}$. Using $G \leq S_{n}$, disjoint cycle data (Chap. $3 \S 7.1$ ) for elements of $\mathrm{C}_{z^{\prime}}$ is sufficient for some applications, though not for the more serious.

Definition 7.10 (Inertia groups). The branch cycle $g_{z^{\prime}, \psi}$ in Lem. 7.9 generates a group, $I_{z^{\prime}, \psi}$ of $G(\hat{L} / \mathbb{C}(z))$. This is the inertia group attached to the embedding $\psi$. The notation $I_{z^{\prime}}$ refers to any choice of the groups conjugate to $I_{z^{\prime}, \psi}$. Points $z^{\prime} \in \mathbb{P}_{z}^{1}$ for which $I_{z^{\prime}}$ is nontrivial are the branch points of $L / \mathbb{C}(z)$.
7.5.4. Two definitions of branch points. There are now two definitions of branch points. Def. 7.10 gives it for the function field $L / \mathbb{C}(z)$ and $\S 6.2$ for the pair $(m, z)$. They are related though they may not be equal [9.11].

Proposition 7.11. Suppose $m(z, f(z)) \equiv 0$ and $L=\mathbb{C}(z, f(z))$. If $z^{\prime} \in \mathbb{C}$ is a branch point of $L / \mathbb{C}(z)$, then it is also a branch point of $(m, z)$.

Proof. Suppose $z^{\prime}$ is a branch point of $L / \mathbb{C}(z)$. Then, there is an embedding $\psi: \mathbb{C}(z, f(z)) \rightarrow \mathcal{P}_{z^{\prime}, e}$ where the image of $f$ is not in $\mathcal{L}_{z^{\prime}}$. In particular, the power series $\psi(f)$ and $g_{z^{\prime}}(\psi(f))$ in $\left(z-z^{\prime}\right)^{1 / e}$ have the same value after substituting 0 for $\left(z-z^{\prime}\right)^{1 / e}$. Since $\left(w-\psi(f(z))\left(w-g_{z^{\prime}}(\psi(f))\right)\right.$ divides $m(z, w)$ (in $\mathcal{P}_{z^{\prime}}[w]$ ), this shows $m\left(z^{\prime}, w\right)$ has multiple zeros.

## 8. Abelian functions from branch of $\log$

A branch of log isn't an algebraic function. Still, it allows explicit construction of all the algebraic functions we call abelian, the topic of this subsection.
8.1. Further notation around extensible functions. Let $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ be the extensible (meromorphic) functions on $U_{\boldsymbol{z}}$ (as in Def. 4.5; given by elements of $\mathcal{L}_{z_{0}}$ ). Denote algebraic elements of $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ (as in Def. 1.1) by $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$.

Definition 8.1. Let $G$ be a finite group having a specific property $P^{*}$. Say an element $f \in \mathcal{E}\left(U_{z}, z_{0}\right)^{\text {alg }}$ has property $P^{*}$ if its monodromy group $G_{f}$ (§4.4.1) has this property. This allows referring to abelian, nilpotent ( $G_{f}$ is a product of its $p$-Sylow subgroups), solvable or primitive functions.

Example: Suppose $[\mathbb{C}(z, f): \mathbb{C}(z)]=n$. Then, $f$ is primitive if $G_{f}$ is a primitive subgroup of $S_{n}$ (Chap. 3 Def. 7.9). Equivalently, by the Galois correspondence, there is no field properly between $\mathbb{C}(z)$ and $\mathbb{C}(z, f)[9.5]$. Later chapters show this is a very important concept. Unfortunately, the word primitive appears in many guises in mathematics (already in this chapter). It has even more meanings in the Webster's dictionary. The closest to our meaning here is this: not derived; as a primitive verb in grammar. So, $\mathbb{C}(z, f)$ is an extension not (even partially) derived
from any other proper extension of $\mathbb{C}(z)$. Note that this is different in English than it being generated by a single element over $\mathbb{C}(z)$ (primitive generator). Denote the abelian (resp. nilpotent) functions in $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$ by $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {ab }}$ (resp. $\left.\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {nil }}\right)$.
8.2. Abelian monodromy. For $e \in \mathbb{Z}$ and $\gamma:[a, b] \rightarrow D$ a path whose range misses all zeros and poles of $f(z),(6.2)$ defines branch of $f(z)^{\frac{1}{e}}$ along $\gamma$.

Here is data for abelian functions of index $e$ :

- distinct points $\boldsymbol{z}=z_{1}, \ldots, z_{r}$ in $\mathbb{P}_{z}^{1}$ : branch points
- $\Delta_{z_{0}}$, a disk neighborhood of $z_{0}$ : base point
- an integer $e$ : index
- a branch $g_{i, j}$ of $\log \left(\frac{z-z_{i}}{z-z_{j}}\right)$ in $\Delta_{z_{0}}, 1 \leq i<j \leq r$

Denote the field $\mathbb{C}\left(z, e^{g_{i, j} / e}, 1 \leq i<j \leq r\right)$ by $L_{e, \boldsymbol{z}}$ : The field of abelian functions (on $\mathbb{P}_{z}^{1}$ ) ramified over $\boldsymbol{z}$ of index dividing $e$. It is a subfield of $\mathcal{L}_{z_{0}}$. Any $f \in L_{e, \boldsymbol{z}}$ defines an analytic $f: \Delta_{z_{0}} \rightarrow \mathbb{P}_{z}^{1}$ according to notation of $\S 4.6$. If some $z_{i}=\infty$ replace $z-z_{i}$ by 1 in the definition. In particular, when $z_{r}=\infty, g_{i, r}$ is a branch of $\log \left(z-z_{i}\right), i=1, \ldots, r-1$. This definition includes all algebraic functions having abelian monodromy group. It will give a valuable comparison in Chap. 4. There is a similar definition of algebraic functions on $D$ with any domain $D$ replacing $\mathbb{P}_{z}^{1}$.
8.2.1. Galois group of $L_{e, \boldsymbol{z}}$. A complete description of $L_{e, \boldsymbol{z}}$ depends only on homology classes of paths in $\Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$.

Corollary 8.2. Assume $\gamma_{1}, \gamma_{2} \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ are homologous and $f$ is an algebraic abelian function on $U_{z}$ corresponding to the data (8.2). Then, the analytic continuations $f_{\gamma_{1}}$ and $f_{\gamma_{2}}$ (back to $z_{0}$ ) are equal. Monodromy from $\Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ induces a faithful action of $H_{1}\left(U_{\boldsymbol{z}}\right) / e H_{1}\left(U_{\boldsymbol{z}}\right)$ on $L_{e, \boldsymbol{z}}$ and therefore on $\mathbb{C}\left(z, \mathcal{A}_{f}\left(U_{\boldsymbol{z}}, z_{0}\right)\right)$ (§4.2.2). In particular, $L_{e, \boldsymbol{z}} / \mathbb{C}(z)$ is Galois with group $H_{1}\left(U_{\boldsymbol{z}}\right) / e H_{1}\left(U_{\boldsymbol{z}}\right)$. For $f \in L_{e, z} / \mathbb{C}(z), \mathbb{C}\left(z, \mathcal{A}_{f}\left(U_{z}, z_{0}\right)\right) / \mathbb{C}(z)$ is Galois with group a quotient of this group.

Proof. For simplicity assume $z_{r}=\infty$. Take $\gamma \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ and

$$
f(z)=m_{1}\left(e^{g_{1, \gamma}(z) / e}, \ldots, e^{g_{r-1, \gamma}(z) / e}\right) / m_{2}\left(e^{g_{1, \gamma}(z) / e}, \ldots, e^{g_{r-1, \gamma}(z) / e}\right)
$$

where $g_{j, \gamma}$ denotes analytic continuation of $g_{j}$ around $\gamma$. Let $m_{j}$ be the winding number of $\gamma$ about $z_{j}$. Analytic continuation of $g_{j}$ around $\gamma$ adds $2 \pi i m_{j}$ to $g_{j}$ (Prop. 3.5). Since $\gamma_{1}$ and $\gamma_{2}$ have the same winding numbers around each $z_{j}$, this proves the effect of their analytic continuations on $f$ are the same.

Note that $L_{e, z} / \mathbb{C}(z)$ is the composite of the field extensions $\mathbb{C}\left(z, e^{g_{j}(z) / e}\right) / \mathbb{C}(z)$. Apply [9.9] using $\left(e^{g_{j}(z) / e}\right)^{e}=z-z_{j}$. Conclude: $\mathbb{C}\left(z, e^{g_{j}(z) / e}\right) \mathbb{C}(z)$ is Galois with group $\mathbb{Z} /(e)$. From $[9.5 \mathrm{~d}]$, the composite of these fields is Galois, with group a subgroup of $\mathbb{Z} /(e) \times \cdots \times \mathbb{Z} /(e)$. The image of $H_{1}\left(U_{\boldsymbol{z}}\right) / e H_{1}\left(U_{\boldsymbol{z}}\right)$ produces field automorphisms of $L_{e, \boldsymbol{z}}$. We know these explicitly. Let a closed path $\lambda$ have respective winding numbers $\left(a_{1}, \ldots, a_{r-1}\right)$ around $\left(z_{1}, \ldots, z_{r-1}\right)$. If $e$ does not divide $a_{j}$, then monodromy action of $\lambda$ on $g_{j}$ is nontrivial. So the automorphism group is all of $(\mathbb{Z} / e)^{r-1}$. This shows the result.
8.3. Deeper into the Monodromy Theorem. Consider $m \in \mathbb{C}[z, w]$ and $D$ a domain in $\mathbb{P}_{z}^{1}$. It is a fundamental to decide when some branch of solutions of $m(z, w)=0$ is a meromorphic function on all of $D$. Riemann's Existence Theorem gives a satisfactory answer to versions of this question.
8.3.1. Simple connectedness. Call a domain in $\mathbb{C}$ simply connected if there is at most one connected component in $\mathbb{P}_{z}^{1} \backslash D$. Chap. 3 has the usual definition of a simply connected topological space. For open subsets of $\mathbb{P}_{z}^{1}$ these definitions describe the same sets. The following is an application of Cauchy's Residue Theorem for later comparison with the general Monodromy Theorem.

Theorem 8.3 (Monodromy Theorem). Suppose $D \subset \mathbb{C} \backslash\left\{z_{1}, \ldots, z_{r}\right\}$ is simply connected. Assume $f$ has no residues in $D$. Then $f(z)$ has a primitive (antiderivative; §2.5) $F(z)$ on $D$. Suppose $\boldsymbol{z}$ contains the zeros and poles of $f(z)$. Apply this to $\frac{d f}{d z} / f$ to conclude there is a branch of $\left.\log (f(z))\right)$ on $D$.
8.3.2. Homological triviality versus simple connectedness. Being simply connected has another characterization: the winding number of any closed path in $D$ relative to any point $z^{\prime}$ outside of $D$ is 0 . That is, $D$ is simply connected if all paths in $D$ are homologous to 0 . Beware! If $D$ is not simply connected, some paths may be homologous to 0 , though not trivial for our applications. For example, any function that isn't abelian has a nontrivial analytic continuation around some path homologous to 0 . For, however, abelian functions, most questions use just the Monodromy Theorem in Prop. 7.4. For example, suppose $m(z, g(z)) \equiv 0$, and $\mathbb{C}(z, g(z)) / \mathbb{C}$ is an abelian extension ( $g$ is abelian). Then, we can characterize those $D$ that aren't simply connected on which $g$ is extendible. It is tougher to be so precise about antiderivatives for even abelian functions $g$ along paths in $D$.
8.4. Primitive tangential base points. Let $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$ and $z^{\prime} \in \boldsymbol{z}$. Suppose $\lambda$ in $U_{\boldsymbol{z}}$ goes from $z_{0}$ to $z_{1}$. Analytic continuation of $f$ produces $f_{\lambda} \in$ $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{1}\right)$. Consider $\lambda$ a restriction map. Applying $\lambda$ restricts $f$ to $f_{\lambda} \in \mathcal{L}_{z_{1}}$.

How about using a path to restrict $f$ to a function around $z^{\prime}$ ? That is, let $\lambda$ be a path with end point close to $z^{\prime}$. Can we consider $f_{\lambda}$ restriction of $g \in \mathcal{P}_{z^{\prime}}$ ? The simple answer is No!, unless $f_{\lambda}$ extends to an analytic function around $z^{\prime}$. It is valuable, however, to add data to $\mathcal{P}_{z^{\prime}}$ so the answer will be Yes!

Choose an open disk $D^{\prime}$ in $U_{\boldsymbol{z}}$, with $z^{\prime}$ on its boundary. Let $g_{e}(z)$ be a branch of $\left(z-z^{\prime}\right)^{1 / e}$ on $D^{\prime}$, one for each positive integer $e$. This always exists from (6.2). Further, we ask the system of these be compatible:
(8.1) For all integers $\left(e, e^{\prime}, e^{\prime \prime}\right)$ satisfying $e e^{\prime}=e^{\prime \prime}, g_{e^{\prime \prime}}(z)^{e^{\prime}}=g_{e}(z)$.

Call this collection $\left\{g_{e}\right\}_{e=1}^{\infty}=\mathcal{G}\left(D^{\prime}, z^{\prime}\right)$ a system of branches on $\left(D^{\prime}, z^{\prime}\right)$. The following is a slight enhancement of Lem. 7.8.

Proposition 8.4. Given $\mathcal{G}\left(D^{\prime}, z^{\prime}\right)$, any system of branches on $\left(D^{\prime}, z^{\prime}\right)$ corresponds one-one with elements of $\hat{\mathbb{Z}}$ (§7.5.1). Precisely: $\left\{m_{e}\right\} \in \hat{\mathbb{Z}} \mapsto\left\{\zeta_{e}^{m_{e}} g_{e}(z)\right\}_{e=1}^{\infty}$.

Let $D^{\prime \prime} \subset D^{\prime}$ be any (open) disk tangent to $z^{\prime}$. Restriction of $\mathcal{G}\left(D^{\prime}, z^{\prime}\right)$ to $D^{\prime \prime}$ defines a system of branches $\mathcal{G}\left(D^{\prime \prime}, z^{\prime}\right)$. Let $\boldsymbol{v}$ be the direction from $z^{\prime}$ along $a$ geodesic on $U_{\boldsymbol{z}}$ toward the center of $D^{\prime}$. (Consider $U_{\boldsymbol{z}}$ a subset of the sphere with its metric; geodesics being great circles.) Containment orders disks tangent to $z^{\prime}$ with $\boldsymbol{v}$ pointed into the disk. There is a maximal element

$$
\mathcal{G}\left(\boldsymbol{v}, z^{\prime}, U_{\boldsymbol{z}}\right)=\mathcal{G}\left(\boldsymbol{v}, z^{\prime}\right)=\mathcal{G}\left(D_{\boldsymbol{v}}, z^{\prime}\right):
$$

Take $D_{\boldsymbol{v}}$ the largest disk in $U_{\boldsymbol{z}}$ having radius along $\boldsymbol{v}$ and tangent to $z^{\prime}$.
So, the set of branch systems satisfying (8.1) is a homogeneous space for $\hat{\mathbb{Z}}$. That is, an action of the group $\hat{\mathbb{Z}}$ on one of them gives all. You still, however, need one choice $\mathcal{G}\left(D^{\prime}, z^{\prime}\right)$ to get the process going.

Definition 8.5. Call $\mathcal{G}\left(\boldsymbol{v}, z^{\prime}\right)=\hat{\boldsymbol{v}}$ a primitive (or naive) tangential base point: $\hat{\boldsymbol{v}}$ has an underlying point $z^{\prime}$, direction $\boldsymbol{v}$ and system of branches on $D_{\boldsymbol{v}}$.

From Cor. 7.5, elements in $\mathcal{P}_{z^{\prime}, e}$ have the form $f^{*}=h\left(\left(z-z^{\prime}\right)^{1 / e}\right)$ with $h \in \mathcal{L}_{z^{\prime}}$. Define $\operatorname{rest}_{\hat{\boldsymbol{v}}}\left(f^{*}\right)$ to be $h\left(g_{e}(z)\right)$. For any simply connected subspace $Y$ of $U_{\boldsymbol{z}}$, denote paths in $U_{\boldsymbol{z}}$ from $z_{0}$ with endpoint in $Y$ by $\Pi_{1}\left(z_{0}, Y\right)$.

Proposition 8.6 (Tangential Base Point Restriction). Assume $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ and $\gamma \in \Pi_{1}\left(z_{0}, D_{\boldsymbol{v}}\right)$. There is a unique $f^{*} \in \mathcal{P}_{z^{\prime}}$ with $\operatorname{rest} \hat{\boldsymbol{v}}\left(f^{*}\right)=f_{\lambda}$.

Proof. Uniqueness of $f^{*}$ is clear. Existence is from Cor. 7.5. Here are details. Let $\delta$ be a clockwise circle bounding a disk $\Delta_{z^{\prime}}$ with center $z^{\prime}$ with $\Delta_{z^{\prime}} \backslash\left\{z^{\prime}\right\} \subset U_{\boldsymbol{z}}$. Assume $\delta$ meets $D_{\boldsymbol{v}}$. Connect the end point of $\lambda$ to some point on $\delta$ by a path lying entirely in $D_{\boldsymbol{v}}$. From Cauchy's Theorem (Prop. 3.6), there is a unique function $g$ defined by a power series on $D_{\boldsymbol{v}}$ that restricts to $f_{\lambda}$. So, any analytic continuation of $f_{\lambda}$ along a path in $D_{\boldsymbol{v}}$ equals $g$. Thus it depends only on the end point of this path. Assume with no loss $\lambda$ ends on $\delta$.

Let $e=e_{f}$ be the order of the monodromy action of $\delta$ on $f_{\lambda}$. Then, Cor. 7.5 says $f_{\lambda}$ is $f^{*}=h\left(g_{e_{f}}(z)\right)$ with $h$ holomorphic in the disk $\delta$ bounds.

Example 8.7 (Deligne tangential base points). Take $z^{\prime}=0$ and $\boldsymbol{v}$ any direction $0 \leq \theta<2 \pi$ on $\mathbb{C}_{z}$ represented by $e^{i \theta}$. Define $g_{e}(z)$ to be $e^{i \theta / e}$ times the unique branch of $\left(e^{-i \theta} z\right)^{1 / e}$ taking positive real values along the direction $\boldsymbol{v}$ from 0: [De89, $\S 15]$ or [Ihar91, p. 103].
8.5. Describing all algebraic abelian functions. Suppose $f(z)$ is algebraic and $\mathbb{C}(z, f) / \mathbb{C}(z)$ is a Galois extension with abelian Galois group $G$. Assume $\boldsymbol{z}$ contains the branch points of $f$ and the ramification indices at all points of $\boldsymbol{z}$ divide some integer $e$. Each $z^{\prime} \in z$ produces an inertia group $I_{z^{\prime}}$ (Def. 7.10). More explicitly it produces a well defined conjugacy class $\mathrm{C}_{z^{\prime}}$ in $G$ (Lem. 7.9). Since, however, $G$ is abelian, this conjugacy class is an element $g_{f, z^{\prime}} \in G$.

THEOREM 8.8. Under the above hypotheses, $g_{f, z^{\prime}}$, as $z^{\prime}$ runs over $\boldsymbol{z}$, determines the field extension $\mathbb{C}(z, f)$. Further, two other properties hold.

- $\left\langle g_{f, z^{\prime}}, z^{\prime} \in \boldsymbol{z}\right\rangle=G$ : generation
- $\prod_{z^{\prime} \in \boldsymbol{z}} g_{f, z^{\prime}}=1$ : product-one condition

Conversely, suppose given $G$ and elements $g_{z^{\prime}} \in G$ for each $z^{\prime} \in z$ satisfying (8.8). Then, there exists algebraic $f$ (given as above by branches of log) satisfying $g_{f, z^{\prime}}=g_{z^{\prime}}$ for $z^{\prime} \in \boldsymbol{z}$. Another algebraic function $f^{*}$ produces the same data if and only if $\mathbb{C}\left(z, f^{*}\right)=\mathbb{C}(z, f)$.

Proof. There is a standard reduction for showing the field is determined by the data $g_{f, z^{\prime}}, z^{\prime} \in z$. Write $G$ as $\prod_{i=1}^{u} G_{i}$ where $G_{i}$ is cyclic of some prime power order. Every finite abelian group has this form ([Isa94, p. 90], see [9.15]). Then, $\mathbb{C}(z, f)$ is the composite of field extensions $L_{i} / \mathbb{C}(z)$ with group $G_{i}, i=1, \ldots, u$. Further, any subextension $\mathbb{C}(z)<M<L_{i}$ is Galois with group a quotient of $G_{i}$. So, it is cyclic of prime power order. So, with no loss assume $\mathbb{C}(z, f) / \mathbb{C}(z)$ is Galois with group isomorphic to $\mathbb{Z} / p^{t}$ for some integer $t$ and prime $p$. List $\boldsymbol{z}$ as $z_{1}, \ldots, z_{r}$, then list the group data as $\left(g_{1}, \ldots, g_{r}\right)$ with $g_{i}=g_{f, i}$ attached to $z_{i}$. Since $G=\mathbb{Z} / p^{t}$, identify $g_{i}$ with an integer $n_{i} \in \mathbb{Z} / p^{t}$.

It is easy to produce a cyclic extension that has exactly this attached data. For simplicity, assume $z_{r}=\infty$. Then, for any $z_{0}$ not in $\boldsymbol{z}$, let $h(z)=\prod_{i=1}^{r-1} h_{i}(z)^{n_{i}}$
with $h_{i}$ a branch of $\left(z-z_{i}\right)^{\frac{1}{p^{t}}}$ in a neighborhood of $z_{0}$. The lemma is done if $\mathbb{C}(z, h(z))=\mathbb{C}(z, f(z))$. Both fields embed in $\mathcal{P}_{z_{i}}$ and the action of $g_{z_{i}}$ restricts to both fields the same way. Any function in the fixed field of all the $g_{i} \mathrm{~s}$ is extensible over the whole Riemann sphere, as in $\S 7.1$. So such a function is a rational function in $z$. Therefore, the fixed field of $\left\langle g_{1}, \ldots, g_{r}\right\rangle$ in $\mathbb{C}(z, f(z))$ is trivial. Apply [9.5d] to the composite of the two fields and conclude they are equal.

Consider the generation condition. Assume $\left\langle g_{f, z^{\prime}}, z^{\prime} \in z\right\rangle=H$ is a proper subgroup of $G$. If $f_{1} \in \mathbb{C}(z, f)$ is in the fixed field of $H$, then $f_{1, \lambda}=f_{1}$ for all $\lambda \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Lem. 7.2 implies $f_{1} \in \mathbb{C}(z)$. So $\mathbb{C}(z)$ is the exact fixed field of $H$ and $H=G$. The product-one condition appears by recognizing $g_{f, z^{\prime}}$ as restriction of the $g_{e, z^{\prime}}$ for the field $L_{e, \boldsymbol{z}}$. Apply the product of the $g_{e, z^{\prime}}$ to generating functions in $L_{e, \boldsymbol{z}}=\mathbb{C}\left(z, e^{g_{i, j} / e}, 1 \leq i<j \leq r\right)$ (from (8.2)). It comes to showing $g_{e, z_{j}} g_{e, z_{i}}\left(e^{g_{i, j} / e}\right)=e^{g_{i, j} / e}$. With no loss take $z_{i}=0$ and $z_{j}=\infty[9.10 \mathrm{a}]$.

The full version of Riemann's Existence Theorem generalizes the generation and product-one conditions (8.8) to $\mathbb{C}(z, f(z))$ where $f$ is any algebraic function. When $G$ is abelian, the product-one condition is independent of the order of the elements $g_{f, z^{\prime}}$. Keep your eye on the analysis that goes into tracking the order of elements appearing in the product-one condition when $G$ is not abelian. This is what produces the significant action of the Hurwitz monodromy group in Chap. 5. Further, the converse holds in generality. Without, however, the abelian condition producing the algebraic function $f$ is more mysterious.

Suppose $G$ and $G^{*}$ are abelian groups and $\boldsymbol{g}_{\boldsymbol{z}}$ and $\boldsymbol{g}_{\boldsymbol{z}^{*}}^{*}$ satisfy the conditions of (8.8). Consider two triples $\mathcal{G}=\left(G, \boldsymbol{z}, \boldsymbol{g}_{\boldsymbol{z}}\right)$ and $\mathcal{G}^{*}=\left(G^{*}, \boldsymbol{z}^{*}, \boldsymbol{g}_{\boldsymbol{z}^{*}}^{*}\right)$ as in Thm. 8.8. Assume $\boldsymbol{z}$ is a subset of $\boldsymbol{z}^{*}$. For this discussion, if $z^{\prime} \in \boldsymbol{z}^{*} \backslash \boldsymbol{z}$ regard $\boldsymbol{g}_{\boldsymbol{z}}$ as having the identity element at $z^{\prime}$. Also, assume there is a homomorphism $\alpha: G^{*} \rightarrow G$ taking $g_{z^{\prime}}^{*}$ to $g_{z^{\prime}}$ for $z^{\prime} \in \boldsymbol{z}^{*}$. Regard $\alpha=\alpha_{\mathcal{G}^{*}, \mathcal{G}}$ as a map from $\mathcal{G}^{*}$ to $\mathcal{G}$.

Corollary 8.9. The projective system $\left\{\mathcal{G}, \alpha_{\mathcal{G}^{*}, \mathcal{G}}\right\}$ of triples with maps has a limit consisting of a group $\mathcal{G}^{\mathrm{ab}}$ and elements $g_{z^{\prime}}^{\mathrm{ab}}$ running over $z^{\prime} \in \mathbb{P}_{z}^{1}$. Then, $\mathcal{G}^{\mathrm{ab}}$ identifies with the maximal abelian quotient of the absolute Galois group of $\mathbb{C}(z)$. Also, $g_{z^{\prime}}^{\mathrm{ab}}$ acts trivially on any abelian algebraic function in $\mathcal{L}_{z^{\prime}}$ and identifies with a generator of the automorphisms of $\mathcal{P}_{z^{\prime}} / \mathcal{L}_{z^{\prime}}$ in its restriction to the abelian algebraic functions in $\mathcal{P}_{z^{\prime}}$ (Cor. 7.5).

Call the group $\mathcal{G}^{\text {ab }}$, the Galois group of the maximal abelian extension of $\mathbb{C}(z)$. A collection $\left\{g_{z^{\prime}}^{\mathrm{ab}}\right\}_{z^{\prime} \in \mathbb{P}_{z}^{1}}$ will be a canonical system of generators of $\mathcal{G}^{\mathrm{ab}}$. Any $g \in \mathcal{G}^{\mathrm{ab}}$ acts on the abelian algebraic functions in $\mathcal{P}_{z^{\prime}}$ for any $z^{\prime}$. This action is also the restriction of an automorphism of $\mathcal{P}_{z^{\prime}} / \mathcal{L}_{z^{\prime}}$. So monodromy action on a branch of $\log \left(z-z^{\prime}\right)$ determines this restriction element as a multiple of $g_{z^{\prime}}^{\mathrm{ab}} \in \hat{\mathbb{Z}}$.

## 9. Exercises

Some exercises remind of basic Galois Theory. Use char $(K)$ to denote the characteristic of a field $K$ : The minimal positive integer $n$ for which $n$ times the identity in $K$ is 0 (if such an integer exists, or 0 otherwise).
9.1. Substitutions and the chain rule. Consider more on (2.7c) as the defining property of analyticity.
(9.1a) For a path $\lambda:[a, b] \rightarrow \mathbb{C}$, compose it with any analytic function $h: \mathbb{C} \rightarrow \mathbb{C}$ to give $h \circ \lambda:[a, b] \rightarrow \mathbb{C}$, another path. If $g$ and $h$ satisfy (2.7c), show

$$
\begin{aligned}
\frac{d}{d t}(g \circ h)\left(\lambda\left(t_{0}\right)\right)=\frac{d}{d t}(g(h(\lambda)))\left(t_{0}\right) & =\frac{d g}{d w} \left\lvert\, w=h\left(\lambda\left(t_{0}\right)\right) \frac{d}{d t}(h \circ \lambda)_{\mid t=t_{0}}\right. \\
= & \frac{d g}{d w}\left(h \circ \lambda_{\mid t=t_{0}}\right) \frac{d h}{d z}\left(\lambda\left(t_{0}\right)\right) \frac{d \lambda}{d t}\left(t_{0}\right) .
\end{aligned}
$$

(9.1b) Show: Existence of $f^{\prime}\left(z_{0}\right)$ requires only checking (2.5) for $\lambda:[-1,1] \rightarrow D$ by $t \mapsto z_{0}+t v$ with $v \neq 0$. That is, check directional derivative rule (2.7b).
(9.1c) Conclude, if in (2.7c) two of $g \circ h, g, h$ are analytic, then so is the third.

With $m(z, w)=w^{k}-h(z)$ and $w(t)$ and $z(t)$ (nonconstant) rational functions with $w(t)^{k} \equiv h(z(t))$ for all $t$, consider indefinite integrals for $I(z)=\int h(z)^{\frac{1}{k}} d z$.
(9.2a) Substitute $z(t)$ for $t$. Rewrite $I(z)$ as an antiderivative for $\frac{d z(t)}{d t} / w(t)$. Apply this with $k=2$ and $h(z)=z^{2}+a z+b$ using [9.3d].
(9.2b) Ex. [9.10f] shows [9.2a] won't work often, not even with $k=2$ and $\operatorname{deg}(h)=3$ having no repeated roots. Show it does work for any $h$ with at most two distinct zeros, but arbitrary degree.
(9.2c) Calculus uses a different substitution: $w(t)$ and $z(t)$ are trigonometric in $t$ with $w(t)^{2}=z(t)^{2}+a z(t)+b$. Result: The square root expression disappears; replaced by a function. Why choose transcendental over rational functions? Hint: Consider the antiderivative as a function of $z$.
9.2. Rational functions and field theory. Suppose $K$ is any field. Consider $u(z)=P_{1}(z) / P_{2}(z)$ in $K(z)$. Follow the notation of $\S 1.2 .1$.
(9.3a) Show $P_{1}(w)-z P_{2}(w)$ is irreducible. Hint: Factor it as $m_{1}(z, w) m_{2}(z, w)$. Then compute the degree in $z$ of each factor.
(9.3b) Suppose $m \in K[z, w], \operatorname{deg}_{z}(m)=1$ and $m(z, f(z)) \equiv 0$ for some $f(z)$ analytic on a domain $D$. Show $K(z, f(z))=K(f(z))$.
(9.3c) If $M \leq L_{1} \leq L_{2}$ is a chain of fields, transitivity for degrees says [ $L_{2}$ : $M]=\left[L_{1}: M\right]\left[L_{2}: L_{1}\right]$. Use it to show $\operatorname{deg}\left(u_{1}\left(u_{2}(z)\right)\right)=\operatorname{deg}\left(u_{1}\right) \operatorname{deg}\left(u_{2}\right)$ for $u_{1}, u_{2} \in K(z) \backslash\{0\}$.
(9.3d) Suppose $M$ is a field and $\operatorname{char}(K) \neq 2$. Assume $m(z, w) \in K[z, w]$ of total degree 2 is irreducible, $z_{0}, w_{0} \in K, m\left(z_{0}, w_{0}\right)=0$ and $w^{\prime}$ is a zero of $m(z, w)$ in $\overline{K(z)}$. Show $K(z)\left(w^{\prime}\right)$ is isomorphic to $K(t)$ for some $t \in$ $K(z)\left(w^{\prime}\right)$. Hint: With $t$ and $s$ variables, let $z_{0}+s=z$ and $w^{\prime}=w_{0}+t s$. Solve for $s$ as a function of $t$ in $m\left(z,{ }^{\prime} w\right)=0$.
(9.3e) Show $z_{0}, w_{0} \in K$ is necessary for the existence of $t$ in (9.3d).
(9.3f) The fundamental theorem of algebra follows from knowing a function $f(z)$ bounded and analytic on $\mathbb{C}$ is constant. How does this imply every analytic function $P: \mathbb{P}_{z}^{1} \rightarrow \mathbb{P}_{w}^{1}$ (§4.6) by $z \mapsto P(z)$ is an element of $\mathbb{C}(z)$ ?
Now consider parametrizations by rational function curves. Use $\S 6.1 .2$ with $f=f_{1} / f_{2} \in \mathbb{C}(w)$ and $\left(f_{1}, f_{2}\right)=1$. Parametrize $X_{f}$ near $\left(z_{0}, w_{0}\right)$ if $w_{0}$ is not a zero of the Wronskian $f_{1}(w) f_{2}^{\prime}(w)-f_{2}(w) f_{1}^{\prime}(w)$ of $f_{1}, f_{2}$ and $f_{2}\left(w_{0}\right) \neq 0$.
(9.4a) Use Def. 4.14 to show this includes when $w_{0}$ is a zero of $f_{2}\left(z_{0}=\infty\right)$.
(9.4b) Extend a) to $w_{0}=\infty$. Show an analytic parametrization of a neighborhood by $(z, g(z))$ exists if and only if $\left|\operatorname{deg}\left(f_{1}\right)-\operatorname{deg}\left(f_{2}\right)\right| \leq 1$.
(9.4c) Suppose $f(g(z)) \equiv z$ for $g(z)$ analytic in a neighborhood of $z_{0}$. With these extensions, show the maximal number of branch points for $\mathbb{C}(z, g(z))$ $(\S 6.2)$ is $2(\operatorname{deg}(f)-1)$ with equality occurring for some rational functions $f$ of degree $n$ for any positive integer $n$.
(9.4d) Suppose $w_{0}$ is a zero of $f_{1}(w) f_{2}^{\prime}(w)-f_{2}(w) f_{1}^{\prime}(w)$ of multiplicity $e_{w_{0}}-1$ and $f\left(w_{0}\right)=z_{0}$. Apply the Cor. 7.5 proof to find $e_{w_{0}}$ distinct functions $g(u)$ analytic around 0 with $g(0)=w_{0}$ and $f(g(u))-z_{0}-u^{e_{w_{0}}} \equiv 0$ ? Extend d) to have either $z_{0}$ or $w_{0}$ is $\infty$. Conclude for $f \in \mathbb{C}(z) \backslash \mathbb{C}$ :

$$
\begin{equation*}
2(\operatorname{deg}(f)-1)=\sum e_{w_{0}}-1 \tag{9.4e}
\end{equation*}
$$

9.3. Galois theory of composite fields and using group theory. Suppose $L_{1} / K$ and $L_{2} / K$ are two field extensions. Given a field $L$ containing both $L_{1}$ and $L_{2}$, there is an immediate minimal field $L_{1} \cdot L_{2}$ in $L$ containing them both [Isa94, Chap. 18].
(9.5a) Suppose $M / K$ is Galois: Its group of automorphisms $G(M / K)=G$ fixed on $K$ has order $[M: K]$. Consider $K<L<M$, a chain of fields. Suppose $L=L_{1}, \ldots, L_{n}$ are the fields conjugate to $L / K$. Show $L_{1} \cdot L_{i}=L_{1}$, $i=1, \ldots, n$, if and only if $L / K$ is Galois $(G(M / L)$ is a normal subgroup; closed under conjugation from $G$ ).
(9.5b) Let $T: G \rightarrow S_{n}$ be the permutation representation of $G$ on cosets of $G(M / L)$ (as in a). Show there is $j \neq 1$ with $L_{1}=L_{1} \cdot L_{j}$ if and only if $(1) T(g)=1 \Leftrightarrow(j) T(g)=j$ for each $g \in G$.
(9.5c) The following notation holds for the next two subexercises. Suppose $M_{i} / K$ is Galois with group $G_{i}, i=1,2$. Consider the group $G$ defined as follows:

$$
\left\{g=\left(g_{1}, g_{2}\right) \in G_{1} \times G_{2} \mid g_{1}(\alpha)=g_{2}(\alpha), \alpha \in M_{1} \cap M_{2}\right\}
$$

Show $G$ acts as automorphisms of $M_{1} \cdot M_{2}$.
(9.5d) Show $|G|=\left[M_{1} \cdot M_{2}: K\right]$, and so $M_{1} \cdot M_{2} / K$ is Galois with group $G$. Hint: Apply the Fundamental Theorem of Galois Theory [Isa94, Thm. 18.21] to the fixed field of $G$.
(9.5e) Conclude $M_{1} \cdot M_{2}$ doesn't depend (up to isomorphism over $K$ ) on what field they both sit inside if both extensions are Galois.
(9.5f) Assume char $(K)$ is $p$ (a prime or 0 ). Suppose $K$ has at most one extension of degree $n$ for any integer $n>0$ (or if $p>0$, prime to $p$ ). Show extensions of $K$ of degree prime to $p$ are Galois with cyclic group.
We warmup in interpreting field theory with group theory. Let $K=\mathbb{C}(z)$. If $f$ is algebraic over $K$ denote $K(f)$ by $L_{f}$, and the Galois closure of $L_{f} / K$ by $\hat{L}_{f}$. Suppose $m_{i} \in \mathbb{C}[z, w]$, of degree $n_{i}$ in $w$, is the irreducible polynomial for a function $f_{i}$ (algebraic according to (1.2)) over $K, i=1,2$. Denote $G\left(\hat{L}_{f_{i}} / K\right)$ by $G_{i}, i=1,2$. As in [9.5d], regard $G \stackrel{\text { def }}{=} G\left(\hat{L}_{f_{1}} \cdot \hat{L}_{f_{2}} / K\right)$ as a subgroup of $S_{n_{1}} \times S_{n_{2}}$. Let $\pi_{i}: G_{1} \times G_{2} \rightarrow G_{i}$ be projection on the $i$ th factor.
(9.6a) For $H$ a subgroup of $G_{1} \times G_{2}$, let $\operatorname{ker}\left(\pi_{i}(H)\right)$ be the kernel of projection of $H$ on $G_{i}$. For $H \leq G_{1} \times G_{2}$ with $\pi_{i}(H)=G_{i}, i=1,2$, let $A_{H}$ be $\left\langle\operatorname{ker}\left(\pi_{1}(H)\right), \operatorname{ker}\left(\pi_{2}(H)\right)\right\rangle$. Show $H=\left\{\left(g_{1}, g_{2}\right) \mid \psi_{1}\left(g_{1}\right)=\psi_{2}\left(g_{2}\right)\right\}$ with $\psi_{i}: G_{i} \rightarrow G_{1} \times G_{2} / A_{H}=G_{H}: H$ is the fiber product of $\psi_{1}$ and $\psi_{2}$.
(9.6b) Consider $L / F$ and $F / M$ algebraic field extensions, with $\psi: F \rightarrow \bar{M}$ an embedding of $F$ in the algebraic closure of $M$. Galois theory depends on the Extension Theorem [Isa94, Thm. 17.30]: There exists an embedding $\psi^{\prime}: L \rightarrow \bar{M}$ extending $\psi$. Explain why this shows $\pi_{i}(G)=G_{i}, i=1,2$.
(9.6c) Let $G_{2}(1)=G\left(\hat{L}_{f_{2}} / L_{f_{2}}\right)$. Consider $\pi_{2}^{-1}\left(G_{2}(1)\right)$, the biggest subgroup of $G$ projecting to $G_{2}(1)$. Show $m_{1}$ is irreducible over $L_{f_{2}}$ if and only if $\pi_{1}\left(\pi_{2}^{-1}\left(G_{2}(1)\right)\right)$ is transitive.
(9.6d) Let $f^{(1)}, \ldots, f^{(n)}$ be the conjugates of $f^{(1)}=f$ with $f$ algebraic over $K$ of degree $n$. Denote $G\left(\hat{L}_{f} / K\left(f^{(i)}\right)\right)$ by $G(i)$. Show: $K\left(f^{(1)}\right)$ contains $f^{(i)}$ if and only if $G(1)=G(i)$.
9.4. Branch of $\log$ and Puiseux expansions. Assume $D \subset \mathbb{C}^{*}$ is a domain.
(9.7a) A classical domain $D$ supporting a branch of $\log$ on $D$ is any (subdomain of a) sector: $S_{\theta_{1}, \theta_{2}}=\left\{r e^{i \theta} \mid \theta_{1}<\theta<\theta_{2}\right\}$ under the condition $\theta_{2}-\theta_{1} \leq 2 \pi$. Give the branches of $\log$ on $S_{\theta_{1}, \theta_{2}}$.
(9.7b) If $H_{1}(z)$ and $H_{2}(z)$ are two branches of $\log$ in $D$ and $H_{1}\left(z_{0}\right)=H_{2}\left(z_{0}\right)$ for $z_{0} \in D$, show $H_{1}(z)=H_{2}(z)$ for $z \in D$.
(9.7c) Prop. 3.2 shows there exists a branch $g_{\lambda}$ of $\log$ along any path in $D$. If for any $\lambda \in \Pi_{1}\left(D, z_{0}\right), g_{\lambda}(1)=g_{\lambda}(0)$, show there is a branch of $\log$ on $D$. Hint: Let $G(z)$ be $g_{\lambda}(b)$ with $\lambda:[a, b] \rightarrow D$ so $\lambda(a)=z_{0}, \lambda(b)=z$ and $g_{\lambda}$ is a branch of $\log$ along $\lambda$ with $g_{\lambda}(a)=w_{0}$ (fixed). Apply Lem. 4.11.
(9.7d) Show there is a branch of $\log$ in a domain $D$ if and only if each closed path in $D$ has winding number 0 about the origin.
(9.7e) Consider $\gamma_{1}, \gamma_{2} ;[0,1] \rightarrow \mathbb{P}_{z}^{1}$ with these properties: $\gamma_{1}(0)=\gamma_{2}(0)=0$, $\gamma_{1}(1)=\gamma_{2}(1)=\infty$, and for $t \in(0,1) \gamma_{1}(t) \neq \gamma_{2}(t)$, and $\gamma_{i}(t) \in \mathbb{C}^{*}$, $i=1,2$. Let $D$ be any component ([9.17e]: there are two) of $\mathbb{C}^{*} \backslash\left\{\gamma_{1}, \gamma_{2}\right\}$. Show there is a branch of $\log$ in $D$.
Assume $f(z)$ is analytic near $z_{0}$ and algebraic according to (1.2): $m(z, f(z)) \equiv 0$ for some nonzero $m \in \mathbb{C}[z, w]$.
(9.8a) Why can we assume $m(z, w)$ is irreducible in the ring $\mathbb{C}[z, w]$ ? How does this same observation show the ring of analytic functions on a (connected) domain $D$ is an integral domain. Hint: $h(z)$ analytic on $D$ and zero at a set with a limit point in $D$ is identically zero [Ahl79, p. 127].
(9.8b) Assume $\left(f, D, z_{0}\right)$ is extensible. As in (1.1), why does $h(z) \in \mathcal{A}_{f}(D)$ also satisfy $m(z, h(z)) \equiv 0$. Conclude: $f(z)$ satisfies (1.1b).
(9.8c) Note in b) for given $D$, the conclusion requires only that $m(z, w)$ has coefficients meromorphic on $D$ (not necessarily on $\mathbb{P}_{z}^{1}$ ).
(9.8d) Use $\S 6.1$ to complete showing $f(z)$ satisfies (1.1).
(9.8e) Suppose $f(z)$ is a branch of $\log$ on $D$. Show it satisfies neither of the properties (1.1a) or (1.1b). Yet, it does satisfy (1.1c).
(9.8f) If $g: D_{1} \rightarrow D$ is analytic and $f(g(z)) \equiv z$, show $g(z)$ satisfies (1.2).
(9.8g) Suppose $f \in \mathcal{H}(\mathbb{C})$. Let $\boldsymbol{z}=\{\infty\}$. Then, $f$ satisfies (1.1a) and (1.1b). Suppose $f$ is not a polynomial function. Show it doesn't satisfy (1.1c). Hint: Apply the Caseroti-Weierstrass theorem [Con78, p. 109].
Consider how branches of log closely tie to Puiseux expansions. Use notation of $\S 1.3$ for the field $\mathcal{L}_{z^{\prime}}$ around $z^{\prime}$. For integer $e>1$ create a copy $\mathcal{P}_{z^{\prime}, e}$ of $\mathcal{L}_{z^{\prime}}$ by replacing $z-z^{\prime}$ by a new variable $u_{e}$. Set $e^{2 \pi i / e}=\zeta_{e}$.
(9.9a) Why is $\mathcal{L}_{z^{\prime}}$ a field?
(9.9b) Suppose $e \mid e^{*}: t=e^{*} / e$. Map $\mathcal{P}_{z^{\prime}, e}$ to $\mathcal{P}_{z^{\prime}, e^{*}}$ by substituting $u_{e^{*}}^{t}$ for $u_{e}$. Show this map extends to a field homomorphism.
(9.9c) Identify $\mathcal{P}_{z^{\prime}, e}$ with its image in $\mathcal{P}_{z^{\prime}, e^{*}}$. Form the union, the ring of Puiseux expansions $\mathcal{P}_{z^{\prime}}$, over all $e$. Why is it a field?
(9.9d) Show $\mathcal{P}_{z^{\prime}, e}$ is a Galois extension of $\mathcal{L}_{z^{\prime}}$ with group $\mathbb{Z} /(e)$. Hint: A generator acts by $u_{e} \mapsto \zeta_{e} u_{e}$.
(9.9e) Suppose $z_{0} \neq z^{\prime}$. Let $h(z)$ be a branch of $\log \left(z-z^{\prime}\right)$ in a neighborhood $D$ of $z_{0}$. Show $f_{e}(z)=e^{h(z) / e}$ is a branch of solutions of $w^{e}=z-z^{\prime}$. So $f(z)$ is an algebraic function.
(9.9f) If $e>1$, show $f_{e}(z)$ is not the analytic continuation of a function in $\mathcal{L}_{z^{\prime}}$.
$(9.9 \mathrm{~g})$ Consider $\varphi: \mathbb{P}_{w}^{1} \rightarrow \mathbb{P}_{z}^{1}$ by $w \mapsto w^{e}+z^{\prime}$. Form $g(w)=f_{e} \circ \varphi$ and show it is an analytic continuation of some function (of $w$ ) around 0 .
We may equally consider Puiseux expansions at $\infty$. Denote the Laurent series around $\infty$ by $\mathcal{L}_{\infty}$ : expressions $(1 / z)^{n} h(1 / z)$ with $n$ an integer and $h(z)$ convergent near $z=0$. As in [9.9], form a copy $\mathcal{P}_{\infty, e}$ of $\mathcal{L}_{\infty}$ by replacing $1 / z$ by $u_{e}$.
(9.10a) Follow [9.9] to form $\mathcal{P}_{\infty}$, the analog of $\mathcal{P}_{z^{\prime}}$. Analytically continue a branch of $z^{1 / e}$ counterclockwise on a circle around $\infty$. Hint: Apply $z \mapsto 1 / z$; it is the same as continuing $z^{-1 / e}$ clockwise around the origin.
(9.10b) For $f(w) \in \mathbb{C}[w]$ of degree $n$ with leading coefficient 1 , write $f(w)=$ $w^{n}+a_{n-1} w^{n-1}+\cdots+a_{0}$, let $m(z, w)=f(w)-z$. Show there is $g(z) \in \mathcal{P}_{\infty}$ of form $z^{\frac{1}{n}}+\sum_{j=0}^{\infty} b_{j} z^{-\frac{j}{n}}$ with $f(g(z)) \equiv z$.
(9.10c) Let $L_{f}$ be $\mathbb{C}(z, g(z)), g$ from b). Let $\hat{L}_{f} / \mathbb{C}(z)$ be the splitting field of $L_{f} / \mathbb{C}(z)$. Show there is $g \in G\left(\hat{L}_{f} / \mathbb{C}(z)\right)$ acting as an $n$-cycle on conjugates of $g(z)$. Hint: Apply $1 / z^{\frac{1}{n}} \mapsto \zeta_{n} 1 / z^{\frac{1}{n}}$.
(9.10d) Consider $f, h \in \mathbb{C}[w]$ with $\operatorname{deg}(h)=m$. Apply [9.6c] to $\hat{L}_{f}$ and $\hat{L}_{h}$. Show the group of $\hat{L}_{f} \cdot \hat{L}_{h} / \mathbb{C}(z)$ contains $\sigma$ of order $n m / \operatorname{gcd}(n, m)$ with restriction of $\sigma$ to $\hat{L}_{f}$ an $n$-cycle and its restriction to $\hat{L}_{h}$ an $m$-cycle.
(9.10e) If $(\operatorname{deg}(f), \operatorname{deg}(h))=1$, show $f(w)-h(u)$ is irreducible. Hint: Irreducibility is equivalent to $[K(w): K]=\operatorname{deg}(w)$ with $K=\mathbb{C}(u)$. Use that d) shows $\left[K^{\prime}(w): K^{\prime}\right]=\operatorname{deg}(w)$ with $K^{\prime}=\mathbb{C}((1 / u))$.
(9.10f) Suppose in d) (with $(\operatorname{deg}(f), \operatorname{deg}(h))=1), L_{f} \cdot L_{h}$ is pure transcendental (equals $\mathbb{C}(t))$. Show for some choice of $t$ there are polynomials $g(t), k(t)$ of respective degrees $m$ and $n$ with $f(g(t))=h(k(t))$.
(9.10g) Apply f) to $f(w)=w^{2}$ and $h(u)=u^{3}-a u-b$ where $h$ has distinct zeros. Show $L_{f} \cdot L_{h}$ is not pure transcendental. Hint: Zeros of $g(t)^{2}$ are multiple.
Critical points over $z \in \mathbb{C}$ appear in (6.8). Now consider $z=\infty$. With $m \in \mathbb{C}[z, w]$ of degree $n$ and $m=h_{0}(z) w^{n}+h_{1}(z) w^{n-1}+\cdots+h_{n}(z)$, assume $h_{0}$ has $z_{0}$ as multiplicity $t$ zero. When $h_{0}$ is constant call $m$ integral (over $z$ ).
(9.11a) Write $t=k n+t_{0}$ with $0 \leq t_{0}<n$. Show there is an integral polynomial $m_{1}(z, w) \in \mathbb{C}[z, w]$ satisfying $m_{1}\left(z,\left(z-z_{0}\right)^{k+1} w\right) \equiv\left(z-z_{0}\right)^{n-t_{0}} m(z, w)$.
(9.11b) Suppose $K$ is a field and $P_{1}, P_{2} \in K[w]$. The Euclidean algorithm gives the greatest common divisor of $P_{1}$ and $P_{2}$. Write $P_{1}=R_{0}, P_{2}=R_{1}$. Form the remainder $R_{2}$ of the division $R_{1} \sqrt{R_{0}}$. Inductively form successive remainders, $R_{3}, \ldots, R_{u}$, until the next stage remainder is 0 . Do an induction to produce $A(w), B(w) \in K[w]$ with $A(w) P_{1}(w)+B(w) P_{2}(w)=R_{u}(w)$.
(9.11c) Continue b): Use that $\mathbb{C}[z]$ has unique factorization to clear denominators on $A(w) P_{1}(w)+B(w) P_{2}(w)=R_{u}(w)$. Suppose $P_{i}=P_{i}(z, w) \in \mathbb{C}[z, w]$, $i=1,2$, have no common factor in $w$. Find $A(z, w), B(z, w) \in \mathbb{C}[z, w]$ and $M(z) \in \mathbb{C}[z] \backslash\{0\}$ with $A(z, w) P_{1}(z, w)+B(z, w) P_{2}(z, w)=M(z)$.
(9.11d) Result c) applies with any unique factorization domain replacing $\mathbb{C}[z]$. Comment on how it applies to $K=\mathcal{L}_{z^{\prime}}$.
(9.11e) We outline examples where critical points of $(m, w)(m(z, f(z)) \equiv 0)$ properly contain critical points of $\mathbb{C}(z, f) / \mathbb{C}(z)$. Let $g_{z_{0}}$ be the conjugacy class of the branch cycle for $m$ at $z_{0}$. Suppose $e=e_{z_{0}}$ is the order of $g_{z_{0}}$. Show, if $m\left(u^{e}+z_{0}, w\right)=m_{1}(u, w)$ is irreducible, then $u=0$ is a branch point of $m_{1}(u, w)$ but not a branch point of $\mathbb{C}\left(u, f\left(u^{e}+z_{0}\right)\right)$.
(9.11f) Apply [9.10e] to give examples of e) by taking $h \in \mathbb{C}[w]$ of degree prime to $e$, so $h(w)-u^{e}$ is irreducible.
9.5. Elementary permutations from $\Pi_{1}\left(D, z_{0}\right)$. Let $\Delta_{z^{\prime}}$ be a disk about $z^{\prime}$ and $\Delta_{z^{\prime}}^{0}=\Delta_{z^{\prime}} \backslash\left\{z^{\prime}\right\}$. Choose $z_{0} \in \Delta_{z^{\prime}}^{0}$.
(9.12a) Suppose $h(t)$ is a branch of $\log \left(z-z^{\prime}\right)$ along $\lambda:[a, b] \rightarrow \mathbb{C}-\backslash\left\{z^{\prime}\right\}$. Then, what path is $h(t)$ a branch of log along?
(9.12b) Suppose $f(z)=\left(z-z^{\prime}\right) h(z)$ is analytic in $\Delta_{z^{\prime}}$ with $h(z) \neq 0$ for any point in $\Delta_{z^{\prime}}$. Show a branch $F(z)$ of $f(z)^{\frac{1}{e}}$ exists at any point in $\Delta_{z^{\prime}}^{*}$. Further, show there is an embedding of the field $\mathbb{C}(z, F(z))$ into $\mathcal{P}_{z^{\prime}, e}$.
(9.12c) Let $g_{j}(z)$ be a branch of $\left(z-z_{j}\right)^{1 / e_{j}}, j=1, \ldots, r$ analytic in a neighborhood of $z_{0}$. With $f(z)=\prod_{j=1}^{r} g_{j}$ and $\lambda:[a, b] \rightarrow \mathbb{C}$ a path with winding number $m_{j}$ around $z_{j}$, explicitly relate $f(z)$ and $f_{\lambda}(z)$.
Consider how analytic continuation easily forces us into groups that are not abelian. Follow Thm. 5.6 notation.
(9.13a) Show the conclusion of the case $\infty \in D$ as in $\S 5.4 .4$ follows.
(9.13b) Recall the semi-direct product $M \times{ }^{s} H$ of groups of $H$ and $M$ with $\psi$ : $H \rightarrow \operatorname{Aut}(M)$ a homomorphism into the automorphisms of $M$. Then, $(m, h) \cdot\left(m^{\prime}, h^{\prime}\right) \stackrel{\text { def }}{=}\left(m \cdot \psi(h)(m), h \cdot h^{\prime}\right)$ defines multiplication on $M \times H$. Consider $M_{0}=\mathbb{Z}^{3}$, and $H_{0}=\mathbb{Z} / 3$ where $1 \in H_{0}$ maps $\left(m_{1}, m_{2}, m_{3}\right) \in M_{0}$ to $\left(m_{2}, m_{3}, m_{1}\right)$. Show $M_{0} \times{ }^{s} H_{0}$ is not abelian.
(9.13c) Let $f(z)$ be a branch of $z^{1 / 3}$ around $z_{0} \neq 0$. For $a \notin\left\{0, \infty, z_{0}\right\}$, consider $h=\frac{1}{f(z)^{2}(f(z)-a)} \in \mathbb{C}(z, f(z))$. Find $\boldsymbol{z} \subset \mathbb{P}_{z}^{1}$ so $\left(h, U_{\boldsymbol{z}}\right)$ is extensible. Find the image of the permutation representation of $\Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ on $\mathcal{A}_{h}\left(U_{\boldsymbol{z}}\right)$.
(9.13d) Let $H(z)$ be a primitive for $h($ in d$)$ ) around $z_{0}$. Show the image of the permutation representation of $\Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ on $\mathcal{A}_{H}\left(U_{\boldsymbol{z}}\right)$ is $M_{0} \times{ }^{s} H_{0}$ from b). Hint: Substitute $w$ with $w^{3}=z$.
9.6. Fractional transformations and the elementary divisor theorem. Recall: For any ring $R$ and integer $n \geq 1, \operatorname{PGL}_{n}(R)$ is $\operatorname{GL}_{n}(R) /\left\langle R^{*} I_{n}\right\rangle$ and $\operatorname{PSL}_{n}(R)=\mathrm{SL}_{n}(R) / \mathrm{SL}_{n}(R) \cap\left\langle R^{*} I_{n}\right\rangle$. Several nonabelian subgroups of $\mathrm{PGL}_{2}(\mathbb{C})$, like $\operatorname{PGL}(\mathbb{R})$ and $\mathrm{PSL}_{2}(\mathbb{Z})$ appear often in complex variables. We contrast their different appearances. Let $\mathcal{T}$ be the translations $\left\{\alpha \in \mathrm{PGL}_{2}(\mathbb{C}) \mid \alpha(z)=z+a, a \in \mathbb{C}\right\}$. Let $\mathcal{M}$ be the multiplications $\left\{\alpha \in \mathrm{PGL}_{2}(\mathbb{C}) \mid \alpha(z)=b z, a \in \mathbb{C}^{*}\right\}$. Finally, consider $\tau: z \mapsto 1 / z$.
(9.14a) Show each $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$ has is one of $a^{\prime}\left(z-z_{1}\right), a^{\prime}\left(z-z_{1}\right) /\left(z-z_{2}\right)=$ $a^{\prime}\left(1+\left(z_{2}-z_{1}\right) /\left(z-z_{2}\right)\right)$, or $a^{\prime} /\left(z-z_{2}\right)$. Why is $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$ a composition of elements from $\mathcal{M}, \mathcal{T}$ and $\tau: \mathcal{M}, \mathcal{T}$ and $\gamma$ generate $\mathrm{PGL}_{2}(\mathbb{C})$.
(9.14b) Give an $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$ mapping $\mathbb{R}$ to the boundary of the unit circle.
(9.14c) Elements of $\mathrm{PGL}_{2}(\mathbb{C})$ mapping $\mathbb{R} \cup\{\infty\}$ to itself are in $\mathrm{PGL}_{2}(\mathbb{R})$. What is the subgroup of these mapping the upper half plane $\mathbb{H}$ (Chap. $3 \S 3.2 .2$ ) into itself? Hint: $z \mapsto 1 / z$ does not.
(9.14d) Combine with b) to describe elements of $\mathrm{PGL}_{2}(\mathbb{C})$ mapping $\mathbb{R} \cup\{\infty\}$ to the unit circle. Which map $\mathbb{H}$ to the inside of the circle?
(9.14e) Which $f \in \mathbb{C}(z)$ map the unit circle into the unit circle. Hint: $f \in \mathbb{C}(z)$ mapping $\mathbb{R} \rightarrow \mathbb{R}$ has zero and pole set closed under complex conjugation.
Let $R$ be a principal ideal domain, $M$ a finitely generated free $R$ module, and $N$ an $R$ submodule of $M$. The Elementary Divisor Theorem (EDT [Jac85, p. 192]): There is a basis $\boldsymbol{v}_{1}, \ldots, \boldsymbol{v}_{m}$ of $M$ and elements $a_{1}, \ldots, a_{m} \in R$ with nonzero elements of $a_{1} \boldsymbol{v}_{1}, \ldots, a_{m} \boldsymbol{v}_{m}$ a basis of $N$. If $a_{1}, \ldots, a_{t}$ are the nonzero $a_{i} \mathrm{~s}$, then we may choose $a_{1}, \ldots, a_{t}$ so $a_{i} \mid a_{i+1}, i=1, \ldots, t$.
(9.15a) Consider an abelian group quotient $A$ of $\mathbb{Z}^{n}$. Apply EDT to show $A$ is isomorphic to $\oplus_{i=1}^{n} \mathbb{Z} /\left(a_{i}\right)$ for some integers $a_{1}, \ldots, a_{n} \in \mathbb{Z}$.
(9.15b) Show in a), if $A$ is a finite group and $a_{1}\left|a_{2}\right| \cdots \mid a_{m}$ are positive integers, then the $a_{1}, \ldots, a_{n}$ are unique.
(9.15c) $\mathrm{SL}_{2}(\mathbb{Z})(2 \times 2$ matrices over $\mathbb{Z}$ of determinant 1$)$ acts on $M_{2}=\mathbb{Z}^{2}$ taking one basis to another. If $N$ is a subgroup of $M_{2}$ of index $n$, then $\mathrm{SL}_{2}(\mathbb{Z})$ maps it in an orbit of index $n$ subgroups. Apply EDT to count $N \leq M_{2}$ of index $n=p^{k}$ ( $p$ a prime). Hint: Start with $N$ for which $M / N$ is cyclic.
(9.15d) Each $N$ from c) defines a subgroup $\Gamma_{N}$ of $\operatorname{PSL}_{2}(\mathbb{Z})$ : the image of the stabilizer in $\mathrm{SL}_{2}(\mathbb{Z})$ of $N$. If $n=p$ is a prime, and $U$ is the biggest normal subgroup of $\mathrm{PSL}_{2}(\mathbb{Z})$ in $\Gamma_{N}$, show $\mathrm{PSL}_{2}(\mathbb{Z}) / U=\mathrm{PSL}_{2}(\mathbb{Z} / p)$.
Let $\Delta$ be the open unit circle. Denote the linear fractional transformations that map $\Delta \rightarrow \Delta$ by $\mathrm{PGL}_{2}(\Delta)$. Form

$$
\left(w_{3}-w_{1}\right)\left(w-w_{2}\right) /\left(w_{2}-w_{1}\right)\left(w-w_{3}\right)=L(w)=L\left(w_{1}, w_{2}, w_{3}, w\right)
$$

for $w_{1}, w_{2}, w_{3} \in \mathbb{C}$. This problem follows a treatment from [Spr57, $\left.\S 9.2\right]$
(9.16a) Use [9.14]. Show $\mathrm{PGL}_{2}(\mathbb{C})$ fixes $L(w)$ :
$L\left(w_{1}, w_{2}, w_{3}, w\right)=L\left(\alpha\left(w_{1}\right), \alpha\left(w_{2}\right), \alpha\left(w_{3}\right), \alpha(w)\right)$, for $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$.
(9.16b) Suppose $w_{1}, \ldots, w_{4} \in \mathbb{C}$ are on a circle in that order. Show: $L\left(w_{4}\right)>1$. Conclude: With $w_{1}, w_{2}, w_{3}$ fixed, $w \mapsto L(w)$ maps the interior of the disk bounded counterclockwise by $w_{1}, w_{2}, w_{3}$ to the upper half plane $\mathbb{H}$.
(9.16c) Suppose $w_{2}, w_{3} \in \Delta$. Let $C_{w_{2}, w_{3}}$ be the unique circle containing $w_{2}$ and $w_{3}$ meeting the unit circle at right angles (at two points). Why is $C_{w_{2}, w_{3}}$ unique? Hint: Use $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$ taking the unit circle to the real line.
(9.16d) Let $w_{1}$ be the point on $C_{w_{2}, w_{3}} \cap \partial \Delta$ closest to $w_{2}$. Similarly, $w_{4}$ is the other point of intersection closest to $w_{3}$. Define the distance $d\left(w_{2}, w_{3}\right)$ to be $\frac{1}{2} \log \left(L\left(w_{1}, w_{2}, w_{3}, w_{4}\right)\right)$. When $w_{2}=0$ and $w_{3}=r e^{i \theta}$ express this as a function of $r$.
(9.16e) Notice $\beta_{w_{2}}(w)=\frac{w-w_{2}}{1-\bar{w}_{2} w}$ is in $\mathrm{PGL}_{2}(\Delta)$ and it maps $w_{2} \mapsto 0$. Use this to $\operatorname{express} d\left(w_{2}, w_{3}\right)$ as $\frac{1}{2} \log \left(\frac{1+\left|\beta_{w_{2}}\left(w_{3}\right)\right|}{1-\left|\beta_{w_{2}}\left(w_{3}\right)\right|}\right)$.
9.7. Metrics on $\mathbb{P}_{z}^{1}, \Delta$ and more generally. The metric topology on $\mathbb{P}_{z}^{1}$ identifies it with the sphere around the origin in $\mathbb{R}^{3}$. Use coordinates $(r, u, v)$ : $z_{0} \in \mathbb{P}_{z}^{1} \mapsto\left(r_{0}, u_{0}, v_{0}\right) \in \mathbb{R}^{3}$. The unit sphere has this analytical description: $\left\{(r, u, v) \mid r^{2}+u^{2}+v^{2}=1\right\}=S$.
(9.17a) From vector calculus, this implicit description of $S$ gives a unit normal direction to $S$ at $\left(r_{0}, u_{0}, v_{0}\right)$. It is a unit vector $\mathbb{N}_{\left(r_{0}, u_{0}, v_{0}\right)}$, (from the origin)
in the direction of the gradient of $f(r, u, v)=r^{2}+u^{2}+v^{2}$. Compute two such vectors. Which suits the definition of outward normal vector?
(9.17b) Let $\mathbb{T}_{\left(r_{0}, u_{0}, v_{0}\right)}$ be points on the plane through $\left(r_{0}, u_{0}, v_{0}\right)$ tangent to the sphere. There are two possible definitions of $\mathbb{T}_{\left(r_{0}, u_{0}, v_{0}\right)}$. Suppose the range of $(x, y) \mapsto(r(x, y), u(x, y), v(x, y))=H(x, y)$ is a neighborhood of $\left(r_{0}, u_{0}, v_{0}\right) ; H$ is differentiable in a neighborhood of the origin and $H(0,0)=\left(r_{0}, u_{0}, v_{0}\right)$, and $\frac{\partial H}{\partial x}(0,0)$ and $\frac{\partial H}{\partial y}(0,0)$ are linearly independent vectors in $\mathbb{R}^{3}$. Apply the chain rule to show
$\mathbb{T}_{\left(r_{0}, u_{0}, v_{0}\right)}^{\dagger} \stackrel{\text { def }}{=}\left\{\left.\left(r_{0}, u_{0}, v_{0}\right)+x \frac{\partial H}{\partial x}(0,0)+y \frac{\partial H}{\partial y}(0,0) \right\rvert\,(x, y) \in \mathbb{R}^{2}\right\}$
is independent of the choice of $H$.
(9.17c) The second definition of $\mathbb{T}_{\left(r_{0}, u_{0}, v_{0}\right)}$ is

$$
\mathbb{T}_{\left(r_{0}, u_{0}, v_{0}\right)}^{\dagger \dagger} \stackrel{\text { def }}{=}\left\{(r, u, v) \mid\left((r, u, v)-\left(r_{0}, u_{0}, v_{0}\right)\right) \cdot \mathbb{N}_{\left(r_{0}, u_{0}, v_{0}\right)}=0\right\}
$$

Use the expression $f(H(x, y)) \equiv 0$ to show $\mathbb{T}_{\left(r_{0}, u_{0}, v_{0}\right)}^{\dagger \dagger}=\mathbb{T}_{\left(r_{0}, u_{0}, v_{0}\right)}^{\dagger}$.
(9.17d) Let $\gamma:[a, b] \rightarrow S$ be a simple closed path. Suppose $\frac{d \gamma}{d t}$ exists and is nonzero at $t_{0} \in[a, b]$. Define the direction to the left of $\gamma$ at $t_{0}$ to be the unit vector $\boldsymbol{u}_{1}$ for which $\operatorname{det}\left(\boldsymbol{u}_{1}\left|\mathbb{N}_{\gamma\left(t_{0}\right)}\right| \frac{d \gamma}{d t}\left(t_{0}\right)\right)$ is positive.
(9.17e) The complement $S \backslash \gamma$ of a simple closed path has two components $U_{1}$ and $U_{2}$ : The Jordan curve Theorem. For simplicial $\gamma$ this is easy (Chap. 4 [11.3]). Assume $t_{0}$ as in d). Give meaning to this: $\gamma$ has positive orientation relative to $U_{1}$. Hint: Interpret $\boldsymbol{u}_{1}$ being parallel to $U_{1}$.
We explore $d\left(w_{2}, w_{3}\right)$ from [9.16], to prove the triangle inequality and to find its differential distance tensor. Use $U(z)=\frac{1+|z|}{1-|z|}$.
(9.18a) Use [9.16e] and find $\beta(w) \in \mathrm{PGL}_{2}(\Delta)$ with $\beta\left(w_{2}\right)=0, \beta\left(w_{1}\right)=a>0$ to reduce $d\left(w_{1}, w_{3}\right) \leq d\left(w_{1}, w_{2}\right)+d\left(w_{2}, w_{3}\right), w_{1}, w_{2}, w_{3} \in \Delta$ to showing $U\left(\frac{z-a}{1-a z}\right) \leq U(a) \cdot U(z)$ with $a \in[0,1)$ and $z \in \Delta$.
(9.18b) Write $z=b e^{i \theta}$. Show $U\left(\frac{z-a}{1-a z}\right)$ is maximum in $\theta$ when $z$ is real. Conclude the inequality of a). Hint: $U(w)$ is increasing in $|w|$ and $\frac{z-a}{1-a z}$ maps the circle of radius $b$ on a circle with real center.
(9.18c) Use [9.16e] to compute the differential distance $S(x, y, d x, d y)$ by considering $w_{1}=x+i y$ close to $w_{2}$. Show $S(x, y, d x, d y)$ to be $\left|\frac{d x+i d y}{1-\left(x^{2}+y^{2}\right)}\right|$.
(9.18d) Apply $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$ mapping the upper half plane $\mathbb{H}$ to $\Delta$. Define a distance on $\mathbb{H}$ by pulling back two points and using the value of the distance on $\Delta$. Show this depend on the particular choice of $\alpha$. Show geodesics on $\mathbb{H}$ are half-circles perpendicular to the real axis.
(9.18e) Use d) to show the metric on $\mathbb{H}$ has differential distance element $\frac{|d x+i d y|}{y}$.

Consider [9.18] from the differential distance tensor view:

$$
F_{\Delta}=\left|\frac{d x+i d y}{1-\left(x^{2}+y^{2}\right)}\right|=h(x, y) \sqrt{d x^{2}+d y^{2}}
$$

with $h(x, y)=\left|1-\left(x^{2}+y^{2}\right)\right|^{-1 / 2}$. Recover this metric's geodesics, circles perpendicular to the boundary of $\Delta$, by applying the Euler-Lagrange variational principle from f). Consider $F^{2}=\boldsymbol{y} \cdot Q(\boldsymbol{x})(\boldsymbol{y})$ in (2.3a): $Q(\boldsymbol{x})$ is an $n \times n$ positive definite symmetric matrix. Tensor notation replaces $\boldsymbol{y}$ by $d x_{1}, \ldots, d x_{n}$. Classically, $F^{2}=\sum_{1 \leq i, j \leq n} q_{i, j}(\boldsymbol{x}) d x_{i} \otimes d x_{j}$ (with $q_{i, j}=q_{j, i}$ ) for a 2-tensor.
(9.19a) Suppose $\gamma$ and $\lambda$ are a pair of paths with $\gamma\left(t_{0}\right)=\lambda_{2}\left(t_{0}\right)=\boldsymbol{x}_{0}$. Define:

$$
F^{2}\left(\frac{d \gamma}{d t}\left(t_{0}\right), \frac{d \lambda}{d t}\left(t_{0}\right)\right)=\sum_{i, j} q_{i, j}\left(\boldsymbol{x}_{0}\right) \frac{d \gamma_{i}}{d t} \frac{d \lambda_{j}}{d t}
$$

Show $\frac{F^{2}\left(\frac{d \gamma}{d t}\left(t_{0}\right), \frac{d \lambda}{d t}\left(t_{0}\right)\right)}{F\left(\frac{d \gamma}{d t}\left(t_{0}\right), \frac{d \gamma}{d t}\left(t_{0}\right)\right) F\left(\frac{d \lambda}{d t}\left(t_{0}\right), \frac{d \lambda}{d t}\left(t_{0}\right)\right)}$ has absolute value at most 1 . So, it has the form $\cos (\theta(\gamma, \lambda))$. Show $\theta(\gamma, \lambda)$, the angle between $\gamma$ and $\lambda$ at $\boldsymbol{x}_{0}$, is independent of their parametrizations.
(9.19b) Apply Ex. 9.1. Show $\sum_{i, j} \int_{a}^{b} \sqrt{q_{i, j}(\gamma(t)) \frac{d \gamma_{i}}{d t} \frac{d \gamma_{j}}{d t}} d t$ is independent of how we parametrize the range of $\gamma$ assuming $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$ is one-one.
(9.19c) Let $H: \mathbb{R}^{2} \rightarrow \mathbb{R}^{n}$ by $\left(u_{1}, u_{2}\right) \mapsto\left(h_{1}\left(u_{1}, u_{2}\right), \ldots, h_{n}\left(u_{1}, u_{2}\right)\right)=\boldsymbol{h}(\boldsymbol{u})$ be a one-one (differentiable) map. Define $H^{*}\left(F^{2}\right)$, pullback of $F^{2}$ on the range of $H$, as $\sum_{1 \leq i, j \leq n} q_{i, j}(\boldsymbol{h}(\boldsymbol{u})) d h_{i} \otimes d h_{j}: d h_{i}(\boldsymbol{u})=\frac{\partial h_{i}}{\partial u_{1}} d u_{1}+\frac{\partial h_{i}}{\partial u_{2}} d u_{2}$. Suppose $\gamma:[a, b] \rightarrow H\left(\mathbb{R}^{2}\right)$. Show $\int_{\gamma} F=\int_{H^{-1} \circ \gamma} \sqrt{H^{*}\left(F^{2}\right)}$ from b).
(9.19d) Consider $H^{*}\left(F^{2}\right)$ in c) when $n=2$. Call $H$ isothermal coordinates if $H^{*}\left(F^{2}\right)$ is $h\left(u_{1}, u_{2}\right)\left(d u_{1} \otimes d u_{1}+d u_{2} \otimes d u_{2}\right)$. Use $n=2$ to factor $F^{2}$ to

$$
\left(A(\boldsymbol{x}) d x_{1}+B(\boldsymbol{x}) d x_{2}\right) \otimes\left(A(\boldsymbol{x}) d x_{1}+\bar{B}(\boldsymbol{x}) d x_{2}\right)
$$

$(\bar{B}(\boldsymbol{x})$ is the complex conjugation of $B(\boldsymbol{x})$ ). Suppose $k(\boldsymbol{x})$ (complex valued) gives $k(\boldsymbol{x})\left(A(\boldsymbol{x}) d x_{1}+B(\boldsymbol{x}) d x_{2}\right)$ with the form $d u_{1}+i d u_{2}$. Show $\left(u_{1}(\boldsymbol{x}), u_{2}(\boldsymbol{x})\right)$ gives isothermal coordinates.
(9.19e) Produce $k(\boldsymbol{x})$ near any $\left(x_{1}^{0}, x_{2}^{0}\right)$, as in c). Outline: Take real and imaginary parts. Rewrite: $d u_{i}=\frac{\partial u_{i}}{\partial x_{1}} d x_{1}+\frac{\partial u_{i}}{\partial x_{2}} d x_{2}$. Finding $k$ comes to this. Suppose $M_{1}(\boldsymbol{x}), M_{2}(\boldsymbol{x})$ are real valued and differentiable. Then, there is $k_{1}(\boldsymbol{x})$ and $M^{*}(\boldsymbol{x})$ with $k_{1}\left(M_{1}(\boldsymbol{x}) d x_{1}+M_{2}(\boldsymbol{x}) d x_{2}\right)$ of form $d M^{*}(\boldsymbol{x})$. Then, $M_{1}(\boldsymbol{x}) d x_{1}+M_{2}(\boldsymbol{x}) d x_{2}=0$ defines $\left\{\left(x_{1}, x_{2} \mid M^{*}\left(x_{1}, x_{2}\right)=0\right\}\right.$, an implicit surface, near $\left(x_{1}^{0}, x_{2}^{0}\right)$. Find $k_{1}$.
(9.19f) We assume the situation of [9.18]. Let $\gamma=\gamma_{1}+i \gamma_{2}:[0,1] \rightarrow \Delta$ be a path from $z_{0}$ to $z_{0}^{\prime}$. Minimize $\int_{\gamma} F_{\Delta}=\int_{0}^{1} S\left(\gamma_{1}(t), \gamma_{2}(t), \frac{d \gamma_{1}}{d t}, \frac{d \gamma_{2}}{d t}\right) d t$ over all such $\gamma$. The Euler-Lagrange variation produces two partial differential equations, one for $x, \frac{d}{d t} \frac{\partial S}{\partial \dot{x}}=\frac{\partial S}{\partial x}$, and a similar one for $y$. Solve to show $F_{\Delta}$ geodesics are circles perpendicular to the boundary of $\Delta$.

## CHAPTER 3

## COMPLEX MANIFOLDS AND COVERS

Chap. 4 replaces the field $\mathbb{C}(z, f(z))$ generated by an algebraic function $f(z)$ over $\mathbb{C}(z)$ by a geometric object, a 1-dimensional complex manifold (Riemann surface) that maps to the Riemann sphere $\mathbb{P}_{z}^{1}$. To prepare for this idea requires building some manifolds, and developing intuition for basic examples. We use fundamental groups to create new 1-dimensional complex manifolds from the space $U_{\boldsymbol{z}}$ with $\boldsymbol{z}$ a finite subset of $\mathbb{P}_{z}^{1}$.

Chap. 5 collects various Riemann surfaces into families. The parameter spaces for these families - one point in the space for each member of the family - are manifolds called moduli spaces. Chap. 4 has a prelude, the moduli space classically called the $j$-line: $\mathbb{P}_{j}^{1} \backslash\{\infty\}$. We use it for more general families than do classical texts on Riemann surfaces. Our moduli spaces may have arbitrarily high complex dimension. Still, their construction uses covering spaces (coming from fundamental groups) of open subsets of projective spaces. This chapter builds an intuition for using group theory to construct these spaces.

## 1. Fiber products and relative topologies

There is so much topology and we have so little space for it despite the need for some special constructions. The treatment is expedient and not completely classical to emphasize some subtle properties of manifolds.
1.1. Set theory constructions. For $X$ and $Y$ sets, the cartesian product of $X$ and $Y$ is the set

$$
X \times Y=\{(x, y) \mid x \in X, y \in Y\}
$$

Let $\left\{X_{\alpha}\right\}_{\alpha \in I}$ be a collection of subsets of the set $X$ indexed by the set $I$. The union of $\left\{X_{\alpha}\right\}_{\alpha \in I}$ is the set of $x \in X$ for which $x \in X_{\alpha}$ for some $\alpha \in I$. Denote this $\bigcup_{\alpha \in I} X_{\alpha}$. The complement of $X_{\alpha}$ in $X, X \backslash X_{\alpha}$, is $\left\{x \in X \mid x \notin \bigcup_{\alpha \in I} X_{\alpha}\right\}$. The intersection of $\left\{X_{\alpha}\right\}_{\alpha \in I}$ is the set of $x \in X$ with $x \in X_{\alpha}$ for each $\alpha \in I$. Denote this $\bigcap_{\alpha \in I} X_{\alpha}$.

Definition 1.1. For $X_{1}$ and $X_{2}$ sets, $Y_{i} \subset X_{i}, i=1,2$, let $f: Y_{1} \rightarrow Y_{2}$ be a one-one onto function. The sum of $X_{1}$ and $X_{2}$ along $f$ is the disjoint union of $X_{1} \backslash Y_{1}, Y_{2}$, and $X_{2} \backslash Y_{2}$. Denote this $X_{1} \bigcup_{f} X_{2}$. Along with this, we have maps $f_{i}: X_{i} \rightarrow X_{1} \bigcup_{f} X_{2}, i=1,2$ : with $f_{2}\left(x_{2}\right)=x_{2}$ for $x_{2} \in X_{2}, f_{1}\left(x_{1}\right)=x_{1}$ if $x_{1} \in X_{1} \backslash Y_{1}$, and $f_{1}\left(x_{1}\right)=f\left(x_{1}\right)$ for $x_{1} \in Y_{1}$. Call $f_{1}$ and $f_{2}$ the canonical maps.

Example 1.2 (The set behind a non-Hausdorff space). Consider

$$
X_{i}=\left\{(t, i) \in \mathbb{R}^{2} \mid-1<t<1\right\}, i=1,2, \text { with }
$$

$Y_{i}=X_{i} \backslash\{(0, i)\}, i=1,2$, and $f: Y_{1} \rightarrow Y_{2}$ by $f(t, 1)=(t, 2)$ for $(t, 1) \in Y_{1}$. Then, $X_{1} \bigcup_{f} X_{2}$ is the disjoint union of $X_{2}$ and the point $(0,1)$ (see Def. 1.4 and Ex. 2.4).

Definition 1.3 (Set theoretic fiber products). Let $f_{i}: X_{i} \rightarrow Z$ be two functions with range $Z, i=1,2$. The fiber product $X_{1} \times{ }_{Z} X_{2}$ consists of

$$
\left\{\left(x_{1}, x_{2}\right) \in X_{1} \times X_{2} \mid f_{1}\left(x_{1}\right)=f_{2}\left(x_{2}\right)\right\}
$$

Denote the natural map back to $Z$ by $f_{1} \times_{Z} f_{2}$. Suppose $X_{i} \subset Z$ and $f_{i}: X_{i} \rightarrow Z$ is inclusion, $i=1,2$. Then, identify $X_{1} \times{ }_{Z} X_{2}$ with $X_{1} \cap X_{2}$.

Suppose $X_{1}=X_{2}=Z=\mathbb{C}$, and $f_{1}$ and $f_{2}$ are polynomials. Then, $X_{1} \times{ }_{Z} X_{2}$ is the subset of $\left(x_{1}, x_{2}\right) \in \mathbb{C}^{2}$ defined by $f_{1}\left(x_{1}\right)=f_{2}\left(x_{2}\right)$. Define the $i$ th projection map, $\operatorname{pr}_{i}: X_{1} \times{ }_{Z} X_{2} \rightarrow X_{i}$ by $\operatorname{pr}_{i}\left(x_{1}, x_{2}\right) \mapsto x_{i}, i=1,2$.
The fiber product is an implicit set: an equation describes it.
The ball of radius $r$ about $\boldsymbol{x}_{0} \in \mathbb{R}^{n}$ is the basic open set $\left\{\boldsymbol{x} \in \mathbb{R}^{n}| | \boldsymbol{x}-\boldsymbol{x}_{0} \mid<r\right\}$. When necessary denote this $B\left(\boldsymbol{x}_{0}, r\right)$. Open sets of $\mathbb{R}^{n}$ are either empty or are (arbitrary) unions of basic open sets. Closed sets are complements (in $\mathbb{R}^{n}$ ) of open sets. Bounded sets are those contained in some basic open set. The collection of open sets, $\mathcal{U}$, in $\mathbb{R}^{n}$ therefore satisfies the axioms for a topology: $\mathcal{U}$ contains the empty set and the whole space, and it is closed under taking arbitrary unions and finite intersections.

Definition 1.4 (Relative topology I). Let $X$ be a subset of $\mathbb{R}^{n}$. Denote the collection of sets $X \cap U$ for $U$ open subset in $\mathbb{R}^{n}$ by $\mathcal{U}_{X}$. Then $\mathcal{U}_{X}$ gives the relative topology on $X$. For $x_{1}, x_{2} \in X$, two distinct points, $B\left(x_{1}, r / 3\right) \cap X$ and $B\left(x_{2}, r / 3\right) \cap X$ are disjoint open neighborhoods of the respective points $x_{1}$ and $x_{2}$ if $r=\left|x_{1}-x_{2}\right|$. Thus, in this relative topology, $X$ is a Hausdorff space.

Suppose $X$ (resp. $Y$ ) is a topological space with open sets $\mathcal{U}_{X}$ (resp., $\mathcal{U}_{Y}$ ). Let $f: X \rightarrow Y$ be a function with domain a subset of $X$. Then $f$ is continuous (for the relative topology) if for each $U \in \mathcal{U}_{Y}$,

$$
f^{-1}(U)=\{x \text { in the domain of } f \mid f(x) \in U\} \text { is in } \mathcal{U}_{X}
$$

For $U$ open in $Y$, denote restriction of $f$ to $f^{-1}(U)$ by $f_{U}: f^{-1}(U) \rightarrow U$. If $f$ is continuous, so is $f_{U}$.

The concept of relative topology generalizes to data $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$ on a set $X$ with the following properties: $\bigcup_{\alpha \in I} X_{\alpha}=X ; \varphi_{\alpha}: X_{\alpha} \rightarrow \mathbb{R}^{n}$ is a one-one map into $\mathbb{R}^{n}$; and $\varphi_{\beta} \circ \varphi_{\alpha}^{-1}: \varphi_{\alpha}\left(X_{\alpha} \cap X_{\beta}\right) \rightarrow \varphi_{\beta}\left(X_{\alpha} \cap X_{\beta}\right)$ is a continuous function for each $\alpha, \beta \in I$. We call the functions $\left\{\varphi_{\beta} \circ \varphi_{\alpha}^{-1}\right\}_{\alpha, \beta \in I}$ transition functions.

Definition 1.5 (Relative topology II). Let $X$ and $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$ be as above. Consider subsets of $X$ that are unions of $\varphi_{\alpha}^{-1}(U)$ with $U$ running over open sets of $\varphi_{\alpha}\left(X_{\alpha}\right), \alpha \in I$. Denote this collection of sets by $\mathcal{U}_{X}$. The topology on $X$ from $\mathcal{U}_{X}$ is the relative topology on $X$ induced from the topologizing data $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$. For $x \in X$ and $U$ an open set containing $x, U$ is a neighborhood of $x$.
1.2. Extending topologies from $\mathbb{R}^{n}$. Two sets of topologizing data on $X$, $\left\{\left(X_{\alpha^{\prime}}^{\prime}, \varphi_{\alpha^{\prime}}^{\prime}\right)\right\}_{\alpha^{\prime} \in I^{\prime}}$ and $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$, are equivalent (the same, or give the same topology) if each defines the same open sets on $X$.

Consider $X$ and $Y$, topological spaces with respective data $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$ and $\left\{\left(Y_{\beta}, \psi_{\beta}\right)\right\}_{\beta \in J}$. A one-one map $f: X \rightarrow Y$ is a (topological) embedding if the topologizing data from $\left\{\left(f^{-1}\left(Y_{\beta}\right), \psi_{\beta} \circ f\right)\right\}_{\beta \in J}$ is equivalent to $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$. Note: Ex. 2.4 has a space with no embedding in $\mathbb{R}^{n}$ (for any $n$ ). It isn't Hausdorff. Yet, each point has a neighborhood embeddable as an open interval in $\mathbb{R}^{1}$.

Associate to each subset $Y$ of a topological space $X$ the closure $\bar{Y}$ of $Y$ in $X$ : $\bar{Y}$ (a closed set) is the points $x \in X$ with each neighborhood of $x$ containing at
least one point of $Y$. If each neighborhood of $x$ contains a point of $Y$ distinct from $x$, then $x$ is a limit point of $Y$.

Compact subsets of $\mathbb{R}^{n}$ are those both closed and bounded. The Heine-Borel covering theorem [Rud76, p. 40] characterizes these sets through the concept of an open covering. A collection $\mathcal{U}=\left\{U_{\alpha}\right\}_{\alpha \in I}$ of open subsets of $\mathbb{R}^{n}$ is an open cover of $Y$ if $Y \subseteq \bigcup_{\alpha \in I} U_{\alpha}$. Then $Y$ has the finite covering property if for each open cover $\mathcal{U}$ there is a finite collection $\left\{U_{\alpha_{i}}\right\}_{i=1}^{t}, \alpha_{1}, \ldots, \alpha_{t} \in I$, covering $Y$.

Theorem 1.6 (Heine-Borel). The finite covering property is equivalent to compactness for subsets of $\mathbb{R}^{n}$.

Thus, for any topological space $X$, without reference to the concept of bounded set, one says a subset $Y$ is compact if it has the finite covering property.

A subset $Y$ of a topological space $X$ is disconnected if there are two nonempty open sets $U_{1}$ and $U_{2}$ of $Y$ (in the relative topology) with $U_{1} \cap U_{2}$ empty and $U_{1} \bigcup U_{2}=Y$. If $Y$ is not disconnected call it connected (in $X$ ). For any $x \in X$, there is a maximal connected set $U_{x}$ containing $x$. So, each topological space decomposes into a union of disjoint connected components. If $f: Y \rightarrow X$ is continuous, the image of any connected subset of $Y$ is a connected subset of $X$.

## 2. Functions on $X$ from functions on $\mathbb{R}^{n}$

There are several points to make about Def. 1.5. First it includes many topologies as our next example illustrates.

Example 2.1. Let $X$ be any set whose points, $x_{\alpha}$, are indexed by $\alpha \in I$. Let $X_{\alpha}=\left\{x_{\alpha}\right\}$ and $\varphi_{\alpha}:\left\{x_{\alpha}\right\} \rightarrow\{\mathbf{0}\}, \alpha \in I$, where $\mathbf{0}$ is the origin of $\mathbb{R}^{n}$. The relative topology on $X$ is the discrete topology.

By using another target space $Y$ with a well-known topology on it (like the $p$-adic numbers $\mathbb{Z}_{p}$, replacing $\mathbb{R}^{n}$ ), we could include $p$-adic topologies, too. Still, it does not include all the topologies significant to modern mathematics even for spaces we consider as manifolds. Later we will extend it to Grothendieck topologies. It is appropriate for that example to notice we don't need a topology on $X$ to start the process (§2.1).

Further, the point of topologizing data is to pull back functions (differentials, and other objects) from $\mathbb{R}^{n}$ so $X$ has local functions (differentials, etc.) just like those of $\mathbb{R}^{n}$. Since $\mathbb{R}^{n}$ also has the notion of real analytic, differentiable and harmonic functions, transition functions also allow us to pull those back, to identify such functions on $X$. For these definitions, however, to be meaningful, they must be locally independent of which function we use for pullback. This requires the transition functions also have these respective properties (§3).

When $n=2 m$ is even, suppose the following two conditions hold.
(2.1a) We have chosen a fixed $\mathbb{R}$ linear map $L=L_{n}: \mathbb{R}^{n} \rightarrow \mathbb{C}^{m}$.
(2.1b) Using $L$, the transition functions are analytic from $\mathbb{C}^{m} \rightarrow \mathbb{C}^{m}$.

These conditions allow identifying a set of functions in a neighborhood of any point on $X$ as analytic (§3.1.2).

Finally, there is a warning. Local function theory immediately challenges us to identify global functions and differentials on $X$ through their local definitions. There is an immediate first problem to assure a simple property we expect from functions in $\mathbb{R}^{n}$. If a function $f$ in a neighborhood of $x \in X$ has good behaviour
as $x^{\prime} \in X$ approaches $x$, then it should have a unique limit value (see $\S 2.2$ on the Hausdorff property).
2.1. Defining a topological space from its atlas. Def. 1.5 shows we don't need $X$ to start with a topology. It inherits one from its topologizing data. So, it is reasonable to ask if we need an a priori space $X$ at all.
2.1.1. Equivalence relations define topological spaces. For example, suppose $\left\{U_{\alpha}\right\}_{\alpha \in I}$ is a collection of open sets in $\mathbb{R}^{n}$, and for some subset $(\beta, \alpha) \in I \times I$, there are invertible continuous maps $\psi_{\beta, \alpha}: V_{\beta}^{\alpha} \rightarrow V_{\alpha}^{\beta}$, with $V_{\alpha}^{\beta}$ open in $U_{\alpha}$ (resp. $V_{\beta}^{\alpha}$ open in $U_{\alpha}$ ). Can we form an $X$ so that $\left\{\psi_{\beta, \alpha}\right\}_{\alpha, \beta \in I}$ are the transition functions for its topological structure? Almost!

Let $X$ be the disjoint union $\dot{U}_{\alpha \in I} U_{\alpha}$ modulo the relation $R_{I}$ on this union defined by $x \in U_{\alpha} \sim x^{\prime} \in U_{\beta}$ if $\psi_{\beta, \alpha}(x)=x^{\prime}$. If $R_{I}$ is an equivalence relation, then the equivalence classes form a set $X$ and on it a topological structure. On this space, of course, the open sets do look like those of $\mathbb{R}^{n}$ (in contrast to Ex. 2.1). The following lemma keeps track of the definitions.

Lemma 2.2. The relation $R_{I}$ is an equivalence relation if and only if the following properties hold:
(2.2a) $\psi_{\alpha, \alpha}$ is the identity map; $\psi_{\alpha, \beta}=\psi_{\beta, \alpha}^{-1}$; and
(2.2b) $\psi_{\gamma, \beta} \circ \psi_{\beta, \alpha}=\psi_{\gamma, \alpha}$ wherever any two of the maps are defined.

Suppose $R_{I}$ is an equivalence relation. Then the inverse of the natural inclusion maps $U_{\alpha} \rightarrow X$ are functions $\varphi_{\alpha}$ giving transition functions $\varphi_{\beta} \circ \varphi_{\alpha}^{-1}=\psi_{\beta, \alpha}$.
2.1.2. Quotient topologies. Suppose $X$ is a topological space with topologizing data $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$. Let $f: X \rightarrow Y$ be any surjective map. Then, there is a topology on $Y$ with open sets $\mathcal{U}_{Y}$ the images by $f$ of all sets in $\mathcal{U}_{X}$. We can't, however, expect topologizing data on $Y$ by pushing down the functions $\varphi_{\alpha}$ without extra conditions. It usually makes sense to write $f$ for restriction of $f$ to any subset $V \subset X$. The argument here, however, requires tracking the domain, and so we write $f_{V}$.

Let $J$ be the subset of $I$ for which $f_{X_{\beta}}: X_{\beta} \rightarrow Y$ is one-one for $\beta \in J$. Let $\mathcal{U}_{X, Y}$ be $\left\{X_{\beta}\right\}_{\beta \in J}$ and assume $\mathcal{U}_{X, Y}$ is a cover of $X$. With no loss assume the coordinate chart for $X$ contains only sets from $\mathcal{U}_{X, Y}$. The hypothesis provides coordinate functions $\psi_{\alpha}: f\left(X_{\alpha}\right) \rightarrow \mathbb{R}^{n}$ by setting $\psi_{\alpha}=\varphi_{\alpha} \circ f_{X_{\alpha}}^{-1}$ on $f\left(X_{\alpha}\right)$.

From Lem. 2.2 we want an equivalence relation on $\dot{U}_{\alpha \in J} \psi_{\alpha}\left(f\left(X_{\alpha}\right)\right)$ that reproduces the set $Y$ as equivalence classes: $y \in \psi_{\alpha}\left(f\left(X_{\alpha}\right)\right) \sim y^{\prime} \in \psi_{\beta}\left(f\left(X_{\beta}\right)\right)$ if $\psi_{\beta, \alpha}(y)=y^{\prime}$. So, the problem is to define $\psi_{\beta, \alpha}$, using that $f_{X_{\alpha}}^{-1}$ is different from $f_{X_{\beta}}^{-1}$ on $f\left(X_{\alpha}\right) \cap f\left(X_{\beta}\right)$. If $f\left(X_{\alpha} \cap X_{\beta}\right)=f\left(X_{\alpha}\right) \cap f\left(X_{\beta}\right)$, then it is consistent to define $\psi_{\beta, \alpha}$ as $\psi_{\beta} \circ \psi_{\alpha}^{-1}=\varphi_{\beta} \circ \varphi_{\alpha}^{-1}$. More generally, an additional hypothesis is essentially necessary and sufficient if we use the full set $\mathcal{U}_{X, Y}$.

Lemma 2.3. Suppose in addition to the above, for each pair $X_{\alpha}, X_{\beta} \in \mathcal{U}_{X, Y}$ with $f\left(X_{\alpha}\right) \cap f\left(X_{\beta}\right) \neq \emptyset$, there exists $X_{\beta^{\prime}} \in \mathcal{U}_{X, Y}$ with
(2.3) $f\left(X_{\beta^{\prime}}\right)=f\left(X_{\beta}\right)$ and $f\left(X_{\alpha} \cap X_{\beta^{\prime}}\right)=f\left(X_{\alpha}\right) \cap f\left(X_{\beta^{\prime}}\right)$.

Then, the topologizing data on $X$ provides topologizing data on $Y$.
Proof. Apply $f$ to $\mathcal{U}_{X, Y}$ get $\mathcal{U}_{Y}$. Suppose $f\left(X_{\alpha}\right) \cap f\left(X_{\beta}\right) \neq \emptyset$. Then, choose $\left(X_{\beta^{\prime}}, \varphi_{\beta^{\prime}}\right)$ and form $\psi_{\beta, \alpha}$ by replacing $f_{\beta}^{-1}$ by $f_{\beta^{\prime}}^{-1}$.
2.2. $\mathbb{R}^{n}$-like behavior requires Hausdorffness. Here is the problem with a space that isn't Hausdorff. Suppose $f:[0,1) \rightarrow X$ is a continuous function, everything of a path except the end point. Manifolds in this book appear as extensions of open subsets of $\mathbb{R}^{n}$. So, the only thing that should prevent us from extending our path (continuously) to $f^{*}:[0,1] \rightarrow X$ is that there is no point $f^{*}(1) \in X$ giving a continuous $f^{*}$. If there are several possible choices $f^{*}(1)$ giving a continuous function $f^{*}$, these extending points would have more exotic neighborhoods than do points in $\mathbb{R}^{n}$. In practice, the use of Hausdorff is to assure in theorems of Chap. Chap. 4 that there is a unique manifold solution to many existence problems.

Example 2.4 (Continuation of Ex. 1.2). As in Ex. 1.2, let $\varphi_{i}: X_{i} \rightarrow \mathbb{R}^{1}$ by $\varphi_{i}(t, i)=t, i=1,2$. The relative topology on $X_{1} \bigcup_{f} X_{2}$ is not Hausdorff [9.1].

Figure 1. An undecided function.


There is a topological formulation of the possibility that we could end a path in two different points. That is, $(f, f):[0,1) \rightarrow X \times X$ has topological closure not in the diagonal $\left.\Delta_{X}=\{(x, x) \mid X] \times X\right\}$. That is, if $f^{*}(1)$ and $f^{\dagger}(1)$ are two different ways to extend $f$ to a path on $[0,1]$, then $\left(f^{*}(1), f^{\dagger}(1)\right)$ is in the closure of $\Delta_{X}$. Conveniently, the exact property that prevents this situation is that $X$ is Hausdorff [9.1b].

Lemma 2.5. $X$ is Hausdorff if and only if $\Delta_{X}$ is closed in $X \times X$ [9.1d].
Here is a classical fact. If $f: X \rightarrow Y$ is continuous and one-one and $Y$ is Hausdorff, then the restriction of $f$ to any compact subset of $X$ is a homeomorphism onto its image. This uses that the image of a compact set is compact; then Hausdorff assures that the image of the compact set (and all closed subsets of it) is closed. It is, however, common to have such an $f$ where the inverse image of some compact sets are not compact. For example, let $f: \mathbb{C}_{z}^{*} \rightarrow \mathbb{C}_{z}$ be the identity map. Then, the inverse image of the unit disk is not compact (compare with [9.1e]). Call a map $f: X \rightarrow Y$ proper if the inverse image of compact sets is compact.

## 3. Manifolds: differentiable and complex

Let $X$ be a topological space with topologizing data $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$ (relative to $\left.\mathbb{R}^{n}\right)$. We add conditions to define differentiable and complex manifolds. Classical cases of the latter include the Riemann sphere, the complex torus and algebraic sets defined by $m \in \mathbb{C}[z, w]$ with nonzero gradient everywhere.

Definition 3.1. Let $X$ be a Hausdorff space with $\left\{\left(X_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$ as topologizing data. Assume $\varphi_{\alpha}$ maps $U_{\alpha}$ to an open connected subset of $\mathbb{R}^{n}$ for each $\alpha \in I$. Call $X$ an $n$-dimensional (topological) manifold.

In this case, replace the open sets $X_{\alpha}$ by the notation $U_{\alpha}$. Call $\left\{\left(U_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$ a coordinate system or atlas. An individual member $\varphi_{\alpha}: U_{\alpha} \rightarrow \mathbb{R}^{n}$ is a (coordinate) chart. Ex. 2.4 shows the Hausdorff condition isn't automatic.
3.1. Manifold structures. Let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}$ be a continuous function defined on an open set $U$. For $\boldsymbol{x}_{0} \in U$ and $\boldsymbol{v} \in \mathbb{R}^{n}$, the directional derivative of $f$ at $\boldsymbol{x}_{0}$ in the direction $\boldsymbol{v}$ is the limit

$$
\lim _{t \rightarrow 0} \frac{\left.f\left(\boldsymbol{x}_{0}+t \boldsymbol{v}\right)-f\left(\boldsymbol{x}_{0}\right)\right)}{t} \stackrel{\text { def }}{=} \frac{\partial f}{\partial \boldsymbol{v}}\left(\boldsymbol{x}_{0}\right)
$$

if it exists. If $\boldsymbol{e}_{i}=\boldsymbol{v}$ is the vector with 1 in the $i$ th coordinate and 0 in the other coordinates, denote the directional derivative by $\frac{\partial f}{\partial x_{i}}\left(\boldsymbol{x}_{0}\right)$. Then

$$
\nabla f\left(\boldsymbol{x}_{0}\right) \stackrel{\text { def }}{=}\left(\frac{\partial f}{\partial x_{1}}\left(\boldsymbol{x}_{0}\right), \ldots, \frac{\partial f}{\partial x_{n}}\left(\boldsymbol{x}_{0}\right)\right)
$$

is the gradient of $f$ at $\boldsymbol{x}_{0}$.
Lemma 3.2. [Rud76, p. 218] Suppose $\frac{\partial f}{\partial x_{i}}$ exists and is continuous near $\boldsymbol{x}_{0}$ for $i=1, \ldots, n$. Then, for each vector $\boldsymbol{v}, \frac{\partial f}{\partial \boldsymbol{v}}\left(\boldsymbol{x}_{0}\right)$ exists and equals $\nabla f\left(\boldsymbol{x}_{0}\right) \cdot \boldsymbol{v}$.

Call a function satisfying the hypotheses of Lemma 3.2 differentiable at $\boldsymbol{x}_{0}$. A function $\boldsymbol{f}=\left(f_{1}(\boldsymbol{x}), \ldots, f_{m}(\boldsymbol{x})\right)$ from $\mathbb{R}^{n}$ to $\mathbb{R}^{m}$ is differentiable at $\boldsymbol{x}_{0}$ if each of the coordinate functions $f_{i}(\boldsymbol{x})$ is differentiable at $\boldsymbol{x}_{0}$. While it is not absolutely necessary, our manifolds often have transition functions with continuous partial derivatives of all orders: smoothly differentiable.

Assume $g: \mathbb{R}^{m} \rightarrow \mathbb{R}$ is a composite of $\mathbb{R}^{m} \xrightarrow{\dot{H}} \mathbb{R}^{n} \xrightarrow{f} \mathbb{R}$. Let $\boldsymbol{y}_{0} \in \mathbb{R}^{m}$. Suppose each coordinate function from $H(\boldsymbol{y})=\left(h_{1}(\boldsymbol{y}), \ldots, h_{n}(\boldsymbol{y})\right)$ of $H$ is differentiable at $\boldsymbol{y}_{0}$ and $f$ is differentiable at $H\left(\boldsymbol{y}_{0}\right)$. Write $J(H)\left(\boldsymbol{y}_{0}\right)$ for the matrix whose $i$ th row is $\nabla h_{i}\left(\boldsymbol{y}_{0}\right)$. As a slight generalization of Lem. 3.2, $\nabla g\left(\boldsymbol{y}_{0}\right)$ exists and equals

$$
\begin{equation*}
=\nabla f\left(H\left(\boldsymbol{y}_{0}\right)\right) \cdot J(H)\left(\boldsymbol{y}_{0}\right) \tag{3.1}
\end{equation*}
$$

3.1.1. Differentiable functions. Let $X$ be an $n$-dimensional manifold. Denote an atlas for it by $\left\{\left(U_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$.

Definition 3.3. Call $X$ a differentiable manifold if each transition function $\varphi_{\beta} \circ \varphi_{\alpha}^{-1}$ is smoothly differentiable on its domain of definition.

For any $x \in U_{\alpha}$ on a chart of a differentiable manifold $X$, define the (smoothly) differentiable functions on $U_{\alpha}$ to be $C^{\infty}\left(U_{\alpha}\right)=\left\{f \circ \varphi_{\alpha} \mid f \in \mathrm{C}^{\infty}\left(\varphi_{\alpha}\left(U_{\alpha}\right)\right)\right\}$. This definition should be independent of the chart: We declare that restricting a differentiable function to an open subset of $U_{\alpha}$ still gives a differentiable function. This, however, must be compatible with the definition of differentiable using any other coordinate chart $\left(U_{\beta}, \varphi_{\beta}\right)$ which also contains $x$.

Lemma 3.4. Suppose $x \in U_{\alpha} \cap U_{\beta}$, and $f \circ \varphi_{\alpha}$ is restriction of a differentiable function to an open neighborhood $W$ of $x$ in $U_{\alpha} \cap U_{\beta}$. Then, $f \circ \varphi_{\alpha}=g \circ \varphi_{\beta}$ for some differentiable function $g$ defined on $\varphi_{\beta}(W)$.

Proof. Write $f \circ \varphi_{\alpha}$ as $f \circ \varphi_{\alpha} \circ \varphi_{\beta}^{-1} \circ \varphi_{\beta}$ and take $g$ as $f \circ \varphi_{\alpha} \circ \varphi_{\beta}^{-1}$. This is defined on $\varphi_{\beta}(W)$. As the composite of two differentiable functions $f$ and $\varphi_{\alpha} \circ \varphi_{\beta}^{-1}$, $g$ is differentiable from (3.1).

Definition 3.5 (Global differentiable functions on $X$ ). If $X$ is a differentiable manifold, then a function $f: X \rightarrow \mathbb{R}$ is differentiable if its restriction to each $U_{\alpha}$ in a coordinate chart is differentiable.
3.1.2. Complex functions. Decompose a complex number $z_{i}$ into its real and complex parts as $x_{i}+i y_{i}$. This produces (as in (2.1)) a natural one-one map:

$$
L=L_{n}: \mathbb{R}^{2 n} \rightarrow \mathbb{C}^{n} \text { by }\left(x_{1}, y_{1}, \ldots, x_{n}, y_{n}\right) \mapsto\left(z_{1}, \ldots, z_{n}\right)
$$

Topologize $\mathbb{C}^{n}$ so $L$ (and its inverse) are continuous. Identify $\mathbb{C}^{n}$ and $\mathbb{R}^{2 n}$ to consider any differentiable function: $g: \mathbb{R}^{2 n} \rightarrow \mathbb{R}$ as a function $g \circ L^{-1}: \mathbb{C}^{n} \rightarrow \mathbb{R}$. Further, a pair $u$ and $v$ of differentiable functions with a common domain $U$ from $\mathbb{R}^{2 n} \rightarrow \mathbb{R}$ produces a differentiable function $f: \mathbb{C}^{n} \rightarrow \mathbb{C}$ on $U$ :

$$
\boldsymbol{z} \mapsto u \circ L^{-1}(\boldsymbol{z})+i v \circ L^{-1}(\boldsymbol{z}) .
$$

Call $f: \mathbb{C}^{n} \rightarrow \mathbb{C}$ analytic at $\boldsymbol{z}_{o}=\left(z_{1,0}, \ldots, z_{n, 0}\right)$ if each complex partial derivative

$$
\frac{\partial f}{\partial z_{i}}\left(\boldsymbol{z}^{\prime}\right)=\lim _{z_{i} \rightarrow z_{i}^{\prime}} \frac{\left(f\left(z_{1}^{\prime}, \ldots, z_{i-1}^{\prime}, z_{i}, z_{i+1}^{\prime}, \ldots, z_{n}^{\prime}\right)-f\left(\boldsymbol{z}^{\prime}\right)\right)}{z_{i}-z_{i}^{\prime}}
$$

exists and is continuous, $i=1, \ldots, n$, with $\boldsymbol{z}^{\prime}$ near $\boldsymbol{z}_{0}$. We say $\boldsymbol{f}=\left(f_{1}(\boldsymbol{z}), \ldots, f_{m}(\boldsymbol{z})\right)$ from $\mathbb{C}^{n}$ to $\mathbb{C}^{m}$ is analytic at $\boldsymbol{z}_{0}$ if each coordinate function $f_{i}(\boldsymbol{z})$ is analytic at $\boldsymbol{z}_{0}$. Analytic functions behave for differentiation (or integration) as if each $z_{i}$ ranging over a 2-dimensional set were a single real variable. [9.4] explores how changing the particular linear identification $L_{n}$ affects this definition. In the first half of the 1800's, researchers realized the geometry underlying this definition could characterize special recurring collections of integrals. A motivating problem (Chap. 4) was whether the integrals of these functions were serious new functions. By, however, defining — as in Def. 3.6 - analytic manifolds, Riemann replaced complicated sets of functions by geometric properties.

To match with previous notation, if $U$ be an open connected subset of $\mathbb{C}^{n}$, denote the analytic functions on $U$ by $\mathcal{H}(U)$. The natural quotient field $\mathcal{M}(U)$ of $\mathcal{H}(U)$ (Lem. 3.9), the field of meromorphic functions on $U$, consists of ratios from $\mathcal{H}(U)$ with nonzero denominators. When $n=1$, at each point of $U$ any meromorphic function takes a well-defined value in $\mathbb{P}_{z}^{1}$. Simple examples like $\frac{z_{1}}{z_{2}}$ at $(0,0)$ show this is not true for $n \geq 2$ [9.11e].

Definition 3.6. Let $X$ be a $2 n$-dimensional manifold with atlas $\left\{\left(U_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$ where $\varphi_{\alpha}: U_{\alpha} \rightarrow \mathbb{C}^{n}$. Call $X$ an analytic (or complex) $n$-dimensional manifold if each transition function $\psi_{\beta, \alpha}=\varphi_{\beta}^{\circ} \circ \varphi_{\alpha}^{-1}$ is analytic on $\varphi_{\alpha}\left(U_{\alpha} \cap U_{\beta}\right)$. So, an analytic manifold is differentiable. A Riemann surface is a 1-dimensional complex manifold.

For any $x \in U_{\alpha}$ on a chart $\mathcal{U}$ of an analytic manifold $X$, define analytic (resp. meromorphic) functions on $U_{\alpha}$ to be $\mathcal{H}_{\mathcal{U}}\left(U_{\alpha}\right)=\left\{f \circ \varphi_{\alpha} \mid f \in \mathcal{H}\left(\varphi_{\alpha}\left(U_{\alpha}\right)\right)\right\}$ (resp. $\mathcal{M}_{\mathcal{U}}\left(U_{\alpha}\right)$ where we replace $f$ analytic by $f$ meromorphic). Exactly as previously, Lem. 3.4 has a version for analytic or meromorphic functions. What changes if we adjust the atlas $\left\{\left(U_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$ in simple ways?

Definition 3.7. Assume $X=X_{\mathcal{U}}$ is an $n$-dimensional analytic manifold, and $h_{\alpha}: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ is one-one, differentiable, but not necessarily analytic, on $\varphi_{\alpha}\left(U_{\alpha}\right)$ for each $\alpha \in I$. Topologies of $X$ from $\left\{\left(U_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}=\mathcal{U}$ or $\left\{\left(U_{\alpha}, h_{\alpha} \circ \varphi_{\alpha}\right)\right\}_{\alpha \in I}=\mathcal{U}_{\boldsymbol{h}}$ are the same. Call $\boldsymbol{h}$ a coordinate adjustment and $\mathcal{U}_{\boldsymbol{h}}$ the adjustment of $\mathcal{U}$ by $\boldsymbol{h}$. Then, $\boldsymbol{h}$ is an analytic adjustment if transition functions for $\mathcal{U}_{\boldsymbol{h}}$ are analytic.

Only special coordinate adjustments are analytic. Even if $\boldsymbol{h}$ is an analytic adjustment, unless all the $h_{\alpha} \mathrm{s}$ are analytic themselves, the functions we call analytic (or meromorphic) on an open set $U_{\alpha}$ of $X_{\mathcal{U}}$ are usually different from those on the same open set of $X_{\mathcal{U}_{\boldsymbol{h}}}$. For example, suppose $I=\{\alpha\}$ and $U_{\alpha}=D$ is an open set in
$\mathbb{C}$. Then, the functions $\mathcal{H}_{(D, h)}(D)=\{f \circ h \mid f \in \mathcal{H}(D)$ we call analytic on $\{(D, h)\}$ are the same as $\mathcal{H}(D)$ if and only if $h$ is analyic.

If $D$ is simply connected (and not all of $\mathbb{C}_{z}$ ), then Riemann's Mapping Theorem says $\mathcal{H}(D, h)$ is isomorphic as a ring to the convergent power series on the unit disk in $\mathbb{C}_{z}$. [Ahl79, p. 230] says this if $h$ is the identity, though composing with $h^{-1}$ for any diffeomorphisms is a ring isomorphism. A nontrivial case of adjustments is where all the $h_{\alpha} \mathrm{s}$ are the same (see $[9.4 \mathrm{c}]$ ). We explore this further in Chap. $4 \S ? ?$. In the next observation (see $\S 5.2 .1$ for the definition of $\frac{\partial}{\partial \bar{z}}$ ) denote range variables for $\varphi_{\alpha}: U_{\alpha} \rightarrow \mathbb{C}^{n}$ by $z_{\alpha, 1}, \ldots, z_{\alpha, n}$.

Lemma 3.8. That $X_{\mathcal{U}_{h}}$ is an analytic manifold is equivalent to
(3.2) $h_{\beta} \circ \varphi_{\beta} \circ \varphi_{\alpha}^{-1} \circ h_{\alpha}^{-1}$ is analytic on $h_{\alpha} \circ \varphi_{\alpha}\left(U_{\alpha} \cap U_{\beta}\right)$ for all $(\alpha, \beta) \in I^{2}$ : $\frac{\partial}{\partial \bar{z}_{\alpha, i}}$ applied to each of its matrix entries is $0, i=1, \ldots, n$.
If the $\left\{h_{\alpha}\right\}_{\alpha \in I}$ are all analytic, then $\mathcal{H}_{\mathcal{U}}\left(U_{\alpha}\right)=\mathcal{H}_{\mathcal{U}_{\boldsymbol{h}}}\left(U_{\alpha}\right)$ for all $\alpha \in I$.
Suppose $X_{\mathcal{U}}$ and $X_{\mathcal{U}_{h}}$ are both analytic manifolds. Lem. 3.8 shows the local analytic functions change unless $\boldsymbol{h}$ consists of analytic functions. We regard the complex structures as the same if and only if both $X_{\mathcal{U}}$ and $X_{\mathcal{U}_{h}}$ have the same analytic functions in a neighborhood of each point. A special case appears often in the theory of complex manifolds. It is when all the functions $h_{\alpha}$ are complex conjugation (Chap. 4 Lem. ??). Notice: Complex conjugation reverses orientation in $\mathbb{C}$ by mapping clockwise paths around the origin to counterclockwise paths.
3.1.3. A tentative definition of algebraic manifold. For complex manifolds, a coordinate chart allows us to define global meromorphic functions as a collection $g_{\alpha} \in \mathcal{M}\left(U_{\alpha}\right)$ for which $g_{\alpha}=g_{\beta}$ on any points of $U_{\alpha} \cap U_{\beta}$ where both make sense. Our major study treats families of compact Riemann surfaces. Often each family member appears explicitly with a finite set of points removed, using Riemann's Existence Theorem to produce such surfaces as covers of $U_{\boldsymbol{z}}$. Meromorphic functions mean for us functions meromorphic on some compactification of this manifold. This includes that the functions are ratios of holomorphic functions at those points that might not be included in the initial presentation. For example, global meromorphic functions on $U_{z}$ refer to elements of $\mathbb{C}(z)$. They are among the ratios of algebraic functions on $U_{\boldsymbol{z}}$, so they have no essential singularities as we approach $\boldsymbol{z}$.

Understanding manifolds which have a coordinate description is important to the goals of this book. When we deal with compact complex manifolds, global coordinate functions live inside the field of global meromorphic functions. Our first tentative definition of algebraic excludes some manifolds that everyone considers algebraic. Still, it is simple, close to the general meaning of algebraic and it leads naturally to that definition.

Lemma 3.9. Suppose $X_{\mathcal{U}}$ is a connected topological space and an analytic manifold. Then, the (global) meromorphic functions on $X=X_{\mathcal{U}}$ form a field, $\mathbb{C}(X)$.

Proof. Add (resp. multiply) functions of form $f_{1}\left(\varphi_{\alpha}\right)$ and $f_{2}\left(\varphi_{\alpha}\right)$ by computing the value at $x \in U_{\alpha}$ as $f_{1}\left(\varphi_{\alpha}(x)\right)+f_{2}\left(\varphi_{\alpha}(x)\right)$ (resp. $f_{1}\left(\varphi_{\alpha}(x)\right) f_{2}\left(\varphi_{\alpha}(x)\right)$ ). Quotients, too, are obvious for they will also be ratios of holomorphic functions at each point. We need only to see that $\mathbb{C}(X)$ is an integral domain. If, however, $f_{1}\left(\varphi_{\alpha}(x)\right) f_{2}\left(\varphi_{\alpha}(x)\right)=0$ for $x \in U_{\alpha}$, then $f_{1}(z) f_{2}(z)=0$ for $z$ on the open set $\varphi_{\alpha}\left(U_{\alpha}\right)$. Chap. 2 [9.8a] shows either $f_{1}\left(\varphi_{\alpha}\right)$ or $f_{2}\left(\varphi_{\alpha}\right)$ is 0 on $U_{\alpha}$.

A goal for compact Riemann surfaces is to understand adjustments well enough to be able to list the isomorphism classes of fields $\mathbb{C}\left(X_{\mathcal{U}_{\boldsymbol{h}}}\right)$, the function field of
$X_{\mathcal{U}_{\boldsymbol{h}}}$, as $\boldsymbol{h}$ varies. How can we describe the complete set of function fields up to isomorphism? This book shows how to apply various answers to many seemingly unrelated problems.

Suppose $x_{1}, x_{2} \in X$ and $f \in \mathbb{C}\left(X_{\mathcal{U}}\right)$ are holomorphic in a neighborhood of $x_{1}$ and $x_{2}$ and takes different values there. We say $f$ separates $x_{1}, x_{2}$. If for each pair of distinct points $x_{1}, x_{2} \in X$ there is an $f \in \mathbb{C}\left(X_{\mathcal{U}}\right)$ separating them, we say $\mathbb{C}\left(X_{\mathcal{U}}\right)$ separates points. Suppose $X_{\mathcal{U}}$ has complex dimension $n, x \in X_{\mathcal{U}}$ is in a coordinate chart $\varphi_{\alpha}: U_{\alpha} \rightarrow \mathbb{C}^{n}$ and there are $n$ functions $f_{1}, \ldots, f_{n} \in \mathbb{C}\left(X_{\mathcal{U}}\right)$ all holomorphic in a neighborhood of $x$. If the Jacobian of $f_{1}, \ldots, f_{n}$ - determinant of the matrix with $(i, j)$-entry of $\frac{\partial f_{i} \circ \varphi_{\alpha}^{-1}}{\partial z_{j}}, i=1, \ldots, n, j=1, \ldots, n$ - is nonzero at $\varphi_{\alpha}(x)$, we say $f_{1}, \ldots, f_{n}$ separate tangents at $x$.

Definition 3.10. An $n$-dimensional compact complex manifold $X$ (with topologizing data $\mathcal{U})$ is $\mathbb{P}^{1}$-algebraic if there is a collection $f_{1}, \ldots, f_{N} \in \mathbb{C}\left(X_{\mathcal{U}}\right)$ so the following conditions hold.
(3.3a) For each $x \in X_{\mathcal{U}}$, there is a collection $\epsilon_{1}, \ldots, \epsilon_{N} \in\{ \pm 1\}$ (dependent on $x$ ) so that $f_{1}^{\epsilon_{1}}, \ldots, f_{N}^{\epsilon_{N}}$ are all holomorphic at $x$.
(3.3b) Among $f_{1}^{\epsilon_{1}}, \ldots, f_{N}^{\epsilon_{N}}$ there are $n$ that separate tangents at $x$.
(3.3c) Given distinct $x_{1}, x_{2} \in X$, one from $f_{1}, \ldots, f_{N}$ separates $x_{1}$ and $x_{2}$.

Note: In (3.3c), if $f_{i}$ is holomorphic at $x$, and $f_{i}(x)=0$, we include $\infty$ as the value of $1 / f_{i}(x)$. Algebraic manifolds are the analytic manifolds $X_{\mathcal{U}}$ most significant to us ( $\mathbb{P}^{1}$-algebraic manifolds are a special case; see $\S 4.1 .2$ ). There are 2-dimensional analytic manifolds with function fields consisting only of constant functions. Our examples will be complex torii. The phrase abelian variety (Chap. 4§??; usually with a extra structure called a polarization) is the name for a complex torus that is algebraic. Chap. 4 analyzes all analytic structures on a dimension one complex torus by corresponding them precisely to the isomorphism class of their function fields. This topic starts in § 3.2.2.

There are two distinct generalizations: To compact Riemann surfaces and to abelian varieties. The former are $\mathbb{P}^{1}$-algebraic while the latter are not in general.
3.2. Classical examples. We discuss two natural first cases of compact complex manifolds.
3.2.1. The Riemann sphere $\mathbb{P}_{z}^{1}$. Let $X$ be the disjoint union of the complex plane $\mathbb{C}$ and a point labeled $\infty$. Here is a coordinate chart:

$$
\begin{aligned}
U_{1}=\mathbb{C}, & \varphi_{1}: U_{1} \rightarrow \mathbb{C} \text { by } \varphi_{1}(z)=z ; \text { and } \\
U_{2}=(\mathbb{C} \backslash\{0\}) \cup\{\infty\}, & \varphi_{2}: U_{2} \rightarrow \mathbb{C} \text { by } \varphi_{2}(\infty)=0 \text { and } \\
& \varphi_{2}(z)=\frac{1}{z} \text { for } z \in \mathbb{C} \backslash\{0\} .
\end{aligned}
$$

Chap. 2 used the Riemann sphere. It embeds in $\mathbb{R}^{3}$. So it is Hausdorff. Then, $X$ is a complex manifold: $\varphi_{2} \circ \varphi_{1}^{-1}(z)=\varphi_{1} \circ \varphi_{2}^{-1}(z)=\frac{1}{z}$ on $\mathbb{C} \backslash\{0\}$ are analytic.

If a complex manifold is compact, some atlas for it contains only finitely many elements. The Riemann sphere required only two (one wouldn't do, would it?).
3.2.2. Complex torus. An atlas for our next example will require four open sets. Let $\omega_{1}$ and $\omega_{2}$ be two nonzero complex numbers satisfying the lattice condition: $\frac{\omega_{2}}{\omega_{1}}$ is not real. Consider the lattice $\omega_{1}$ and $\omega_{2}$ generate:

$$
\begin{equation*}
L\left(\omega_{1}, \omega_{2}\right)=\left\{m_{1} \omega_{1}+m_{2} \omega_{2} \mid m_{1}, m_{2} \in \mathbb{Z}\right\} \tag{3.4}
\end{equation*}
$$

The lattice condition guarantees the natural quotient map $\mathbb{C} \rightarrow \mathbb{C} / L\left(\omega_{1}, \omega_{2}\right)$ has open sets that are like open sets in $\mathbb{C}[9.6 \mathrm{c}]$. According to Lem. 2.3, the manifold
structure on $\mathbb{C}$ automatically gives the manifold structure on $\mathbb{C} / L\left(\omega_{1}, \omega_{2}\right)$. Use the chart $\left\{\left(U_{i}^{\prime}, \varphi_{i}^{\prime}\right)\right\}_{i \in\{0,1,2,3\}}$ of Fig. 3 with $\varphi_{i}^{\prime}$ the inclusion of $U_{i}^{\prime}$ in $\mathbb{C}$. This assures satisfying the Lem. 2.3 condition: Each $z \in \mathbb{C}$ has an $i=i_{z}$ for which $z \in U_{i}^{\prime}$ and the natural map $\mathbb{C} \rightarrow \mathbb{C} / L\left(\omega_{1}, \omega_{2}\right)$ is one-one on $U_{i}$.

The resulting complex manifold $\mathbb{C} / L\left(\omega_{1}, \omega_{2}\right)$ depends only on $L\left(\omega_{1}, \omega_{2}\right)$. Among the many choices we can make of $\omega_{1}, \omega_{2}$ generating this lattice, it is traditional to choose them satisfying special conditions. Elements of the group $\mathrm{SL}_{2}(\mathbb{Z})$ act on $\omega_{1}, \omega_{2}$ to give all pairs of basis elements Chap. 2 [9.15c]. Further, for $a \in \mathbb{C}^{*}$ the scaling $\mathbb{C} \rightarrow \mathbb{C}$ by $z \mapsto a z$ induces a homomorphism $\mathbb{C} / L\left(\omega_{1}, \omega_{2}\right) \rightarrow \mathbb{C} / L\left(a \omega_{1}, a \omega_{2}\right)$ of abelian groups. At the level of coordinate charts, the same scaling gives the map. So, it induces an analytic isomorphism (for precision use Def. 4.1). With no loss take $a=1 / \omega_{1}$, to change the basis of the lattice to $1, \omega_{2} / \omega_{1}$. The ratio $\omega_{2} / \omega_{1}=\tau$ aptly indicates the shape of the parallelogram (3.6). This starts a typical normalization for the complex structure. If we could uniquely indicate the complex structure by $\tau$, that would be an excellent way to parametrize them. The problem is that the complex structure depends only on the lattice $L(1, \tau)$ generated by 1 and $\tau$. Many values of $\tau$ giving the same $L(1, \tau)$. For example, here are three obvious changes:
(3.5a) If necessary, replace $\{1, \tau\}$ by $\{1,-\tau\}$ to assume $\Im(\tau)$ is in the upper half plane $\mathbb{H} \stackrel{\text { def }}{=}\{\tau \in \mathbb{C} \mid \Im(\tau)>0\}$; or
(3.5b) replace $\{1, \tau\}$ by $\{1, \tau+n\}$ for some integer $n$ to assume $0 \leq \Re(\tau)<1$; or (3.5c) scale by $-1 / \tau$ to replace $\{1, \tau\}$ by $\{1,-1 / \tau\}$.

Changes from (3.5) generate a group, $\mathrm{PSL}_{2}(\mathbb{Z})\left(<\mathrm{PSL}_{2}(\mathbb{R}) ; \S 8.2\right)$, acting on $\tau \in \mathbb{H}$.
Lemma 3.11. Together, (3.5) permits restricting a $\tau$ representing a given complex torus (up to isomorphism) to the narrow strip in $\mathbb{H}$ over the closed interval $[0,1) \subset \mathbb{R}$ lying within the closed unit circle around the origin.

Transition functions restrict on each connected component of an intersection of charts to be translation in the complex plane. Topologically this is the same as a torus in $\mathbb{R}^{3}$. Topologists deal with torii, too, though they concentrate especially on the topological space in which the torii sit (see [9.5] for the point of Fig. 2). We care most about this additional complex structure, while they rarely distinguish between one complex torus and another. See $\S 7.2 .3$ for additional comments on attempts to draw pictures in $\mathbb{R}^{3}$.

Figure 2. These two torii could unknot in $\mathbb{R}^{4}$.


Here is the set behind the manifold:

$$
\begin{equation*}
X=\left\{t_{1} \omega_{1}+t_{2} \omega_{2} \mid 0 \leq t_{i}<1, i=1,2\right\} . \tag{3.6}
\end{equation*}
$$

Standard open parallelograms in $\mathbb{C}$ represent each of four coordinate charts in Fig. 3, $U_{i}, i=0,1,2,3$, that do lie in $X$.

Figure 3. Four open sets sort of covering a torus


Let $U_{0}=\left\{t_{1} \omega_{1}+t_{2} \omega_{2} \mid 0<t_{i}<1, i=1,2\right\}$, with $\varphi_{0}: U_{0} \rightarrow \mathbb{C}$ the identity map. The corresponding $U_{0}^{\prime}$ is equal to $U_{0}$ in Fig. 3. On the other hand, consider

$$
U_{1}=\left\{t_{1} \omega_{1}+t_{2} \omega_{2} \left\lvert\, \frac{1}{3}<t_{2}<\frac{2}{3}\right. \text { and either } 0 \leq t_{1}<\frac{1}{3} \text { or } \frac{2}{3}<t_{1}<1\right\}
$$

and $\varphi_{1}: U_{1} \rightarrow \mathbb{C}$ by

$$
\varphi_{1}\left(t_{1} \omega_{1}+t_{2} \omega_{2}\right)=\left\{\begin{array}{ll}
t_{1} \omega_{1}+t_{2} \omega_{2} & \text { for } 0 \leq t_{1}<\frac{1}{3} \\
\left(t_{1}-1\right) \omega_{1}+t_{2} \omega_{2} & \text { for } \frac{2}{3}<t_{1}<1 .
\end{array} .\right.
$$

Form the corresponding $U_{1}^{\prime}$ by translating a pieces of the range of $\varphi_{1}$.
The remaining charts are similar (though slightly more complicated):

$$
\begin{gathered}
U_{2}=\left\{t_{1} \omega_{1}+t_{2} \omega_{2} \left\lvert\, \frac{1}{3}<t_{1}<\frac{2}{3}\right. \text { and either } 0 \leq t_{2}<\frac{1}{3} \text { or } \frac{2}{3}<t_{1}<1\right\}, \\
\varphi_{2}\left(t_{1} \omega_{1}+t_{2} \omega_{2}\right)= \begin{cases}t_{1} \omega_{1}+t_{2} \omega_{2} & \text { for } 0 \leq t_{2}<\frac{1}{3} \\
t_{1} \omega_{1}+\left(t_{2}-1\right) \omega_{2} & \text { for } \frac{2}{3}<t_{2}<1\end{cases} \\
U_{3}=\left\{t_{1} \omega_{1}+t_{2} \omega_{2} \left\lvert\, 0 \leq t_{1}<\frac{1}{2}\right. \text { or } \frac{1}{2}<t_{1}<1,0 \leq t_{2}<\frac{1}{2} \text { or } \frac{1}{2}<t_{2}<1\right\}, \text { and } \\
\varphi_{3}\left(t_{1} \omega_{1}+t_{2} \omega_{2}\right)= \begin{cases}t_{1} \omega_{1}+t_{2} \omega_{2} & \text { for } 0 \leq t_{1}, t_{2}<\frac{1}{2} \\
\left(t_{1}-1\right) \omega_{1}+t_{2} \omega_{2} & \text { for } \frac{1}{2}<t_{1}<1,0 \leq t_{2}<\frac{1}{2} \\
t_{1} \omega_{1}+\left(t_{2}-1\right) \omega_{2} & \text { for } 0 \leq t_{1}<\frac{1}{2}, \frac{1}{2}<t_{2}<1 \\
\left(t_{1}-1\right) \omega_{1}+\left(t_{2}-1\right) \omega_{2} & \text { for } \frac{1}{2}<t_{1}, t_{2}<1\end{cases}
\end{gathered}
$$

To see $X$ is a 1 -dimensional complex manifold check the transition functions $\varphi_{j} \circ \varphi_{i}^{-1}: \varphi_{i}\left(U_{i} \cap U_{j}\right) \rightarrow \varphi_{j}\left(U_{i} \cap U_{j}\right)$. For each $i$ and $j, \varphi_{i}\left(U_{i} \cap U_{j}\right)$ is the union of a finite number of connected open sets. For example,

$$
\varphi_{0}\left(U_{0} \cap U_{1}\right)=U_{1}^{\prime} \backslash\left\{t_{2} \omega_{2} \left\lvert\, \frac{1}{3}<t_{2}<\frac{2}{3}\right.\right\} .
$$

On each connected component of $\varphi_{i}\left(U_{i} \cap U_{j}\right), \varphi_{j} \circ \varphi_{i}^{-1}$ is translation by one of the complex numbers $\delta_{1} \omega_{1}+\delta_{2} \omega_{2}$ where $\delta_{k}$ is 0 or $\pm 1, k=1,2$.

With this manifold structure, $X$ is the complex torus with periods $\omega_{1}$ and $\omega_{2}$.
3.3. Manifolds from algebraic functions. Let $m \in \mathbb{C}[z, w]$ be an irreducible polynomial. Denote the branch points of $m$ by $\boldsymbol{z}$ with $z_{0} \in U_{\boldsymbol{z}}=\mathbb{P}_{z}^{1} \backslash \boldsymbol{z}$ as in Chap. 2 Def. 6.3. Assume $f(z)$ is analytic in a neighborhood of $z_{0}$ and it satisfies $m(z, f(z)) \equiv 0$. Chap. 2 started with two definitions of algebraic functions Def. 1.1 and Def. 1.2. They characterize the same set of functions (Chap. 2 Prop. 7.3).

Riemann's Existence Theorem starts by attaching to each algebraic function a unique (up to analytic isomorphism) compact complex manifold of dimension 1. The next two examples are the first step in that construction, producing an
open subset of the final manifold. We introduce some algebraic geometry using as an excuse showing how to construct explicit manifold compactifications in special cases. We expect coordinates for the abstract compactification of a general Riemann surface to be somewhat mysterious.
3.3.1. An unramified cover of $U_{\boldsymbol{z}}$. Consider first the set

$$
X^{[0]}=X_{f}^{0}=\{(z, w) \in \mathbb{C} \times \mathbb{C} \mid z \notin z, m(z, w)=0\}
$$

Proposition 3.12. The projection map $\mathrm{pr}_{z}: X^{[0]} \rightarrow U_{\boldsymbol{z}}$ by $(z, w) \mapsto z$ produces a natural atlas on $X^{[0]}$ making it a connected complex manifold. For $\lambda \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}, z_{1}\right)$ (Chap. $2 \S 1.1$ ), naturally identify the manifolds $X_{f}^{0}$ and $X_{f_{\lambda}}^{0}$.

Proof. To simplify the construction, assume $\infty \in \boldsymbol{z}$. As usual, apply an element of $\mathrm{PGL}_{2}(\mathbb{C})$ to $\boldsymbol{z}$ to arrange that situation (Chap. 2 §5.2.1; see Lem. 4.3).

Use the implicit function theorem (Chap. $2 \S 6.2$ ) as follows. For $\left(z^{\prime}, w^{\prime}\right) \in X^{[0]}$, let $\Delta_{z^{\prime}}$ be the open disk centered at $z^{\prime}$ of radius the minimum distance from $z^{\prime}$ to a point of $\boldsymbol{z}$. Then, for some one-one analytic function $f_{z^{\prime}, w^{\prime}}(z)$ the following holds.
(3.7) The points $\left(z, f_{z^{\prime}, w^{\prime}}(z)\right)$ are on $X^{[0]}$ and $f_{z^{\prime}, w^{\prime}}\left(z^{\prime}\right)=w^{\prime}$.

For each $\left(z^{\prime}, w^{\prime}\right)$ let $U_{z^{\prime}, w^{\prime}}$ be the range of $z \mapsto F_{z^{\prime}}(z) \stackrel{\text { def }}{=}\left(z, f_{z^{\prime}, w^{\prime}}(z)\right)$ on $\Delta_{z^{\prime}}$. The inverse of $F_{z^{\prime}}$ is $\operatorname{pr}_{z}$, projection of a pair $(z, w)$ onto its $z$-coordinate. Compatible with the definition of manifold, here denote $\mathrm{pr}_{z}$ by $\varphi_{z^{\prime}, w^{\prime}}$. Then, $F_{z^{\prime}}$ parametrizes the neighborhood $U_{z^{\prime}, w^{\prime}}$ of $\left(z^{\prime}, w^{\prime}\right)$ and $\varphi_{z^{\prime}, w^{\prime}}$ maps it into $\mathbb{C}_{z}$. If $V=U_{z^{\prime}, w^{\prime}} \cap U_{z^{\prime \prime}, w^{\prime \prime}}$ is nonempty, then $\varphi_{z^{\prime \prime}, w^{\prime \prime}} \circ \varphi_{z^{\prime}, w^{\prime}}^{-1}$ is the identity map on the overlap of $\Delta_{z^{\prime}} \cap \Delta_{z^{\prime \prime}}$.

That gives an atlas. As it is a subspace of the Hausdorff space $\mathbb{C} \times \mathbb{C}, X^{[0]}$ is Hausdorff. So, it is a connected (from Chap. $2 \S 6.4$ ) complex manifold. Let $\lambda$ be a path as in the statement of the proposition. The point set of $X_{f}^{0}$ consists of pairs $\left(z^{\prime}, x^{\prime}\right) \in \mathbb{C} \times \mathbb{C}$ of the form $\left(z^{\prime}, f_{\gamma}\left(z^{\prime}\right)\right)$ with $\gamma:[a, b] \rightarrow U_{z}$ with $\gamma(a)=z_{0}$ and $\gamma(b)=z^{\prime}$. As $X_{f_{\lambda}}^{0}$ is connected, we can write any point on it as the endpoint of $\left(z, f_{\lambda \cdot \gamma}\right)$ for some $\lambda$. So, $X_{f_{\lambda}}^{0}$ is the same subset of points in $\mathbb{C} \times \mathbb{C}$.

Note: Each $z^{\prime} \in U_{\boldsymbol{z}}$ has a neighborhood $\Delta_{z^{\prime}}$ with this property.
(3.8) $\mathrm{pr}_{z}$ restricted to each connected component $U_{z^{\prime}, w^{\prime}}{\operatorname{of~} \mathrm{pr}_{z}^{-1}\left(\Delta_{z^{\prime}}\right) \text { is a home- }}^{\text {a }}$ omorphism with $\Delta_{z^{\prime}}$.
This is a stronger property than $\mathrm{pr}_{z}$ being an immersion. It means $\mathrm{pr}_{z}: X^{[0]} \rightarrow U_{z}$ is an (unramified) cover according to Def. 7.12. The inverse image by $\mathrm{pr}_{z}$ of small closed disks around $z^{\prime}$ are closed disks around points lying over $z^{\prime}$. That is, the preimage of a compact set is compact, and $\mathrm{pr}_{z}$ is a proper map [9.1d].

REmark 3.13 (Finite atlas). The atlas of Prop. 3.12 contains an infinite number of elements. For a manifold that adds one complication $\S 3.2 .1$ and $\S 3.2 .2$ don't have. This came about to include a deleted neighborhood of $\left(z_{i}, w^{\prime}\right)$ with $z_{i} \in \boldsymbol{z}$ and $w^{\prime}$ a solution of $m\left(z_{i}, w^{\prime}\right)$. That's because we chose disks on $\mathbb{C}_{z}$ as the domain for the $F_{z^{\prime}}$ parametrization. To remedy this choose other simply connected sets, including traditional slit disks given by scaling, translating and rotating

$$
\{z \in \mathbb{C}||z|<1\} \backslash\{0 \leq \Re(z)<1\} .
$$

Chap. $4 \S 2.4$ has further justification for these charts.
3.3.2. Further compactification and use of equations. Chap. 4 Thm. 2.6 shows there is a unique compact complex manifold, up to analytic isomorphism (Def. 4.1), extending $X^{[0]}$ (and also the analytic map to $\mathbb{P}_{z}^{1}$ ). To get it we must compatibly add points and analytic disk neighborhoods to match with the analytic structure on $X^{[0]}$. Using the equation $m(z, f(z)) \equiv 0$ often allows adding further points $\left(z_{i}, w^{\prime}\right)$ to $X^{[0]}$ and their local analytic functions to extend the complex manifold structure. The simplest such extension includes those points $\left(z_{i}, w^{\prime}\right)$ where, even though $z_{i}$ is a branch point, $\frac{\partial m}{\partial w}\left(z_{i}, w^{\prime}\right) \neq 0$. That is, consider

$$
X^{[1]}=\left\{(z, w) \in \mathbb{C} \times \mathbb{C} \mid m(z, w)=0, \frac{\partial m}{\partial w}(z, w) \neq 0\right\} .
$$

The variable for a local chart around $w^{\prime}$ is $w$. Prop. 3.15 gives the details.
Example 3.14. Suppose $h \in \mathbb{C}[w]$ of degree $n>1$ produces $h: \mathbb{P}_{w}^{1} \rightarrow \mathbb{P}_{z}^{1}$. Let $z_{i}$ be a branch point of $m(z, w)=h(w)-z$ and let $g_{z_{i}} \in S_{n}$ be a representative of the conjugacy class attached to $z_{i}$ (Chap. 2 Lem. 7.9). Then, there is a one-one correspondence between the following sets. Chap. 2 [9.4]:
(3.9a) Points $\left(z_{i}, w^{\prime}\right)$ over $z_{i}$ for which $z \mapsto\left(z, f_{z_{i}, w^{\prime}}(z)\right)$ (3.7) parametrizes a neighborhood of $\left(z_{i}, w^{\prime}\right)$.
(3.9b) Disjoint cycles of length 1 in $g_{z_{i}}$.

Example: Consider $h_{1}(w)=w(w-1)(w-2)$. Use notation from Chap. 2 Lem. 7.9. The group attached to an algebraic $f_{1}(z)$ satisfying $h_{1}\left(f_{1}(z)\right)-z \equiv 0$ is $S_{3}$.

Branch cycles $g_{z_{1}}$ and $g_{z_{2}}$ at the two branch points $z_{1}, z_{2}$ have the shape (1)(2) (§7.1.1): disjoint cycles of length 1 and 2. So each branch point has two points above it. Then, for each $z_{i}$ there are two solutions $w_{i, 1}$ and $w_{i, 2}$ of $h(w)-z_{i}$. Select $w_{i, 1}$ so that $\frac{d h}{d w}\left(w_{i, 1}\right) \neq 0$ and $\frac{d h}{d w}\left(w_{i, 2}\right)=0, i=1,2$. Adding $\left(z_{i}, w_{i, 1}\right)$ to $X^{[0]}$ produces an open set on which $\mathrm{pr}_{z}$ maps one-one to $\mathbb{P}_{z}^{1}$. This does not hold for the point $\left(z_{i}, w_{i, 2}\right)$. So, $X^{[1]}$ has exactly one point on it over each of $z_{1}$ and $z_{2}$.

For any $h(w)$ in Ex. 3.14, $X^{[1]}$ will have missing points in that the map $\mathrm{pr}_{z}$ is not proper over some points $z_{i} \in \boldsymbol{z}(\S 2.2)$. For $f$ analytic in several variables $z_{1}, \ldots, z_{n}$ in a neighborhood of a point $\boldsymbol{z}_{0}$, we call

$$
\nabla f\left(\boldsymbol{z}_{0}\right) \stackrel{\text { def }}{=}\left(\frac{\partial f}{\partial z_{1}}\left(\boldsymbol{z}_{0}\right), \ldots, \frac{\partial f}{\partial z_{n}}\left(\boldsymbol{z}_{0}\right)\right)
$$

the complex gradient of $f$ at $\boldsymbol{z}_{0}$. Now consider a set (usually) larger than $X^{[1]}$ :

$$
X^{[2]}=\{(z, w) \in \mathbb{C} \times \mathbb{C} \mid m(z, w)=0, \quad \nabla(m)(z, w) \neq 0\}
$$

Proposition 3.15. A natural atlas makes $X^{[2]}$ into a complex manifold.
Proof. Since $X^{[2]}$ is a subspace of $\mathbb{C} \times \mathbb{C}$ it is Hausdorff. From Prop. 3.12 we have only to add $\left(z_{i}, w^{\prime}\right)$ lying over $z_{i} \in z$ sitting in $X^{[2]}$ to their neighborhoods in $X^{[1]}$. Change the $w^{\prime}$ coordinate by an element of $\mathrm{PGL}_{2}(\mathbb{C})$ to assume none of the finitely many $w^{\prime} \mathrm{s}$ is $\infty$.

By assumption $\nabla(m)\left(z_{i}, w^{\prime}\right) \neq 0$, though by definition $\frac{\partial m}{\partial w}\left(z_{i}, w^{\prime}\right)=0$. Therefore, $\frac{\partial m}{\partial z}\left(z_{i}, w^{\prime}\right) \neq 0$. Apply the implicit function theorem to find a disk $\Delta_{w^{\prime}} \subset \mathbb{C}_{w}$ and $h_{z_{i}, w^{\prime}}(w)$ analytic on $\Delta_{w^{\prime}}$ with the following properties.
(3.10a) The points $\left(h_{z_{i}, w^{\prime}}(w), w\right)$ are on $X^{[1]}$.
(3.10b) The radius of $\Delta_{w^{\prime}}$ is the minimum distance from $w^{\prime}$ to any branch point of $m^{*}(w, z) \stackrel{\text { def }}{=} m(z, w)$ (switch the variables $z$ and $w$ ).

Similar to the proof of Prop. 3.12, let $V_{z_{i}, w^{\prime}}$ be the range of $w \mapsto\left(h_{z_{i}, w^{\prime}}(w), w\right)$ on $\Delta_{w^{\prime}}$. Then the coordinate map at $\left(z_{i}, w^{\prime}\right)$ is $\mathrm{pr}_{w}$ by $(z, w) \mapsto w$.

The essence of producing the manifold structure is to check the transition functions. The key check occurs when the intersection a neighborhood of $\left(z_{i}, w^{\prime}\right)$ meets a neighborhood of $\left(z^{\prime \prime}, w^{\prime \prime}\right)$ with $z^{\prime \prime} \notin \boldsymbol{z}$. For example:

$$
\operatorname{pr}_{w} \circ \operatorname{pr}_{z}^{-1}: z \mapsto\left(z, f_{z^{\prime \prime}, w^{\prime \prime}}(z)\right) \mapsto f_{z^{\prime \prime}, w^{\prime \prime}}(z)
$$

is analytic. Similarly, so is

$$
\operatorname{pr}_{z} \circ \operatorname{pr}_{w}^{-1}: w \mapsto\left(h_{z_{i}, w^{\prime}}(w), w\right) \mapsto h_{z_{i}, w^{\prime}}(w)
$$

That concludes the proof of the lemma.

## 4. Coordinates and meromorphic functions

Here we define analytic maps between complex manifolds. In many areas of mathematics, being able to compare all objects of study with a core of special cases can help. For example, it is helpful to know that all finite groups have a Jordan-Hölder series of finite simple groups and that this collection of finite simple groups (including their multiplicities) is an invariant of the group. Still, even an expert on the classification of finite simple groups can't be confident of a complete understanding of the finite group from knowing its Jordan-Hölder series.

For certain compact complex manifolds, knowing how to use their meromorphic functions can help decide how such a manifold fits among all related manifolds. That is a rough statement of how we use coordinates on compact complex manifolds. This subsection uses explicit (though only partial) compactification of Riemann surfaces of algebraic functions to illustrate how coordinates give defining equations.
4.1. Comparing analytic spaces. We define maps between analytic spaces, and then emphasize the significance of such maps to $\mathbb{P}^{1}$.
4.1.1. Maps between spaces. Let $X_{i}$ be a differentiable (resp., complex) manifold of dimension $n_{i}$ with topologizing data $\left\{\left(U_{\alpha_{i}}, \varphi_{\alpha_{i}}\right)\right\}_{\alpha_{i} \in I_{i}}$. Consider a function $f: X_{1} \rightarrow X_{2}$ and the functions

$$
\begin{equation*}
\varphi_{\alpha_{2}} \circ f \circ \varphi_{\alpha_{1}}^{-1}: \varphi_{\alpha_{1}}\left(U_{\alpha_{1}} \cap f^{-1}\left(U_{\alpha_{2}}\right)\right) \rightarrow \varphi_{\alpha_{2}}\left(f\left(U_{\alpha_{1}}\right) \cap U_{\alpha_{2}}\right) \tag{4.1}
\end{equation*}
$$

for $\left(\alpha_{1}, \alpha_{2}\right) \in I_{1} \times I_{2}$.
Definition 4.1 (Analytic map). Call $f$ differentiable (resp. analytic) if the functions of (4.1) are differentiable (resp. analytic) on their domains. For $X_{1} \subseteq \mathbb{R}^{n}$ and $X_{2} \subseteq \mathbb{R}^{m}$, this is equivalent to $f$ being differentiable as usual. If $f$ is one-one and onto, call $f$ a differentiable (resp. analytic) isomorphism between $X_{1}$ and $X_{2}$.

The phrase isomorphism in Def. 4.1 implies there is a differentiable (resp. analytic) $g: X_{2} \rightarrow X_{1}$ inverse to $f$. That is the gist of our next statement.

Lemma 4.2. Let $X$ and $Y$ be differentiable manifolds. Assume $f: Y \rightarrow X$ is a differentiable map, and in a neighborhood $U_{y}$ of some point $y \in Y$, one-one. Then, there exists differentiable $g: f\left(U_{y}\right) \rightarrow U_{y}$ that is an inverse to $f$. So, if $f$ is one-one and onto, it has differentiable inverse. If we replace the word differentiable by analytic, there is an analogous result.

Proof. Both statements are consequences of the inverse function theorem. This says that a local inverse exists and is differentiable. There is an inverse function to a one-one onto map ( $\S 2.2$ ), so the differentiability is all we need. The definition
of differentiable (or analytic) function reverts this result to one about $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ (or for $f: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ ) for some integer $n$.

Chap. $2 \S 6.1$ discusses the inverse function theorem for one complex variable. The full inverse function theorem is an inductive procedure for several complex variables. See [C89, p. 72] or [Rud76, p. 224] for the general case. For differentiable functions, equation (3.1) says that an inverse $g$ to $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ would have Jacobian matrix $J(g)(\boldsymbol{y})=J(f(\boldsymbol{x}))^{-1}$ at $\boldsymbol{y}=f(\boldsymbol{x})$. This is a differential equation for $g=$ $\left(g_{1}(\boldsymbol{y}), \ldots, g_{n}(\boldsymbol{y})\right)$, given $f$. The case when $f$ is real analytic is much more likely for our use, and that has easier proofs in the literature.
4.1.2. $\mathbb{P}^{1}$-algebraic spaces. Let $\varphi: X \rightarrow Y$ be an analytic map of complex manifolds. If $U$ is an open subset of $Y$, denote the restriction of $\varphi$ over $U$ by $\varphi_{U}: \varphi^{-1}(U) \rightarrow U$. Then, composing holomorphic functions on an open set $U \subset Y$ with $\varphi$ produces a $\operatorname{map} \varphi^{*}: \mathcal{H}(U) \rightarrow \mathcal{H}\left(f^{-1}(U)\right)$. In particular, if both spaces are connected, and $\varphi$ is onto, this induces an injection $\varphi^{*}: \mathbb{C}(X) \rightarrow \mathbb{C}(Y)$, an embedding of the function field of $Y$ into that of $X$.

Chap. 2 Def. 4.14 includes the definition of analytic maps from a domain on $\mathbb{P}_{w}^{1}$ to $\mathbb{P}_{z}^{1}$, a special case of Def. 4.1. More generally, for any complex manifold $X$, a nonconstant analytic map $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is a meromorphic function on $X$ (represented by $z$ ). Chap. 2 Lem. 2.1 guarantees a nonconstant map of compact Riemann surfaces is surjective. This also applies to $\varphi$, even if $X$ (compact) has larger dimension, for again these functions come locally from power series expressions and so give an open map. Further, if $X$ is a compact Riemann surface, Chap. 4 Thm. 2.6 shows any meromorphic function on $X$ extends to give an analytic map from $X$ to projective 1 -space. Chap. 4 Lem .2 .1 shows the following points. If $X$ is compact (and $\varphi$ is nonconstant), then $\varphi$ has a degree, $\left|\varphi^{-1}\left(z^{\prime}\right)\right|$ for $z^{\prime} \in \mathbb{P}_{z}^{1}$ not in a finite set of values where this cardinality is a smaller number. Further, if we count points in $\varphi^{-1}\left(z^{\prime}\right)$ with appropriate multiplicity for their appearance in the fiber, the degree is independent of $z^{\prime} \in \mathbb{P}_{z}^{1}$.

Many compact complex manifolds of dimension at least 2 (example: $\mathbb{P}^{n}, n \geq 2$, [9.11e]), have the following property. Though they have many nonconstant meromorphic functions, none are represented by an analytic map to $\mathbb{P}_{z}^{1}$. The compact complex manifolds that are $\mathbb{P}^{1}$-algebraic are exactly those that embed in $\left(\mathbb{P}^{1}\right)^{N}$ for some integer $N$. That is, they have sufficiently many functions represented by an analytic map to $\mathbb{P}^{1}$, the gist of condition (3.3a).

A virtue of the definition $\mathbb{P}^{1}$-algebraic is its simplicity, this use of special elements of the function field giving maps to $\mathbb{P}^{1}$. Still, Chap. $4 \S 6.1 .1$ extends this, as is traditional, to say a manifold is algebraic if it embeds in $\mathbb{P}^{N}$ for some $N$. The effect of that is to show why a set of basic principles forces extending $\mathbb{P}^{1}$-algebraic manifolds to include $\mathbb{P}^{N}$ as algebraic. We hope this adds historical perspective on what was less than a century ago a complicated issue. Witness this [Mu66, p. 15] quote on going directly from affine space to projective space:

Among others, Poncelet realized that an immense simplication could be introduced in many questions by by considering "projective" algebraic sets (cf. Felix Klein, Die Enwicklund der Mathematik, Part I, p. 80-82). Even to this day, . . . projective algebraic sets play a central role in algebro-geometric questions: therfore we shall define them as soon as possible.

Mumford's quote, and the total acceptance of it in [Har77], shouldn't deny the natural way that $\mathbb{P}^{1}$-algebraic spaces and fiber products illuminate special meromorphic functions arise in providing coordinates.

In practice, on many intensely studied algebraic manifolds, you can choose a finite set, $f_{1}, \ldots, f_{m}$, of global meromorphic functions to construct the manifold, whose points we can then see as given by the values of $f_{1}, \ldots, f_{m}$ at the given point. From these, it is theoretically possible to construct anything else you would expect attached to the manifold from $f_{1}, \ldots, f_{m}$. Still, much classical algebraic geometry spends great time on using coordinates (embeddings in projective space) of special types to make these constructions. For many applications, however, this is a too-detailed reliance on specific use of coordinates. We hope discussions in this chapter help the reader see why coordinates are necessary, though one shouldn't insist on seeing them explicitly at all stages.

We especially study families of compact Riemann surfaces with each family member appearing with an attached equivalence class of maps to $\mathbb{P}_{z}^{1}$. What, however, is the analogy, so important to individual measurements, for comparing different function fields (Lem. 3.9) associated to different complex manifolds? Where would we expect such comparisons to arise? Comparing Riemann surfaces is possible if there is an efficient labeling of function field generators. The easiest event is if all these Riemann surfaces embed naturally in a space with global coordinates that restrict to give coordinates on the individual surfaces. $\S 4.2$ gives examples of how coordinates can help compactify some Riemann surfaces.

An easy way to get new analytic maps from old appears if $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is a meromorphic function. Let $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$. Then $\alpha \circ \varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is a new meromorphic function.

For Ex. 3.14, Prop. 3.15, produces $X^{[2]}$ analytically isomorphic to $\mathbb{C}_{w}$. We already knew this was a manifold. The proof of Props. 3.12 and 3.15 simplifies because $\infty \in z$. The following lemma removes that assumption [9.1b].

Lemma 4.3. Let $U_{i} \subset \mathbb{P}_{z}^{1}, i=1,2$ be domains. Let $\varphi: X \rightarrow U_{1} \cup U_{2}$ denote projection of a manifold for an algebraic function onto the $z$ coordinate. With $\alpha_{i} \in \mathrm{PGL}_{2}(\mathbb{C}), i=1,2$, assume $\alpha_{i} \circ \varphi_{U_{i}}: \varphi^{-1}\left(U_{i}\right) \rightarrow \alpha_{i}^{-1}\left(U_{i}\right)$ is a manifold from the construction of Prop. 3.15, $i=1,2$. Then, $X$ is a complex manifold extending the manifold structure on $\varphi^{-1}\left(U_{i}\right)$.

Assume $X$ is a manifold from Prop. 3.15. Let $\varphi: X \rightarrow U \subset \mathbb{P}_{z}^{1}$ be the algebraic function giving projection onto the $z$ coordinate. Riemann's Existence Theorem (Chap. 4) produces a unique compact complex manifold $\bar{X}$ containing $X$ as an open subset. We do this by extending $\varphi$ to an analytic map $\bar{\varphi}: \bar{X} \rightarrow \mathbb{P}_{z}^{1}$. This is an abstract approach to compactification. It will help to see preliminary examples that relate compactifications and coordinates. In $\S 4.2$ we give these.
4.2. Compactifications and fiber products. Continue the notation for $m$ and its branch points $\boldsymbol{z}$ from $\S 3.3$. Denote

$$
\left\{w^{\prime} \in \mathbb{P}_{w}^{1} \mid\left(z^{\prime}, w^{\prime}\right) \in X^{[0]}, z^{\prime} \notin \boldsymbol{z}\right\} \text { by } U_{\operatorname{pr}_{z}^{-1}(z)}
$$

To further compactify we might embed the subset $X^{[0]}$ of $U_{\boldsymbol{z}} \times U_{\mathrm{pr}_{z}^{-1}(\boldsymbol{z})}$ into a compact space $Z$; then take the closure $X$ of $X^{[0]}$ in $Z$. (Or apply to the already extended spaces $X^{[1]}$ or $X^{[2]}$.) As a closed subspace of compact space, $X$ is compact.
4.2.1. Local holomorphic functions from equations. We note especially that equations give more than an (implicit) description of a point set. Using the implicit function theorem, they often give local parametrizing functions. In this section we use spaces $Z$ to compactify that give natural local equations around points of the closure of $X^{[0]}$. Such equations help decide which points of the closure have extensions to the analytic structure on $X^{[0]}$ (or just manifold structure). This is an aspect of saying such $Z$ provide global coordinates.

We need a notation for holomorphic functions compatible with $\S 1.3$ for the Laurent field $\mathcal{L}_{z^{\prime}}$. We use $\mathcal{L}_{z^{\prime}}^{h}$ for the ring of functions, with each holomorphic in some disk (dependent on the function) about $z^{\prime}$ : power series $\sum_{n=0}^{\infty} a_{n}\left(z-z^{\prime}\right)^{n}$, convergent in some neighborhood of $z^{\prime}$. For a general space $X$ and point $x \in X$, the notation would be $\mathcal{L}_{X, x}^{h}$. For the holomorphic elements of $\mathcal{P}_{z^{\prime}, e}$ use $\mathcal{P}_{z^{\prime}, e}^{h}$.

We've been giving examples of point sets $\{(z, w) \mid m(z, w)=0\}$ in $\mathbb{C} \times \mathbb{C}$ using just one equation. Defining algebraic functions $f\left(z_{1}, \ldots, z_{n}\right)$ in several variables is easy: Consider $X_{m}=\left\{\left(z_{1}, \ldots, z_{n}, w\right) \mid m\left(z_{1}, \ldots, z_{n}, w\right)=0\right\}$, and we say $m$ algebraically defines $f\left(z_{1}, \ldots, z_{n}\right)$, holomorphic in the variables $z_{1}, \ldots, z_{n}$, if $m\left(z_{1}, \ldots, z_{n}, f\left(z_{1}, \ldots, z_{n}\right)\right) \equiv 0$. Also, the notation above extends to consider $\mathcal{L}_{z_{1}^{\prime}, \ldots, z_{n}^{\prime}}^{h}$. Suppose $m\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right)=0$, for $\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right) \in \mathbb{C}^{n+1}$. Assume also that $m$ defines $f\left(z_{1}, \ldots, z_{n}\right)$ algebraically, and $f\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}\right)=w^{\prime}$. Then, we say the local holomorphic (or analytic) functions around $\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right)$ consists of elements of the ring $\mathcal{L}_{z_{1}^{\prime}, \ldots, z_{n}^{\prime}}^{h}$. This ring is invariant under analytic change of variables.

The next definition extends this to consider local holomorphic functions even with no a priori algebraic function $f$ satisfying $m$. Recall the residue class map $\mathrm{rc}_{z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}}: \mathbb{C}\left[z_{1}, \ldots, z_{n}, w\right] \rightarrow \mathbb{C}$ by $\left(z_{1}, \ldots, z_{n}, w\right) \mapsto\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right)$. This is a ring homomorphism, and we record this in the form of the following. The completion of the ring $\mathbb{C}\left[z_{1}, \ldots, z_{n}, w\right] /(m)$ at $\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}\right)$ is

$$
\mathcal{L}_{z_{1}^{\prime}, \ldots, z_{n}^{\prime}}^{h}\left[z_{1}, \ldots, z_{n}, w\right] /\left(m\left(z_{1}, \ldots, z_{n}, w\right)\right) \stackrel{\text { def }}{=} \mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}}^{h} .
$$

Definition 4.4. Analytic functions on $X_{m}$ around $\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right)$ are elements of the localization of $\mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}}^{h}$ at $w=w^{\prime}$ :

$$
\mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}} \stackrel{\text { def }}{=}\left\{u / v \mid u \in \mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}}^{h}, \quad v \in \mathbb{C}\left[z_{1}, \ldots, z_{n}, w\right]\right.
$$ with $u\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right) \neq 0$.

Lemma 4.5. Assume $z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}$ is on $X_{m}$. Then $\mathrm{rc}_{z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}}$ factors through $\mathbb{C}\left[z_{1}, \ldots, z_{n}, w\right] /(m)$; even through $\mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}}^{h}$. This defines the value, $\mathrm{rc}_{( }(s)$, of $s \in \mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}}$ at $\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right)$. Assume the lead coefficient of $m(z, w)$ is invertible in $\mathcal{L}_{z_{1}^{\prime}, \ldots, z_{n}^{\prime}}^{h}$. Use $W_{m}$ for the set of distinct $w^{\prime}=w \operatorname{solving~} m\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w\right)=$ 0 (allowing multiple zeros). There is a natural injective homomorphism

$$
\mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}}^{h} \rightarrow \oplus_{w^{\prime} \in W_{m}} \mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}}^{h}
$$

DEfinition 4.6 (Local holomorphic functions). Suppose $m\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right)=$ 0 , for $\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right) \in \mathbb{C}^{n+1}$ and there are but finitely many solutions $w$ to $m\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right)=0$. Then, the local holomorphic (or analytic) functions that $m$ defines consist of elements of $\mathcal{L}_{X_{m}, z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}}^{h}\left[z_{1}, \ldots, z_{n}, w\right] /\left(m\left(z_{1}, \ldots, z_{n}, w\right)\right)=R$. We say this defines a manifold neighborhood if $R$ is isomorphic to the convergent power series around a point of $\mathbb{C}^{n}$.

It is appropriate to say $R$ is the restriction of local holomorphic functions on $\mathbb{C}^{n+1}$ to the set $X_{m}$ around $\left(z_{1}^{\prime}, \ldots, z_{n}^{\prime}, w^{\prime}\right)$. Further, the definition works as well if several equations, $m_{1}, \ldots, m_{u}$, instead of just one, define the set.
4.2.2. $\mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$ compactification. Since $Z=\mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$ is a product of compact spaces, it is compact. Further, the compactification of $X^{[0]}$, if it is a manifold, suits the definition for $\mathbb{P}^{1}$-algebraic in (3.3).

The natural manifold structure on $Z$ has four open sets in its atlas following Ex. 3.2.1. Label these $U_{i, z} \times U_{j, w}, 1 \leq i, j \leq 2: U_{1, z}=\mathbb{C}_{z}$ and $U_{2, z}=\mathbb{C}_{z}^{*} \cup\{\infty\}$, etc. The atlas gives an isomorphism of each of the four opens sets $U_{i, z} \times U_{j, w}$ with $\mathbb{C} \times \mathbb{C}$, by a map we call $\varphi_{i, j}$. Let $\bar{X}$ be the closure of $X^{[0]}$ in $Z$. We describe the part of $\bar{X}$ lying inside $U_{i, z} \times U_{j, w}$ by an algebraic equation. Then a previous procedure allows checking points at which $X$ has a manifold structure.

Start with $\bar{X} \cap U_{2, z} \times U_{2, w}$, and leave the other open sets as analogous. On the open subset $\mathbb{C}^{*} \times \mathbb{C}^{*} \subset U_{2, z} \times U_{2, w}, \varphi_{2,2}$ acts as

$$
(z, w) \mapsto(1 / z, 1 / w)=\left(z^{\prime}, w^{\prime}\right)
$$

An equation in $\left(z^{\prime}, w^{\prime}\right)$ describes $\varphi_{2,2}$ applied to $X \cap\left(U_{2, z} \times U_{2, w}\right)=X_{2,2}: \varphi_{2,2}\left(X_{2,2}\right)$ is the closure of $\left\{\left(z^{\prime}, w^{\prime}\right) \mid m\left(1 / z^{\prime}, 1 / w^{\prime}\right)=0\right\}$ in $\mathbb{C}_{z^{\prime}} \times \mathbb{C}_{w^{\prime}}$. Get the closure points by allowing $z^{\prime}$ or $w^{\prime}$ to go to 0 . To include those limit values, multiply $m\left(1 / z^{\prime}, 1 / w^{\prime}\right)$ by the minimal powers of $z^{\prime}$ and $w^{\prime}$ to clear the denominators.

Example 4.7 (Continuation of Ex. 3.14). Continue with $m(z, w)=h(w)-z$ and $\operatorname{deg}(h)=n$. The set $\varphi_{2,2}\left(X_{2,2}\right)$ is $\left\{\left(z^{\prime}, w^{\prime}\right) \mid z^{\prime} h^{*}\left(w^{\prime}\right)-\left(w^{\prime}\right)^{n}=0\right\}$ where $h^{*}\left(w^{\prime}\right)=h\left(1 / w^{\prime}\right)\left(w^{\prime}\right)^{n}$. Check that $X_{1,2}$ and $X_{2,1}$ have no new points beyond those already in $X_{1,1}$. Still, $X_{2,2}$ has a new point, corresponding to $\left(z^{\prime}, w^{\prime}\right)=(0,0)$. The gradient of $z^{\prime} h^{*}\left(w^{\prime}\right)-\left(w^{\prime}\right)^{n}$ at zero is $\left(h^{*}(0), 0\right) \neq(0,0)$. So, there is a manifold neighborhood of this point [9.10a].
4.2.3. Tensor products and fiber products of $\mathbb{P}^{1}$ covers. We combine two cases of Ex. 4.7. Suppose $m(z, w)=h(w)-g(z)$, a variables separated equation. Rename $z$ to a variable $w^{\prime}$, and use $z$ for the value $h(w)$. Rewrite $m(z, w)$ as $m\left(w^{\prime}, w\right)$.

Consider $\left(w^{\prime}, w\right) \in \mathbb{C}_{w^{\prime}} \times \mathbb{C}_{w}$ satisfying $m\left(w^{\prime}, w\right)=0$. Call this $X_{m}$. Denote the Riemann surface for a function $w^{\prime}(z)$ (resp. $w(z)$, as in Ex. 4.7) of $z$ satisfying $h\left(w^{\prime}(z)\right) \equiv z($ resp. $g(w)=z)$ by $X_{w^{\prime}}\left(\right.$ resp. $\left.X_{w}\right)$. There is a map $\varphi_{w^{\prime}}: X_{w^{\prime}} \rightarrow \mathbb{P}_{z}^{1}$ by $w^{\prime} \mapsto h\left(w^{\prime}\right)=z$. Similarly for a $\operatorname{map} \varphi_{w}$.

Compare with Def. 1.3: $X_{m}$ as a set is the same as the fiber product of these two maps. Now apply the $\mathbb{P}_{w^{\prime}}^{1} \times \mathbb{P}_{w}^{1}$ compactification to $m\left(w^{\prime}, w\right)$. The resulting set is $\bar{X}_{w^{\prime}} \times_{\mathbb{P}_{z}^{1}} \bar{X}_{w}=\bar{X}_{m}$. (In our example, $\bar{X}_{w^{\prime}}=\mathbb{P}_{w^{\prime}}^{1}$ and $\bar{X}_{w}=\mathbb{P}_{w}^{1}$. .) This is the fiber product (over $\mathbb{P}_{z}^{1}$ ) of the compactifications of $X_{w^{\prime}}$ and $X_{w}$ from Ex. 4.7.

Now consider points of $\bar{X}_{m}$ to decide what are the natural local analytic functions in a neighborhood within one of the four charts for $\mathbb{P}^{1} \times \mathbb{P}^{1}$ :

$$
\begin{equation*}
X_{i, j}=U_{i, z} \times U_{j, w}, 1 \leq i, j \leq 2 \tag{4.2}
\end{equation*}
$$

For $\left(w_{0}^{\prime}, w_{0}\right) \in \bar{X}_{m}$. Let $e_{w_{0}^{\prime}}\left(\right.$ resp. $\left.e_{w_{0}}\right)$ be the ramification index (Chap. 2 Def. 7.6) of $w_{0}^{\prime}$ over $h\left(w_{0}^{\prime}\right)=z_{0}$ (resp. $w_{0}$ over $g\left(w_{0}\right)=z_{0}$ ). New cases are with $e_{w_{0}^{\prime}}=e^{\prime}>1$ and $e_{w_{0}}=e>1$.

Local holomorphic functions in a neighborhood of $\left(w_{0}^{\prime}, w_{0}\right)$ that come from the coordinates $w^{\prime}$ and $w$ are analytic in the solutions $w^{\prime}$ of $h\left(w^{\prime}\right)=z$ expanded about $w_{0}^{\prime}$ and in the solutions $w$ of $g(w)=z$ expanded about $w_{0}$. As usual, use $\zeta_{d}$ for the complex number $e^{2 \pi i / d}$. Assume $R$ is a ring, and $S_{1}$ and $S_{2}$ are two $R$ algebras.

Then the tensor product $S_{1} \otimes_{R} S_{2}$ is the natural direct sum of $R$ algebras. That is, it is an $R$ algebra $T$ with $R$ algebra homomorphisms $\psi_{i}: S_{i} \rightarrow T, i=1,2$ ( $\psi_{1}: s_{1} \in S_{1} \mapsto s_{1} \otimes 1$, etc.) and any such homomorphism will naturally factor through the map to $S_{1} \otimes_{R} S_{2}$. As in Chap. 2 Cor. 7.5: $\left[e_{1}, e_{2}\right]$ is the least common multiple of $e_{1}$ and $e_{2} ; u(z)=\left(z-z^{\prime}\right)^{1 /\left[e_{1}, e_{2}\right]}$ is a choice of $\left[e_{1}, e_{2}\right]$ th root of $z-z^{\prime}$ (a generator of $\left.\mathcal{P}_{z^{\prime},\left[e_{1}, e_{2}\right]}\right)$; and $\zeta_{d}=e^{2 \pi i / d}$. Our first lemma is a famous consequence of the Euclidean algorithm.

Lemma 4.8. Assume $K$ is a characteristic 0 field and $f \in K[x]$ is $\prod_{i=1}^{u} g_{i}(x)^{r_{i}}$ with $g_{1}, \ldots, g_{u}$ irreducible and distinct monic polynomials over $K$. Then the natural $\operatorname{map} \mu: K[x] /(f(x)) \rightarrow \oplus_{i=1}^{u} K[x] /\left(g_{i}^{e_{i}}\right)$ by $h(x) \mapsto\left(h \bmod \left(g_{i}^{e_{1}}\right), \ldots, h \bmod \left(g_{i}^{e_{1}}\right)\right.$ is an isomorphism.

Proof. Check that the kernel of $\mu$ trivial. So this linear vector space map, injects a space of dimension $\operatorname{deg}(f)$ into one of the same dimension $\sum_{i=1}^{u} e_{i} \operatorname{deg}\left(g_{i}\right)$. Conclude: $\mu$ is onto.

Proposition 4.9. Suppose $U$ is an open subset of $\mathbb{P}_{z}^{1}$, and $\varphi: X \rightarrow U$ is an analytic map of Riemann surfaces. For $x^{\prime}$ over $\varphi\left(x^{\prime}\right)=z^{\prime}$ with ramification index $e_{x^{\prime} / z^{\prime}}=e, \mathcal{L}_{X, x^{\prime}}^{h}$ is a natural $\mathcal{L}_{\mathbb{P}_{z}^{1}, z^{\prime}}^{h}$ algebra that identifies with $\mathcal{P}_{z^{\prime}, e}^{h}$.

Let $\varphi_{i}: X_{i} \rightarrow U$ be two such maps, with $x_{i}^{\prime} \in X_{i}$ over $z^{\prime}$ having ramification index $e_{i}, i=1,2$. Let $d=\left(e_{1}, e_{2}\right)$. Then, the ring of local holomorphic functions about $\left(x_{1}, x_{2}\right)$ on $X_{1} \times_{\mathbb{P}_{z}^{1}} X_{2}=Y$ is $\mathcal{L}_{X_{1}, x_{1}^{\prime}}^{h} \otimes_{\mathcal{L}_{\mathbb{P}_{z}^{1}}^{h} z^{\prime}} \mathcal{L}_{X_{2}, x_{2}^{\prime}}^{h}$. So $u^{e_{1} / d}=u_{2}$ (resp. $u^{e_{2} / d}=u_{1}$ ) is an $e_{2}$ th ( $e_{1} t h$ ) root of $\left(z-z^{\prime}\right)$. Then, $\mathcal{L}_{Y,\left(x_{1}, x_{2}\right)}^{h}$ naturally identifies with $\mathcal{L}_{\mathbb{P}_{z}^{1}, z^{\prime}}^{h}\left[u_{1} \otimes 1,1 \otimes u_{2}\right]=R\left(\right.$ with $\left(u_{1} \otimes 1\right)^{e_{1}}=z \otimes 1=\left(1 \otimes u_{2}\right)^{e_{2}}$ according to the rules of tensoring over $\mathcal{L}_{\mathbb{P}_{z}^{1}}^{h} z^{\prime}$ ). This ring has a single maximal ideal. There is an injective homomorphism

$$
\mu: \mathcal{L}_{\mathbb{P}_{z}^{1}, z^{\prime}}^{h}\left[u_{1} \otimes 1,1 \otimes u_{2}\right] \rightarrow \oplus_{j+1}^{d} \mathcal{L}_{\mathbb{P}_{z}^{1}, z^{\prime}}^{h}[x, y] /\left(x^{e_{1} / d}-\zeta_{d}^{j} y^{e_{2} / d}, z=y^{e_{2}}\right)
$$

by $u_{1} \otimes 1 \mapsto x$ and $1 \otimes u_{2} \mapsto y$ in each coordinate. Each summand on the right of (4.9) is an integral domain whose quotient field naturally identifies with $\mathcal{P}_{z^{\prime},\left[e_{1}, e_{2}\right]}$.

Then, $R$ is an integral domain if and only if $d=1$, and the image of $\mu$ in each summand is a proper subring of the summand unless one of $e_{i} / d$ is 1 . Conclude: Restricting local holomorphic functions on $X_{1} \times X_{2}$ defines an analytic manifold structure around $\left(x_{1}^{\prime}, x_{2}^{\prime}\right)$ if and only if one of the $e_{i} s$ is 1 . Yet, the image of $\mu$ generates the quotient field of each summand.

Proof. According to Def. 4.1, by rewriting $\varphi$ using local analytic coordinates $z_{x^{\prime}}$ and $z_{z^{\prime}}$ around $x^{\prime}$ and $z^{\prime}$, we get a very simple normal form. A local analytic change of variables identifies $z_{x^{\prime}}$ with one of the solutions of $u^{e}=z_{z^{\prime}}$. Chap. 2 Cor. 7.5 shows this when $\varphi$ is given by an algebraic function. Chap. 4 (proof of Lem. 2.1) shows it is not dependent on a priori knowing $\varphi$ is algebraic. That gives the first paragraph in the lemma.

Now consider $\varphi_{i}, i=1,2$, in the statement of the lemma. From above, identify an analytic coordinate around $x_{i}(z)$ around $x_{i}^{\prime}$ with $\left(z-z^{\prime}\right)^{1 / e_{i}}$ and the map $\varphi_{i}$ with the $e_{i}$ th power map, $i=1,2$. The only relations among $u_{1} \otimes 1$ and $1 \otimes u_{2}$ are generated by $\left(u_{1} \otimes 1\right)^{e_{1}}=z \otimes 1=\left(1 \otimes u_{2}\right)^{e_{2}}$ and the kernel of the map $\mu$ is in the ideal generated by this relation.

If $d>1$, then $\left(u_{1} \otimes 1\right)^{e_{1} / d}-\zeta_{d}^{j}\left(1 \otimes u_{2}\right)^{e_{2} / d}$ divides $\left(u_{1} \otimes 1\right)^{e_{1}}-\left(1 \otimes u_{2}\right)^{e_{2}}=z$.

Replace $\mathcal{L}_{\mathbb{P}_{z}^{1}, z^{\prime}}^{h}$ by $\mathcal{L}_{z^{\prime}}(y)=K$, a field (leaving $x$ as a variable). Then applying Lem. 4.8 to $\mu$ actually gives an isomorphism. The corresponding summands on the right side of (4.9) would be fields identified with the quotient fields of the summands on the right side of the actual (4.9). So, to finish the result we have only to show the quotient field of the summand $\mathcal{L}_{\mathbb{P}_{z}^{1}, z^{\prime}}^{h}[x, y] /\left(x^{e_{1} / d}-\zeta_{d}^{j} y^{e_{2} / d}, z=y^{e_{2}}\right)$ identifies with $\mathcal{P}_{z^{\prime},\left[e_{1}, e_{2}\right]}$, though the summand itself is a proper subring of the locally holomorphic functions in $\left(z-z^{\prime}\right)^{1 /\left[e_{1}, e_{2}\right]}$ [9.11b].

Now apply Prop. 4.9 to (4.2).
Corollary 4.10. Restricting local holomorphic functions on $\mathbb{P}_{w^{\prime}}^{1} \times \mathbb{P}_{w}^{1}$ to the fiber product $\mathbb{P}_{w^{\prime}}^{1} \times_{\mathbb{P}_{z}^{1}} \mathbb{P}_{w}^{1}$ compactification gives an analytic manifold structure around $\left(w_{0}^{\prime}, w_{0}\right)$ if and only if $\left(e_{w_{0}}^{\prime}, e_{w_{0}}\right)=1$.

Remark 4.11 (simplifying the use of Prop. 4.9). Riemann's Existence Theorem gives a unique compact manifold by completing a cover of $U_{\boldsymbol{z}}$. In so doing, it computes precisely what to expect when you take the fiber product of two ramified covers of $\mathbb{P}_{z}^{1}$ (over of any other Riemann surface). Chap. $4 \S 3.3 .2$ shows the combinatorial result of getting $d$ distinct points on the correctly compactified fiber product (ramified of order $\left[e_{1}, e_{2}\right]$ over $z^{\prime}$ ) over the pair $\left(x_{1}^{\prime}, x_{2}^{\prime}\right)$ is built transparently into the use of branch cycles. Since, however, fiber products (and tensor products) are so important, Prop. 4.9 gives a relatively simple example readers may return to for help with other examples.
4.3. $\mathbb{P}^{n}$ compactifications. Denote the origin in $\mathbb{C}^{n+1}$ by 0 . There is an action of $\mathbb{C}^{*}$ on $\mathbb{C}^{n+1} \backslash\{0\}$. Given a nonzero vector $\boldsymbol{v}=\left(v_{0}, \ldots, v_{n}\right) \in \mathbb{C}^{n+1}$ and $\alpha \in \mathbb{C}^{*}$ form the result of scalar multiplication $\alpha \cdot \boldsymbol{v}=\left(\alpha v_{0}, \ldots, \alpha v_{n}\right)$. Projective $n$-space is a quotient definition like that of a complex torus: $\mathbb{P}^{n}=\mathbb{C}^{n+1} \backslash\{0\} / \mathbb{C}^{*}$. Mapping $\boldsymbol{v}$ to the set equivalent to $\boldsymbol{v}$ gives $\Gamma_{n}: \mathbb{C}^{n+1} \backslash\{0\} \rightarrow \mathbb{P}^{n}$.
4.3.1. An atlas on $\mathbb{P}^{n}$. In this form, it can be convenient (though cumbersome) to label $\mathbb{P}^{n}$ as either $\mathbb{P}_{v_{1} / v_{0}, \ldots, v_{n} / v_{0}}^{n}$ (inhomogeneous coordinates) or $\mathbb{P}_{v_{0}, \ldots, v_{n}}^{n}$ (homogenous coordinates). The extra notation means we have added data for a standard set of coordinate functions for $\mathbb{P}^{n}$. Algebraic geometry texts might refer to a manifold analytically isomorphic to this manifold as $\mathbb{P}^{n}$. Still, there is a significance to adding specific coordinates as Chap. 5 does. To practice this distinction try [9.11e]. Taking $n=1$ and $v_{1} / v_{0}=z$ gives the notation for $\mathbb{P}_{z}^{1}$ from Chap. 2.

Standard coordinates on $\mathbb{P}^{n}$ produce standard transition functions for its manifold structure. Typical of forming an object by an equivalence relation, each point of $\mathbb{P}^{n}$ is a set in $\mathbb{C}^{n+1}$. As some coordinate is not 0 , such a point has a representative with some coordinate equal 1. If you tell which coordinate that is, the representative will be unique.

Let $U_{i}$ be the points with representative having 1 in the $i$ th position. Each point of $\mathbb{P}^{n}$ has a representative in $U_{i}$ for some $i$. Projecting $U_{i}$ onto coordinates different from the $i$ th gives a coordinate chart $\varphi_{i}: U_{i} \rightarrow \mathbb{C}^{n}, i=0, \ldots, n$. If $\boldsymbol{v}$ is any other representative of a point in $U_{i}$, first scale it by $1 / v_{i}$ before this projection.

Lemma 4.12. The atlas $\left\{U_{i}, \varphi_{i}\right\}_{i=1}^{n}$ makes $\mathbb{P}^{n}$ a compact dimension $n$ complex manifold. The map $\mathbb{C}^{n+1} \backslash\{\mathbf{0}\} \rightarrow \mathbb{P}^{n}$ is a map of analytic manifolds.

Proof. An explicit computation of the transition function

$$
\varphi_{j} \circ \varphi_{i}^{-1}: \mathbb{C}_{v_{0}, v_{1}, \ldots, v_{i-1}, v_{i+1}, \ldots, v_{n}} \rightarrow \mathbb{C}_{v_{0}, v_{1}, \ldots, v_{j-1}, v_{j+1}, \ldots, v_{n}}^{n}
$$

is easy. If $i=j$ it is the indentity. Otherwise, it maps $\left(v_{0}, v_{1}, \ldots, v_{i-1}, v_{i+1}, \ldots, v_{n}\right)$ to $1 / v_{j}\left(v_{0}, v_{1}, \ldots, v_{j-1}, v_{j+1}, \ldots, v_{n}\right)$ (with $v_{i}=1$ ). It is analytic on $\varphi_{i}\left(U_{i} \cap U_{j}\right)$.

To see $\mathbb{P}^{n}$ is compact, use the standard absolute value $|v|$ on $\mathbb{C}$. Let $\mathbb{C}_{c}^{n+1}$ be the vectors $\boldsymbol{v}$ with $\max _{i=0}^{n}\left(\left|v_{i}\right|\right) \leq 1$. This is a closed bounded subset of $\mathbb{C}^{n+1}$. So, by the Heine-Borel compactness theorem, it is compact. Every point of $\mathbb{P}^{n}$ has a representative in $\mathbb{C}_{c}^{n+1}$ : Scale it by the largest nonzero entry. Now use that the image of a compact set under a continuous map is compact. An alternate could use this characterization of compactness: Infinite sequences of points in a separable metric space have convergent subsequences [9.10b].

The diagonal in $\mathbb{P}^{n} \times \mathbb{P}^{n}$ is the image of a compact subset of the diagonal in $\mathbb{C}^{n+1} \times \mathbb{C}^{n+1}$. Though the image is compact, until we know $\mathbb{P}^{n}$ is Hausdorff we can't invoke Lem. 2.5 to see the image is closed. Here, however, a direct argument can establish that $\mathbb{P}^{n}$ is Hausdorff. Suppose two points are in one of the $U_{i} s$, a copy of $\mathbb{C}^{n}$. As this is Hausdorff, separate the two points by open sets. So, given any two points it suffices to change coordinates to assure, in the new coordinates, these are both in one of the $U_{i} \mathrm{~s}$. Do that choosing a linear combination $L_{\boldsymbol{a}}=\sum_{i=0}^{n} a_{i} v_{i}$ so neither point lies on the zero set of $L_{\boldsymbol{a}}$. Use $L_{\boldsymbol{a}}$ in place of $v_{j}$ as one of the new coordinates for any $j$ for which $a_{j} \neq 0$.

Use $\Gamma_{n}^{-1}\left(U_{i}\right)=V_{i} \subset \mathbb{C}^{n+1}$ and the same transition functions for a coordinate chart on $\mathbb{C}^{n+1}$. This shows $\Gamma_{n}$ is a map of complex manifolds.
4.3.2. $\mathbb{P}_{z, w, u}^{2}$ compactifications. As in $\S 4.2 .2$, let $Z^{\prime}=\mathbb{P}_{z, w, u}^{2}$. Embed $\mathbb{C}_{z} \times \mathbb{C}_{w}$ in this by $\varphi_{u}^{-1}:(z, w) \mapsto(z, w, 1) \bmod \mathbb{C}^{*} \in Z^{\prime}$. Call the image $U_{u}$. Similarly, let $U_{w}$ be points of $\mathbb{P}_{z, w, u}^{2}$ with a representative of form $(z, 1, u)$ and $U_{z}$ points with a representative of form $(1, w, u)$. Take $X^{\prime}$ to be the closure of $\{(z, w) \mid m(z, w)=0\}$ in the compact space $Z^{\prime}$. To check points of $X^{\prime}$ for a manifold neighborhood requires an equation around each point of $X^{\prime}$. It suffices to define this equation for points of $X^{\prime} \cap U_{z}$ and $X^{\prime} \cap U_{w}$. We do the former; the latter is similar.

Since $\varphi_{z}$ identifies $U_{z}$ with $\mathbb{C}_{w} \times \mathbb{C}_{u}$, it suffices to define the image of $X^{\prime} \cap U_{z}$ under $\varphi_{z}$. With $n^{\prime}$ the total degree of $m$, it is

$$
X_{z}^{\prime}=\left\{(w, u) \mid u^{n^{\prime}} m(1 / u, w / u) .\right.
$$

4.3.3. Hyperelliptic curves. Suppose $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is a degree 2 map of compact Riemann surfaces. Let $z$ be the finite set of branch points (as in Chap. 4 Lem. 2.1). The theme of Chap. $2 \S 8$ is that we already know, from branches of $\log$, what are the abelian covers of $U_{\boldsymbol{z}}=U$ (see Chap. 4 Prop. 2.11). That is, $\pi_{U}: X_{U} \rightarrow U_{\boldsymbol{z}}$ is equivalent to the cover defined by a branch of square root of $h(z) \in \mathbb{C}(z)$. Also, $h$ has multiplicity one zeros and poles contained in $\boldsymbol{z}$ (Chap. 2 (6.2)): $\varphi$ is a cover from a branch of solutions $f(z)$ of $m(z, w)=w^{2}-h(z)$ with $h(z)=\frac{\prod_{i=1}^{t}\left(z-z_{i}\right)}{\prod_{j=t+1}^{r}\left(z-z_{j}\right)}$.

Suppose the $z_{i}$ s are distinct, and all different from 0 or $\infty(r=2 t$ so the degrees of the numerator and denominator are the same). Then, according to Prop. 4.9, this is an if and only if condition that for a manifold compactification given by the fiber product embedding in $\operatorname{pr}_{z}^{1} \times \mathbb{P}_{w}^{1}$. This is good, yet the standard normalization of hyperelliptic curves changes the variables so that $h$ is a polynomial. Do this by multiplying both sides by the square of the denominator, then change the variable $w$ to $w \prod_{j=t+1}^{r}\left(z-z_{j}\right)$. For simplicity we keep the name of the variables the same. So, now consider the equation $w r=h(z)$ where $h=\prod_{i=1}^{r}\left(z-z_{i}\right)$. Here $r$ is even, and we assume it is at least 4. Another common normalization is make the changes
$z \mapsto z_{1}+1 / z$ and $w \mapsto w / z$, thereby replacing $h$ by a polynomial having odd degree $r \geq 3$. As it stands let us consider the $\mathbb{P}_{z, w, u}^{2}$ compactification.

Then, $X_{u}^{\prime}=\left\{(z, w) \mid w^{2}-h(z)=0\right\}$ has a manifold neighborhood around each point: $\nabla(m)=0$ implies $w=0$ and $\frac{d h}{d z}=0(z$ is a repeated root of $h)$. From above,

$$
\begin{align*}
X_{w}^{\prime} & =\left\{(z, u) \mid u^{n-2}-u^{n} h(z / u)=m^{(w)}(z, u)=0\right\} \text { and } \\
X_{z}^{\prime} & =\left\{(w, u) \mid u^{n-2} w^{2}-u^{n} h(1 / u)=m^{(z)}(w, u)=0\right\} . \tag{4.3}
\end{align*}
$$

On $X_{w}^{\prime}$ new points (not already represented on $X_{u}^{\prime}$ ) have $u=0$ and $z=0$. For $r>3$, $\nabla\left(m^{(w)}\right)(0,0)=0$. So, it has no manifold neighborhood. Note this contrasts with the $\mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$ compactification of $m$, in which all points have manifold neighborhoods when you use the right algebraic change of coordinates [9.11c]. For $r=3$, however, the point $(0,0)$ has a manifold neighborhood in $\mathbb{P}^{2}$. There are no new points on $X_{z}^{\prime} ; u=0$ gives no solution in $w$ to $m^{(z)}(w, u)=0$.
4.3.4. Coordinates give meromorphic functions. Let $\bar{X}$ be the $\mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$ compactification (§4.2.2) of $X=\{(z, w) \mid m(z, w)=0\}$ with $m \in \mathbb{C}[z, w]$. Assume every point of $\bar{X}$ has a manifold neighborhood in this compactification. Then, every point of $\bar{X}$ has the form $(z, w) \in \mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$. Thus, projection of $(z, w)$ onto $z$ (or onto $w$ ) provides a meromorphic function on $\bar{X}$.

Similarly, suppose $\bar{X}$ is the $\mathbb{P}_{z, w, u}^{2}$ compactification (§4.3.2) of $X$ and every point of $\bar{X}$ has manifold neighborhood. Then, many meromorphic functions come from this compactification. A linear form in $(z, w, u)$ is a nonzero linear combination of $z, w, u$ (like $L_{\boldsymbol{a}}$, used in the proof of Lem. 4.12). Assume $\bar{X}$ is not in the zero set of any linear form. For example, suppose $m(z, w)$ is irreducible and has total degree $n>1$.

Proposition 4.13. Let $L_{1}$ and $L_{2}$ be linear forms in $(z, w, u)$, not multiples of one another. Let $\left(z_{0}, w_{0}, u_{0}\right)$ represent the unique point of intersection of the zero sets of $L_{1}$ and $L_{2}$. Then, with $z^{\prime}=L_{1}(z, w, u) / L_{2}(z, w, u)$, there is a natural (nonconstant) meromorphic function $\bar{\varphi}: \bar{X} \rightarrow \mathbb{P}_{z^{\prime}}^{1}$. The degree of $\bar{\varphi}$ is $n$ if $\left(z_{0}, w_{0}, u_{0}\right) \notin \bar{X}$ and $n-1$ otherwise.

Proof. Give the map by $(z, w, u) \in \bar{X} \mapsto L_{1}(z, w, u) / L_{2}(z, w, u)$. We verify this map is well-defined. If $\left(z_{0}, w_{0}, u_{0}\right) \notin \bar{X}$, then meaningfully assign a value $z^{\prime} \in \mathbb{C} \cup\{\infty\}$ to the evaluation of $L_{1} / L_{2}$ at any point of $\bar{X}$. Let $H_{z_{0}^{\prime}}$ be the line in $\mathbb{P}^{2}$ given as the zero set of $L_{z_{0}^{\prime}}=L_{1}-z_{0}^{\prime} L_{2}$. To see the degree, check the number of points in the intersection of $H_{z_{0}^{\prime}}$ and $\bar{X}$ if $z_{0}^{\prime}$ is suitably general. This is $n$. These are exactly the points that go to $z_{0}^{\prime}$.

On the other hand, suppose $\left(z_{0}, w_{0}, u_{0}\right) \in \bar{X}$. Then each $H_{z_{0}^{\prime}}$ goes through $\left(z_{0}, w_{0}, u_{0}\right)$. If $z_{0}^{\prime}$ is general, $L_{1}(z, w, u) / L_{2}(z, w, u)$ has a clear ratio value at the $n-1$ points other than $\left(z_{0}, w_{0}, u_{0}\right)$. So, this gives a map of degree $n-1$ of $\bar{X} \rightarrow \mathbb{P}_{z^{\prime}}^{1}$. Check: For only one value $z_{0}^{\prime}$ is $H_{z_{0}^{\prime}}$ tangent to $\bar{X}$ at $\left(z_{0}, w_{0}, u_{0}\right)$ because we assumed $\bar{X}$ is nonsingular [9.11f]. Interpret such a $z_{0}^{\prime}$ as having $\left(z_{0}, w_{0}, u_{0}\right)$ above it.

## 5. Paths, vectors and forms

Notation for paths started in Chap. $2 \S 2.2$. Let $X$ be a topological space. A path in $X$ is a continuous $\gamma:[a, b] \rightarrow X$ for some choice of $a$ and $b$ with $a<b$. The points $\gamma(a)$ and $\gamma(b)$ are, respectively, the initial and end points of the path. The path $\gamma$ is closed if $\gamma(a)=\gamma(b)$.

The idea a path being piecewise differentiable (simplicial) works if $X$ is an $n$-dimensional differentiable manifold (or, more generally, a finite union of differentiable manifolds), with topologizing data $\left\{\left(U_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$. Then, $\gamma$ is differentiable if $\frac{d}{d t}\left(\varphi_{\alpha} \circ \gamma(t)\right)=\boldsymbol{v}_{\alpha}(t)$ exists for each $t \in[a, b]$ (use one-sided limits at the endpoints) and each $\alpha \in I$ with $\gamma(t) \in U_{\alpha}$. The vector $\boldsymbol{v}_{\alpha}(t)$ is the tangent vector to $\gamma$ at $t$ with respect to $\left(U_{\alpha}, \varphi_{\alpha}\right)$. It depends only on $\gamma$ close to $t$.

As in Chap. 2, simplicial paths support applications to integration, and to forming convenient analytic continuations of functions. Still, it is awkward to analyze homotopy classes of paths without allowing paths that are only continuous in the homotopy (see Prop. 6.10).
5.1. Tangent vectors. The above formulation presents a tangent vector as something attached to a path. We recognize a tangent vector at a point $x_{0}$ without having a path through the point. Let $\mathcal{C}_{x_{0}}^{\infty}=\mathcal{C}_{x_{0}, X}$ be functions, differentiable and complex valued, defined in some neighborhood of $x_{0}$.

Definition 5.1. A (complex valued) tangent vector to a differentiable manifold $X$ at a point $x_{0}$ is a linear map $v: C_{x_{0}}^{\infty} \rightarrow C_{x_{0}}^{\infty}$ satisfying Leiznitz's rule:

$$
\begin{equation*}
\boldsymbol{v}\left(f_{1} f_{2}\right)\left(x_{0}\right)=\boldsymbol{v}\left(f_{1}\right)\left(x_{0}\right) f_{2}\left(x_{0}\right)+\left(f_{1}\right)\left(x_{0}\right) \boldsymbol{v}\left(f_{2}\right)\left(x_{0}\right) \tag{5.1}
\end{equation*}
$$

That is, $\boldsymbol{v}$ is a derivation of $\mathcal{C}_{x_{0}}$ defined at $x_{0}$.
5.1.1. Tangent vectors and paths. To relate to tangent vectors attached to a path, assume $x_{0} \in U_{\alpha}$. A function $f$ in a neighborhood of $x_{0}$ defines a function $f \circ \varphi_{\alpha}^{-1}$ on a neighborhood of $\varphi_{\alpha}\left(x_{0}\right) \in \mathbb{R}^{n}$. Denote the variables of $\mathbb{R}^{n}$ here by $\boldsymbol{y}=\left(y_{1}, \ldots, y_{n}\right)$. Consider $F: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ by $\boldsymbol{y} \mapsto\left(F_{1}(\boldsymbol{y}), \ldots, F_{n}(\boldsymbol{y})\right)$. Suppose each coordinate function $F_{i}(\boldsymbol{y})$ has continuous partial derivatives. The Jacobian matrix $J(F)$ of $F$ is the $n \times n$ matrix with $(i, j)$-entry $\frac{\partial F_{i}}{\partial y_{j}}$ at the point $\boldsymbol{y}$.

Lemma 5.2. [Rud76, p. 214] Identify derivations of functions $f \in \mathcal{C}_{\varphi_{\alpha}\left(x_{0}\right), \mathbb{R}^{n}}$ with linear combinations $T_{\boldsymbol{v}}=\sum_{i=1}^{n} v_{i} \frac{\partial}{\partial y_{i}}, v_{1}, \ldots, v_{n} \in \mathcal{C}_{\varphi_{\alpha}\left(x_{0}\right), \mathbb{R}^{n}}$.

So, $T_{\boldsymbol{v}}(f)\left(\varphi_{\alpha}\left(x_{0}\right)\right)$ is the directional derivative of $f$ in the direction $\boldsymbol{v}\left(\varphi_{\alpha}\left(x_{0}\right)\right)$.
For $\gamma(t) \in U_{\alpha} \cap U_{\beta}$, the chain rule relates $\boldsymbol{v}_{\alpha}(t)$ and $\boldsymbol{v}_{\beta}(t)$ :

$$
\begin{equation*}
\left(J\left(\varphi_{\beta} \circ \varphi_{\alpha}^{-1}\right)_{\left(\varphi_{\alpha} \circ \gamma\right)(t)}\right)\left(\boldsymbol{v}_{\alpha}(t)\right)=\boldsymbol{v}_{\beta}(t) . \tag{5.2}
\end{equation*}
$$

So, $\boldsymbol{v}_{\alpha}(t)$ is nonzero if and only if $\boldsymbol{v}_{\beta}(t)$ is nonzero. To check if $\gamma$ has a nonzero tangent vector doesn't depend on the choice of $\left(U_{\alpha}, \varphi_{\alpha}\right)$.
5.1.2. Vector fields. A vector field $T_{U}$ on an open set $U$ in a (differentiable) manifold $X$ is a differentiable assignment of derivations at each point of $U$. A formal definition shows the effect of transition functions from an atlas. Sometimes it is confusing to use $\boldsymbol{y}$ for variables of all copies of $\mathbb{R}^{n}$. So, we use $\boldsymbol{y}_{\alpha}=\left(y_{\alpha, 1}, \ldots, y_{\alpha, n}\right)$ for variables in the range of $\varphi_{\alpha}$.

Definition 5.3. Assume $\left\{U_{\alpha}, \varphi_{\alpha}\right\}_{\alpha \in I}$ is an atlas for the differentiable manifold $X$. Then, $T_{U}$ consists of giving $T_{\alpha}=\sum_{i=1}^{n} f_{\alpha, i} \frac{\partial}{\partial y_{\alpha, i}}$ with the $f_{\alpha, i}$ s differentiable functions on $V_{\alpha}=\varphi_{\alpha}\left(U_{\alpha}\right)$, for each $\alpha \in I$, subject to the following rule. Assume $U_{\alpha} \cap U_{\beta}$ is nonempty. Consider any differentiable function $f: U_{\alpha} \rightarrow \mathbb{R}^{n}$. Use the same notation $T_{\alpha}$ for the restriction of $T_{\alpha}$ to $\varphi_{\alpha}\left(U_{\alpha} \cap U_{\beta}\right)$. Here is a relation between $T_{\alpha}$ and $T_{\beta}$ on $U_{\alpha} \cap U_{\beta}$ :

$$
\begin{equation*}
T_{\alpha}\left(f \circ \varphi_{\alpha}^{-1}\left(y_{\alpha, 1}, \ldots, y_{\alpha, n}\right)\right)=T_{\beta}\left(f \circ \varphi_{\beta}^{-1}\left(y_{\beta, 1}, \ldots, y_{\beta, n}\right)\right) . \tag{5.3}
\end{equation*}
$$

Apply $\left(\frac{\partial}{\partial y_{\beta, 1}}, \ldots, \frac{\partial}{\partial y_{\beta, n}}\right)$ to $f \circ \varphi_{\alpha}^{-1}\left(\boldsymbol{y}_{\beta}\right)=f_{\alpha}\left(\boldsymbol{y}_{\alpha}\right)$ to get a gradient vector of $\left(f_{\alpha_{1}}, \ldots, f_{\alpha_{n}}\right)\left(\boldsymbol{y}_{\alpha}\right)$ functions. A traditional expression rewrites (5.3) as

$$
\begin{equation*}
J\left(\psi_{\boldsymbol{y}_{\beta}, \boldsymbol{y}_{\alpha}}\right)^{-1}\left(\frac{\partial}{\partial y_{\alpha, 1}}, \ldots, \frac{\partial}{\partial y_{\alpha, n}}\right)=\left(\frac{\partial}{\partial y_{\beta, 1}}, \ldots, \frac{\partial}{\partial y_{\beta, n}}\right) \tag{5.4}
\end{equation*}
$$

applied to $f\left(\psi_{\beta, \alpha}\left(\boldsymbol{y}_{\alpha}\right)\right)$ [9.14c]. Thus, (5.3) translates to a linear relation between $\left.\left(f_{\alpha, 1}, \ldots, f_{\alpha, n}\right)\left(\boldsymbol{y}_{\alpha}\right)\right)$ and $\left(f_{\beta, i}, \ldots, f_{\beta, n}\right)\left(\psi_{\beta, \alpha}\left(\boldsymbol{y}_{\alpha}\right)\right)$ [9.14].

So, a chart produces a preferred basis for vector fields and a preferred basis for differential 1-forms from the coordinate functions for the chart.

Definition 5.4. As in Chap. 2 Def. 2.1, $\gamma:[a, b] \rightarrow X$ is simplicial if there is an integer $n$ and $t_{0}=a<t_{1}<\cdots<t_{n-1}<t_{n}=b$ with $\gamma_{\left[t_{i}, t_{i+1}\right]}$ differentiable, $i=0, \ldots, n-1$. Also, $\gamma$ is special simplicial if either $\frac{d}{d t}(\gamma(t))$ is identically zero for $t \in\left(t_{i}, t_{i+1}\right)$ or it is nonzero for each $t \in\left(t_{i}, t_{i+1}\right), i=0, \ldots, n-1$. A space $X$ is simplicially connected if, for each pair $x_{0}, x_{1} \in X$, there is a simplicial path $\gamma:[a, b] \rightarrow X$ with $\gamma(a)=x_{0}, \gamma(b)=x_{1}$.

Lemma 5.5 (Integrating vector fields). Let $T_{U}$ be a vector field on the open set $U$ of the differentiable manifold $X$. For each $u_{0} \in U$ there exists $\epsilon>0$ and $a$ unique differentiable path $\gamma:[-\epsilon, \epsilon] \rightarrow U$, with $\gamma(0)=u_{0}$, so the following holds. The derivation $T_{U, \gamma(t)}$ at $\gamma(t)$ is the directional derivative of $\gamma$ at $t \in[-\epsilon, \epsilon]$.

Proof. With no loss, assume $u_{0}$ is in an atlas element $U_{\alpha}$. We summarize the meaning of the lemma using the previous notation $\varphi_{\alpha}: U_{\alpha} \rightarrow \mathbb{C}^{n}$.

Let $\boldsymbol{y}$ be coordinates on $\mathbb{R}^{n} \supset \varphi_{\alpha}\left(U_{\alpha}\right)$. Use the path $t \mapsto \varphi_{\alpha} \circ \gamma(t)=\gamma^{*}(t)$. By definition, $T_{U_{\alpha}}$ is an expression $\sum_{i=1}^{n} f_{\alpha, i} \frac{\partial}{\partial y_{i}}$. The lemma says there is $\gamma^{*}(t)$ so $\frac{d \gamma_{i}^{*}}{d t}(t)=f_{\alpha, i}\left(\gamma^{*}(t)\right), i=1, \ldots, n$.

Many books quote this result ([Hi65, p. 12], for example) by referring to the existence and uniqueness of solutions to ordinary differential equations. The path in $U_{\alpha}$ is then $\varphi_{\alpha}^{-1}\left(\gamma^{*}(t)\right)$. All general proofs we've seen use fixed point arguments and involve considerable detail, as in the exercises of [Rud76, p. 118, \#25-29, p. $170, \# 25-26]$ giving uniquess and existence under all conditions that would come up for us. Analytic dependence of the solutions on $u_{0}$ is considered more difficult (see [Bo86, p. 171-174, Thm. 4.1]).

Suppose $T_{U}$ is a vector field on $U$ and $\gamma:[a, b] \rightarrow U$ is a differentiable path. Then, call $\gamma$ an integral curve of $T_{U}$. With some assumptions there is a useful converse producing $T_{U}$ from a path. [9.13].
5.2. Holomorphic vector fields and differential forms. Analogs of differentiable vector fields reflect the complex structure on a manifold $X$. The main example from Def. 5.3 has $V_{\alpha}$ as $T_{\alpha}=\sum_{i=1}^{n} f_{\alpha, i} \frac{\partial}{\partial z_{\alpha, i}}$ with the $f_{\alpha, i}$ s holomorphic in the complex coordinates $z_{\alpha, i}, i=1, \ldots, n$. Though $T_{\alpha}$ initially only applies to functions analytic in $\left(z_{\alpha, 1}, \ldots, z_{\alpha, n}\right)$, we may extend it to all differentiable functions taking complex values.
5.2.1. Extend $T$ differentiably. Let $\boldsymbol{z}=\left(z_{1}, \ldots, z_{n}\right)$ be the coordinate functions on $\mathbb{C}^{n}$. Write $z_{j}=x_{j}+i y_{j}$ and $\bar{z}_{j}=x_{j}-i y_{j}$, breaking the coordinates into their real and imaginary parts. Then, $x_{j}=\frac{1}{2} z_{j}+\bar{z}_{j}$ and $y_{j}=\frac{1}{2 i} z_{j}-\bar{z}_{j}$. Define $\frac{\partial}{\partial z_{j}}$ on holomorphic functions $f\left(z_{1}, \ldots, z_{n}\right)$ as the $j$ th partial derivative with respect to the
variables $z_{1}, \ldots, z_{n}$. The partials $\frac{\partial}{\partial x_{j}}$ and $\frac{\partial}{\partial y_{j}}$ act on any differentiable functions of the variables $x_{1}, \ldots, x_{n}, y_{1}, \ldots, y_{n}$ (see Chap. 2 Lem. 2.6).

Lemma 5.6. The operator $\frac{1}{2}\left(\frac{\partial}{\partial x_{j}}-i \frac{\partial}{\partial y_{j}}\right)$ maps $z_{j}$ to $1, z_{k}$ to 0 for $k \neq j$. Further, it maps $\bar{z}_{l}$ to 0 for all l. So, it extends $\frac{\partial}{\partial z_{j}}$ to act as previously on holomorphic functions, and to kill anti-holomorphic functions. Similarly, $\frac{1}{2}\left(\frac{\partial}{\partial x_{j}}+i \frac{\partial}{\partial y_{j}}\right)$ extends $\frac{\partial}{\partial \bar{z}_{j}}$ from anti-holomorphic functions to all differentiable functions.
5.2.2. Vector fields in local coordinates. Suppose $T_{\alpha}$ and $T_{\beta}$ are the expressions for a holomorphic vector field on two coordinate charts. Interpret the relation between the $f_{\alpha, i} \mathrm{~s}$ and $f_{\beta, j} \mathrm{~s}$ given by the complex version of the Jacobian of the transition functions. So, for $X$ a 1-dimensional complex manifold, the equation relating $f_{\alpha}\left(z_{\alpha}\right) \frac{\partial}{\partial z_{\alpha}}$ and $f_{\beta}\left(z_{\beta}\right) \frac{\partial}{\partial z_{\beta}}$ comes from expecting the same value upon application of both to $z_{\beta}=\psi_{\beta, \alpha}\left(z_{\alpha}\right)$ :

$$
\begin{equation*}
f_{\beta}\left(\psi_{\beta, \alpha}\left(z_{\alpha}\right)\right)=f_{\alpha}\left(z_{\alpha}\right) \frac{\partial \psi_{\beta, \alpha}}{\partial z_{\alpha}} \tag{5.5}
\end{equation*}
$$

5.2.3. Differential 1-forms. Now consider the collection of differential 1-forms $\Omega_{U}$ defined on an open set $U$ in a differentiable manifold $X$. Use notation of $\S 5.1 .2$ analogous to that for vector fields. As in $\S$ Chap. 22.3 our motivation is to form integrals of $\omega_{U} \in \Omega_{U}$ along any piecewise differentiable path in $U$.

DEFINITION 5.7. Such an $\omega_{U}$ comes by giving $\omega_{\alpha}=\sum_{i=1}^{n} g_{\alpha, i} d y_{\alpha, i}$ with the $g_{\alpha, i}$ s differentiable functions on $V_{\alpha}=\varphi_{\alpha}\left(U_{\alpha} \cap U\right)$, for each $\alpha \in I$, subject to the following rule. If $V_{\alpha} \cap V_{\beta}$ is nonempty, denote restriction of $\omega_{\alpha}$ to $\varphi_{\alpha}\left(V_{\alpha} \cap V_{\beta}\right)$ also by $\omega_{\alpha}$ and let $\gamma:[a, b] \rightarrow V_{\alpha} \cap V_{\beta}$ be a differentiable path. Then,

$$
\begin{equation*}
\int_{\varphi_{\alpha} \circ \gamma} \omega_{\alpha}=\int_{\varphi_{\beta} \circ \gamma} \omega_{\beta} . \tag{5.6}
\end{equation*}
$$

Equation (5.6) translates to a linear relation between $\left(g_{\alpha, 1}, \ldots, g_{\alpha, n}\right)\left(\boldsymbol{y}_{\alpha}\right)$ and $\left(g_{\beta, i}, \ldots, g_{\beta, n}\right)\left(\psi_{\beta, \alpha}\left(\boldsymbol{y}_{\alpha}\right)\right)$. This formula applies with $\gamma_{[t, t+\epsilon]}$ (restriction of $\gamma$ to $[t, t+\epsilon]$ ) replacing $\gamma$ for any value of $t \in[a, b]$ and $\epsilon>0$. So, it gives equality of the integrands as a function of $t$.

Definition 5.8 (Contraction). Suppose $T_{U}$ is a vector field defined on $U$. Use the previous notation for expressing $T_{U}$ on $V_{\alpha}: T_{\alpha}=\sum_{i=1}^{n} f_{\alpha, i} \frac{\partial}{\partial y_{\alpha, i}}$. The contraction of $T_{\alpha}$ and $\omega_{\alpha}$ is the function $\sum_{i=1}^{n} f_{\alpha, i} g_{\alpha, i}$. Denote it by $\left\langle T_{\alpha}, \omega_{\alpha}\right\rangle$. More generally, the contraction $\left\langle T_{U}, \omega_{U}\right\rangle$ of $T_{U}$ and $\omega_{U}$ is $F \in C_{U}^{\infty}$ with this property.
(5.7) $F \circ \varphi_{\alpha}^{-1}\left(\boldsymbol{y}_{\alpha}\right)=\left\langle T_{\alpha}, \omega_{\alpha}\right\rangle$ on $\varphi_{\alpha}\left(V_{\alpha}\right)$, for each $\alpha \in I$.

Lemma 5.9. As above, $F \circ \varphi_{\alpha}^{-1}$ at $\varphi_{\alpha}(x)$ does not depend on $\alpha$ and the contraction $\left\langle T_{U}, \omega_{U}\right\rangle$ is a differentiable function on $U$. Further, the vector of differentials $\left(d y_{\beta, 1}, \ldots, d y_{\beta, n}\right)$ evaluated at $\psi_{\beta, \alpha}\left(\boldsymbol{y}_{\alpha}\right)$ is $J\left(\psi_{\beta, \alpha}\right)\left(d y_{\alpha, 1}, \ldots, d y_{\alpha, 1}\right)$.

Proof. By explicit computation using Lemma 5.2, $f \circ \varphi_{\alpha}^{-1}$ is the integrand of the left of (5.6). The comment following (5.6) shows this equals the contraction for $\beta$ evaluated at $\psi_{\beta, \alpha}\left(\boldsymbol{y}_{\alpha}\right)$. To conclude the proof use the vector field formula [9.14c]. Contract each side with the differentials $d y_{\beta, j}$ to see the transformation formula for differentials is inverse to that for vector fields.
5.2.4. Tensors. Suppose $\left\{U_{\alpha}, \varphi_{\alpha}\right\}_{\alpha \in I}$ is an atlas for a differentiable manifold $X$. On each $U_{\alpha}$ let $\mathbb{T}_{U_{\alpha}}^{0}$ (resp. $\mathbb{D}_{U_{\alpha}}^{0}$ ) be the tensor algebra over $C^{\infty}\left(U_{\alpha}\right)$ generated by tangent vectors (resp. differential 1-forms) on $U_{\alpha}$. By definition that means elements of $\mathbb{T}_{U_{\alpha}}^{0}$ are finite sums of terms $g T_{1} \otimes T_{2} \otimes \cdots \otimes T_{k}$ with $k$ any nonnegative integer, $g \in C^{\infty}\left(U_{\alpha}\right)$ and $T_{1}, \ldots, T_{k}$ tangent vectors on $U_{\alpha}$. If $k=0$, the element is just the function $g$.

Suppose $h_{1}, h_{2} \in C^{\infty}$ and $T_{i}^{(1)}$ and $T_{i}^{(2)}$ are tangent vectors on $U_{\alpha}$. Further, interpret the tensor sign $\otimes$ to be a formal symbol modulo the following relations. Replacing $T_{i}$ by $h_{1} T_{i}^{(1)}+h_{2} T_{i}^{(2)}$ replaces $g T_{1} \otimes \cdots \otimes T_{i} \otimes \cdots \otimes T_{k}$ by the sum

$$
g h_{1} T_{1} \otimes \cdots \otimes T_{i}^{(1)} \otimes \cdots \otimes T_{k}+g h_{2} T_{1} \otimes \cdots \otimes T_{i}^{(2)} \otimes \cdots \otimes T_{k}
$$

There are two things to note:
(5.8a) Unless it follows from these allowed relations, we do not expect $T_{1} \otimes T_{2}$ to equal $T_{2} \otimes T_{1}$.
(5.8b) Declaring $T_{1} \otimes \cdots \otimes T_{k}$ times $T_{1}^{\prime} \otimes \cdots \otimes T_{k^{\prime}}^{\prime}$ (in that order) to be $T_{1} \otimes \cdots \otimes$ $T_{k} \otimes T_{1}^{\prime} \otimes \cdots \otimes T_{k^{\prime}}^{\prime}$ generates an associative ring multiplication on $\mathbb{T}_{U_{\alpha}}^{0}$.
Similarly for $\mathbb{D}^{0}\left(U_{\alpha}\right)$. Both have $C^{\infty}\left(U_{\alpha}\right)$ as a subring acting by multiplication on each element of $\mathbb{T}_{U_{\alpha}}^{0}\left(\right.$ or $\left.\mathbb{D}_{U_{\alpha}}^{0}\right)$ : These are associate algebras over $C^{\infty}\left(U_{\alpha}\right)$. We may even tensor together elements of $\mathbb{T}_{U_{\alpha}}^{0}$ and $\mathbb{D}_{U_{\alpha}}^{0}$ for a bigger algebra $\mathbb{T}_{U_{\alpha}}^{0} \otimes$ $\mathbb{D}_{U_{\alpha}}^{0}$. In this convention, however, we can distinguish between tangent vectors and differential forms, and typically we pass all the tangent vectors to the left.

A subtlety occurs in comparing elements $\omega_{\alpha} \in \mathbb{T}_{U_{\alpha}}^{0} \otimes \mathbb{D}_{U_{\alpha}}^{0}$ and $\omega_{\beta} \in \mathbb{T}_{U_{\alpha}}^{0} \otimes \mathbb{D}_{U_{\alpha}}^{0}$ on the intersection $U_{\alpha} \cap U_{\beta}$. Use the transition function $\varphi_{\beta} \circ \varphi_{\alpha}^{-1}$ to reexpress $\omega_{\beta}$ in the variables $y_{\alpha, 1}, \ldots, y_{\alpha, n}$ for $\left(U_{\alpha}, \varphi_{\alpha}\right)$ as previously for 1-forms (and vectors). Then, using the formal rules for $\otimes$, compare $\omega_{\alpha}$ and $\omega_{\beta}$ upon their restriction to $U_{\alpha} \cap U_{\beta}$. Suppose the restriction of $\omega_{\alpha}$ and $\omega_{\beta}$ (using the variables $y_{\alpha, 1}, \ldots, y_{\alpha, n}$ ) are the same on $U_{\alpha} \cap U_{\beta}$. Then, we declare them together as forming a general element $\omega$ of the tensor algebra on $U_{\alpha} \cup U_{\beta}$. The subtlety is that $\omega$ likely will not be in $\mathbb{T}_{U_{\alpha} \cup U_{\beta}}^{0} \otimes \mathbb{D}_{U_{\alpha} \cup U_{\beta}}^{0}$. Drop the 0 superscript for a more general algebra.

Definition 5.10. The (mixed) tensor algebra $\mathbb{T}_{X} \otimes \mathbb{D}_{X}$ on $X$ consists of collections $\omega_{\alpha_{i}} \in \mathbb{T}_{U_{\alpha_{i}}}^{0} \otimes \mathbb{D}_{U_{\alpha_{i}}}^{0}, i=1, \ldots, t$, with $\cup_{i=1}^{t} U_{\alpha_{i}}=X$ and $\omega_{\alpha_{i}}$ and $\omega_{\alpha_{j}}$ restricting to equal elements in $\mathbb{T}_{U_{\alpha_{i}} \cap U_{\alpha_{j}}}^{0} \otimes \mathbb{D}_{U_{\alpha_{i}} \cap U_{\alpha_{j}}}^{0}$ for all allowed $i$ and $j$.

Elements of $\mathbb{D}_{X}$ are covariant tensors. If everywhere locally $\omega \in \mathbb{D}_{X}$ is a sum of terms with each a tensor of exactly $k$ differential 1-forms, then it is a $k$-covariant tensor. Generalize contraction (Def. 5.8) to define $\omega$ paired with $k$ ordered tangent vectors $\left(T_{1}, \ldots, T_{k}\right)$. Notice how this requires local expressions of $\omega$ as a sum of terms like $g \omega_{1} \otimes \cdots \otimes \omega_{k}$, with each $\omega_{i}$ a local differential 1-form. This contraction, $\left\langle\left(T_{1}, \ldots, T_{k}\right), \omega\right\rangle$, is a global $C^{\infty}$ function on $X$. For $\omega=g \omega_{1} \otimes \cdots \otimes \omega_{k}$ write it as as $g \prod_{i=1}^{k}\left\langle T_{i}, \omega_{i}\right\rangle$. Such an $\omega$ is symmetric if $\left\langle\left(T_{1}, \ldots, T_{k}\right), \omega\right\rangle=\left\langle\left(T_{(1) \pi}, \ldots, T_{(k) \pi}\right), \omega\right\rangle$ for any permutation $\pi \in S_{k}$. It is alternating (or a differential $k$-form) if

$$
\left\langle\left(T_{1}, \ldots, T_{k}\right), \omega\right\rangle=\operatorname{Det}(\pi)\left\langle\left(T_{(1) \pi}, \ldots, T_{(k) \pi}\right), \omega\right\rangle \pi \in S_{k} \text { (§7.1.4). }
$$

5.2.5. Orientation of a differentiable manifold. A traditional and fuller treatment of the tensor algebra appears in texts on Riemannian geometry like $[\mathbf{H i 6 5}$, Chap. 4]. Riemannian geometry starts with a differentiable manifold and a given symmetric 2-tensor furnished for measuring distances and angles [9.19]. From that tensor appear others for measuring other quantities on the manifold. For example,
if on a differentiable 2-manifold we can measure distances along parametrized paths, then we should also be able to define the area of an open subset. The problem here is that you aren't likely to find a single parametrization by $\mathbb{R}^{2}$ of the whole area, and you must parametrize it in pieces, then add up the resulting areas. This forces the notion of orientation. The only 2 -manifolds that have a well-defined area are orientable, which does include all Riemann surfaces Chap. 4 [11.11].

An orientation on a 2-dimensional differentiable manifold $X$ consists of a rule for continuously assigning a left and right direction at the transversal meeting of two paths on the manifold. Precisely: Suppose given $\gamma^{i}:[-1,1] \rightarrow X, i=1,2$, differentiable paths for which $x \gamma_{i}(0)=x \in X, i=1,2$, and $\left(U_{\alpha}, \varphi_{\alpha}\right)$ is a coordinate chart containing $x$. So, we start with oriented 1-dimensional differential manifolds meeting at a point. Assume also that $\frac{\varphi_{i} \circ \gamma_{i}}{d t}(0)=\boldsymbol{v}_{i}, i=1,2$, are distinct nonzero vectors. View a traveler as moving along $\varphi_{\alpha} \circ \gamma_{1}(t)$, facing at time $t=0$ the direction $\boldsymbol{v}_{1}$ in $\mathbb{R}^{2}$ regarded as the $(x, y)$ plane in $\mathbb{R}^{3}$. Then, the parametric line $L_{0}=\left\{\varphi_{\alpha} \circ \gamma_{1}(0)+s \boldsymbol{v}_{1} \mid s \in \mathbb{R}^{1}\right\}$ cuts the plane so that $\boldsymbol{v}_{2}$ points in the direction of the left half or the right half.

Definition 5.11. Suppose there is a new $\left\{\left(V_{\beta}, \psi_{\beta}\right)\right\}_{\beta \in J}$ on $X$, compatible with the original atlas (usually taken as a subcollection of its coordinate charts) with this property. Independently of the choice of a coordinate chart in the new atlas containing $x$, the vector $\boldsymbol{v}_{2}$ lies consistently in the same half plane (left or right) defined by the corresponding $L_{0}$. Then, we say the new atlas defines an orientation at $x$. The atlas defines an orientation on $X$ if it gives an orientation at each $x \in X$. Riemann surfaces are examples of oriented manifolds.

A generalizing definition inductively allows discussing an orientation of $X$ defined by the oriented meeting of an oriented $n-1$ dimensional manifold meeting an oriented 1-dimensional manifold Chap. 4 [11.5c].
5.3. Meromorphic vector fields and differentials. The definition of vector fields and differential forms is formal. So for each chart, $\left(U_{\alpha}, \varphi_{\alpha}\right)$, it extends to objects of form $T_{\alpha}=\sum_{i=1}^{n} f_{\alpha, i} \frac{\partial}{\partial z_{\alpha, i}}$ or $\omega_{\alpha}=\sum_{i=1}^{n} f_{\alpha, i} d z_{\alpha, i}$ with the $f_{\alpha, i}$ s meromorphic in the complex coordinates $z_{\alpha, i}, i=1, \ldots, n$. Then, since the jacobian of transition functions (and its inverse) have holomorphic function entries, this assures it maps a vector of meromorphic functions to a vector of meromorphic functions.

Example 5.12 (Differential of a meromorphic function). Suppose $X$ is a Riemann surface (not necessarily compact) and $\psi: X \rightarrow \mathbb{P}_{z}^{1}$ is a (nonconstant) meromorphic function on $X$. We produce a meromorphic differential from $\psi$ and an atlas $\mathcal{U}_{X}=\left\{U_{\alpha}, \varphi_{\alpha}\right\}_{\alpha \in I}$ for $X$. Define $d \psi_{\alpha}$ to be $\frac{d \psi \circ \varphi_{\alpha}^{-1}}{d z_{\alpha}} d z_{\alpha}$. Check: This is a differential form satisfying transformation formula (5.6).

Finally, let $\omega$ be a meromorphic differential 1-form on the Riemann surface $X$. Let $x_{0} \in X$ lie in $U_{\alpha}$ where $\omega$ has the expression $f_{\alpha}\left(z_{\alpha}\right) d z_{\alpha}$. Suppose $\varphi_{\alpha}\left(x_{0}\right)=0$. Then, the order $m_{x_{0}}$ of $\omega$ at $x_{0}$ is the order of $f_{\alpha}$ at 0 . Transition functions have neither zeros nor poles. So this order doesn't change if we compute it from another coordinate chart through $\varphi_{\beta}$ with $x_{0} \in U_{\beta}$.
5.3.1. Divisors. Conclude: For a given $\omega$, the formal sum $\sum_{x \in X} m_{x} x$ has meaning. Denote it $(\omega)$ or $D_{\omega}$ depending on the notational context. It is the divisor of $\omega$. Similarly, for any meromorphic function and meromorphic tangent vector on $X$ we may define its divisor $(f)$ or $D_{f}$. Call any formal sum $D=\sum_{x \in X} m_{x} x$ a divisor, and $m_{x}$ is its support multiplicity at $x$.

Lemma 5.13. On a connected Riemann surface $X$, let $D$ be the divisor of a nonconstant meromorphic differential, function or tangent vector. Then, the points of nonzero support multiplicity for $D$ have no accumulation point. So, if $X$ is also compact, divisors of nonconstant meromorphic differentials, functions or tangent vectors have only a finite number of nonzero support multiplicities.

Proof. We do the case for differentials. The others are similar. Suppose $(\omega)=\sum_{x \in X} m_{x} x$ is the divisor of a differential and infinitely many of the $m_{x}$ are nonzero. Then, this set of $x \mathrm{~s}$ has an accumulation point, $x_{0}$. Let $\left(U_{\alpha}, \varphi_{\alpha}\right)$ be a coordinate chart containing $x_{0}$, so the statement is that on $\varphi_{\alpha}\left(U_{\alpha}\right)$ we have a meromorphic differential $f_{\alpha}\left(z_{\alpha}\right) d z_{\alpha}$ having an accumulation of zeros or poles at $\varphi_{\alpha}\left(x_{0}\right)=z_{\alpha}^{\prime}$. As in Chap. 2 [9.8a], this implies $f_{\alpha}$ is identically zero (or $\infty$ ) and using connectedness, that the same holds for the differential, contrary to our assumption (for extra help, see the argument of Chap. 4 Lem. 2.1).

If $X$ is not compact, divisors as in Lem. 5.13 may have infinitely many nonzero support terms (as with a holomorphic nonpolynomial function in the complex plane $\mathbb{C}_{z}$ ). In fact, the next general result in the complex plane has a similar version for any noncompact Riemann surface attached to an algebraic function [Ahl79, p. 195].

Proposition 5.14 (Weierstrass factorization). Suppose $\left\{m_{x_{i}}\right\}_{i \in I}$ is any collection of nonzero integers attached to a sequence of distinct points $\left\{x_{i} \in \mathbb{C}_{z}\right\}_{i \in I}$ with no accumulation point in $\mathbb{C}_{z}$. Then, there is a holomorphic function $f(z)$ with $(f)=\sum_{i \in I} m_{x_{i}} x_{i}$. Also, $f(z) d z$ (resp. $f(z) \frac{\partial}{\partial z}$ ) is a holomorphic differential (resp. vector field) with exactly the same divisor.

Still, our tool will be the investigation of differentials, functions, etc., that extend meromorphically to a natural compactification of $X$. So, we typically assume (unless otherwise said) that $m_{x}=0$ except for finitely many $x \in X$. For such a divisor $D$, the sum $\sum_{x \in X} m_{x}$ is the degree $\operatorname{deg}(D)$ of $D$. A divisor $D$ is positive (or $D \geq 0$ ) if all its support multiplicities are nonnegative. This definition gives a partial ordering on divisors: With $D=\sum_{x \in X} m_{x} x$ and $D^{\prime}=\sum_{x \in X} m_{x}^{\prime} x, D \geq D^{\prime}$ if $m_{x} \geq m_{x}^{\prime}$ for each $x \in X$. Equivalently, with the obvious subtraction of divisors, $D-D^{\prime}$ is positive.

Multiplying two functions or a function and a differential gives an object with divisor having the sum of the constituent multiplicities: $(f \omega)=(f)+(\omega)$.

Definition 5.15. Suppose $X$ is a compact Riemann surface. We say two divisors $D_{1}$ and $D_{2}$ on $X$ are linearly equivalent if $D_{2}-D_{1}=(f)$ for some meromorphic function $f: X \rightarrow \mathbb{P}_{z}^{1}$. This is an equivalence relation between divisors.

Our notation for the linear equivalence class of a divisor $D$ on a compact Riemann surface will be $[D]$. On a compact Riemann surface, the divisor of a meromorphic function has degree 0 (Chap. 4 Lem. 2.1; see Ex. 5.17). Anticipating that, conclude there is a well-defined degree attached to a linear equivalence class of divisors. Finally, we have a crucial definition attached to a divisor for which the reader should practice the notation.

Definition 5.16. For any divisor $D$ on a Riemann surface, the linear system of $D, L(D)$, is the collection of meromorphic functions $f$ for which $(f)+D \geq 0$.
5.3.2. Relation between functions and differentials. As in Ex. 5.12, any (nonconstant) meromorphic function on a Riemann surface $X$ provides us a nontrivial
meromorphic differential form. Further, assume $\omega_{1}, \omega_{2}$ are meromorphic differentials and $\omega_{1}$ is not a constant multiple of $\omega_{2}$. This produces a nonconstant meromorphic function $\psi: X \rightarrow \mathbb{P}_{z}^{1}$ by the formula

$$
\begin{equation*}
\psi \circ \varphi_{\alpha}^{-1}\left(z_{\alpha}\right)=\omega_{\alpha, 1} / \omega_{\alpha, 2} . \tag{5.9}
\end{equation*}
$$

So, all nonconstant differentials are linearly equivalent, and (see Def. 5.15), on a compact Riemann surface, all have the same degree.

Example 5.17 . Consider the identity map $z: \mathbb{P}_{z}^{1} \rightarrow \mathbb{P}_{z}^{1}$ by $z \mapsto z$. Carefully consider what is $d z=\omega$ using Ex. 3.2.1. To clarify notation, denote $\varphi_{1}$ by $\varphi_{\alpha}$ and $\varphi_{2}$ by $\varphi_{\alpha^{\prime}}$. Then, $\varphi_{\alpha}: \mathbb{C}_{z} \rightarrow \mathbb{C}_{z_{\alpha}}$ by $z \mapsto z$, and so $\omega_{\alpha}=\frac{d z_{\alpha}}{d z_{\alpha}} d z_{\alpha}=d z_{\alpha}$. Also, $\varphi_{\alpha^{\prime}}: \mathbb{C}_{z}^{*} \cup\{\infty\} \rightarrow \mathbb{C}_{z_{\alpha^{\prime}}}$ by $z \mapsto z^{-1}$. So,

$$
\begin{equation*}
\omega_{\alpha^{\prime}}=\frac{d z_{\alpha^{\prime}}^{-1}}{d z_{\alpha^{\prime}}} d z_{\alpha^{\prime}}=-z_{\alpha^{\prime}}^{-2} d z_{\alpha^{\prime}} \tag{5.10}
\end{equation*}
$$

The differential $d z$ is meromorphic, not holomorphic, and it has degree -2. To see there are no nonconstant holomorphic differentials on $\mathbb{P}_{z}^{1}$, write such a differential as $g(z) d z$ with $g$ a meromorphic function on $\mathbb{P}_{z}^{1}$. Liouville's Theorem says $g$ has as many zeros as poles [Ahl79, p. 122]. So the degree of $g(z) d z$ also is -2 , and $(g(z) d z)$ cannot be positive. A similar computation shows the vector space of holomorphic differentials on a complex torus has dimension 1 [9.8].
5.3.3. Pulling back differentials. Let $f: X_{1} \rightarrow X_{2}$ be an analytic and surjective map between complex manifolds. Then, a meromorphic function $\psi: Y \rightarrow \mathbb{P}_{z}^{1}$ produces a meromorphic function $\psi \circ f \stackrel{\text { def }}{=} f^{*}(\psi): X \rightarrow \mathbb{P}_{z}^{1}$ giving an embedding $\mathbb{C}(Y) \subset \mathbb{C}(X)$ (§4.1.2).

Lemma 5.18. We may extend $f^{*}$ to embed meromorphic differentials $\mathcal{M}^{1}(Y)$ on $Y$ into meromorphic differentials $\mathcal{M}^{1}(X)$. Further, this maps holomorphic differents $\Omega^{1}(Y)$ on $Y$ into holomorphic differentials on $X$. Then $\varphi^{*}$ has the following property. For $\omega \in \mathcal{M}^{1}(Y)$, suppose $\gamma \in \Pi_{1}\left(X, x_{0}\right)$ does not go through a pole of $\varphi^{*}(\omega)$. Then, $\int_{\gamma} \varphi^{*}(\omega)=\int_{\varphi_{*}(\gamma)} \omega$.

Proof. Use the notation of (4.1). To simplify we do this for the case of 1-dimensional complex manifolds, though the many variable case is just a slight addition to the notation. This is truely a local statement. Write $\omega$ as $h_{\alpha_{2}}\left(z_{\alpha_{2}}\right) d z_{\alpha_{2}}$ on $\varphi_{\alpha_{2}}\left(f\left(U_{\alpha_{1}}\right) \cap U_{\alpha_{2}}\right)$. Then, define $f^{*}(\omega)$ by

$$
h_{\alpha_{2}}\left(\varphi_{\alpha_{2}} \circ f \circ \varphi_{\alpha_{1}}^{-1}\left(z_{\alpha_{1}}\right)\right) d\left(\varphi_{\alpha_{2}} \circ f \circ \varphi_{\alpha_{1}}^{-1}\left(z_{\alpha_{1}}\right)\right) \text { on } U_{\alpha_{1}} \cap f^{-1}\left(U_{\alpha_{2}}\right) .
$$

The equality of the integrals is nothing more, after substituting for the coordinates of the path $\gamma$, than the change of variables formula Chap. 2 Lem. 2.3.
5.4. Half-canonical differentials. Square-roots of differentials appear on a Riemann surface $X$ when we seek a canonical choice of $\theta$ function attached to the surface. The case when $X$ has genus 1 (Chap. $4 \S 7.5$ ) will be our guide.

Riemann's $\theta$ functions often allow us to put coordinates (as in the initial discussion of $\S 4$ ) on such total familes. Whenever possible, we would like the construction of such coordinates to be canonical. Usually, however, constructing $\theta$ functions depends on choices. So, we are careful to note, for curves in families, how the construction varies with the points parametrizing the family members.

Riemann used $\theta$ functions to give coordinates for constructing objects, like differentials and functions on a Riemann surface. When the Riemann surface has
genus 1 (or 0) there are natural choices for working with Riemann's coordinates. When, however, the genus exceeds 1 , and the surface is not special, there are several $\left(2^{2 g-1}-2^{2 g-2}\right)$ potential choices of the odd $\theta$ function Riemann required to generalize Abel's Theorem. We will see that half-canonical differentials precisely differentiate between these choices.
5.4.1. Cocycles. For $X$ an $n$-dimensional complex manifold, let $\left\{U_{\alpha}, \varphi_{\alpha}\right\}_{\alpha \in I}$ be the coordinate chart, and $\left\{\psi_{\beta, \alpha}=\varphi_{\beta} \circ \varphi_{\alpha}^{-1}\right\}_{\alpha, \beta \in I}$ the corresponding collection of transition functions (as in Def. 3.6). Each $\psi_{\beta, \alpha}$ then is a one-one analytic function on an open subset of $\mathbb{C}^{n}$ whose coordinates we label $z_{\alpha, 1}, \ldots, z_{\alpha, n}$. Denote the $n \times n$ complex Jacobian matrix for $\psi_{\beta, \alpha}$ by $J\left(\psi_{\beta, \alpha}\right)$. Call the matrices $\left\{J\left(\psi_{\beta, \alpha}\right)\right\}_{\alpha, \beta \in I}$ the (transformation) cocycle attached to meromorphic differentials. Similarly $\left\{J\left(\psi_{\beta, \alpha}\right)^{-1}\right\}_{\alpha, \beta \in I}$ is the cocycle attached to meromorphic tangent vectors. Recall the notation for $n \times n$ matrices, $\mathbb{M}_{n}(R)$ with entries in an integral domain $R$ and for the invertible matrices $\mathrm{GL}_{n}(R)$ with entries in $R$ under multiplication. Cramer's rule says for each $A \in \mathbb{M}_{n}(R)$ there is an adjoint matrix $A^{*}$ so that $A A^{*}$ is the scalar matrix $\operatorname{det}(A) I_{n}$ given by the determinant of $A$. This shows the invertibility of $A \in \mathbb{M}_{n}(R)$ is equivalent to $\operatorname{det}(A)$ being a unit (in the multiplicatively invertible elements $R^{*}$ ) of $R$. Denote the $n \times n$ identity matrix (resp. zero matrix) in $\mathrm{GL}_{n}(R)$ by $I_{n}\left(\right.$ resp. $\left.\mathbf{0}_{n}\right)$.

Definition 5.19 (1-cycocle). Suppose $g_{\beta, \alpha} \in \operatorname{GL}_{n}\left(\mathcal{H}\left(U_{\alpha} \cap U_{\beta}\right)\right), \alpha, \beta \in I$. Assume also that $g_{\gamma, \beta} g_{\beta, \alpha}=g_{\gamma, \alpha}$ for all $\alpha, \beta, \gamma \in I$ on $U_{\alpha} \cap U_{\beta} \cap U_{\gamma}$ (if this is nonempty). Then, $\left\{g_{\beta, \alpha}\right\}_{\alpha, \beta \in I}$ is a multiplicative 1-cocycle with values in $\mathcal{G} \mathcal{L}_{n, X}$. Similarly, suppose $g_{\beta, \alpha} \in \mathbb{M}_{n}\left(\mathcal{H}\left(U_{\alpha} \cap U_{\beta}\right)\right)$, $\alpha, \beta \in I$. Suppose $g_{\gamma, \beta}+g_{\beta, \alpha}=g_{\gamma, \alpha}$ for all $\alpha, \beta, \gamma \in I$ on $U_{\alpha} \cap U_{\beta} \cap U_{\gamma}$. Then, $\left\{g_{\beta, \alpha}\right\}_{\alpha, \beta \in I}$ is an additive 1-cocycle with values in $\mathcal{G} \mathcal{L}_{n, X}$.

We also name (1-)cocycles for collections of subgroups in $\mathcal{G} \mathcal{L}_{n, X}$ (resp. $\mathcal{M}_{n, X}$ ) for which it makes sense to multiply (resp. add) $g_{\gamma, \beta}$ and $g_{\beta, \alpha}$. So, for example, we may consider a multiplicative cocycle with values in $\left\{ \pm I_{n}\right\}$ or an additive cocycle with values in $\mathbb{Z} I_{n}$. When there are 1-cocycles, there are also 0 -chains and their associated 1-boundaries. We write the definition for $\mathrm{GL}_{n}$, recognizing there are analogous versions for all other types of cocycles.

DEFINITION 5.20 (1-boundary). With $u_{\alpha} \in \operatorname{GL}_{n}\left(\mathcal{H}\left(U_{\alpha}\right)\right), \alpha \in I$, suppose $g_{\beta, \alpha}=u_{\beta}\left(u_{\alpha}\right)^{-1}$ for all $\alpha, \beta, \gamma \in I$ in $U_{\alpha} \cap U_{\beta}$ (if nonempty). Then, $\left\{g_{\beta, \alpha}\right\}_{\alpha, \beta \in I}$ is a 1 -cocycle, called a 1 -boundary with values in $\mathcal{G} \mathcal{L}_{n, X}$. Call the set $\left\{u_{\alpha}\right\}_{\alpha \in I}$ a 0 -chain with values in $\mathcal{G} \mathcal{L}_{n, X}$.
5.4.2. Half-canonical divisors. Suppose $\omega$ is a meromorphic differential on a Riemann surface $X$, written locally as $f_{\alpha}\left(z_{\alpha}\right) d z_{\alpha}$ on simply connected domains $U_{\alpha}$ (Chap. 2 §8.3). Assume also the square hypothesis:
(5.11) The divisor of $f_{\alpha}\left(z_{\alpha}\right)$ has the form $2 D_{\alpha}$ for $U_{\alpha}$ running over a subchart covering $X$.
Then, there is a branch $h_{\alpha}\left(z_{\alpha}\right)$ of square root (of $f_{\alpha}\left(z_{\alpha}\right)$ ) on $U_{\alpha}$ (Chap. 2 (6.2)). Of course, there are two of these; our notation means we have chosen one. Call the symbol $\tau_{\alpha}=h_{\alpha}\left(z_{\alpha}\right) \sqrt{d z_{\alpha}}$, a half-canonical divisor on $U_{\alpha}$. The squares of these form a global differential on $X$. Denote the collection $\left\{h_{\alpha}\left(z_{\alpha}\right)\right\}_{\alpha \in I}$, by $\boldsymbol{h}$ and refer to it as a square-root of $\omega$.

Lemma 5.21 (Half-canonical divisor). The collection of divisors $\left\{\left(h_{\alpha}\left(z_{\alpha}\right)\right)\right\}_{\alpha \in I}$ from a square root of $\omega$ give a well-defined divisor: a half-canonical divisor on $X$.

Proof. Let $D=(\omega)$ be the divisor of $\omega$. Since, $h_{\alpha}^{2}=f_{i, \alpha}$, the support multiplicities of $D$ are all even integers. So, a square-root of $\omega$ defines $D_{1 / 2}=(\omega) / 2$, a divisor uniquely given by the zeros and poles of the $h_{\alpha} \mathrm{s}$.

Now consider how to decide, based on a square-root of $\omega$, if there is an object $\omega_{1 / 2}$ with values at points on $X$ whose divisor is $D_{1 / 2}=(\omega) / 2$. Continue the transition function notation $\psi_{\beta, \alpha}$ from $\S 5.4 .1$. This requires us to make sense, on $U_{\alpha} \cap U_{\beta}$, of equality between

$$
\begin{equation*}
\tau_{\alpha}\left(z_{\alpha}\right)=h_{\alpha}\left(z_{\alpha}\right) \sqrt{d z_{\alpha}} \text { and } \tau_{\beta}\left(\psi_{\beta, \alpha}\left(z_{\alpha}\right)\right)=h_{\beta}\left(\psi_{\beta, \alpha}\left(z_{\alpha}\right)\right) \sqrt{d \psi_{\beta, \alpha}\left(z_{\alpha}\right)} . \tag{5.12}
\end{equation*}
$$

Proposition 5.22. Assume each component of $U_{\alpha} \cap U_{\beta},(\alpha, \beta) \in I \times I$ is simply connected and for such, we have made a choice of $\sqrt{J\left(\psi_{\beta, \alpha}\right)}=g_{\beta, \alpha}$ on $U_{\alpha} \cap U_{\beta}$. Then, independent of $\alpha$ with $x^{\prime} \in U_{\alpha}$, setting the value of $\tau_{\alpha}$ to $h_{\alpha}\left(\varphi_{\alpha}\left(x^{\prime}\right)\right)$ is welldefined if and only if $\left\{g_{\beta, \alpha}\right\}_{(\alpha, \beta) \in I \times I}=\boldsymbol{g}$ is a 1-cocycle. If there is a $\boldsymbol{g}$ that is a 1 -cocycle, call the resulting half-canonical differential $\omega_{1 / 2, \boldsymbol{h}, \boldsymbol{g}}$. Then, with $\boldsymbol{g}$ fixed, but $\boldsymbol{h}^{\prime}$ varying over square-roots of $\omega$, any pair of $\omega_{1 / 2, \boldsymbol{h}^{\prime}, \boldsymbol{g}}$ differ by a 1-boundary with values in $\{ \pm 1\}$.

Proof. We need only add that the cocycle condition on $\boldsymbol{g}$ is necessary and sufficient for (5.12). For this check that if $x^{\prime} \in U_{\alpha} \cap U_{\beta} \cap U_{\gamma}$, then all the values $h_{\alpha}\left(\varphi_{\alpha}\left(x^{\prime}\right)\right), h_{\beta}\left(\varphi_{\alpha}\left(x^{\prime}\right)\right)$ and $h_{\gamma}\left(\varphi_{\gamma}\left(x^{\prime}\right)\right)$ at $x^{\prime}$ match up using $\boldsymbol{g}$. Comparing (5.12) for each of the pairs $(\alpha, \beta),(\beta, \gamma)$ and $(\alpha, \gamma)$ gives the cocycle condition.
5.4.3. Square-hypothesis for hyperelliptic curves. Suppose the affine part of a hyperelliptic curve $X$, with compactification from Ex. 4.2.3, is $\left\{(z, w) \mid w^{2}=h(z)\right\}$. We explicitly display differentials $\omega$ satisfying the square hypothesis of (5.11). For simplicity, assume $h$ has odd degree and distinct zeros $z_{1}, \ldots, z_{r-1}$ (with $z_{r}=\infty$ ). Denote the point on $X$ over $z_{i}$ by $x_{i}$, with $x_{\infty}$ lying over $z=\infty$. As in [Mum76, p. 7], form the differentials

$$
\omega_{i}=\frac{\left(z-z_{i}\right)^{\frac{1}{2}}}{\left(\prod_{j \neq i} z-z_{j}\right)^{\frac{1}{2}}} d z, i=1, \ldots, r-1
$$

Since $w=\sqrt{h(z)}$, the factor in front of the $d z$ in $\omega_{i}$ is just $\frac{z-z_{i}}{w}$, a meromorphic function on $X$. The divisor of $\omega_{i}$ is therefore $2 x_{i}-2 x_{\infty}=D_{i}$. For the check at a neighborhood of $x_{\infty}$ over $z=\infty$, use $t=1 / \sqrt{z}$ as the uniformizing parameter on $X$. Consider the case $\operatorname{deg}(h)=3$. Then, $\left(t^{-1}-z_{i}\right)\left(-2 w t^{3}\right) d t$ has $t=0$ as a pole of order 2. So, $D_{i}$ is the same divisor as $\left(z-z_{i}\right)$.

Now consider the case $\operatorname{deg}(h)=r-1, r \geq 6$ an even integer. Similarly, $\left(\omega_{i}\right)=2 x_{i}+2(r / 2-3) x_{\infty}$, as $\frac{z-z_{i}}{-2 w t^{3}} d t$ has $t=0$ as a zero of multiplicity $2(r / 2-3)$.

## 6. Homotopy, monodromy and fundamental groups

Complex structure provides the notion of analytic continuation. We detect the effects of analytic continuation through monodromy action, a representation of some fundamental group. In practice this can be a permutation representation, a representation as automorphisms of a vector space or a representation into automorphisms of a more general group. The prototype use of monodromy is Riemann's Existence Theorem: We replace constructing a compact Riemann surface using charts with permutation representations of a fundamental group. For example, using classical generators (Chap. 4 Fig. 3) for the fundamental group of
$U_{z}=\mathbb{P}_{z}^{1} \backslash\{\boldsymbol{z}\}$ gives an effective listing of Riemann surface covers (and their corresponding algebraic functions; Chap. 4 Cor. 2.9).
6.1. Homotopy of paths. Let $\gamma_{i}:\left[a_{i}, b_{i}\right] \rightarrow X, i=1,2$, be two one-one simplicial paths in $X$ with the same range, initial, and end points. The function $f(t)=\gamma_{2}^{-1} \circ \gamma_{1}$ is a simplicial path $f:\left[a_{1}, b_{1}\right] \rightarrow\left[a_{2}, b_{2}\right]$ for which $\frac{d}{d t}(f(t)) \geq 0$ (where the derivative is defined) and $\gamma_{2}(f(t))=\gamma_{1}$. (Use the chain rule.) We give a more general statement.

Definition 6.1 (Image equivalent paths). Let $\gamma:\left[a_{1}, b_{1}\right] \rightarrow X$ be a simplicial path in $X$, and let $f_{1}:\left[a_{2}, b_{2}\right] \rightarrow\left[a_{1}, b_{1}\right]$ and $f_{2}:\left[a_{1}, b_{1}\right] \rightarrow\left[a_{2}, b_{2}\right]$ be simplicial paths with $\frac{d}{d t}\left(f_{i}(t)\right) \geq 0$ where it is defined, $i=1,2$. Assume also $\gamma \circ f_{1} \circ f_{2}(t)=\gamma(t)$ for $t \in\left[a_{1}, b_{1}\right]$. Call $\gamma$ and $\gamma \circ f_{1}$ image equivalent paths. It is a simple exercise to show each path is image equivalent to a path $\gamma:[0,1] \rightarrow X$.

Definition 6.2 (Homotopically equivalent paths). Consider a continuous map $F:[a, b] \times[0,1] \rightarrow X$, and points $x_{a}, x_{b} \in X$, with the following properties: $F(t, s)=\gamma_{s}(t)$ is a path for each $s \in[0,1]$ with initial point $x_{a}$ and end point $x_{b}$. Call $F$ a homotopy between $\gamma_{0}$ and $\gamma_{1}$ (or $\gamma_{0}$ and $\gamma_{1}$ are homotopic).

REmARK 6.3 (Warnings!). The end points of the paths $\gamma_{s}$ remain fixed throughout a homotopy, or else all paths in a connected space would be homotopic.

Even if $\gamma_{0}$ and $\gamma_{1}$ are simplicial paths, we do not initially assume $\gamma_{s}$ is also simplicial. Still, the argument of Chap. 2 Lem. 4.3 generalizes easily to any (union of) differentiable manifold(s) to say that any continuous path is homotopic to a simplicial path. Further, it is then image equivalent to a product of simplicial paths that are either constant or have nonzero derivative, and if it is a nonconstant path, you can toss out - up to equivalence - the constant paths. We use this statement freely [9.12]. It is common to think of both $s$ and $t$ as time parameters. It is compatible to consider the range of $\gamma_{0}$ as a physical object layed down parametrically. As a function of time, each point $\gamma_{0}(t)$ of the range of $\gamma_{0}$ moves to a different position $\gamma_{s}(t)$. So, $F$ represents deforming an initial path, perhaps along which it is more efficient to accrue similar information from traversing $\gamma_{0}$.

In Fig. 4 the space $X$ is the same as Fig. 3. Note: $\gamma_{1}$ and $\gamma_{2}$ are closed, beginning and ending at $0 \bmod L \in \mathbb{C} / L$.

Figure 4. $\gamma_{1}$ can't deform to $\gamma_{2}$ on $X$


Definition 6.4. Extend the definition of homotopic paths. We say two paths $\gamma_{i}:\left[a_{i}, b_{i}\right] \rightarrow X, i=1,2$, with $\gamma_{1}\left(a_{1}\right)=\gamma_{2}\left(a_{2}\right)$ and $\gamma_{1}\left(b_{1}\right)=\gamma_{2}\left(b_{2}\right)$ are equivalent (or homotopic) if $\gamma_{1}$ and $\gamma_{2}$ are image equivalent, respectively, to homotopic paths $\gamma_{i}^{*}:[a, b] \rightarrow X, i=1,2$, for some $a<b$. This is an equivalence relation.
6.2. Analytic continuation on a manifold. Suppose $f \in \mathcal{E}\left(D, z_{0}\right)$ is extensible in a domain $D$ and $\gamma[a, b] \rightarrow D$ is a path. Chap. 2 Rem. 4.4 notes the production of a simplicial path $\gamma^{*}$ in $D$ for which the analytic continuations $f_{\gamma}$ and $f_{\gamma^{*}}$ are the same. Further, assume $f$ is extensible as a holomorphic (rather than just meromorphic) function in $D$. Then, define $F_{\gamma}$ for any antiderivative $F$ of $f$ (around $z_{0}$ ) as the analytic continuation $F_{\gamma^{*}}$. Chap. 2 Lem. 4.3 produces $\gamma^{*}$ from $\gamma$ by a succession of homotopies, between a piece of path on $\gamma$ contained in a disk and a line segment joining two points on the boundary of the disk. Disks are a crucial case of the following definition. The simple lemma following it, hidden in the construction of $\gamma^{*}$, appears in most arguments about homotopy classes.

Definition 6.5. Call a topological space $X$ contractible (to $x_{0} \in X$ ) if there is a continuous function $f: X \times[0,1] \rightarrow X$ satisfying $f(x, 0)=x$ and $f(x, 1)=x_{0}$ for each $x \in X$.

Lemma 6.6. A closed or open ball (or anything homeomorphic to such) in $\mathbb{R}^{n}$ is contractible. If $X$ is contractible, then any two paths with the same endpoints are homotopic [9.12b].

Analytic continuation of a meromorphic function (Chap. 2 Def. 4.1) extends to manifolds by imitating the other extensions to manifolds. Suppose $X$ is a complex manifold with coordinate chart $\left\{\left(U_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}$. Consider any path $\gamma:[a, b] \rightarrow X$. Our notation follows the case for a dimension 1 complex manifold, though it extends easily to the general case.

By a disk (or ball) $D$ on $X$ we mean an open set in $X$ which lies in one coordinate neighborhood $U_{\alpha}$ where $\varphi_{\alpha}(D)$ is a disk (or ball) in $\varphi_{\alpha}\left(U_{\alpha}\right)=V_{\alpha}$.
6.2.1. Extensible functions on $X$. Follow Chap. $2 \S 4.1$ to extend analytic continuation of a function along a path to where the path is in a complex manifold.

Definition 6.7 (Analytic continuation along a path). Suppose $f$ is meromorphic in a neighborhood $U_{x_{0}} \subset X$ of $x_{0} \in X$ and $\gamma:[a, b] \rightarrow X$ is a path based at $x_{0}$. Let $f^{*}:[a, b] \rightarrow \mathbb{P}_{z}^{1}$ be a continuous function with the following properties.
(6.1a) $f^{*}(t)=f(\gamma(t))$ for $t$ close to $a$ (in $[a, b]$ ).
(6.1b) For each $t^{\prime} \in[a, b]$, there is a neighborhood $U_{\gamma\left(t^{\prime}\right)}$ of $\gamma\left(t^{\prime}\right)$ and an analytic function $h_{t^{\prime}}: U_{\gamma\left(t^{\prime}\right)} \rightarrow \mathbb{P}_{z}^{1}$ with $h_{t^{\prime}}(\gamma(t))=f^{*}(t)$ for $t$ near $t^{\prime}$ (in $[a, b]$ ).
As before, $h_{t^{\prime}}$ is the analytic continuation of $f$ to $t^{\prime}$. It is an analytic function in some neighborhood of $\gamma\left(t^{\prime}\right)$. Reference is usually to the end function $h_{b}=f_{\gamma}$, analytic in a neighborhood of $\gamma(b)$. This is the analytic continuation of $f$ (along $\gamma$ ). As with analytic continuation along a path in $\mathbb{P}_{z}^{1}, f^{*}(t)$ determines all data for an analytic continuation. Also, it is unique: its difference from another function suiting (6.1) must be constant (restrict to coordinate neighborhoods of points of the path and apply Chap. 2 [9.8a]). Again, there is a related definition.
6.2.2. Algebraic functions on $X$. An analytic function $\hat{f}: X \rightarrow \mathbb{P}_{z}^{1}$ satisfying $\hat{f}(x)=f(x)$ for all $x \in U_{x_{0}}$ is an analytic continuation or extension of $f$ to $X$.

Definition 6.8. Denote by $\mathcal{E}\left(X, x_{0}\right)$ all functions meromorphic in a neighborhood of $x_{0}$ that analytically continue along every path in $X$ based at $x_{0}$.

Further, suppose there is compact Riemann surface $\bar{X}$ with $X=\bar{X} \backslash \boldsymbol{x}$ where $\boldsymbol{x}$ is a finite set of points on $\bar{X}$. Chap. 4 shows, if such a $\bar{X}$ exists, it is unique up to analytic isomorphism. If $\boldsymbol{x}$ consists of $r$ points, call such an $X$ an $r$-punctured Riemann surface. Dropping reference to $r$, call it just a punctured Riemann surface. This tacitly assumes $r$ is a finite number.

Definition 6.9. Suppose $X$ is a punctured Riemann surface. Then, $\mathcal{E}\left(X, x_{0}\right)^{\text {alg }}$ consists of the $f \in \mathcal{E}\left(X, x_{0}\right)$ for which both the following sets are finite.
(6.2a) All analytic continuations, $\mathcal{A}_{f}(X)=\left\{f_{\gamma}\right\}_{\gamma \in \Pi_{1}\left(X, x_{0}\right)}$ of $f$ in $X$.
(6.2b) For $x^{\prime} \in \boldsymbol{x}$, the limit endpoint values of $f_{\gamma}$ along all $\gamma \in \Pi_{1}\left(X, x_{0}, x^{\prime}\right)$.

Proposition 6.10. Let $D$ be a disk on $X$, and suppose $f: D \rightarrow \mathbb{P}_{z}^{1}$ is analytic. There is a partition $a=t_{0}<t_{0}^{*}<t_{1}<t_{1}^{*}<\cdots<t_{n-1}^{*}<t_{n}=b$ of $[a, b]$, coordinate neighborhoods $\left(U_{i}, \varphi_{i}\right)$, a disk $D_{i}$ centered about $\gamma\left(t_{i}\right)$ in $U_{i}$ and $f_{i} \in \mathcal{H}\left(D_{i}\right), i=$ $1, \ldots, n-1$, with these properties.
(6.3a) $D_{i} \cap D_{i+1} \neq \emptyset$ and $f_{i}(z)=f_{i+1}(z)$ for $z \in D_{i} \cap D_{i+1}$.
(6.3b) $\gamma(t) \in D_{i}$ for $t \in\left[t_{i}, t_{i}^{*}\right], \gamma(t) \in D_{i+1}$ for $t \in\left[t_{i}^{*}, t_{i+1}\right], i=0, \ldots, n-1$.
(6.3c) $f_{0}(z)=f(z)$ for $z \in U_{z_{0}}$.

Further, let $\gamma^{*}$ be the path along the consecutive line segments $\gamma\left(t_{i}\right)$ to $\gamma\left(t_{i}^{*}\right)$, then $\gamma\left(t_{i}^{*}\right)$ to $\gamma\left(t_{i+1}\right), i=0, \ldots, n-1$. Then, $f_{\gamma^{*}}=f_{\gamma}$.

Proof. The proof reduces to that of Chap. 2 Lem. 4.3 by using the definition of function and coordinate charts on a complex manifold.

Proposition 6.11 (The general monodromy theorem). Let $\gamma_{1}, \gamma_{2}:[a, b] \rightarrow X$ be two paths with $\gamma_{1}(a)=\gamma_{2}(a)=x_{0}$ and $\gamma_{1}(b)=\gamma_{2}(b)=x_{1}$. Suppose $\gamma_{1}$ and $\gamma_{2}$ are homotopic on $X$. Let $U_{x_{0}}$ be a neighborhood of $x_{0}$ and $f: U_{x_{0}} \rightarrow \mathbb{P}_{z}^{1}$ Then, $f_{\gamma_{1}}=f_{\gamma_{2}}([\mathbf{A h l 7 9}$, p. 295] and [Con78, p. 219] $)$.

Proof. Let $F:[a, b] \times[0,1] \rightarrow X$ be a homotopy between $\gamma_{1}$ and $\gamma_{2}$ fixing points $x_{a}=x_{0}, x_{b}=x_{1} \in X$. A continuous function on a compact space is absolutely continuous. From absolute continuity of $F$ there are partitions

$$
a=s_{0}<s_{1}<\cdots<s_{n}=b \text { of }[a, b] \text { and } 0=t_{0}<t_{1}<\cdots<t_{m}=1 \text { of }[0,1]
$$

so that $F:\left[s_{i}, s_{i+1}\right] \times\left[t_{j}, t_{j+1}\right] \rightarrow X$ has range in a coordinate chart $U_{i, j}$ on $X$ and $\varphi_{i, j}: U_{i, j} \rightarrow \mathbb{C}$ has range in a disk.

Suppose $h$ is meromorphic in a neighborhood of $F\left(s_{i}, t_{j}\right)$ and extensible on the range of $F$ on $\left[s_{i}, s_{i+1}\right] \times\left[t_{j}, t_{j+1}\right]$. Denote the product of the paths

$$
s \mapsto F\left(s, t_{j}\right)=F_{i j, 1}, s \in\left[s_{i}, s_{i+1}\right] \text { and } t \mapsto F\left(s_{i+1}, t\right)=F_{i+1 j, 2}, t \in\left[t_{j}, t_{j+1}\right]
$$

by $\mu_{i j}^{+}$. Similarly, let $\mu_{i j}^{-}$be the product of paths $t \mapsto F\left(s_{i}, t\right)=F_{i j, 2}, t \in\left[t_{j}, t_{j+1}\right]$ and $s \mapsto F\left(s, t_{j+1}\right)=F_{i j+1,1}, s \in\left[s_{i}, s_{i+1}\right]$. From Chap. 2 Lem. 4.6, $h_{\mu_{i j}^{+}}=h_{\mu_{i_{j}}^{-}}$.

Write the path $\gamma_{1}$ as the product of the paths $F_{i 0,1}, i=0, \ldots, m$. Similarly, $\gamma_{2}$ is the product of the paths $F_{i n, 1}, i=0, \ldots, m$. We give a sequence of paths (with the same endpoints) that starts with $\gamma_{1}$, and ends with $\gamma_{2}$. The terms of the sequence differ from path-to-path in the chain by a product of paths of form $\left(\mu_{i j}^{+}\right)^{-1} \mu_{i j}^{-}$or of form $\gamma \gamma^{-1}$. This shows $f_{\gamma_{1}}=f_{\gamma_{2}}$. Simply replace $F_{i 0,1}$ by

$$
F_{i 0,1} F_{i+10,2} F_{i+10,2}^{-1}\left(\mu_{i 0}^{+}\right)^{-1} \mu_{i 0}^{-}
$$

for each $i=1, \ldots, m$. These substitutions lead from $\gamma_{1}$ to the path that is the product of $F_{i 1,1}, i=0, \ldots, m$. Continue inductively to the path $\gamma_{2}$, which is the product of $F_{i 1, n}, i=0, \ldots, m$.

Chap. $2 \S 4.4$ defines the product of two paths $\gamma_{i}:\left[a_{i}, b_{i}\right] \rightarrow X, i=1,2$, for which the end point of $\gamma_{1}$ is the initial point of $\gamma_{2}$. Many treatments on fundamental groups (like [Ma; Chap. 2]) restrict the domain interval for a path to $[0,1]$. The treatment here aids computation of the Artin braid group (Chap. 4 [11.8] and

Chap. 5). It has other virtues: If $\gamma_{i}:\left[a_{i}, b_{i}\right] \rightarrow X, i=1,2,3$, are three paths with $\gamma_{i}\left(b_{i}\right)=\gamma_{i+1}\left(a_{i+1}\right), i=1,2$, then $\gamma_{1}\left(\gamma_{2} \gamma_{3}\right)$ and $\left(\gamma_{1} \gamma_{2}\right) \gamma_{3}$ are identical rather than just equivalent as in [Ma67, p. 59]. Thus, forming products is trivially associative.
6.3. Path equivalence classes form a group. We say $\gamma:[a, b] \rightarrow X$, a closed path with initial (and end) point $x_{0} \in X$, is based at $x_{0}$. The set of paths based at $x_{0}$ is closed under taking products. Denote the (homotopy) equivalence class of $\gamma$ by $[\gamma]$. Note: $\left[\gamma_{1}^{*} \gamma_{2}^{*}\right]$ is independent of the choice of $\gamma_{i}^{*} \in\left[\gamma_{i}\right], i=1,2$. The function $\gamma:[a, b] \rightarrow X$ by $\gamma(t)=x_{0}$ is called a constant path; denote $[\gamma]$ by $\epsilon_{x_{0}}$. The set of equivalence classes of paths in $X$ based at $x_{0}$ is the fundamental group of $X$ based at $x_{0}$.

Theorem 6.12. Equivalence classes of paths into $X$ based at $x_{0}$ form a group, denoted $\pi_{1}\left(X, x_{0}\right)$, under the multiplication given by $\left[\gamma_{1}\right]\left[\gamma_{2}\right] \stackrel{\text { def }}{=}\left[\gamma_{1} \gamma_{2}\right]$. The identity element is $\epsilon_{x_{0}}$. The inverse of $[\gamma]$ is the class $\left[\gamma^{-1}\right]$ (Chap. 2 §4.4).

Proof. Consider $\gamma:[a, b] \rightarrow X$ and $\gamma^{-1}$ as above. Let $s^{\prime}=a+s(b-a)$ and consider the function $F:[a, 2 b-a] \times[0,1] \rightarrow X$ defined by

$$
F(t, s)= \begin{cases}\gamma(t) & \text { for } t \in\left[a, s^{\prime}\right]  \tag{6.4}\\ \gamma\left(s^{\prime}\right) & \text { for } t \in\left[s^{\prime}, 2 b-s^{\prime}\right] \\ \gamma(2 b-t) & \text { for } t \in\left[2 b-s^{\prime}, 2 b-a\right]\end{cases}
$$

So, $F$ is a homotopy between $\gamma \gamma^{-1}$ and the constant path from $[a, 2 b-a]$ into $\left\{x_{0}\right\}$.
From [9.12b], for $\gamma_{0}:\left[a_{0}, b_{0}\right] \rightarrow\left\{x_{0}\right\}$, the paths $\gamma_{0} \gamma$ and $\gamma \gamma_{0}$ are equivalent to $\gamma$. Thus, $[\gamma]\left[\gamma^{-1}\right]=\epsilon_{x_{0}},[\gamma] \epsilon_{x_{0}}=[\gamma]=\epsilon_{x_{0}}[\gamma]$. This shows $\pi_{1}\left(X, x_{0}\right)$ is a group.

The fundamental group does depend on the base point $x_{0}$, though its isomorphism class does not. Indeed, for $x_{0}, x_{1} \in X$, let $\alpha:[a, b] \rightarrow X$ be a path with initial point $x_{0}$ and end point $x_{1}$. Define $\psi\left(x_{0}, x_{1}\right): \pi_{1}\left(X, x_{1}\right) \rightarrow \pi_{1}\left(X, x_{0}\right)$ by

$$
\psi\left(x_{0}, x_{1}\right)([\gamma])=\left[\alpha \gamma \alpha^{-1}\right] \text { for each }[\gamma] \in \pi_{1}\left(X, x_{1}\right)
$$

Check that $\psi\left(x_{0}, x_{1}\right)$ is a homomorphism of groups inverse to the homomor$\operatorname{phism} \psi\left(x_{1}, x_{0}\right):[\gamma] \in \pi_{1}\left(X, x_{0}\right) \mapsto\left[\alpha^{-1} \gamma \alpha\right] \in \pi_{1}\left(X, x_{0}\right)$. Note: The isomorphism $\pi\left(x_{0}, x_{1}\right)$ depends on the choice of $\alpha$ if $\pi_{1}\left(X, x_{0}\right)$ is not an abelian group.

Corollary 6.13. For $x_{0}, x_{1} \in X, \pi_{1}\left(X, x_{0}\right)$ and $\pi_{1}\left(X, x_{1}\right)$ are isomorphic.
Still, we eventually come to fundamental groups of members of a family of topological spaces (Chap. 5), where all members have the same fundamental group. Our most profound (the braid and Hurwitz monodromy) groups appear to account for different identifications among these fundamental groups.
6.4. Fundamental group of a circle. For any differentiable manifold $X$, there is a natural map from the fundamental group $\pi_{1}\left(X, x_{0}\right)$ computed with piecewise differentiable paths to the fundamental group computed with continuous paths, $\pi_{1}\left(X, x_{0}\right)^{\text {cont }}$. This induces an isomorphism (though we don't exploit this seriously) from Rem. 6.3. This point shows in a comparison of the two fundamental groups when $X=S^{1}$, a circle which we take to be the unit circle in $\mathbb{C}_{z}$. We give two proofs that it is isomorphic to $\mathbb{Z}$. The first explicitly uses simplicial paths. The other uses the universal covering space (Lem. 8.4).

Consider the path $\gamma_{[a, b]}^{*}:[a, b] \rightarrow S^{1}$ by $t \mapsto \cos (2 \pi t)+i \sin (2 \pi t), t \in[a, b]$. For $n \geq 0$ an integer, denote $\gamma_{[0, n]}^{*}$ by $\gamma_{n}^{*}$, and let $S^{1}$ be the image of $\gamma_{1}^{*}$. Denote
the inverse of $\gamma_{\left.\right|_{[0,1]} ^{*}}^{*}$ by $\left(\gamma^{*}\right)_{\left.\right|_{[0,1]}}^{-1}$. Since $\left(\gamma_{1}^{*}\right)^{n}=\gamma_{n}^{*}$ it is consistent to define $\gamma_{-n}^{*}$ to be $\left(\gamma_{1}^{*-1}\right)^{n}$. For $n=0$ let $\gamma_{0}^{*}$ be the constant path mapping to 1 .

Figure 5. Homotopically speaking, a path going nowhere. Traversal for $t \in\left[\frac{i}{6}, \frac{i+1}{6}\right], i=0,1,2,3,4,5$


THEOREM 6.14. The group $\pi_{1}\left(S^{1}, 1\right)$ is infinite cyclic with generator $\left[\gamma_{1}\right]$.
Proof. From Rem. 6.3 any nonconstant path $\gamma:[a, b] \rightarrow S^{1}$ is equivalent (Def. 6.4) to a product of paths with nonzero derivative. Each such is then image equivalent to $\left(\gamma^{*}\right)_{\mid[r, s]}^{\epsilon}$ for some $r<s$ and $\epsilon \in\{ \pm 1\}$. So, we can write the path as $\prod_{i=1}^{\ell}\left(\left.\gamma^{*}\right|_{\left[r_{i}, s_{i}\right]} ^{\epsilon_{i}}\right.$ with $s_{i}=r_{i+1}$. Suppose $\epsilon_{i}$ and $\epsilon_{i+1}$ have opposite sign. Further subdivide one of paths corresponding to $i$ or to $i+1$ to assume $\left[r_{i}, s_{i}\right]$ and $\left[r_{i+1}, s_{i+1}\right]$ have the same length. From (6.4),

$$
\left(\gamma^{*}\right)_{\left.\right|_{\left[r_{i}, s_{i}\right]} ^{\epsilon_{i}}}^{\epsilon_{i}}\left(\gamma^{*}\right)_{\left.\right|_{\left[r_{i+1}, s_{i+1}\right]} ^{\epsilon_{i+1}}}
$$

is equivalent to the constant path with image $\left(\gamma^{*}\right)^{\epsilon_{i}}\left(r_{i}\right)$ [9.12a]. Thus the whole path is equivalent to a path with a smaller $\ell$. An induction on the integer $\sum_{i=1}^{\ell}\left|\epsilon_{i+1}-\epsilon_{i}\right|$ shows $\gamma$ is equivalent to $\gamma_{n}^{*}$ for some integer $n$.

The proof is complete if $\gamma_{n}^{*}$ is inequivalent to $\gamma_{m}^{*}$ for $m \neq n$. Decompose $\gamma:[a, b] \rightarrow S^{1}$ into its real and imaginary parts: $\gamma=\gamma_{1}+i \gamma_{2}$ where $\gamma_{i}:[a, b] \rightarrow \mathbb{R}$, $i=1,2$. Define $\operatorname{deg}(\gamma)$ through the formula

$$
\begin{aligned}
2 \pi i \operatorname{deg}(\gamma)= & \int_{a}^{b}( \\
& \left.\gamma_{1}(t), \gamma_{2}(t)\right) \cdot\left(\frac{d \gamma_{1}}{d t}(t), \frac{d \gamma_{2}}{d t}(t)\right) d t \\
& +i \int_{a}^{b}\left(-\gamma_{2}(t), \gamma_{1}(t)\right) \cdot\left(\frac{d \gamma_{1}}{d t}(t), \frac{d \gamma_{2}}{d t}(t)\right) d t
\end{aligned}
$$

(as in Chap. 2 Lem. 2.3). By direct computation $\operatorname{deg}\left(\gamma_{n}^{*}\right)=n$.
If $\gamma$ is homotopic to $\gamma_{n}^{*}$, then Chap. 2 Lem. 2.3 shows $\operatorname{deg}(\gamma)=n$. As $\operatorname{deg}(\gamma)$ depends only on $[\gamma][9.12 \mathrm{~d}]$, $\left[\gamma_{n}^{*}\right]$ is distinct from $\left[\gamma_{m}^{*}\right]$ for $n \neq m$.

Chap. 4 computes fundamental groups of many spaces from Thm. 6.14.
Let $\gamma:[a, b] \rightarrow X_{1}$ be a (simplicial) path. Consider $f \circ \gamma:[a, b] \rightarrow X_{2}$, and for $x_{1} \in X_{1}$, denote $f\left(x_{1}\right)$ by $x_{2}$. For $[\gamma] \in \pi_{1}\left(X_{1}, x_{1}\right)$, $[f \circ \gamma] \in \pi_{1}\left(X_{2}, f\left(x_{1}\right)\right)$ is independent of the choice of $\gamma$ representing $[\gamma]$. To a product of paths $\gamma_{1} \gamma_{2}$ in $X_{1}$, apply the formula $f \circ\left(\gamma_{1} \gamma_{2}\right)=\left(f \circ \gamma_{1}\right)\left(f \circ \gamma_{2}\right)$. This shows $\left[f \circ \gamma_{1}\right]\left[f \circ \gamma_{2}\right]=\left[f \circ\left(\gamma_{1} \gamma_{2}\right)\right]$.

Lemma 6.15. Conclude: $f$ induces a homomorphism of groups

$$
f_{*}: \pi_{1}\left(X_{1}, x_{1}\right) \rightarrow \pi_{1}\left(X_{2}, x_{2}\right)
$$

If $f$ is one-one and onto then $f_{*}$ is an isomorphism of groups.
Example 6.16. Let $X_{1}=X_{2}=S^{1}$ and consider $\cos (2 \pi t)+i \sin (2 \pi t)=z(t)$. For a fixed positive integer $n$ define a function $f$ by the formula $f(z(t))=z(n t)=$ $\cos (2 \pi n t)+i \sin (2 \pi n t)$. Thus $f_{*}: \pi_{1}\left(S^{1}, 1\right) \rightarrow \pi_{1}\left(S^{1}, 1\right)$. Also, for $\gamma_{1}^{*}$, the generating path for $\pi_{1}\left(S^{1}, 1\right), f \circ \gamma_{1}^{*}(t)=f(z(t))$. Therefore $f \circ \gamma_{1}^{*}$ is image equivalent to $\gamma_{n}^{*}$. Identify $\pi_{1}\left(S^{1}, 1\right)$ with $\mathbb{Z}$, the group of integers, by identifying the integer 1 with [ $\gamma_{1}^{*}$ ]. Then, $f_{*}: \pi_{1}\left(S^{1}, 1\right) \rightarrow \pi_{1}\left(S^{1}, 1\right)$ sends the integer $m$ to $f_{*}(m)=n m$. The image of $f_{*}$ is the subgroup of $\pi_{1}\left(S^{1}, 1\right)=\mathbb{Z}$ that $n$ generates.
6.5. Fundamental group of a product. Let $\left(X, x_{0}\right)$ and $\left(Y, y_{0}\right)$ be two differentiable manifolds with a base point. The projections onto each factor, $\mathrm{pr}_{X}$ : $X \times Y \rightarrow X$ and $\mathrm{pr}_{Y}: X \times Y \rightarrow Y$, induce homomorphisms
$\operatorname{pr}_{X_{*}}: \pi_{1}\left(X \times Y,\left(x_{0}, y_{0}\right)\right) \rightarrow \pi_{1}\left(X, x_{0}\right)$ and $\operatorname{pr}_{Y *}: \pi_{1}\left(X \times Y,\left(x_{0}, y_{0}\right)\right) \rightarrow \pi_{1}\left(Y, y_{0}\right)$.
So, there is a homomorphism

$$
\begin{equation*}
\left(\operatorname{pr}_{X *}, \operatorname{pr}_{Y *}\right): \pi_{1}\left(X \times Y,\left(x_{0}, y_{0}\right)\right) \rightarrow \pi_{1}\left(X, x_{0}\right) \times \pi_{1}\left(Y, y_{0}\right) \tag{6.5}
\end{equation*}
$$

The right side is the product group with factors $\pi_{1}\left(X, x_{0}\right)$ and $\pi_{1}\left(Y, y_{0}\right)$.
Theorem 6.17. $\pi_{1}\left(X \times Y,\left(x_{0}, y_{0}\right)\right)$ and $\pi_{1}\left(X, x_{0}\right) \times \pi_{1}\left(Y, y_{0}\right)$ are isomorphic.
Proof. Let $f^{X}$ (resp. $f^{Y}$ ) map $X \rightarrow X \times Y$ by $f^{X}(x)=\left(x, y_{0}\right)$ (resp. map $Y \rightarrow X \times Y$ by $\left.f^{Y}(y)=\left(x_{0}, y\right)\right)$. For $\gamma:[a, b] \rightarrow X \times Y$ consider the paths $\left(f^{X} \circ \operatorname{pr}_{X} \circ \gamma\right)=\psi^{X}:[a, b] \rightarrow X \times Y$ and $\left(f^{Y} \circ \operatorname{pr}_{Y} \circ \gamma\right)=\psi^{Y}:[a, b] \rightarrow X \times Y$.

We show the map taking $\left(\left[\gamma_{1}\right],\left[\gamma_{2}\right]\right) \in \pi_{1}\left(X, x_{0}\right) \times \pi_{1}\left(Y, y_{0}\right)$ to $f_{*}^{X}\left(\left[\gamma_{1}\right]\right) f_{*}^{Y}\left(\left[\gamma_{2}\right]\right)$ in $\pi_{1}\left(X \times Y,\left(x_{0}, y_{0}\right)\right)$ is inverse to $\left(\operatorname{pr}_{X *}, \operatorname{pr}_{Y *}\right)$. This only requires showing $\gamma$ is equivalent to $\psi^{X} \psi^{Y}$. Fig. 6.5 illustrates this when $X=Y=S^{1}$ and $X \times Y$ is the complex torus of Fig. 3 with $\omega_{1}=1$ and $\omega_{2}=i[9.5 \mathrm{~b}]$.

Figure 6. The diagonal recomposes itself


Write $\gamma(t)=\left(\gamma^{X}(t), \gamma^{Y}(t)\right)$ for $t \in[a, b]$ and assume $[a, b]=[0,1]$. Then $\gamma$ is image equivalent to the path $\left(\gamma^{X}\left(\frac{t}{2}\right), \gamma^{Y}\left(\frac{t}{2}\right)\right)$ for $t \in[0,2]$. Also, $\psi^{X}$ is the path $t \mapsto\left(\gamma^{X}(t), y_{0}\right)$ for $t \in[0,1]$ and $\left(x_{0}, \gamma^{Y}(t-1)\right)$ for $t \in[1,2]$. Here is a homotopy between these paths running over $s \in[0,1]$ :

$$
\gamma_{s}(t)= \begin{cases}\left(\gamma^{X}\left(\frac{t}{2-s}\right), y_{0}\right) & \text { for } t \in[0, s] \\ \left.\gamma^{X}\left(\frac{t}{2-s}\right), \gamma^{Y}\left(\frac{t-s}{2-s}\right)\right) & \text { for } t \in[s, 2-s] \\ \left(x_{0}, \gamma^{Y}\left(\frac{t-s}{2-s}\right)\right) & \text { for } t \in[2-s, 2]\end{cases}
$$

Example 6.18 (Continuation of $\S 3.2 .2$ ). Here $X^{i}=\mathbb{C} / L\left(\omega_{1}^{i}, \omega_{2}^{i}\right)$ is

$$
\left\{t_{1} \omega_{1}^{i}+t_{2} \omega_{2}^{i} \mid 0 \leq t_{i}<1, i=1,2\right\}
$$

where $\omega_{1}^{i} / \omega_{2}^{i} \in \mathbb{C} \backslash \mathbb{R}, i=1,2$. For the lattice $\left\{m_{1} \omega_{1}^{i}+m_{2} \omega_{2}^{i} \mid m_{1}, m_{2} \in \mathbb{Z}\right\}$ use the letter $L_{i}, i=1,2$. For $z \in \mathbb{C}$, there is a unique $\omega \in L_{i}$ with $z-\omega \in X^{i}$. Then $z-\omega$ represents the coset $z \bmod L_{i} \stackrel{\text { def }}{=}\left\{z+u \mid u \in L_{i}\right\}($ as in $\S 7.1)$. Let $\pi^{i}: \mathbb{C} \rightarrow \mathbb{C} / L_{i}$ be the map that takes $z$ to $z \bmod L_{i}$. Then $\pi^{i}$ is an analytic map. It becomes a homomorphism of groups if we make $X^{i}$ into a group using this addition formula:

$$
z_{1} \bmod L_{i}+z_{2} \bmod L_{i} \stackrel{\text { def }}{=} z_{1}+z_{2} \bmod L_{i}[9.9 d]
$$

Suppose $L_{1} \subseteq L_{2}$. Then, for $z \in \mathbb{C}$, the set

$$
\left(\pi^{1}\right)^{-1}\left(z \bmod L_{1}\right)=\left\{z+\omega \mid \omega \in L_{1}\right\}
$$

is in $\left(\pi^{2}\right)^{-1}\left(z \bmod L_{2}\right)$. So, the map $f$ taking $z \bmod L_{1}$ to $z \bmod L_{2}$ depends only on $z \bmod L_{1}$, not on $z$. Identify $\pi_{1}\left(X^{i}, 0\right)$ with $L_{i}($ as in $[9.9 \mathrm{~g}])$. The induced map $f_{*}$ is the inclusion $L_{1}$ into $L_{2}$. For each $x_{2} \in X_{2}$ the cardinality of the set $f^{-1}\left(x_{2}\right)$ is the order of the quotient group $L_{2} / L_{1}[9.7 \mathrm{~d}]$.
Note: These concepts work equally well for finite unions of manifolds.

## 7. Permutation representations and covers

Two types of group theory arise in analyzing algebraic functions from Riemann's viewpoint. One is the presentation of fundamental groups, as free groups on generators with relations. Elementary examples of that do appear in many topology books (here too, starting with Chap. $4 \S 1.1$ ). The second type is less common: Analyzing homomorphisms of fundamental groups into other groups. Motivating problems and sufficient group theory show how finite and profinite group theory apply to the study of moduli of Riemann surfaces. The group theory starts with permutation representations and their associated group representations.
7.1. Permutation representations. Denote by $\{\boldsymbol{x}\}=\left\{x_{1}, \ldots, x_{n}\right\}$ any set of $n$ distinct elements. Let $S_{n}$ be the collection of permutations of $\{\boldsymbol{x}\}$, and regard $S_{n}$ as a group in the usual way. Multiplication of permutations corresponds to functional composition of maps on $\{\boldsymbol{x}\}$. Reminder: As the introduction states, we typically act with $S_{n}$ on the right of elements from $\boldsymbol{x}$, though sometimes the presence of a second action forces us to act on the left.
7.1.1. Permutation notation and actions. Denote the identity element of $S_{n}$ by 1. Here is an inefficient, though clear way to express the effect of $g \in S_{n}$ :

$$
\left(\begin{array}{cccc}
1 & 2 & \cdots & n \\
(1) g & (2) g & \cdots & (n) g
\end{array}\right)
$$

where $k=(j) g$ is the integer subscript of the image of $x_{j}$ under $g$.
Example 7.1. Suppose $n=16$, and the display of $g$ is

$$
\left(\begin{array}{cccccccccccccccc}
1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 & 10 & 11 & 12 & 13 & 14 & 15 & 16 \\
16 & 12 & 9 & 8 & 1 & 3 & 2 & 5 & 6 & 10 & 11 & 7 & 4 & 13 & 14 & 15
\end{array}\right) .
$$

The notation indicates $g$ maps $x_{9}$ to $x_{6}$. Disjoint cycle notation for $g$ represents it as a product of disjoint cycles of integers. It requires fewer symbols than the
complete permutation notation. Also, it shortens computations in $S_{n}$ by parsing the group action into memorable pieces. The disjoint cycle representation for $g$ :
$(116151413485)(2127)(963)$.
The order of the disjoint cycles is unimportant; $(i) g$ goes to the right of $i$. That is, $(1) g=16$ is right of 1 , and the cycle closes at 5 because (5) $g$ is 1 , back to the beginning. Exclude cycles of length $1((10) g=10$ gives a cycle (10)) for efficiency. An element of $S_{n}$ is a $k$-cycle, $k>1$ if it has one and only one cycle - of length $k$ - of length bigger than 1 .

For another unique, less orthodox way to write permutations see [9.17a].
Let $G$ be any group. A degree $n$ permutation representation of $G$ is a homomorphism $T: G \rightarrow S_{n}$. Such a $T$ is the same as giving an action of $G$ on the set $S=\left\{x_{1}, \ldots, x_{n}\right\}$.

With $G$ a group and $S$ a set, a right action is a function: $A=A_{R}: S \times G \rightarrow S$ : $A(s, g) \mapsto(s) g$ with two action properties:
(7.1a) $(s) g_{1} g_{2}=\left((s) g_{1}\right) g_{2}$ for $s \in S, g_{1}, g_{2} \in G$. Using $A$ we would write this

$$
A\left(A\left(s, g_{1}\right), g_{2}\right)=A\left(s, g_{1} g_{2}\right)
$$

(7.1b) $(s) 1_{G}=s$ for $s \in S$ (the identity in $G$ leaves $s \in S$ fixed).

A left action is from a function $A_{L}: G \times S \rightarrow S$ with the action composite

$$
A_{L}\left(g_{1}, A_{L}\left(g_{2}, s\right)\right)=A_{L}\left(g_{1} g_{2}, s\right)
$$

An orbit of an action is the range of the set $s \times G$, under $A$, for some $s \in S$. The kernel of the action $\operatorname{ker}(A)$ consists of those $g \in G$ that act like the identity on $S$. The most important example is where $G$ acts on the right cosets of a subgroup $H$ of $G$. The set $H g=\{h g\}_{h \in H}$ is a right coset of $H$ in $G$. Two right cosets $H g$ and $H g^{\prime}$ are either equal or have no elements in common. Assume there are exactly $n$ distinct right cosets of $H$ in $G: H, H g_{2}, \ldots, H g_{n}$. Call $n$ the index $(G: H)$ of $H$ in $G$. Finding good representatives for cosets is an art (try [9.17c]).

The archetype of a right action: $A:\left(H g^{\prime}, g\right) \mapsto H g^{\prime} g$, or $g \in G$ maps a right coset $H g^{\prime}$ to $\left(H g^{\prime}\right) g=H g^{\prime} g$. For any subgroup $H$ there is both a set of right cosets of $H$ and a set of left cosets of $H$. Only if $H$ is normal in $G$ are all right cosets also left cosets. The map $\left(g, g^{\prime} H\right) \mapsto g g^{\prime} H$ is a left action on left cosets. There are further actions of groups in [9.16]. We emphasize a right action because this is the natural action of fundamental groups acting on points as in Lem. 7.13.

Definition 7.2. Suppose $G$ is a group with a normal subgroup $H$ and another subgroup $W$. Assume $\langle H, W\rangle=G$ and $H \cap W=\{1\}$. We say $G$ is the semi-direct product of $H$ and $W$, written $H \times{ }^{s} W$.

If $G=H \times{ }^{s} W$, then elements of $G$ act as automorphisms of $H$ by conjugation. This is an action $A$ : For $g \in G, A(g): h \in H \mapsto g^{-1} h g \stackrel{\text { def }}{=} h^{g}$. This is a right action. The following lemma, in a left or right action form is in almost all graduate texts in algebra.

Lemma 7.3. Each element of $H \times^{s} W$ has a unique expression as $w h, h \in H$, and $w \in W$. Suppose $A: W \rightarrow \operatorname{Aut}(H)$ is a homomorphism giving a right action of $W$ on $H$. Then, there is a group $G$ given as a semi-direct product of $H$ and $W$. Multiplication in this group satisfies the formula $w_{1} h_{1} w_{2} h_{2}=w_{1} w_{2}\left(h_{1}\right) A\left(w_{2}\right) h_{2}$.

REMARK 7.4 (Affine action). There is a memorable notation for multiplication by imitating matrix multiplication of lower triangular $2 \times 2$ matrices. Associate $w_{i} h_{i}$ with $\left(\begin{array}{cc}w_{i} & 0 \\ h_{i} & 1\end{array}\right), i=1,2$. Then, the multiplication in $H \times{ }^{s} W$ imitates an expected matrix calculation:

$$
\left(\begin{array}{cc}
w_{1} & 0 \\
h_{1} & 1
\end{array}\right)\left(\begin{array}{cc}
w_{2} & 0 \\
h_{2} & 1
\end{array}\right)=\left(\begin{array}{cc}
w_{1} w_{2} & 0 \\
\left(h_{1}\right) A\left(w_{2}\right) h_{2} & 1
\end{array}\right) .
$$

Further, $H \times{ }^{s} W$ acts as permutations of $H$. For its matrix form, replace $h^{\prime} \in H$ by the vector $\left(h^{\prime}, 1\right):\left(h^{\prime}, 1\right)\left(\begin{array}{cc}w & 0 \\ h & 1\end{array}\right)=\left(\left(h^{\prime}\right) A(w) h, 1\right)$ or $h^{\prime} \mapsto\left(h^{\prime}\right) A(w) h$. The left action version with upper triangular matrices has a little glitch in it, unless $H$ is an abelian group. That, however, comes up often in important examples (see [9.18]).
7.1.2. Transitive and intransitive representations. We discuss concepts that use coset representations. Lem. 7.7 shows how to go from the definition of action to the language of homomorphisms. When using groups acting on manifolds we often translate from actions into representations.

Definition 7.5. The right coset representation $T_{H}: G \rightarrow S_{n}$, defined by the subgroup $H \leq G$, comes from the formula
(7.2) for $g \in G, i \in\{1,2, \ldots, n\},(i) T_{H}(g)=j$ with $H g_{j}$ the right coset equal to $H g_{i} g$.
Denote the subgroup of elements $g \in G$ for which $T(g)$ fixes the integer $j$ by $G(T, j)=G(j)$. For $T$ a permutation representation, $\operatorname{ker}(T)$ is $\left\{g \in G \mid T(g)=1_{G}\right\}$, the kernel of the action of $G$. Call $T$ faithful if $\operatorname{ker}(T)$ consists only of $1_{G}$. Also, $T$ is transitive ( $G$ under $T$ has one orbit) if for each $i \in\{1,2, \ldots, n\}$, there is $g_{i} \in G$ with (1) $T\left(g_{i}\right)=i$. Then, $G(1) g_{i}$ is the set of $g \in G$ taking 1 to $i$. By definition, $\operatorname{ker}(T)$ is $\bigcap_{i=1}^{n} G(i)$. Assume $T$ is transitive and (1)T( $\left.g_{i}\right)=i, i=1, \ldots, n$. Then, $g_{i}^{-1} G(1) g_{i}$, the conjugate of each element of $G(1)$ by $g_{i}$, equals $G(i)$. So, $G(1) \ldots, G(n)$ is a complete list of conjugates of $G(1)$ in $S_{n}$.

Definition 7.6. Let $T_{i}$ be a degree $n$ permutation representation of $G, i=1,2$. Suppose there is $h \in S_{n}$ with $h^{-1} T_{1}(g) h=T_{2}(g)$ for each $g \in G$. Then $T_{1}$ is permutation equivalent to $T_{2}: T_{1}$ and $T_{2}$ are equivalent as permutation representations.

LEMMA 7.7. In notation above, $G$ acts on (right) cosets of $H \leq G$, permuting them, and $T_{H}: G \rightarrow S_{n}$ is a homomorphism. The kernel is those $g \in G$ that fix each coset. This is the same as the elements of $\cap_{g \in G} g^{-1} H g$. Reordering cosets of $H$ in $G$ changes the representation $T_{H}$ only up to permutation equivalence.

Suppose $A_{S}$ (resp. $A_{S^{\prime}}$ ) is an action of $G$ on $S$ (resp. $S^{\prime}$ ) with $S$ and $S^{\prime}$ disjoint sets. Then, there is an action of $G$ on $S \times S^{\prime}$, the direct product action: $A \times A^{\prime}$ : $\left(S \times S^{\prime}\right) \times G \rightarrow S \times S^{\prime}$ by $g \in G:\left(s, s^{\prime}\right) \in S \times S^{\prime} \mapsto\left((s) g,\left(s^{\prime}\right) g\right)$. There is also an action of $G$ on $S \cup S^{\prime}$, the direct sum action: $A \oplus A^{\prime}:\left(S \dot{\cup} S^{\prime}\right) \times G \rightarrow S \dot{\cup} S^{\prime}$ by $g \in G: s \in S \dot{\cup} S^{\prime} \mapsto(s) g$ given by $A$ if $s \in S$, and by $A^{\prime}$ if $s \in S^{\prime}$. For $T: G \rightarrow S_{n}$ an arbitrary permutation representation, partition $\{1, \ldots, n\}$ into a disjoint union $X_{1} \cup X_{2} \cup \cdots X_{t}$ of the $G$ orbits. Suppose $n_{i}=\left|X_{i}\right|, i=1, \ldots, n$.

Theorem 7.8. Let $T_{H}: G \rightarrow S_{n}$ be the right coset representation associated to the subgroup $H$ of $G$. Then $T_{H}$ is a transitive representation with $\operatorname{ker}\left(T_{H}\right)$ equal to $\bigcap_{g \in G} g^{-1} H g$. Conversely, if $T: G \rightarrow S_{n}$ is a transitive representation of $G$, then $T$ is permutation equivalent to $T_{H}$ with $H=G(1)$. Generally, in the notation
above for $T, T=\oplus_{i=1}^{t} T_{i}: G \rightarrow \oplus_{i=1}^{t} S_{n_{i}}$ presents $T$ as the direct sum of right coset representations corresponding to subgroups of $G$.

Proof. For each $i \in\{1,2, \ldots, n\}$, formula (7.2) shows (1) $T_{H}\left(\sigma_{i}\right)=i$. So $T_{H}$ is transitive. The subgroup $\operatorname{ker}\left(T_{H}\right)$ consists of the $g^{\prime} \in G$ such that $H g_{i} g^{\prime}=H g_{i}$, $i=1, \ldots, n: g^{\prime} \in g_{i}^{-1} H g_{i}, i=1, \ldots, n$. Each element in $G$ has the form $h g_{i}$ for some $h \in H$ and $i \in\{1,2, \ldots, n\}$. So, $g^{\prime} \in \operatorname{ker}\left(T_{H}\right)$ if and only if $g^{\prime} \in \bigcap_{g \in G} g^{-1} H g$.

Let $T: G \rightarrow S_{n}$ be an arbitrary transitive permutation representation. Choose $g_{1}, \ldots, g_{n}$ so that (1)T( $\left.g_{i}\right)=i, i=1, \ldots, n$. Thus, the cosets $G(1) g_{1}, \ldots, G(1) g_{n}$ are distinct. Conclude that (7.2), with $G(1)$ replacing $H$, gives $T_{G(1)}$. As

$$
\{g \in G \mid(i) T(g)=j\}=g_{j}^{-1} G(1) g_{i}
$$

${ }_{(i)} T(g)=j$ exactly if $(i) T_{G(1)}(g)=j$. This means $T_{G(1)}$ and $T$ are the same permutation representation. We made choices in selecting the $g_{j} \mathrm{~s}$. So, independent of choices, the representations are permutation equivalent.

Now suppose the representation is not transitive. Since the orbits are all distinct, there is a natural map from the representation to the direct sum representation on the collection of orbits.
7.1.3. Primitive representations and equivariant maps. A subgroup $H \leq G$ is normal if $g^{-1} \mathrm{Hg}=H$ for each $g \in G$. Only then is the set of pairwise products $H g H g^{\prime}$ of two cosets a single coset, equal to $H g g^{\prime}$. So, the cosets have a natural group multiplication. Denote this set by $G / H$ : Each element $\bar{g}=g \bmod H \in G / H$ denotes the coset $H g$. For $H$ any subgroup of $G$, the normalizer of $H$ in $G$ is $N_{G}(H)=\left\{g \in G \mid g^{-1} H g=H\right\}$. Similarly, define the centralizer of $H$ in $G$ :

$$
\operatorname{Cen}_{G}(H)=\left\{g \in G \mid g^{-1} h g=h \text { for each } h \in H\right\}[9.15] .
$$

Definition 7.9. Consider a transitive permutation representation $T: G \rightarrow S_{n}$ of $G$. Call $T$ primitive if there are no groups properly between $G(1)$ and $G$. Let $G(1)$ be the subgroup of $G$ that fixes 1 . If $T$ is transitive, then it is $T$ is doubly transitive if for each $j \in\{2, \ldots, n\}$ there is a $g \in G(1)$ with $(2) T(g)=j: G(1)$ is transitive on $\{2, \ldots, n\}$.

When the notation shows $G$ is in $S_{n}$, we drop the $T$ notation for permutation representations. The transitivity formula for a chain of subgroups $K \leq H \leq G$ says that $(G: K)=(G: H)(H: K)$.

Lemma 7.10. Doubly transitive permutation representations are primitive.
Proof. Suppose $G \leq S_{n}$ is doubly transitive. Let $H$ be a subgroup of $G$ properly containing $G(1)$. Choose $h \in H \backslash G(1)$. Then (1) $h=j \in\{2, \ldots, n\}$. For any $j^{\prime} \in\{2, \ldots, n\}$, use double transitivity to produce $g^{\prime}$ with $(1) g^{\prime}=1$ and $(j) g^{\prime}=j^{\prime}: h g^{\prime} \in H$ takes 1 to $j^{\prime}$. So, the number of cosets of $G(1)$ in $H$ is the same as the number of cosets of $G(1)$ in $G$. Apply the transitivity formula to the chain $G(1)<H \leq G$ to conclude the index of $H$ in $G$ is 1 and $T$ is primitive.

Assume group $G$ acts on two sets: It has an action $A_{S}$ (resp. $A_{S^{\prime}}$ ) on $S$ (resp. $S^{\prime}$ ) with $S$ and $S^{\prime}$ related by a function $f: S \rightarrow S^{\prime}$. We say $f$ commutes with (is equivariant for) these actions if $f\left((s, g) A_{S}\right)=(f(s), g) A_{S^{\prime}}$ for $s \in S, g \in G$.

Example 7.11 (Compatible permutation representations). For $G$ a group and $M$ a normal subgroup, let $u_{M}: G \rightarrow G / M$ be the natural homomorphism with kernel $H$. Suppose $H_{1}$ is a subgroup of $G$ and $H_{2}$ is a subgroup of $G / M$ for which
$f_{M}\left(H_{1}\right) \leq H_{2}$. Then $u_{M}$ induces a map $f_{M}:\left\{H_{1} g \mid g \in G\right\} \rightarrow\left\{H_{2} g \mid g \in G\right\}$. This map commutes with $G$ acting on the cosets of $H_{1}$ and on the cosets of $H_{2}$.
7.1.4. Representations from permutation representations. [9.18] gives many examples of primitive groups that are not doubly transitive. For $g \in G$, some authors abuse notation to write $T(g)=\left(s_{1}\right) \cdots\left(s_{t}\right)$ where $s_{1}, \cdots, s_{t}$ are the integer lengths of the disjoint cycles of $T(g)$ (we usually omit cycles of length one) to indicate a cycle type (conjugacy class) in $S_{n}$. Denote the count of length one cycles in $T(g)$ by $\mathrm{t}(T(g))$, the trace of $T(g)$. For example, the permutation example of $\S 7.1 .1$ has trace 2 and its cube has trace 5 . We remind why $T(g)$ it is a trace.

Regard the formal symbols $\left\{x_{1}, \ldots, x_{n}\right\}$ as basis vectors for a vector space $V$ over a field $F$. Then each permutation $g \in S_{n}$ extends linearly to act on $V$. That is, applying $g \in G$ to $v=\sum_{i=1}^{n} a_{i} x_{i} \in V$ gives $\sum_{i=1}^{n} a_{i} x_{(i) g}$. Write the result of $g$ on $x_{i}$ to be $\sum_{j=1}^{n} a_{i, j} x_{i}$ with coefficients denoting what would appear in the $i$ th position of a matrix $M_{g}$ acting on the right of (row) vectors. When $F$ has characteristic 0 , the matrix $M_{g}$ has trace $\sum_{i=1}^{n} a_{i, i}$, the count of the number of $x_{i}$ s that $g$ fixes. In each row and column the matrix $M_{g}$ has exactly one non-zero entry and that is a 1 . So, $M_{g}$ is an element of the orthogonal group $O_{n}: M_{g}$ times its transpose is the identity matrix. The determinant function is multiplicative on $n \times n$ matrices. Conclude that $M_{g}$ has determinant $\operatorname{Det}\left(M_{g}\right) \stackrel{\text { def }}{=} \operatorname{Det}(g)$ equal to $\pm 1$. When the field $F$ has characteristic $p$, the count of the integers fixed by $g$ is the trace $\bmod p$. We may revert, when acting with matrices to a traditional left-hand action.

The result is that a degree $n$ permutation representation $T$ of a group $G$ produces a homomorphism $\rho_{T}: G \rightarrow \mathrm{GL}_{n}(F)$. If $T$ is a faithful permutation representation, then $\rho_{T}$ is a faithful group representation: Its kernel is trivial. Any homomorphism $\rho: G \rightarrow \mathrm{GL}_{n}(F)$ is called a representation of $G$ over the field $F$. With $V=F^{n}$, we often write $V_{T}$ to indicate we mean $V$ with the action through $T$. Then, for any representation, extend this notation to use $V_{\rho}$. In fact, group theory doesn't restrict to just finite dimensional representations, though we will. Most situations regard permutation representations as the same if they are equivalent. If $M \in \mathrm{GL}_{n}(F)$, then the two permutation representations $g \mapsto \rho(g)$ and $g \mapsto M^{-1} \rho(g) M$ are (representation) equivalent. Though two permutation representations may be inequivalent, their corresponding representations might be equivalent (§8.6.2 and [9.20]).

The group representation attached to the sum of permutation representations is the action on the direct sum of the vector spaces. When $F$ has characteristic 0 , every permutation representation of degree exceeding 1 is the direct sum of the identity representation and another representation. These are the only summands if and only if the permutation representation is doubly transitive [9.19d]. Further, the group representation of the direct product of two permutation representations is their tensor product; the trace is the product of the constituent traces [9.19a]. The group ring of $G$ over $F$ has the notation $F[G]$. The product of $\sum_{g \in G} a_{g} g$ and $\sum_{g \in G} b_{g} g$ (with $a_{g}, b_{g} \in F$ ) is given by convolution: $\sum_{g \in G} c_{g} g$ with $c_{g}=$ $\sum_{h \in G} a_{h} b_{h^{-1} g}, g \in G$. A representation $\rho$ then produces a homomorphism of associative rings: $\sum_{g \in G} a_{g} g \mapsto \sum_{g \in G} a_{g} \rho(g) \in \mathbb{M}_{\operatorname{deg}(\rho)}(F)$. Call an idempotent $I$ in this ring $G$ invariant if it commutes with multiplication by elements of $G$. That means the range of $I$ is a $G$ invariant space: $I$ is a $G$ invariant projection [9.19h].
7.2. Covering spaces. Let $X$ and $Y$ be differentiable (resp. analytic) manifolds. Assume $f: Y \rightarrow X$ is a differentiable (resp. analytic) map. We will often use that if $f$ is one-one, and onto in a neighborhood of a point, then it has a differentiable (resp. analytic) inverse (Lem. 4.2). Suppose $\varphi: X \rightarrow X^{\prime}$ is any map between spaces, and $x_{0}$ maps to $x_{0}^{\prime}$ under $\varphi$. As in Lem. 6.15, this induces a homomorphism on fundamental groups $\varphi_{*}: \pi_{1}\left(X, x_{0}\right) \rightarrow \pi_{1}\left(X^{\prime}, x_{0}^{\prime}\right)$ by mapping a closed path $\gamma:[a, b] \rightarrow X$ to $[\varphi \circ \gamma] \in \pi_{1}\left(X^{\prime}, x_{0}^{\prime}\right)$. This makes sense because composing with $\varphi$ preserves homotopy classes of paths into $X$. Though obvious, it doesn't trivialize computing the image of $\pi_{1}\left(X, x_{0}\right)$ under $\varphi_{*}$.

Definition 7.12 (Covering space). The pair $(Y, f)$ (or just $Y$ if there is no confusion) is a covering space (or cover) of $X$ if each point $x \in X$ has a connected neighborhood (Chap. $2 \S 2.2 .2$ ) $U_{x}$ with this property: for each connected component $V$ of $f^{-1}\left(U_{x}\right)$, restricting $f$ to $V$ is a one-one and onto map $V \rightarrow U_{x}$.
7.2.1. Degree of a cover. Assume $X$ is connected, and $f: Y \rightarrow X$ is a cover. Then, the cardinality of the fibers $\left|f^{-1}(x)\right|, x \in X$, being locally constant, must actually be constant. This is the degree $\operatorname{deg}(f)$ of $f$. We say $(Y, f)$ is finite, or that $f$ is a finite cover if $\operatorname{deg}(f)<\infty$.

Two covers $f_{i}: Y_{i} \rightarrow X, i=1,2$ are equivalent (as covers of $X$ ) if there is a one-one and onto continuous map $\psi: Y_{1} \rightarrow Y_{2}$ with $f_{2} \circ \psi=f_{1}$ [9.21]. Note: For any covering space $(Y, f)$ of $X, U$ an open subset of $X$, and $V$ a union of connected components of $f^{-1}(U)$, the restriction of $f$ to $V$ gives a cover $\left(V, f_{\left.\right|_{V}}\right)$ of $U$.

A framework for considering equivalence classes of finite covers of a manifold $X$ is the goal remaining to this subsection. This immediately reduces to considering connected finite covers $(Y, f)$; we assume $Y$ is a connected space. The classification hinges on producing an equivalence class, $T(Y, f)$, of permutation representations (§7.1) from an equivalence classes of covers $(Y, f)$. We do that now.

Note: Covers in this section are what topologists call covers. In algebraic geometry the word cover includes complex analytic maps of manifolds having some fibers that ramify (their cardinality is smaller than the degree). The phrase then includes, for example, any nonconstant analytic map $f: Y \rightarrow \mathbb{P}_{z}^{1}$, with $Y$ a compact Riemann surface and $\operatorname{deg}(f) \geq 2$. As the fundamental group of $\mathbb{P}_{z}^{1}$ is trivial, such an $f$ must ramify (Chap. 4 Thm. 1.8). By the end of Chap. 4, a cover will include any surjective analytic map between compact complex manifolds with finite (point sets in their) fibers. Reference back to this chapter will speak of the unramified covers corresponding to subgroups of fundamental groups as in Thm. 7.16.
7.2.2. Covers and permutation representations. Let $f: Y \rightarrow X$ be a cover with $\gamma:[a, b] \rightarrow X$ a path having initial point $x_{0}$ and end point $x_{1}$.

Lemma 7.13 (Action of path lifting). For $y^{\prime} \in Y$ with $f\left(y^{\prime}\right)=x_{0}$, there is a unique path $\tilde{\gamma}:[a, b] \rightarrow Y$ with $f \circ \tilde{\gamma}=\gamma$ : the lift of $\gamma$ with initial point $y^{\prime}$.

So, $\gamma$ produces a unique map $\gamma_{*}: f^{-1}\left(x_{0}\right) \rightarrow f^{-1}\left(x_{1}\right)$ depending only on the image of $\gamma$ in $\pi_{1}\left(X, x_{0}, x_{1}\right)$. In particular, consider paths $\gamma_{i}:\left[a_{i}, b_{i}\right] \rightarrow X, i=1,2$, with $\gamma_{1}\left(b_{1}\right)=\gamma_{2}\left(a_{1}\right)$ and $\gamma_{1}\left(a_{1}\right), \gamma_{2}\left(b_{1}\right), \gamma_{2}\left(b_{2}\right)$ respectively $x_{0}, x_{1}, x_{2}$. Then, there is a transitivity formula:

$$
\begin{equation*}
\left(\gamma_{1} \cdot \gamma_{2}\right)_{*}=\left(\gamma_{1}\right)_{*} \circ\left(\gamma_{2}\right)_{*}: f^{-1}\left(x_{0}\right) \rightarrow f^{-1}\left(x_{2}\right) \tag{7.3}
\end{equation*}
$$

Proof. Each $\gamma(t)$ has a neighborhood $U_{t}$ with $f$ one-one on the connected components of $f^{-1}\left(U_{t}\right)$. The argument of Chap. $2 \S 3.3 .2$ works here as it did there, by assuming you have extended the path lifting $\tilde{\gamma}$ to an interval $\left[a, t^{\prime}\right]$ with $t^{\prime}<b$.

Let $[r, s]$ be a closed nontrivial interval for which $t^{\prime} \in[r, s]$ and there is neighborhood $U_{t^{\prime}}$ of $\left.\tilde{( } t^{\prime}\right)$ containing $\gamma([r, s])$ with $U^{\prime} \subset f^{-1}\left(U_{t^{\prime}}\right)$ a connected component on which $f$ is one-one and $\gamma^{*}\left(t^{\prime}\right) \in U^{\prime}$. For each $t \in[r, s]$ define $\tilde{\gamma}(t)$ to be the unique point of $U^{\prime}$ lying over $\gamma(t)$. Finish exactly as in Chap. $2 \S 3.3 .2$.

Now considering (7.3) Since the map $\gamma_{*}$ is clearly continuous and varies continuously in a homotopy family, as a map on a finite set, it is a homotopy class invariant. So, $\gamma_{*}$ depends only on the image of $\gamma$ in $\pi_{1}\left(X, x_{0}, x_{1}\right)$. The path $\widetilde{\gamma_{1} \cdot \gamma_{2}}$ starting at $y^{\prime}$ is the same as the path $\tilde{\gamma}_{1} \cdot \tilde{\gamma}_{2}$ where $\tilde{\gamma}_{2}$ is the unique path starting at the end point of $\tilde{\gamma}_{1}$. The formula (7.3) just says the endpoint of both of these paths are the same.

Label the points of $f^{-1}\left(x_{0}\right)$ as $\boldsymbol{y}=\left\{y_{1}, \ldots, y_{n}\right\}$. Consider a path $\gamma:[a, b] \rightarrow X$ based at $x_{0}$. Then, the end point of the lift of $\gamma$ with initial point $y_{j}, j=1, \ldots, n$ associates to $\gamma$ and $\boldsymbol{y}$ a unique labeling of $f^{-1}(\gamma(b))$. A closed path $\gamma$ gives an element of $S_{n}, T_{\boldsymbol{y}}(\gamma)$, as follows:
(7.4) $(i) T_{\boldsymbol{y}}(\gamma)=j$ with $y_{j}$ the end point of the lift of $\gamma$ with initial point $y_{i}$.

For $\gamma_{1}, \gamma_{2} \in \Pi_{1}\left(X, x_{0}\right)$ (closed paths based at $\left.x_{0}\right)(7.3)$ gives

$$
T_{\boldsymbol{y}}\left(\gamma_{1} \gamma_{2}\right)=T_{\boldsymbol{y}}\left(\gamma_{1}\right) T_{\boldsymbol{y}}\left(\gamma_{2}\right)
$$

The right side consists of elements multiplied in $S_{n}$. So, $T_{\boldsymbol{y}}$ defines a permutation representation of $\pi_{1}\left(X, x_{0}\right)$ whose equivalence class we denote by $T(Y, f)$.

In Fig. 7 , for example, $w \mapsto w^{n}=z$ gives the map $f: \mathbb{C}_{w}^{*} \rightarrow \mathbb{C}_{z}^{*}\left(\mathbb{C}^{*}=\mathbb{C} \backslash\{0\}\right)$. A lift of $\gamma$ (a clockwise circle, compatible with our choices in Chap. 4) is $\tilde{\gamma}$ going $\frac{1}{n}$ of the way around a clockwise circle. The associated permutation is an $n$-cycle of $S_{n}$ representing that $\tilde{\gamma}$ goes from the lift $y^{\prime}=2^{1 / n}$ of $\gamma(0)=2$ to $y^{\prime \prime}=2^{\frac{1}{n}} e^{\frac{-2 \pi i}{n}}$, the point on $\tilde{\gamma}$ lying $\frac{1}{n}$ of the way around from $y^{\prime}$. $\S 7.2 .3$ discusses a traditional picture representing the $n$th power map as if it were the projection on a real coordinate.

Figure 7. An $n$-cycle of path liftings

7.2.3. Impossible pictures. We discuss the problem of representing covers by pictures in $\mathbb{R}^{3}$. Consider the ramified cover $f: U_{w: 0, \infty} \rightarrow U_{z: 0, \infty}$ by $w \mapsto w^{n}$ in Fig. 7. Points of $U_{w: 0, \infty}$ over $z \in U_{z: 0, \infty}$ correspond on the graph of $f$ to $\mathbb{C} \times \mathbb{C}$ points on the line with constant second coordinate $z$. You can't draw pictures in $\mathbb{C} \times \mathbb{C}=\mathbb{R}^{4}$. So first year complex variables texts try to represent $U_{w: 0, \infty}$ and $U_{z: 0, \infty}$ as subsets of $\mathbb{R}^{3}$.

Let $\left(x_{1}, x_{2}, x_{3}\right)$ be coordinates for $\mathbb{R}^{3}$, and let $x_{3}=0$ represent $U_{z: 0, \infty}$ sitting in $\mathbb{R}^{3} \backslash\{(0,0,0)\}$. Pictures try to represent an annulus around the origin in $U_{w: 0, \infty}$ as a set $M$ in $\mathbb{R}^{3}$ over an annulus $D_{0}$ in $U_{z: 0, \infty}$. Then, points of $M$ over $\left(x_{1}, x_{2}, 0\right) \in D_{0}$ are on the line in $\mathbb{R}^{3}$ whose points have first coordinates $x_{1}$ and $x_{2}$. That is, $f$ appears as a coordinate projection. There is, however, no topological subspace $M$ of $\mathbb{R}^{3}$ that can work! If there were, then a cylinder perpendicular to the plane $x_{3}=0$, with $(0,0,0)$ on its axis, would intersect $M$ in a simple closed path winding $n$ times around the cylinder. Represent such a path by $\gamma:[0,1] \rightarrow \mathbb{R}^{3}$ where $t \in[0,1]$ maps to

$$
\gamma(t)=\left(\cos (2 \pi n t), \sin (2 \pi n t), x_{3}(2 \pi n t)\right) \text { and } x_{3}(2 \pi n)=x_{3}(0)
$$

Conclude: $w(t)=x_{3}(2 \pi n t)-x_{3}(2 \pi n t+2 \pi)$ is 0 for some value of $t$ between 0 and $(n-1) / n$. So, the path isn't simple. The author has never seen such a picture attempt in the literature for any noncyclic cover, much less for more demanding nonsolvable groups. Still, we discuss this more in Chap. $4 \S 2.4$ which also uses symbolic representations that assume we understand cyclic covers from their description in Chap. 2.
7.3. Pointed covers and a Galois correspondence. Let $f: Y \rightarrow X$ be a cover. Call the triple $\left(Y, f, y^{\prime}\right)$ a pointed cover if $y^{\prime} \in Y$. Then, we regard $f\left(x^{\prime}\right)=x_{0}$ as the base point for $X$, and $\left(Y, f, y^{\prime}\right)$ is a pointed cover of $\left(X, x_{0}\right)$.

Definition 7.14. Suppose $\left(Y, f_{i}, y_{i}^{\prime}\right), i=1,2$, are two pointed and connected covers of $X$. We say they are compatibly pointed (or compatible) if whenever we have covers $h: Z \rightarrow X$ and $h_{j}: Y_{j} \rightarrow Z$, with $h \circ h_{j}=f_{j}, j=1,2$, then $h_{1}\left(y_{1}\right)=h_{2}\left(y_{2}\right)$.

If it is clear a cover is pointed, we may refer just to the covering maps $f_{1}$ and $f_{2}$ to say these are compatible. Extension Lem. 8.1 shows the difference between a pointed cover on one hand, and a cover without a point on the other. Group theoretically this interprets as the difference between giving a subgroup of a group and giving a conjugacy class of subgroups.
7.3.1. Fiber products of covers. The basic theorems of Galois theory, including the construction of the Galois closure of a cover (§8.3), that translates geometrically using fiber products.

Lemma 7.15. Given connected covers $f_{j}: Y_{j} \rightarrow X, j=1,2$, of $X$, any connected component of $Y_{1} \times_{X} Y_{2}$ is minimal among among connected covers $(Y, f)$ of $X$ factoring through each $f_{j}$. If the covers are compatibly pointed with $y_{j}^{\prime} \in Y_{j}$, $j=1,2$, then a unique pointed component of $Y_{1} \times_{X} Y_{2}$, $\left(Y,\left(f_{1}, f_{2}\right),\left(y_{1}^{\prime}, y_{2}^{\prime}\right)\right)$ is compatible with both $\left(Y_{i}, f_{i}, y_{i}^{\prime}\right), i=1,2$.

Proof. Let $Y$ be a connected component of $Y_{1} \times{ }_{X} Y_{2}$. Denote projection of $Y$ on $Y_{j}$ by $\mathrm{pr}_{j}$. Consider any $\left(y_{1}, y_{2}\right) \in Y_{1} \times_{X} Y_{2}$ lying over $x \in X$. Choose a neighborhood $U_{x}$ of $x$ for which there is a neighborhood $U_{y_{j}} \subset Y_{j}$ on which $f_{j}$ maps one-one to $U_{x}$. Then, restricting $\left(f_{1}, f_{2}\right)$ to $U_{y_{1}} \times_{U_{x}} U_{y_{2}}$ gives a one-one map that shows $Y$ is a cover of $X$.

Now assume the covers are compatibly pointed. Let $x_{0} \in X$ be $f_{1}\left(y_{1}^{\prime}\right)=f_{2}\left(y_{2}^{\prime}\right)$. Then, a unique component of $Y_{1} \times_{X} Y_{2}$ contains $\left(y_{1}^{\prime}, y_{2}^{\prime}\right)$.

Thm. 7.16 produces covers of any path-connected, locally path-connected space. For, however, our main applications where $X$ is a (complex) manifold, it shows any cover of $X$ is a (complex) manifold with a natural coordinate chart. It also says
one cover of a space $X$ dominates all others. This is the universal covering space $\tilde{X}$ corresponding to $H=\{1\} \leq \pi_{1}\left(X, x_{0}\right)$.

Theorem 7.16 (Unramified Galois correspondence). Let $\left(Y, f, y^{\prime}\right)$ be a pointed cover of $\left(X, x_{0}\right)$. This canonically corresponds to a subgroup $H_{Y, f, y^{\prime}} \leq \pi_{1}\left(X, x_{0}\right)$ which we identify with $\pi_{1}\left(Y, y^{\prime}\right)$. The index $\left(\pi_{1}\left(X, x_{0}\right): \pi\left(Y, y^{\prime}\right)\right)$ is $n=\operatorname{deg}(f)$. Any ordering $\boldsymbol{y}=\left\{y_{1}, \ldots, y_{n}\right\}$ on the fiber $f^{-1}\left(x_{0}\right)$ with $y_{1}=y^{\prime}$ corresponds to $a$ transitive permutation representation $T_{Y, f, \boldsymbol{y}}$ in which the stabilizer of 1 is $H_{Y, f, y^{\prime}}$. If $y^{\prime \prime} \in f^{-1}\left(x_{0}\right)$, then $H_{Y, f, y^{\prime}}$ and $H_{Y, f, y^{\prime \prime}}$ are conjugate subgroups of $\pi_{1}\left(X, x_{0}\right)$ and we identify $y^{\prime \prime}$ with a coset of $H$ in $\pi_{1}\left(X, x_{0}\right)$.

Conversely, each subgroup $H \leq \pi_{1}\left(X, X_{0}\right)$ of index $n$ (possibly $\infty$ ) produces a canonical pointed (connected) degree $n$ cover $\left(Y_{H}, f_{H}, y_{H}^{\prime}\right)$ of $X$. We regard $y_{H}^{\prime}$ as the $H$ coset of the identity in $\pi_{1}\left(X, X_{0}\right)$. The fundamental group of $Y_{H}$ maps one-one onto $H$ under $\left(f_{H}\right)_{*}$.

Suppose $H_{1}$ and $H_{2}$ are two subgroups of $\pi_{1}\left(X, x_{0}\right)$. Then, the unique connected component of $Y_{H_{1}} \times_{X} Y_{H_{2}}$ containing $\left(y_{H_{1}}^{\prime}, y_{H_{2}}^{\prime}\right)$ corresponds to the subgroup $H_{1} \cap H_{2}$. The maximal pointed cover of $X$ through which both $f_{1}$ and $f_{2}$ factor is $\left(Y_{\left\langle H_{1}, H_{2}\right\rangle}, f_{\left\langle H_{1}, H_{2}\right\rangle}, y_{\left\langle H_{1}, H_{2}\right\rangle}^{\prime}\right)$.
$\S 7.3 .2$ consists of a proof of Thm. 7.16 and $\S 8.1$ has corollaries appropriate for covers that aren't pointed.
7.3.2. Proof of Thm. 7.16. Start with $\left(Y, f, y^{\prime}\right)$. Apply (7.4) to a closed path $\gamma:[a, b] \rightarrow X$ based at $x_{0}$. Use a specific ordering of $f^{-1}\left(x_{0}\right)$ with $y_{1}=y^{\prime}$. The lift of $\gamma$ to a path with initial point $y_{1}$ is a closed path in $Y$ based at $y_{1}$ if and only if $(1) T_{\boldsymbol{y}}=1$. So we identify $\pi_{1}\left(Y, y_{1}\right)$ with $H\left(f, y_{1}\right)$, the subgroup of $\pi_{1}\left(X, x_{0}\right)$ stabilizing 1 under the map $f_{*}$.

Now consider how a subgroup $H$ of $\pi_{1}\left(X, x_{0}\right)$ of index $n$ canonically produces a degree $n$ pointed cover of $X$. First: $H$ produces an equivalence class of permutation representations of $\pi_{1}\left(X, x_{0}\right)$ of degree $n$ (Thm. 7.8), with the coset of the identity corresponding to the integer 1 in the permutation representation.

Define $Y_{\infty}$ : As a set it is the collection of all equivalence classes of paths in $X$ — not necessarily closed — with initial point $x_{0}$. For $\gamma \in Y_{\infty}$ let $f_{\infty}([\gamma])$ be the endpoint of $\gamma$. Define $Y_{H}$ to be $Y_{\infty}$ modulo the relation that equivalences

$$
\left[\gamma_{1}\right] \text { and }\left[\gamma_{2}\right] \text { if } f_{\infty}\left(\left[\gamma_{1}\right]\right)=f_{\infty}\left(\left[\gamma_{2}\right]\right) \text { and }\left[\gamma_{1} \gamma_{2}^{-1}\right] \in H
$$

Let $f_{H}: Y_{H} \rightarrow X$ be the map induced by $f_{\infty}$ on the set $Y_{H}$. Now use that $X$ is a connected manifold. For each $x \in X$ choose a path $\gamma$ with initial point $x_{0}$ and endpoint $x$. A ball neighborhood $U_{x}$ of $x$ has this property: For $\gamma_{1}, \gamma_{2}:\left[a^{\prime}, b^{\prime}\right] \rightarrow U_{x}$, two paths with the same initial and endpoints, $\gamma_{1} \gamma_{2}^{-1}$ is equivalent to the constant path in $U_{x}$.

For each such pair $\left(\gamma, U_{x}\right)$ consider the subset of $Y_{H}$ represented by paths $\gamma \gamma_{1}$ with $\gamma_{1}$ a path in $U_{x}$ with initial point $x$. Denote this subset by $V_{\gamma, U_{x}}$. We declare the topology on $Y_{H}$ to have as a basis of open sets these $V_{\gamma, U_{x}} \mathrm{~s}$ running over all pairs $\left(x, U_{x}\right)$. For $y \in Y_{H}$ with $f_{H}(y)=x, f_{H}^{-1}\left(U_{x}\right)$ has $n$ connected components, $V_{\gamma_{i}, U_{x}}, i=1, \ldots, n$, where $\left[\gamma_{1} \gamma_{i}^{-1}\right]$ runs over distinct coset representatives of $H$ in $\pi_{1}\left(X, x_{0}\right)$. With this topology $\left(Y_{H}, f_{H}\right)$ satisfies Def. 7.12. It also has an atlas of open sets inherited from $X$. If we show $Y_{H}$ is Hausdorff, then $\left(Y_{H}, f_{H}\right)$ is a cover of $X$. As usual, since $X$ is Hausdorff, we have only to find disjoint open sets around two points over the same point of $X$. We have done exactly that above.

To complete classifying pointed covers of $X$, we show the following. Given $\left(Y, f, y^{\prime}\right)$ a connected cover and $H\left(f, y^{\prime}\right)$ the corresponding subgroup of $\pi_{1}\left(X, x_{0}\right)$, and $\left(Y_{H\left(f, y^{\prime}\right)}, f_{H\left(f, y^{\prime}\right)}, y_{H\left(f, y^{\prime}\right)}^{\prime}\right)$ the cover of $X$ associated to $H\left(f, y^{\prime}\right)$, then
(7.5) $\left(Y, f, y^{\prime}\right)$ is equivalent to $\left(Y_{H\left(f, y^{\prime}\right)}, f_{H\left(f, y^{\prime}\right)}\right)$.

For $y \in Y$ let $\gamma^{*}:[a, b] \rightarrow Y$ be a path from $y^{\prime}$ to $y$, and let $\psi(y)=f_{H}\left(\gamma^{*}\right)$. Follow the defined maps to see $\psi: Y \rightarrow Y_{H(f, \boldsymbol{y})}$ is a one-one map giving (7.5).

Suppose $\left(Y_{H}, f_{H}, y_{H}\right)$ is the canonical cover defined by $H \leq \pi_{1}\left(X, x_{0}\right)$. Let $\left(Y_{H}, f_{H}, y^{\prime \prime}\right)$ by the same cover, those with a different point, $y^{\prime \prime} \in f_{H}^{-1}\left(x_{0}\right)$. Any $\gamma \in \pi_{1}\left(Y, y_{H}, y^{\prime \prime}\right)$ defines a coset $H[\gamma]$ of $H$ in $\pi_{1}\left(X, x_{0}\right)$. Conversely, the elements of $\pi_{1}\left(X, x_{0}\right)$ that stabilize $H[\gamma]$ are exactly the elements of the conjugate subgroup $\left[\gamma^{-1}\right] H[\gamma]$. That shows that using different points attached to a fixed cover correspond to subgroups conjugate to $H$.

Now suppose $H_{1}$ and $H_{2}$ are two subgroups of $\pi_{1}\left(X, x_{0}\right)$. We must show properties attached to the equivalence of two categories: Pointed covers of ( $X, x_{0}$ ) and subgroups of $\pi_{1}\left(X, x_{0}\right)$. The notion of fiber product is a categorical construction. So, the association between $H_{1} \cap H_{2}$ and $\left(Y_{\left\langle H_{1}, H_{2}\right\rangle}, f_{\left\langle H_{1}, H_{2}\right\rangle}, y_{\left\langle H_{1}, H_{2}\right\rangle}^{\prime}\right)$ is that they are the fiber products of the two givens in their respective categories. Def. 1.3 notes the fiber product for subsets of a set is just their intersection. As the intersection of two subgroups is a subgroup, the fiber product from subgroups of a group is just their intersection. For saying fiber product is categorical, see [9.3a]. Similarly, the correspondence between $\left\langle H_{1}, H_{2}\right\rangle$ and $\left(Y_{\left\langle H_{1}, H_{2}\right\rangle}, f_{\left\langle H_{1}, H_{2}\right\rangle}, y_{\left\langle H_{1}, H_{2}\right\rangle}^{\prime}\right)$ is that these are the pushouts of the two givens in their respective categories [9.3c].

## 8. Group theory and covering spaces

We won't be able to make explicit computations with covers until Chap. 4. Still, the topics of this section come from practical experience with covers. Following a discussion of algebraic functions (§8.2) and a geometric approach to the Galois closure of a cover (§8.3), we consider the decomposing covers ( $\S 8.4$ ) and the relation between covers and locally constant bundles (§8.5). A problem from this on computing components of covers shows the power of an elementary piece from finite group representations (§8.6)
8.1. Corollaries of Thm. 7.16. Suppose $\left(Y_{i}, f_{i}, y_{i}^{\prime}\right), i=1,2$, are any two pointed covers of $\left(X, x_{0}\right)$. By an isomorphism $g:\left(Y_{1}, f_{1}, y_{1}^{\prime}\right) \rightarrow\left(Y_{2}, f_{2}, y_{2}^{\prime}\right)$ between them, we mean an isomorphism between $Y_{1}$ and $Y_{2}$ with these properties:
(8.1a) $g\left(y_{1}^{\prime}\right)=y_{2}^{\prime}(g$ preserves basepoints $)$; and
(8.1b) $f_{2} \circ g=f_{1}$ ( $g$ commutes with projections).

The crucial point is that if two pointed covers are isomorphic, this isomorphism is unique. Suppose, however, we don't assume $g$ preserves basepoints?

Lemma 8.1 (Extension Lemma). Consider a pair of covers $\left(Y_{i}, f_{i}\right), i=1,2$, without their basepoints, and any isomorphism $g$ between them. Then, $g$ maps the fiber $f_{1}^{-1}\left(x_{0}\right)$ one-one to $f_{2}^{-1}\left(x_{0}\right)$, and what $g$ does to any one element of $f_{1}^{-1}\left(x_{0}\right)$ determines $g$. Further, isomorphisms between $\left(Y_{1}, f_{1}\right)$ and $\left(Y_{2}, f_{2}\right)$ correspond oneone with automorphisms $\operatorname{Aut}\left(Y_{i}, f_{i}\right)$ of $\left(Y_{i}, f_{i}\right)$ (for either $i=1$ or 2).

Any automorphism of a cover $(Y, f)$ of $X$ lifts to an automorphism of the universal cover $(\tilde{X}, \tilde{f})$ of $X$. If $X$ is a complex manifold, then $\operatorname{Aut}(Y, f)$ is a group of complex analytic isomorphisms.

Proof. Assume $g$ that maps $y_{1}^{\prime} \in f_{1}^{-1}\left(x_{0}\right)$ to $y_{2}^{\prime} \in f_{1}^{-1}\left(x_{0}\right)$. Then, $g$ is an isomorphism between $\left(Y_{1}, f_{1}, y_{1}^{\prime}\right)$ and $\left(Y_{2}, f_{2}, y_{2}^{\prime}\right)$, and so it is unique. Let $A_{1,2}$ be the set of isomorphisms between $\left(Y_{1}, f_{1}\right)$ and $\left.Y_{2}, f_{2}\right)$. Then, we have an action of $\operatorname{Aut}\left(Y_{1}, f_{1}\right)$ (resp. $\left.\operatorname{Aut}\left(Y_{2}, f_{2}\right)\right)$ on the right (resp. left) of $A_{1,2}$ :

$$
\begin{array}{ll}
A_{1}: A_{1,2} \times \operatorname{Aut}\left(Y_{1}, f_{1}\right) \rightarrow A_{1,2} \text { by } & (g, \alpha) \mapsto g \circ \alpha ; \text { and } \\
A_{2}: \operatorname{Aut}\left(Y_{2}, f_{2}\right) \times A_{1,2} \rightarrow A_{1,2} \text { by } & (\beta, g) \mapsto \beta \circ g .
\end{array}
$$

For $g^{\prime}, g \in A_{1,2}, g^{-1} g^{\prime}=\alpha$ is in $\operatorname{Aut}\left(Y_{1}, f_{1}\right)$. This shows $g \circ \alpha=g^{\prime}$, and $A_{1}$ is transitive on $A_{1,2}$ (as in §7.1). Similarly, $A_{2}$ is transitive on $A_{1,2}$.

Now consider an automorphism $\alpha$ of $(Y, f)$. Again, let $\left(Y, f, y^{\prime}\right)$ with $y^{\prime}$ over $x_{0}$ be a corresponding pointed cover. Then, $\left(Y, f, y^{\prime}\right)$ and $\left(Y, f, \alpha\left(y^{\prime}\right)\right)$ are pointed covers of $\left(X, x_{0}\right)$. So, Thm. 7.16 shows they correspond to conjugate subgroups $H$ and $H_{\alpha}: H_{\alpha}=\left[\gamma^{-1}\right] H[\gamma]$ for some $[\gamma] \in \pi_{1}\left(X, x_{0}\right)$. A natural analytic isomorphism between $\left(Y_{H}, f_{H}, y_{H}^{\prime}\right)$ and $\left(Y_{H_{\alpha}}, f_{H_{\alpha}}, y_{H_{\alpha}}^{\prime}\right)$ comes by mapping the homotopy class of $\left[\gamma^{\prime}\right]$ defining a point of $Y_{H_{\alpha}}$ (in $\S 7.3 .2$ ) to $[\gamma]\left[\gamma^{\prime}\right]$. The new base point (the coset of $[\gamma]$ ) has stabilizer $\left[\gamma^{-1}\right] H[\gamma]$. This automorphism lifts to the universal covering space, because premultiplying by $[\gamma]$ also defines it there.

Definition 8.2. Let $T_{\boldsymbol{y}}: \pi_{1}\left(X, x_{0}\right) \rightarrow S_{n}$ be the representation of (7.4) associated to $(Y, f)$. The image of $\pi_{1}\left(X, x_{0}\right)$ is called the (geometric) monodromy group, $G(Y, f)$, of the cover. It is isomorphic to $\pi_{1}\left(X, x_{0}\right) / \bigcap_{i=1}^{n} \pi_{1}\left(Y, y_{i}\right)$ (Thm. 7.8).

Covers $(Y, f)$ of a manifold $\left(X, x_{0}\right)$ have two extremes. For most, $\operatorname{Aut}(Y, f)$ consists only of the identity element: We say $(Y, f)$ has no automorphisms. The other extreme is in this definition.

Definition 8.3. If $\operatorname{Aut}(Y, f)$ is transitive on the the fiber $f^{-1}\left(x_{0}\right)$, we say $(Y, f)$ is Galois.

The Galois situation is our main tool, though what constantly arises in practice is the situation with no automorphisms. $\oint 8.3$ has the details for distinguishing these and all the cases in between. An example of the Galois situation is the universal cover of $\left(X, x_{0}\right)$ where the automorphism group is isomorphic to the whole fundamental group of $\left(X, x_{0}\right)$. The fiber $f^{-1}\left(x_{0}\right)$ in this case corresponds to the elements of $\pi_{1}\left(X, x_{0}\right)$, and by translation these give a permutation of the points. Automorphisms also give a permutation of $f^{-1}\left(x_{0}\right)$. Still, from Lem. 8.8, only when $\pi_{1}\left(X, x_{0}\right)$ is abelian can we expect to canonically identify these two groups of permutations. The next lemma revisits Chap. 2 Prop. 3.2. As previously, use the notation $\tilde{f}: \tilde{X} \rightarrow X$ for the universal cover of $X$ with paths starting at $x_{0}$ representing its points.

Lemma 8.4. In the notation above, let $[\gamma] \in \pi_{1}\left(X, x_{0}\right)$ and let $\left[\gamma^{\prime}\right]$ represent a homotopy class of paths on $X$ with $\gamma^{\prime}:[a, b] \rightarrow X, \gamma^{\prime}(a)=x_{0}$ and $\gamma^{\prime}(b)=$ $x$. Then, multiplication by $[\gamma]^{-1}$ on the left of $\gamma^{\prime}$ induces an automorphism of $\tilde{X}$ giving an action $A_{L}: \pi_{1}\left(X, x_{0}\right) \times \tilde{X} \rightarrow \tilde{X}$. Regard the fiber $\tilde{f}^{-1}\left(x_{0}\right)$ as elements of $\pi_{1}\left(X, x_{0}\right)$. Then, the usual right action of $\pi_{1}\left(X, x_{0}\right)$ gives the group structure identifying $\pi_{1}\left(X, x_{0}\right)$ with the monodromy group of $\tilde{f}$.

The exponential map $\exp : \mathbb{R} \rightarrow S^{1}$ by $\theta \mapsto e^{2 \pi i \theta}$ presents $\mathbb{R}$ as the universal cover of $S^{1}$ with $\mathbb{Z}$ as its fundamental group. The path $\gamma_{n}^{*}$ corresponds to $n \in \mathbb{Z}$ and the automorphisms of $(\mathbb{R}, \exp )$ identify with $\mathbb{Z}$ acting by translation. Similarly, the fundamental group of a complex torus $\mathbb{C}^{n} / L$ identifies with the lattice $L$.

Proof. The universal covering space is unique up to homeomorphisms commuting with the map to $X$. One way to identify the fundamental group of a space $X$ is to find any space $\tilde{X}$ with trivial fundamental group and a covering map $\tilde{f}: \tilde{X} \rightarrow X$. Given $x_{0} \in X$, any other cover of $X$ that has trivial fundamental group must be isomorphic to ( $\tilde{X}, \tilde{f})$, and this isomorphism is unique up to composition on the left with an element of $(\tilde{X}, \tilde{f})$. Since $\mathbb{R}$ and $\mathbb{C}^{n}$ are contractible, they have trivial fundamental group (Lem. 6.6). The $\operatorname{map} \theta \in \mathbb{R} \mapsto e^{2 \pi i \theta}$ is a covering map with the elements of $\mathbb{R}$ over 1 given by the integers. The permutation of the fiber over 1 given by the path $\gamma_{n}^{*}$ is translation by $n$. The argument is similar for a complex torus.

The next corollary tells when a map between spaces extends to a map between covers of the spaces.

Corollary 8.5. Suppose $\varphi: X \rightarrow X^{\prime}$ is a differentiable map between complex manifolds mapping a point $x_{0} \in X$ to $x_{0}^{\prime} \in X^{\prime}$. Let $\varphi_{H^{\prime}}: Y_{H^{\prime}}^{\prime} \rightarrow X^{\prime}$ be the cover defined by a subgroup $H^{\prime} \leq \pi_{1}\left(X^{\prime}, x_{0}^{\prime}\right)$. Then, there is a continuous (and so automatically differentiable) map $\psi: X \rightarrow Y_{H^{\prime}}^{\prime}$ with $\varphi_{H^{\prime}} \circ \psi=\varphi$ if and only if the induced map $\varphi_{*}: \pi_{1}\left(X, x_{0}\right) \rightarrow \pi_{1}\left(X^{\prime}, x_{0}^{\prime}\right)$ has image in a conjugate of $H^{\prime}$.

Proof. Suppose the induced map $\varphi_{*}$ has image in a conjugate $m^{-1} H^{\prime} m$ of $H^{\prime}$. Let $\gamma^{*}$ be a representative path in $X^{\prime}$ for which $\left[\gamma^{*}\right]=m$. Then, let $\gamma$ : $[a, b] \rightarrow X$ start at $x_{0}$ and end at $x$. Define $\psi_{m, H^{\prime}}: X \rightarrow Y_{H^{\prime}}^{\prime}$ by $\psi(x)$ is the class $m \cdot[\varphi \circ \gamma] \in Y_{H^{\prime}}^{\prime}$ : the product of $m$ and the image under $\psi$ of $\gamma$. To show the map doesn't depend on $\gamma$, we consider another closed path $\gamma^{\prime}$ from $x_{0}$ to $x$. We are done if the closed path $\left(\gamma^{*}\right)^{-1} \cdot \psi\left(\gamma \cdot\left(\gamma^{\prime}\right)^{-1}\right) \cdot \gamma^{*}$ in $X^{\prime}$ defines a closed path in $Y_{H^{\prime}}^{\prime}$. Since, however, $\gamma \cdot\left(\gamma^{\prime}\right)^{-1}$ is a closed path in $X$, its image under $p h i_{*}$ is some $\rho \in m^{-1} H^{\prime} m$ by hypothesis and the image of $\left(\gamma^{*}\right)^{-1} \cdot \psi\left(\gamma \cdot\left(\gamma^{\prime}\right)^{-1}\right) \cdot \gamma^{*}$ is therefore $m \rho m^{-1} \in H^{\prime}$. From the definition of $Y_{H^{\prime}}^{\prime}$ this exactly says the image path is closed.

Conversely, suppose there is such a $\psi: X \rightarrow Y_{H^{\prime}}^{\prime}$. Then, closed paths in $X$ have image under $\psi$ in $X^{\prime}$ that lift to closed paths in $Y_{H^{\prime}}^{\prime}$. So, the image group $\psi_{*}\left(\pi_{1}\left(X, x_{0}\right)\right)=H^{*}$ is a subgroup of $\pi_{1}\left(X^{\prime}, x_{0}^{\prime}\right)$ whose corresponding cover $Y_{H^{*}}^{\prime}$ factors through $\psi_{H^{\prime}}: Y_{H^{*}}^{\prime} \rightarrow X^{\prime}$.

Suppose $X$ is a connected complex manifold (like $U_{\boldsymbol{z}}=\mathbb{P}_{z}^{1} \backslash\{\boldsymbol{z}\}$ ). Define analytic continuation along a path from Def. 6.7. Consider the extensible functions $\mathcal{E}\left(X, x_{0}\right)$ : complex analytic functions defined in a neighborhood of $x_{0}$ that have an analytic continuation along every path in $X$ (as in Chap. 2 Def. 4.5). Let $\varphi: Y \rightarrow X$ be a cover with $y_{0} \in Y$ lying over $x_{0}$. Let $\gamma:[a, b] \rightarrow X$ be a path starting at $x_{0}$ with $\gamma^{\dagger}:[a, b] \rightarrow Y$ its unique path lift starting at $y_{0}$ (Lem. 7.13).

Proposition 8.6. There is an isomorphism (of rings) between $\mathcal{E}\left(Y, y_{0}\right)$ and $\mathcal{E}\left(X, x_{0}\right)$. In particular, for $\left(\tilde{X}, \tilde{x}_{0}\right)$ the universal cover of $\left(X, x_{0}\right)$, holomorphic functions on $\tilde{X}$ form a ring isomorphic to $\mathcal{E}\left(X, x_{0}\right)$. If $\varphi$ is a finite cover of punctured Riemann surfaces, this induces an analytic isomorphism between $\mathcal{E}\left(Y, y_{0}\right)^{\text {alg }}$ and $\mathcal{E}\left(X, x_{0}\right)^{\text {alg }}$. These results hold with extensible meromorphic replacing extensible holomorphic functions.

Proof. Since $\varphi$ is a cover, there is a disk neighborhood $U_{x_{0}}$ of $x_{0}$ and a component $U_{y_{0}}$ of $\varphi^{-1}\left(U_{x_{0}}\right)$ with $y_{0} \in U_{y_{0}}$ on which $\varphi$ maps one-one. So, restriction of a function $f \in \mathcal{E}\left(X, x_{0}\right)$ to $U_{x_{0}}$ transports by $\varphi^{-1}$ to a function $f \in U_{y_{0}}$. There
is no harm in using the same notation to extend $f$ along $\gamma^{\dagger}:[a, b] \rightarrow Y$ starting at $y_{0}$. Let $\gamma$ be $\varphi \circ \gamma^{\dagger}$, and let $f^{*}:[a, b] \rightarrow \mathbb{P}_{z}^{1}$ be the continuous function defining the analytic continuation along $\gamma$. Define the analytic continuation of $f$ along $\gamma^{\dagger}$ to be the same, $f^{*}$. This shows $f$ is extensible in $Y$. Clearly, if $f$ is algebraic (on $X$ ) it will also be algebraic on $Y$.
8.2. The problem of identifying algebraic functions explicitly. Suppose $\tilde{\varphi}: \tilde{X} \rightarrow X$ is the universal covering space of a complex manifold $X$ and $\tilde{x}$ lies over $x_{0} \in X$. Then, similar to formation of complex torii and other quotient manifolds, it is natural to regard points of $X$ as the orbits of the action of $\pi_{1}\left(X, x_{0}\right)$ on $\tilde{X}$. Riemann's approach was to identify the universal covering space of a Riemann surface as a simply connnected domain on the Riemann sphere. Consider the case of Prop. 8.6 when $Y=\tilde{X}$ and $X=U_{\boldsymbol{z}}$, with $|\boldsymbol{z}| \geq 3$. Riemann's Uniformization Theorem says $\tilde{X}$ is analytically isomorphic to a disk $\Delta$ in such a way that the map extends continuously to the boundaries (Chap. 4 Def. 7.9 for an elementary proof, or $[\mathbf{S p r} 57, \mathrm{Thm} .9 .6]$ for the more general case). So, $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ is ring isomorphic to convergent functions in a disk. We find it convenient to replace a disk by the analytically isomorphic upper half plane $\mathbb{H}$. This is the same exact space independent of $\left(z_{0}, \boldsymbol{z}\right)$. What changes, however, with $\boldsymbol{z}$ is the identification of algebraic functions $F_{\boldsymbol{z}}$. Suppose $\varphi_{\boldsymbol{z}}: \mathbb{H} \rightarrow U_{\boldsymbol{z}}$ is this uniformization.

Elements of $\mathrm{PGL}_{2}(\mathbb{R})$ with positive determinant (Chap. $2[9.14 \mathrm{~d}]$; this identifies with $\mathrm{PSL}_{2}(\mathbb{R})$ ) represent the action of complex analytic isomorphisms of $\mathbb{H}$. As $\boldsymbol{z}$ varies, a different subgroup $\Gamma_{\boldsymbol{z}}$ (though abstractly isomorphic as a group) of $\mathrm{PSL}_{2}(\mathbb{R})$ defines $U_{\boldsymbol{z}}$ as a quotient of $\mathbb{H}$.

Prop. 8.6 identifies extensible (meromorphic) algebraic functions on $U_{z}$ with certain meromorphic functions $\mathcal{F}_{z}$ on $\mathbb{H}$. Though, which ones? Given $g^{*}$ meromorphic on $\mathbb{H}$, composing it with an analytic isomorphism of $\mathbb{H}$ produces a new meromorphic function on $\mathbb{H}$. We call the compositions of $g^{*}$ with elements of $\Gamma_{\boldsymbol{z}}$ transforms by $\Gamma_{z}$.

Proposition 8.7. Suppose $f$, meromorphic on $\mathbb{H}$, has only finitely many transforms under the action of $\Gamma_{\boldsymbol{z}}$ and a unique limit value as it approaches any point in $\mathbb{R} \cup\{\infty\}$. Then, $f$ defines an algebraic element of $\mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ and conversely.

Outline. Let $\tilde{x} \in \mathbb{H}$ lie over $z_{0} \in U_{\boldsymbol{z}}$. From Prop. 8.6, any meromorphic extensible function $g$ on $U_{z}$ identifies with a meromorphic function $g^{*}$ on $\mathbb{H}$. Further, the analytic continuation of $g$ around $[\gamma] \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ produces $g_{\gamma}^{*}$, the result of composing $g^{*}$ with the analytic isomorphism of $\mathbb{H}$ associated to $\gamma$. If $g$ is algebraic, then it has only finitely many analytic continuations, so the different transforms $g_{\gamma}^{*}$, running over $\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ are finite in number. Conversely, if the number of transforms of a meromorphic function $g^{*}$ on $\mathbb{H}$ are finite in number, then the identification of $g^{*}$ with $g \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ gives a function with only finitely many analytic continuations.
8.3. Galois theory and covering spaces. Use notation from Lem. 8.1: $(Y, f)$ is a cover of $X$.
8.3.1. Identifying automorphisms of a cover. Having $\operatorname{Aut}(Y, f)$ act on a fiber $\left\{y_{1}, \ldots, y_{n}\right\}=f^{-1}\left(x_{0}\right)$ induces a homomorphism $\Lambda_{\boldsymbol{y}}: \operatorname{Aut}(Y, f) \rightarrow S_{n}$.

It is a mistake to confuse the Galois (geometric monodromy) group of a cover with its automorphism group, even if the cover is Galois. The next lemma efficiently
differentiates $\operatorname{Aut}(Y, f)$ from $G(Y, f)$. It shows that having chosen a right action for $G(Y, f)$ forces using a left action of $\operatorname{Aut}(Y, f)$ on the set $\{1, \ldots, n\}$.

Lemma 8.8. Let $(Y, f)$ be a connected cover of $X$. The homomorphism $\Lambda_{\boldsymbol{y}}$ injects $\operatorname{Aut}(Y, f)$ onto the centralizer $\operatorname{Cen}_{S_{n}}(G(Y, f))$ of $G(Y, f)$ in $S_{n}$. This is isomorphic to $N_{\pi_{1}\left(X, x_{0}\right)}\left(\pi_{1}\left(Y, y_{1}\right)\right) / \pi_{1}\left(Y, y_{1}\right)$ (§7.1) and $|\operatorname{Aut}(Y, f)| \leq n$ with equality if and only if $\pi_{1}\left(Y, y_{1}\right)$ is normal in $\pi_{1}\left(X, x_{0}\right)$.

Proof. For $y \in Y$ let $\gamma^{*}:[a, b] \rightarrow Y$ be a path with initial point $y_{i}$ and endpoint $y$. Consider $\psi \in \operatorname{Aut}(Y, f)$. Then $\psi \circ \gamma^{*}:[a, b] \rightarrow Y$ is the (unique) lift of $f \circ \gamma^{*}$ with initial point $\psi\left(y_{i}\right)$. So, if $i=1$ and $\psi\left(y_{1}\right)=y_{1}$, then $\psi \circ \gamma^{*}=\gamma^{*}$. Thus $\psi(y)=y$ for each $y \in Y$, and $\Lambda_{y}$ is injective. This alone shows $|\operatorname{Aut}(Y, f)| \leq n$.

In the above, assume $\gamma=f \circ \gamma^{*}$ is a closed path. If the endpoint of $\gamma^{*}$ is $y_{j}$, then the endpoint of $\psi \circ \gamma^{*}$ is $\psi\left(y_{j}\right)$. Thus

$$
(i) \Lambda_{\boldsymbol{y}}(\psi)^{-1} \circ T_{\boldsymbol{y}}(\gamma) \circ \Lambda_{\boldsymbol{y}}(\psi)=(i) T_{\boldsymbol{y}}(\gamma)
$$

Equivalently, $\Lambda_{\boldsymbol{y}}(\psi) \in \operatorname{Cen}_{S_{n}}(G(Y, f))$. Conversely, for $\alpha \in \operatorname{Cen}_{S_{n}}(G(Y, f))$ define $\alpha$ to be a permutation of the points $\left\{y_{1}, \ldots, y_{n}\right\}$ from its action on $\{1, \ldots, n\}$. Still, use an action on the left: If $(i) \alpha=j$, write $\left.\alpha\left(y_{i}\right)=y_{j}\right)$. Our goal is to create an automorphism - also called $\alpha$ - on $Y$ that extends this action on the fiber over $x_{0}$.

Take $i=1$ and $\gamma^{*}$ as in the first paragraph above. Define $\psi_{\alpha, \gamma^{*}}$ :
(8.2) $\psi_{\alpha, \gamma^{*}}(y)$ is the endpoint of the lift of $f \circ \gamma^{*}$ with initial point $\alpha\left(y_{1}\right)$.

If we show $\psi_{\alpha, \gamma^{*}}(y)$ is independent of $\gamma^{*}$ having endpoint $y$, then $\psi_{\alpha, \gamma^{*}}$ defines an element $\psi_{\alpha} \in \operatorname{Aut}(Y, f)$. For this purpose let $\gamma^{1}$ (resp., $\gamma^{2}$ ) be a path in $Y$ with initial (resp., end) point $y$ and end (resp., initial) point $y_{1}$. If $\psi_{\alpha, \gamma^{*}}(y) \neq \psi_{\alpha, \gamma^{1}}(y)$, then $\psi_{\alpha, \gamma^{*} \gamma^{2}}\left(y_{1}\right) \neq \psi_{\alpha, \gamma^{1} \gamma^{2}}\left(y_{1}\right)$. Therefore, $\psi_{\alpha, \gamma^{*}}(y)$ is independent of $\gamma^{*}$ if and only if $\psi_{\alpha, \gamma^{*}}\left(y_{1}\right)$ is independent of $\gamma^{*}$ for $\gamma^{*} \in \pi_{1}\left(Y, y_{1}\right)$. That is, we must show $\alpha\left(y_{1}\right)$ is the endpoint of the lift of $f \circ \gamma$ with initial point $\alpha\left(y_{1}\right)$ for each $\gamma \in \pi_{1}\left(Y, y_{1}\right)$.

With $\alpha\left(y_{1}\right)=y_{j}$, this is equivalent to $((1) \alpha) T(Y, f)(f \circ \gamma)=j$. (The right action of $\alpha$ on 1 is intentional- $\alpha$ did come from $S_{n}$.) For $\gamma$ a closed path on $Y$ with initial point $y_{1},(1) T(Y, f)(f \circ \gamma)=1$ is automatic. Apply $\alpha$ to the right side of this and use that $\alpha$ commutes with $T(Y, f)(f \circ \gamma)$ to conclude from [9.15b]. Recall: $G(1)$ is the subgroup of $G(Y, f)$ leaving 1 fixed. Thm. 7.16 identifies $N_{\pi_{1}\left(X, x_{0}\right)}\left(\pi_{1}\left(Y, y_{1}\right)\right) / \pi_{1}\left(Y, y_{1}\right)$ with $N_{G(Y, f)}(G(1)) / G(1)$.
8.3.2. Fiber products and Galois closure. We say a connected cover $(Y, f)$ of $X$ is a Galois cover (or is Galois) if $|\operatorname{Aut}(Y, f)|$ equals $n=\operatorname{deg}(f)$. By Lem. 8.8 this holds if and only if $\pi_{1}\left(Y, y_{1}\right)$ is a normal subgroup of $\left.\pi_{1}\left(X, x_{0}\right)\right)$. Each cover $(Y, f)$ produces a Galois cover $(\hat{Y}, \hat{f})$ of $X$ called the Galois closure of $(Y, f)$. If $H \leq \pi_{1}\left(X, x_{0}\right)$ corresponds to $Y$, then $\cap g^{-1} H g$ corresponds to $(\hat{Y}, \hat{f})$. We use fiber products to give an alternate construction of it (Def. 1.3). It correctly displays the automorphism group action. We again warn: Don't confuse it with the geometric monodromy group, though they are isomorphic for a Galois cover.

Denote the fiber product of $Y \rightarrow X$ taken $n=\operatorname{deg}(f)$ times by

$$
Y_{X}^{n} \stackrel{\text { def }}{=} Y \times_{X} \times \cdots \times_{X} Y
$$

Points of $Y_{X}^{n}$ are $n$-tuples $\left(y_{1}^{\prime}, \ldots, y_{n}^{\prime}\right) \in Y^{n}$ for which $f\left(y_{i}\right)=f\left(y_{j}\right)$ for all $i$ and $j$. The fat diagonal, $\Delta_{Y, f, n}$, is the subset of $n$-tuples of $Y_{X}^{n}$ with at least two equal coordinate entries. Remove it to form $Y_{X}^{n} \backslash \Delta_{Y, f, n}=U_{Y, f, n}$. We use a copy of $S_{n}$ acting on the left of $\{1, \ldots, n\}$ to give an action of automorphisms on this set:
(8.3)

$$
\text { for } \sigma \in S_{n} \text { and } \boldsymbol{y}^{\prime}=\left(y_{1}^{\prime}, \ldots, y_{n}^{\prime}\right) \in U_{Y, f, n}, \alpha_{\sigma} \operatorname{maps} \boldsymbol{y}^{\prime} \text { to }
$$

$$
\left(y_{\sigma(1)}^{\prime}, \ldots, y_{\sigma(n)}^{\prime}\right)=\alpha_{\sigma}\left(\boldsymbol{y}^{\prime}\right)
$$

Restrict the natural map of $Y_{X}^{n}$ to $X$ to $U_{Y, f, n}$ to present $U_{Y, f, n}$ as a degree $n$ ! cover of $X$ with automorphism group containing $S_{n}$. The action of $S_{n}$ is transitive on points mapping to $x_{0}$. Yet, $U_{Y, f, n}$ may not be connected. (We don't consider it a Galois cover of $X$.) Decompose $U_{Y, f, n}$ into connected components $\hat{Y}_{1}, \ldots, \hat{Y}_{t}$. Let $\hat{f}_{i}$ be the restriction to $\hat{Y}_{i}$ of the projection map $U_{Y, f, n} \rightarrow X, i=1, \ldots, t$. A computation shows $\operatorname{deg}\left(\hat{f}_{i}\right)=|G(Y, f)|$ [9.22].

THEOREM 8.9. The covers $\left(\hat{Y}_{i}, \hat{f}_{i}\right)$ are equivalent as covers of $X, i=1, \ldots, t$. Characterize members $(\hat{Y}, \hat{f})$ of this equivalence class from these properties.
(8.4a) $(\hat{Y}, \hat{f})$ is a Galois cover of $X$, with its group a transitive subgroup of $S_{n}$.
(8.4b) There is a commutative diagram of covers of $X$ :

(8.4c) For any Galois cover $\hat{g}: \hat{Z} \rightarrow X$ factoring through $Y$ by $g_{Y}: \hat{Z} \rightarrow Y$, there is commutative diagram of covers of $X$ :


Proof. Choose $y_{1} \in Y$ lying over $x_{0} \in X$. Thm. 7.3.2 identifies the subgroup of $\pi_{1}\left(X, x_{0}\right)$ corresponding to $(Y, f)$ with $\pi_{1}\left(Y, y_{1}\right)$. It also identifies its conjugates (in $\left.\pi_{1}\left(X, x_{0}\right)\right) \pi_{1}\left(Y, y_{i}\right)$ with $y_{i}$ running over $f^{-1}\left(x_{0}\right)$. A Galois cover corresponds to a normal subgroup of $\pi_{1}\left(X, x_{0}\right)$. So, the smallest Galois cover mapping through $(Y, f)$ corresponds to the largest normal subgroup, $H=\bigcap_{i=1}^{n} \pi_{1}\left(Y, y_{i}\right)$, of $\pi_{1}\left(X, x_{0}\right)$ contained in $\pi_{1}\left(Y, y_{1}\right)$. So, there is a cover with property (8.4c).

Let $\mathrm{pr}_{1}: Y_{X}^{n} \rightarrow Y$ be projection onto the first factor, and let $f_{Y, i}$ be the restriction of $\operatorname{pr}_{1}$ to $\hat{Y}_{i}$. Then, with $(\hat{Y}, \hat{f})$ (resp., $f_{Y}$ ) replaced by $\left(\hat{Y}_{i}, \hat{f}_{i}\right)$ (resp., $\left.f_{Y, i}\right)$ properties (8.4a) and (8.4b) hold, $i=1, \ldots, t$. This shows the map $h: \hat{Y}_{i} \rightarrow Y$ has degree 1: $\left(\hat{Y}_{i}, \hat{f}_{i}\right)$ and $(\hat{Y}, \hat{f})$ are equivalent covers of $X$. The proof is complete.

Fig. 8 shows four discs on a degree 4 cover of $U_{\boldsymbol{z}}$ lying over a disk $U_{z_{0}}$ around the base point. Assume the cover has monodromy group $S_{4}$. (Like that from a general degree 4 polynomial $f \in \mathbb{C}[w]$.) We visibly can see the action of any element $\alpha \in S_{4}$ on the four points of $f^{-1}\left(z_{0}\right)$ extend to the four disjoint disks over $U_{z_{0}}$. Yet, there is no continuous extension of any nonidentity $\alpha$ to $f^{-1}\left(U_{z_{0}}\right)$. Lem. 8.8 says such extending $\alpha$ s must centralize the monodromy. We stipulated, however, this is $S_{4}$, a group with trivial center.

Figure 8. $\alpha=(12)(34) \in S_{4}$ tries, but fails, to be an automorphism of $Y$ : The four discs on the left constitute $f^{-1}\left(U_{x_{0}}\right)$

8.3.3. Galois closure orbits. Chap. 2 [9.5] has Galois exercises based on using fields. We now explain how these have analogs where we replace field extensions of a given field by covers of a given space. One tricky point: Composite of two fields makes sense only if there is given a priori a field $L$ containing them both. As with the comments from $\S 4.2 .3$ on local holomorphic functions, the next lemma shows fiber product of covers is dual to tensor product of fields. This analogy will come through even more when we deal with the field of meromorphic functions on a cover in Chap. 4 Prop. 2.11.

Lemma 8.10. Let $K_{i}, i=1,2$, be two finite extensions of a field $K$ (having 0 characteristic). The ring $K_{1} \otimes_{K} K_{2}$ is the direct sum of field extension of $K$. These summands are, up to isomorphism of extensions of $K$, in one-one correspondence with all compositions of $K_{1}$ and $K_{2}$.

Proof. Since the characteristic is 0 (only need separable extensions), the primitive element theorem says $K_{2}=K(\alpha)$ for some $\alpha \in K_{2}$. Up to isomorphism of extensions, $K_{2} / K$ is $K[x] /\left(f_{2}(x)\right)$ with $f_{2}$ the irreducible polynomial for $\alpha$ over $K$. Factor $f_{2}$ as $\prod_{i=1}^{u} g_{i}(x)$ over $K_{1}$, with the $g_{i}$ s monic and distinct. (Again use characteristic 0 , or just that irreducible polynomials have no repeated roots.) Now apply Lem. 4.8 to write $K_{1} \otimes K_{2}=K_{1}[x] /\left(f_{2}(x)\right)$ as $\oplus_{i=1}^{u} K_{1}[x] /\left(g_{i}(x)\right)$. Since each of the $g_{i}$ s is irreducible over $K_{1}$, each of the summands is a field. So each summand is a field generated by extensions of $K$ isomorphic to $K_{1}$ and $K_{2}$.

Conversely, suppose $L$ is a field containing $K_{1}$ and generated by $K_{1}$ and $K^{\prime}=$ $K\left(\alpha^{\prime}\right) / K$ with $\alpha^{\prime}$ the image of $\alpha$ in an isomorphism of $K_{2} / K$ with it. Then, $L$ is isomorphic to one of the summands of $K_{1} \otimes K_{2}$. This concludes the proof.

Suppose $L_{i} / K$ (resp. $f_{i}: Y_{i} \rightarrow X$ ) is a field extension (resp. connected cover) of finite degree $n_{i}$, with $G_{i}$ its Galois closure group and $\hat{L}_{i} / K$ (resp. $\hat{f}_{i}: \hat{Y}_{i} \rightarrow X$ its Galois closure field (resp. cover), $i=1,2$. As in Chap. 2 [9.6a], consider the fiber product $H_{f}$ of $G_{1}$ and $G_{2}$ over the Galois group of the well-defined field extension $\hat{L}_{1} \cap \hat{L}_{2}$. Then, $G\left(\hat{L}_{1} \cdot \hat{L}_{2} / K\right)$ is $H_{f}$. The restriction of elements of $H_{f}$ to $\hat{L}_{i}$ produces a permutation representation $T_{i}, i=1,2$. Now consider the direct product representation $T_{f}$ of $H_{f}$ induced from $T_{1}$ and $T_{2}$ (§7.1.2). The next lemma, in this analogy, shows different composites of field extensions correspond
to the different components of the fiber product of the covers over $X$. The proof shows also that inequivalent composite extensions $L_{1} \cdot L_{2}$ correspond one-one to orbits of $T_{f}$ (compare with Chap. 2 [9.6c]).

Lemma 8.11. Let $g: Y \rightarrow X$ be the maximal cover through which $\hat{f}_{i}, i=1,2$, both factor. Then, $g$ is a Galois cover. If $M$ is its group, this induces homomorphisms $f_{i_{*}}: G_{i} \rightarrow M$. Denote the fiber product of these group homomorphisms by $H_{c}$. Then, any connected component $\hat{Y}_{1,2}$ of $\hat{Y}_{1} \times_{X} \hat{Y}_{2}$ (as a cover of $X$ ) is the minimal Galois cover of $X$ factoring through $\hat{f}_{i}, i=1,2$. The group of this cover is $H_{c}$, a subgroup of $S_{n_{1}} \times S_{n_{2}}$ (acting on pairs $(i, j), 1 \leq i \leq n_{1}, 1 \leq j \leq n_{2}$ ). Orbits of $T_{c}$ correspond one-one to the components of $Y_{1} \times{ }_{X} Y_{2}$.

Proof. Let $\mathcal{C}_{\text {Gal }}$ be the category of Galois covers of $X$ up to isomorphism commuting with the map to $X$. Similarly, let $\mathcal{C}_{\text {Nor }}$ be the category of normal subgroups of $\pi_{1}\left(X, x_{0}\right)$. The first part of the lemma is an equivalencing of fiber products in each of these categories (as at the end of the proof of Thm. 7.16). The fiber product for two normal subgroups of $\pi_{1}\left(X, x_{0}\right)$ is their intersection, which identifies the quotient as $H_{c}$ in this case. Since the fiber product $\hat{Y}_{1} \times{ }_{X} \hat{Y}_{2}$ may not be connected, and therefore not Galois, this cannot be the fiber product in the category of Galois covers of $X$. A connected component, however, of it defines an equivalence class of connected and Galois covers. It is this that is the fiber product in the category $\mathcal{C}_{\text {Gal }}$.

Now consider the statement on orbits of $T_{c}$. Since $H_{c}$ factors through $G_{i}$, with its representation $T_{i}, i=1,2$, it makes sense to form the direct (tensor) product $T_{c}$ of $T_{1}$ and $T_{2}$. Direct summands in the category of permutation representations correspond to components of covers in the category of covers of $X$. Since permutation representations correspond to equivalence classes of covers, to show the statement on orbits we have only to show that the direct product permutation representation $T_{c}$ corresponds to the fiber product $Y_{1} \times_{X} Y_{2}$. This is the equivalence of direct product in their respective categories.
8.4. Imprimitive covers and wreath products. Suppose $f: Y \rightarrow X$ is a (connected) cover, and $f_{f}$ factors through another cover $f_{1}: Y_{1} \rightarrow X$. That gives a series of covers $Y \xrightarrow{f_{2}} Y_{1} \xrightarrow{f_{1}} X$. We say $f_{1} \circ f_{2}$ is a decomposition of $f$ if $\operatorname{deg}\left(f_{i}\right)>1$, $i=1,2$. If there is no such decomposition of $f$, we say it is indecomposable or primitive. Equivalence two decompositions if the their corresponding covers $f_{1}: Y_{1} \rightarrow X$ are equivalent to give equivalence classes of decompositions. As $G(Y, f) \leq S_{n}$, denote the subgroup stabilizing 1 by $G(Y, f)(1)$.

LEMMA 8.12. The monodromy group $G(Y, f)$ is a primitive subgroup of $S_{n}$ if and only if $f$ is primitive (Def. 7.9). Equivalence classes of decompositions of $f$ correspond one-one with subgroups properly between $G(Y, f)$ and $G(Y, f)(1)$.

Proof. Choose a basepoint $y_{1} \in Y$ to apply Thm. 7.16. Groups between $G(Y, f)$ and $H_{1}=\{g \in G(Y, f) \mid(1) g=1\}$ correspond one-one to decompositions of $f$. In particular, $f$ is primitive if and only if there no decomposition of $f$.

Suppose $G$ and $H$ are groups, with $G_{1} \leq G$ and $H_{1} \leq H$. Let $T_{G_{1}}: G \rightarrow S_{n}$ and $T_{H_{1}}: H \rightarrow S_{m}$ be corresponding coset representations. Use $T_{G_{1}}$ to have $G$ act on $H^{n}$, the product of $n$ copies of $H$ :
(8.5) $g \in G$ acts by $\left(h_{1}, \ldots, h_{n}\right) \mapsto\left(h_{(1) T_{G_{1}}(g)}, \ldots, h_{(n) T_{G_{1}}(g)}\right)$.

This gives a natural permutation representation $T_{H \imath G}: H \imath G \stackrel{\text { def }}{=} H^{n} \times{ }^{s} G \rightarrow S_{n m}$ acting on a set $L=\left\{1_{1}, \ldots, 1_{m}, 2_{1}, \ldots, 2_{m}, \ldots, n_{1}, \ldots, n_{m}\right\}$ by this formula：

$$
\left(i_{j}\right) T_{H \imath G}\left(h_{1}, \ldots, h_{n}, g\right)=(i) T_{G}(g)_{(j) T_{H}\left(h_{i}\right)}
$$

Call $T_{H 乙 G}$ the wreath product representation of $T_{G}$ and $T_{H}$ ．Then，$H 乙 G$ is the wreath product of $G$ and $H$ ，though this assumes we know the corresponding per－ mutations representations．Now consider how the wreath occurs in covering theory．

Definition 8．13．Suppose $\psi: \hat{G} \rightarrow G$ is a cover of groups．Let $T_{G_{1}}$（resp．$T_{\hat{G}_{1}}$ ） be a faithful permutation representation of $G$（resp．$\hat{G}$ ）．Call $T_{\hat{G}_{1}}$ an extension of $T_{G_{1}}$ if $\psi$ maps some conjugate of $\hat{G}_{1}$ maps surjectively to $G_{1}: T_{\hat{G}_{1}}$ extends $T_{G_{1}}$ ．

Lemma 8．14．Suppose $f: Y \rightarrow X$ is a（connected）cover，and $f$ factors as a series of covers $Y \xrightarrow{f_{2}} Y_{1} \xrightarrow{f_{1}} X$ ．Let $G_{f_{i}}$ be the group of the Galois closure of $f_{i}$ ，with $T_{f_{i}}$ the corresponding permutation representations，$i=1,2$ ．Use similar notation for $f$ ．Then，$T_{f}$ extends $T_{f_{1}}, G_{f}$ is a transitive subgroup of $G_{f_{2}}$ 乙 $G_{f_{1}}$ and $G_{f_{1}}(1)$ maps surjectively to the group $G_{f_{2}}$ ．Further，$G_{f}=G_{f_{2}} \imath G_{f_{1}}$ if and only if the kernel of $G_{f} \rightarrow G_{f_{1}}$ is isomorphic to $G_{f_{2}}^{\operatorname{deg}\left(f_{1}\right)}$ ．

Proof．Choose a base point $y_{0} \in Y$ and so image base points in $Y_{1}$ and $X$ ．Apply Thm． 7.16 to identify $G_{f}$（resp．$G_{f_{2}}, G_{f_{1}}$ ）with permutation representa－ tions of $\pi_{1}\left(X, f\left(y_{0}\right)\right)$（resp．$\pi_{1}\left(Y_{1}, f_{2}\left(y_{0}\right)\right), \pi_{1}\left(X, f\left(y_{0}\right)\right)$ ）from the cosets of $\pi_{1}\left(Y, y_{0}\right)$ （resp．$\left.\pi_{1}\left(Y, y_{0}\right), \pi_{1}\left(Y_{1}, f_{2}\left(y_{0}\right)\right)\right)$ ．The permutation representation of $G_{f}$ comes from the image $G_{f}(1)$ of $\pi_{1}\left(Y, y_{0}\right)$ in $G_{f}$ ．Similarly，the permutation representation of $G_{f_{1}}$ comes from the image of $G_{f_{1}}(1)$ of $\pi_{1}\left(Y_{1}, f_{2}\left(y_{0}\right)\right)$ in $G_{f_{1}}$ ．As $\pi_{1}\left(Y_{1}, f_{2}\left(y_{0}\right)\right)$ contains $\pi_{1}\left(Y, y_{0}\right), T_{f}$ extends $T_{f_{1}}$ ．All coset permutation representations are transitive．

With $x_{0}=f\left(y_{0}\right)$ ，let $W=y_{1}, \ldots, y_{\operatorname{deg}\left(f_{1}\right)}$ be the points of $Y_{1}$ lying over $x_{0}$ ． Similarly，let $W_{i}=\left\{y_{i, j_{i}}\right\}_{j_{i}=1, \ldots, \operatorname{deg}\left(f_{2}\right)}$ be the points of $Y$ lying over $y_{i}$ ．Intersecting the conjugates of $G_{f}(1)$ gives $K=\operatorname{ker}\left(G_{f} \rightarrow G_{f_{1}}\right)$ ．So，$K$ acts as permutations on each $W_{i}, i=1, \ldots, \operatorname{deg}\left(f_{1}\right)$ ．Restrict $G_{f}(1)$ to $W_{1}$ for the group $G_{f_{2}}$ in the representation $T_{f_{2}}$ ．Similarly，identifying all sets $W_{i}$ ，embeds $K$ as a subgroup of $G_{f_{2}}^{\operatorname{deg}\left(f_{1}\right)}$ ．This identifies $G_{f}$ with a subgroup of the wreath product．Since the order of $G_{f}$ is $\left|G_{f_{1}}\right||K|$ ，the index of $G_{f}$ in $G_{f_{2}} 乙 G_{f_{1}}$ equals $\left(G_{f_{2}}^{\operatorname{deg}\left(f_{1}\right)}: K\right)$ ．This gives the last statement of the lemma．

8．5．Representations and groupoids．Rather than define groupoid gen－ erally，we present a classical case for later use．The idea is that of Deligne and Grothendieck．Deligne has a notion of（fundamental group）realizations．We think of these as ways a space declares its presence through types of analytic continua－ tion．This helps us to explain the profinite fundamental group of a complex manifold （Chap． 4 §8．2）．Mastering the Hurwitz monodromy group in Chap． 5 simplifies if we understand how a fundamental group depends on a base point．That leads to generalizing what will serve as a base point．Tangential base points（Chap． 2 §8．4） are an example．We get much mileage from a particularly significant parameter space，the classical $j$－line（Chap． $4 \S 7.8$ ）．This follows［De89，§10］which used the related $\lambda$－line．

8．5．1．A law of composition．Suppose $\mathcal{C}_{X}$ is the category of unramified covers of an complex manifold $X$ ．For $\varphi: Y \rightarrow X$ an unramified cover and $\psi: X^{\prime} \rightarrow X$ any map of complex manifolds，there is a natural contravariant map $\psi^{*}: \mathcal{C}_{X} \rightarrow \mathcal{C}_{X^{\prime}}$ through fiber products：$\psi^{*}(\varphi)=: X^{\prime} \times_{X} Y \rightarrow X^{\prime}$ ．

Lemma 8.15. The map $\psi^{*}$ preserves fiber products. For $\varphi_{1}, \varphi_{2} \in \mathcal{E}_{X}$ :

$$
\psi^{*}\left(\varphi_{1} \times_{X} \varphi_{2}\right)=\psi^{*}\left(\varphi_{1}\right) \times_{X^{\prime}} \psi^{*}\left(\varphi_{2}\right)
$$

Proof. Seeing this set theoretically makes it clear the cover structures are compatible. First: Identify $\left(Y_{1} \times_{X} Y_{2}\right) \times_{X} X^{\prime}$ with $\left(Y_{1} \times_{X} X^{\prime}\right) \times_{X^{\prime}}\left(Y_{2} \times_{X} X^{\prime}\right)$ by mapping $\left(y_{1}, y_{2}, x^{\prime}\right)$ all lying over a given $x \in X$ to $\left(\left(y_{1}, x^{\prime}\right),\left(y_{2}, x^{\prime}\right)\right)$. Then, both maps send this element to $x^{\prime}$.

Let $\hat{\varphi}: \hat{Y} \rightarrow X$ be the Galois closure of this cover. Suppose this has group $G$. Then $G$ acts faithfully and transitively on the fibers of $\hat{\varphi}$. On $\hat{Y} \times \hat{Y} \rightarrow X \times X$ let $G$ act diagonally: $\left(\hat{y}_{1}, \hat{y}_{2}\right) g \stackrel{\text { def }}{=}\left(\left(\hat{y}_{1}\right) g,\left(\hat{y}_{2}\right) g\right)$.

Denote $\hat{Y} \times \hat{Y} / G$, the orbits of the action of $G$, by $\mathcal{G}$. Let $\mathcal{G}_{i, j}$ be the pullback of $\mathcal{G}$ to $X \times X \times X$ induced from the projection of $X \times X \times X$ on its $(i, j)$ factors. For example, $\mathcal{G}_{1,2}$ consists of triples $\left(\hat{y}_{1}, \hat{y}_{2}, x_{3}\right)$ with $\hat{y}_{i} \in \hat{Y}, i=1,2$, and $x_{3} \in X$.

This gives a composition law $\mathcal{G}_{1,2} \times \mathcal{G}_{2,3} \rightarrow \mathcal{G}_{1,3}$ respecting fibers over $X \times X \times X$. Here is what that means. For $\left(x_{1}, x_{2}, x_{3}\right) \in X \times X \times X$, let ( $\hat{y}_{1}, \hat{y}_{2}, \hat{x}_{3}$ ) (resp. $\left.\left(x_{1}, \hat{y}_{2}^{\prime}, \hat{y}_{3}^{\prime}\right)\right)$ represent a point of $\mathcal{G}_{x_{1}, x_{2}}$ the fiber of $\mathcal{G}_{1,2}$ (resp. $\mathcal{G}_{2,3}$ ) over $\left(x_{1}, x_{2}\right)$ (resp. $\left(x_{2}, x_{3}\right)$ ). The composition law $\mathcal{G}_{x_{1}, x_{2}} \times \mathcal{G}_{x_{2}, x_{3}} \rightarrow \mathcal{G}_{x_{1}, x_{3}}$ uses the following formula. There is a unique $g \in G$ taking $\hat{y}_{2}^{\prime}$ to $\hat{y}_{2}$. Define the product of ( $\hat{y}_{1}, \hat{y}_{2}, x_{3}$ ) and $\left(x_{1}, \hat{y}_{2}^{\prime}, \hat{y}_{3}^{\prime}\right)$ to be $\left(\hat{y}_{1}, x_{2},\left(\hat{y}_{3}^{\prime}\right) g\right)$.

We say $\mathcal{G}=\hat{Y} \times \hat{Y} / G \rightarrow X \times X$ is a groupoid. Most significant is that it induces a groupoid in $\mathcal{F}_{X^{\prime}}$ by pullback, for each $\psi: X^{\prime} \rightarrow X$.
8.5.2. Fundamental groupoid. There is a fundamental groupoid that dominates all (classical) groupoids over $X$. We define this directly, as it will appear in Chap. 5.

Consider this data: $x_{1}, x_{2} \in X$, and $D_{i}$ a simply connected (path-connected) neighborhood of $x_{i}$ on $X, i=1,2$. Suppose $x_{i}^{\prime} \in D_{i}, i=1,2$. To read the next lemma correctly, emphasize the word canonical.

Lemma 8.16. There is a canonical isomorphism (dependent on $\left(D_{1}, D_{2}\right)$ ):

$$
\psi_{D_{1}, D_{2}}: \pi_{1}\left(X, x_{1}, x_{2}\right) \rightarrow \pi_{1}\left(X, x_{1}^{\prime}, x_{2}^{\prime}\right)
$$

Proof. For $\gamma_{i}$ any path from $x_{i}$ to $x_{i}^{\prime}$ in $D_{i}, i=1,2$, map $\gamma \in \pi_{1}\left(X, x_{1}, x_{2}\right)$ to $\left[\gamma_{1}^{-1} \cdot \gamma \cdot \gamma_{2}\right]=\left[\gamma_{1}^{-1}\right][\gamma]\left[\gamma_{2}\right] \in \pi_{1}\left(X, x_{1}^{\prime}, x_{2}^{\prime}\right)$. Under the hypotheses, $\left[\gamma_{i}\right]$ depends only on $x_{i}, x_{i}^{\prime}, D_{i}$ and not the particular choice of path. That shows the lemma. We will, however, confront repeatedly the dependence of $\psi_{D_{1}, D_{2}}$ on $\left(D_{1}, D_{2}\right)$.

Definition 8.17. The fundamental groupoid $\mathcal{P}_{X}$ of $X$ consists of the disjoint union $\dot{U}_{x_{1}, x_{2} \in X} \pi_{1}\left(X, x_{1}, x_{2}\right)$. The composition law for $\pi_{1}\left(X, x_{1}, x_{2}\right) \times \pi_{1}\left(X, x_{2}, x_{3}\right)$ is the usual path multipication: $\left[\gamma_{1,2}\right] \in \pi_{1}\left(X, x_{1}, x_{2}\right)$ times $\left[\gamma_{2,3}\right] \in \pi_{1}\left(X, x_{2}, x_{3}\right)$ is $\left[\gamma_{1,2}\right]\left[\gamma_{2,3}\right] \in \pi_{1}\left(X, x_{1}, x_{3}\right)$.

Restriction of $\mathcal{P}_{X}$ to the diagonal of $X \times X$ is the local system of fundamental groups $\dot{\cup}_{x_{1} \in X} \pi_{1}\left(X, x_{1}\right)$. For $x \in X$, restrict $\mathcal{P}_{X}$ to $X \times\{x\} \subset X \times X$ to get the universal cover of $(X, x)$. Now we trace through an action of a groupoid on various locally constant sets.
8.5.3. Action of a groupoid. We recognized already that the category $\mathcal{C}_{X}$ consists of locally constant finite sets on $X$. That means, given $f: Y \rightarrow X$ an unramified cover, the topology on $Y$ comes from an open cover $\mathcal{U}$ of $X$ so that $f_{U}: Y_{U} \rightarrow U$ makes of $Y_{U}$ a finite collection of disjoint copies of $U$. Generalizing the notion of covers allows defining related locally constant structures. We concentrate here on $\mathbb{V}_{X}$, the category of locally constant - or flat - vector bundles
on $X$. Suppose $V$ is a vector space over $\mathbb{C}\left(\right.$ say, $\left.\mathbb{C}^{n}\right)$. Then, there is a natural fiber preserving addition and scalar multiplication with the expected properties on $V \times U$. An object $\mathcal{V} \in \mathbb{V}_{X}$ consists of an analytic map $L: \mathcal{V} \rightarrow X$ of manifolds with an open cover $\mathcal{U}$ having the following properties.
(8.6a) For $U \in \mathcal{U}$, there is an analytic isomorphism $\psi_{U}: \mathcal{V}_{U_{i}} \rightarrow V \times U_{\gamma\left(t_{i}\right)}$ so that $L_{U}: \mathcal{V}_{U} \rightarrow U$ and $\mathrm{pr}_{U} \circ \psi_{U}: \mathcal{V}_{U} \rightarrow U$ are the same.
(8.6b) Local constancy: For $U, U^{\prime} \in \mathcal{U}$, with $U \cap U^{\prime}$, an element of $\mathrm{GL}_{n}(\mathbb{C})$ gives $\psi_{U}^{-1} \circ \psi_{U^{\prime}}$ restricted to $V \times\left(U \cap U^{\prime}\right)$ along each fiber.
(8.6c) A fiber preserving complex analytic addition and multiplication by $\mathbb{C}$ on $\mathcal{V}$ restricts over each $U \in \mathcal{U}$ to that structure on $V \times U$.
Note the right action in (8.6b). We say $\mathcal{V}$ is a rank $n$ (locally constant, or flat) bundle. Two flat bundles $\mathcal{V}_{1}$ and $\mathcal{V}_{2}$ are bundle isomorphic if there is a compatible open cover $\mathcal{U}$ for both and a fiber preserving analytic isomorphism $\psi: \mathcal{V}_{1} \rightarrow \mathcal{V}_{2}$. Suppose $\psi$ intertwines (8.6) for $\mathcal{V}_{2}$ relative to $\mathcal{U}$ to that for $\mathcal{V}_{1}$ so that for each $U \in \mathcal{U}$, an element $g_{U} \in \mathrm{GL}_{n}(\mathbb{C})$ gives $\psi_{1, U}^{-1} \circ \psi \circ \psi_{2, U}$. Then, $\psi$ is a flat isomorphism. Warning: Some bundle isomorphisms have no corresponding flat isomorphism.

ExAMPLE 8.18 (Flat bundle from a cover). Let $f: Y \rightarrow X$ be a degree $n$ cover (element of $\mathcal{C}_{X}$ ). For each $x \in X$, denote the space spanned over $\mathbb{C}$ by the points of $f^{-1}(x)$ by $V_{x}$. We explain why $\mathcal{V}_{f} \stackrel{\text { def }}{=} \dot{U}_{x \in X} V_{x}$ is a locally constant vector bundle on $X$ by taking $L_{f}$ to be the natural projection. Suppose $U \leq X$ is open, $x^{\prime} \in U$ and $f_{U}$ identifies $Y_{U}$ with $\dot{U}_{y^{\prime} \in f^{-1}\left(x^{\prime}\right)} U_{y^{\prime}}$ where $U_{y} \leq Y$ maps one-one onto $U$. This means we have $n$ sections to the map $f_{U}$. We also call these $y_{1}^{\prime}, \ldots, y_{n}^{\prime}$. So, for each $x \in U,\left\{y_{i}^{\prime}(x)\right\}_{i=1}^{n}$ is a basis for $V_{x}$. Then, we have a natural analytic manifold topology on $\mathcal{V}_{f}$ by identifying $\mathcal{V}_{f, U}$ with $\mathbb{C}^{n} \times U$ by mapping the standard basis of $\mathbb{C}^{n}$ to $y_{1}^{\prime}(x), \ldots, y_{n}^{\prime}(x)$ running over $x \in U$.

Suppose $\mathcal{P}$ is a groupoid on $X$ and $\mathcal{V} \in \mathbb{V}$. Regard $\mathcal{P}$ as a locally constant bundle of sets over $X \times X$. Consider the fiber products $\mathrm{pr}_{i}^{*}(\mathcal{V}) \stackrel{\text { def }}{=} \mathcal{V} \times_{X}(X \times X)$, using $\mathrm{pr}_{i}: X \times X \rightarrow X$, projection on the $i$ th factor, $i=1,2$. We say $\mathcal{P}$ acts on $\mathcal{V}$ if there is a fiber preserving analytic map

$$
\begin{equation*}
A_{X}: \operatorname{pr}_{1}^{*}(\mathcal{V}) \times_{X \times X} \mathcal{P} \rightarrow \operatorname{pr}_{2}^{*}(\mathcal{V}) \tag{8.7}
\end{equation*}
$$

Regard each term $\mathcal{P}, \operatorname{pr}_{1}^{*}(\mathcal{V})$ and $\operatorname{pr}_{2}^{*}(\mathcal{V})$ as a locally constant bundle over $X \times X$.
Denote the vector space $\mathbb{C}^{n}$ (with its canonical basis understood) by $V$, so that there is an action of $\mathrm{GL}_{n}(\mathbb{C})$ on the right of $V$. (To adjust to a left action on $\mathrm{GL}_{n}(\mathbb{C})$, see Ex . $[9.16 \mathrm{f}]$.) For $x_{0} \in X$, and $n$ a positive integer, consider pairs $\left(\mathcal{V}, m_{x_{0}}\right)$ with $\mathcal{V}$ a flat bundle of rank $n$, and $m_{x_{0}}$ a fixed vector space isomorphism of $\mathcal{V}_{x_{0}}$ with $V$, by $\mathbb{V}_{x_{0}, n}$. Compose $m_{x_{0}}$ with any element of $\mathrm{GL}_{n}(\mathbb{C})$ gives a natural action of $\mathrm{GL}_{n}(\mathbb{C})$ on the pairs $\left(\mathcal{V}, m_{x_{0}}\right)$.

Proposition 8.19. The fundamental groupoid $\mathcal{P}_{X}$ acts on every $\mathcal{V} \in \mathbb{V}_{X}$. Each $\left(\mathcal{V}, m_{x_{0}}\right) \in \mathbb{V}_{x_{0}, n}$ produces $\alpha_{\mathcal{V}, m_{x_{0}}} \in \operatorname{Hom}\left(\pi_{1}\left(X, x_{0}\right), \mathrm{GL}_{n}(\mathbb{C})\right)$ and the map $\left(\mathcal{V}, m_{x_{0}}\right) \mapsto \alpha_{\mathcal{V}, m_{x_{0}}}$ is one-one and onto. Flat rank $n$ bundles up to flat isomorphism correspond to elements of $\operatorname{Hom}\left(\pi_{1}\left(X, x_{0}\right), \mathrm{GL}_{n}(\mathbb{C})\right) / G$.

Proof. Consider $\left(x, x^{\prime}\right) \in X \times X$ and $[\gamma] \in \pi_{1}\left(X, x, x^{\prime}\right)$. We give an action: $A_{X}\left(\left(v, x, x^{\prime}\right),[\gamma]\right)=\left(v^{\prime}, x, x^{\prime}\right)$ with $v \in \mathcal{V}_{x}$ and $v^{\prime} \in \mathcal{V}_{x^{\prime}}$. The construction, for a given path $\gamma$, is exactly as in the proof of Lem. 7.13. We set appropriate notation.

If $\gamma:[a, b] \rightarrow X$, then there is a partition $t_{0}=a<t_{1}<\cdots<t_{n}=b$ and contractible open subsets $U_{\gamma\left(t_{i}\right)}, i=0, \ldots, n$, with $U_{\gamma\left(t_{i}\right)} \cap U_{\gamma\left(t_{i+1}\right)}$ contractible, $i=0, \ldots, n-1$, so the following holds.
(8.8a) $\psi_{U_{i}}: \mathcal{V}_{U_{i}} \rightarrow V \times U_{\gamma\left(t_{i}\right)}$ is one of the maps given by (8.6).
(8.8b) $\gamma_{\left[t_{i-1}, t_{i+1}\right]} \leq U_{\gamma\left(t_{i}\right)}, i=0, \ldots, n$, with the provisos $t_{-1}=a$ and $t_{n+1}=b$. Since the path $\gamma$ has the information about the endpoints in it, we may simplify notation by rewriting our expression for $A_{X}$ as $A_{X}(v,[\gamma])=v^{\prime}$ with $v$ (resp. $v^{\prime}$ ) in the beginning (resp. end) point of $\gamma$. Inductively define $A_{X}\left(v, \gamma_{\left[t_{0}, t_{k+1}\right]}\right)=v_{k+1}$ :

$$
A_{X}\left(A_{X}\left(v,\left[\gamma_{\left[t_{0}, t_{k}\right]}\right]\right),\left[\gamma_{\left[t_{k}, t_{k+1}\right]}\right]\right)=A_{X}\left(v_{k},\left[\gamma_{\left[t_{k}, t_{k+1}\right.}\right]\right)=\left(v_{k}\right)\left(\psi_{U_{k}}\right)^{-1} \circ \psi_{U_{k+1}}
$$

That defines the action for a particular path. We need to know the result doesn't depend on the partition, nor on the homotopy class of $\gamma$. Starting from the definition of the action on $\gamma$ with a partition, apply the General Monodromy Theorem 6.11 proof. (Our contractibility assumptions on the $U_{i}$ s allow us to use this proof.) Line-for-line this shows $A_{X}$ depends only on the homotopy class $[\gamma]$ and not on $\gamma$.

Define $\alpha_{\mathcal{V}}$ as $\prod_{k=0}^{n-1}\left(\psi_{U_{k}}\right)^{-1} \circ \psi_{U_{k+1}}$. We use that the constituent elements are in $\mathrm{GL}_{n}(\mathbb{C})$ (locally constant as a function of $x \in X$ ), and that the result is independent of the homotopy class of the path to see it is a homomorphism. Now consider when two flat bundles are flat isomorphic.

Notice that the collection of isomorphisms $\psi_{U}: \mathcal{V}_{U} \rightarrow V \times U$ gives a cocycle condition: For $U, U^{\prime}, U^{\prime \prime}$ intersecting nontrivially,

$$
\left(\psi_{U}^{-1} \circ \psi_{U^{\prime}}\right) \circ\left(\psi_{U^{\prime}}^{-1} \circ \psi_{U^{\prime \prime}}\right)=\psi_{U}^{-1} \circ \psi_{U^{\prime \prime}} .
$$

Apply Lem. 2.2 to see that $\mathcal{V}$ identifies with the disjoint union of $\cup_{U \in \mathcal{U}} V \times U$ modulo the equivalence of points on $V \times U$ with $V \times U^{\prime}$ on the overlap of $U \cap U^{\prime}$ by $\psi_{U}^{-1} \circ \psi_{U^{\prime}}$. Using this, a flat isomorphism between $\mathcal{V}_{1}$ and $\mathcal{V}_{2}$ interprets as the existence of $g_{U} \in \mathrm{GL}_{n}(\mathbb{C})$ for which

$$
g_{U}^{-1} \circ \psi_{1, U}^{-1} \circ \psi_{1, U^{\prime}} \circ g_{U^{\prime}}=\psi_{2, U}^{-1} \circ \psi_{2, U^{\prime}}
$$

In running around any path given by a sequence of $U_{i} \mathrm{~s}$, the conclusion is that $\alpha_{\mathcal{V}_{1}}$ differs from $\alpha_{\mathcal{V}_{2}}$ on this path by conjugation by $g_{U_{0}}$. That effect is determined by its effect on $m_{x_{0}}$. This concludes the proof of the theorem.
8.6. Complete reducibility and covers with equivalent flat bundles. Flat bundles appear in a few well-known papers long ago. [Gun67, p. 97], from which the author first heard of these subjects many years ago, cites [We38] and [At57]. Riemann knew of the distinction between holomorphic vector bundles and flat bundles through his investigation general ordinary differential equations versus differential equations with ordinary singular points. This topic appears in Chap. 4. An advanced reader will note we have yet to define general holomorphic bundles.
8.6.1. Decomposing the representations of a cover. A cover $f: Y \rightarrow X$ has a flat bundle on $X$ associated with it (Ex. 8.18). Let $\rho_{X} \in \operatorname{Hom}\left(\pi_{1}\left(X, x_{0}\right), \mathrm{GL}_{n}(\mathbb{C})\right)$ be the associated homomorphism. We explore the natural map $\mathcal{E}_{X} \rightarrow \mathcal{V}_{X}$, especially noting it is not injective. [Sch70] and [Fri73] are sources for practical problems in which this becomes significant. In particular, Chap. 4 [11.12] uses Riemann's Existence Theorem on the groups of [9.20] to produce primitive, inequivalent covers whose fibers products are reducible. This is a chance to introduce the significant topic of complete reducibility for fundamental groups representations.

Definition 8.20. Let $G$ be a group and $F$ a field. Suppose $\rho: G \rightarrow G L_{n}(F)$ is a representation of $G$. Then, $\rho$ has an invariant subspace $V \leq F^{n}$ if $\rho(g)$ maps $V$ into $V$ for each $g \in G$. A representation is irreducible if it has no invariant subspace. Two invariant subspaces $V$ and $W$ (for $\varphi$ ) are complements if $V$ and $W$ span $F^{n}$, and $V \cap W=\{0\}$. Call $\rho$ completely reducible if every $\rho$ invariant subspace $V$ has a complement.

Recall: With $R$ a ring, $r \in R$ is an idempotent if $r^{2}=r$. Idempotents in $\mathbb{M}_{n}(F)$ are the matrices of projection onto subspaces of $F^{n}$.

Lemma 8.21. Suppose $V$ is a $\rho$ invariant subspace. If $F$ has characteristic 0, then $V$ has a complement.

Proof. Let $P: F^{n} \rightarrow V$ be any projection onto $V$ : Choose a basis $v_{1}, \ldots, v_{k}$ of $V$, extend to a basis $v_{1}, \ldots, v_{n}$ of $V$, and define $P$ by $\sum_{i=1}^{n} a_{i} v_{i} \mapsto \sum_{i=1}^{k} a_{i} v_{i}$. Then, $P^{2}=P$ and $P$ is an idempotent. So, too is $I_{n}-P=\mathrm{P}^{\prime}$, and it defines a complementary space by projection. If $P$ commutes with the action of $G$, then $I_{n}-P$ would also be a $G$ invariant subspace. To get this, average over $G$ : Replace $P$ with $P_{G}=\frac{1}{|G|} \sum_{g \in G} \rho(g)^{-1} P \rho(g)$. Since each term $\rho(g)^{-1} P \rho(g)$ acts like the identity on $V$, for $v \in V,(v) P_{G}=\frac{1}{|G|} \sum_{g \in G}(v) \rho(g)^{-1} P \rho(g)=v$.
[9.19] applies the complete reducibility of finite group representations when $F$ has zero characteristic. Complete reducibility does in general if either $G$ is infinite or $F$ has positive characteristic [9.17]. If a representation $\rho$ is completely reducible, then we may write $F^{n}$ as $\oplus_{i=1}^{k} V_{i}$, a direct sum of invariant and irreducible subspaces for the action of $G$. Another notation for this is $\rho=\left(\rho_{1}, \ldots, \rho_{k}\right)$ with $\rho_{i}$ restriction of $\rho$ to the space $V_{i}: \rho$ is the direct sum of the actions of the $\rho_{i}, i=1, \ldots, k$.

The notation $\mathbf{1}_{G}$ is for the one-dimensional representation of $G$ where the action of $G$ leaves each vector fixed. Given any representation $\rho$ there is natural conjugate representation $\bar{\rho}: g \mapsto \bar{\rho}(g)$ by applying ${ }^{-}$to each entry of $\rho(g)$.
8.6.2. Components of fiber products. Suppose $f_{i}: Y_{i} \rightarrow X$ is a connected cover of degree $n_{i}$, with $\rho_{f_{i}} \in \operatorname{Hom}\left(\pi_{1}\left(X, x_{0}\right), \mathrm{GL}_{n_{i}}(\mathbb{C})\right)$ the corresponding element from Prop. 8.19, $i=1,2$. Then, $\rho_{f_{1}}$ and $\rho_{f_{2}}$ induce the tensor product representation $\rho_{f_{1}} \otimes \rho_{f_{2}} \in \operatorname{Hom}\left(\pi_{1}\left(X, x_{0}\right), \mathrm{GL}_{n_{1} n_{2}}(\mathbb{C})\right)$. Let $G_{i}$ be the group of a Galois closure $\hat{Y}_{i} \rightarrow X$ of $f_{i}, i=1,2$. Lem. 8.11 shows each of these representations factors through a faithful representation of $G=G_{1} \times{ }_{H} G_{2}$ for some group $H$ that is a quotient of both $G_{1}$ and $G_{2}$. Here $G$ is the group of the minimal Galois cover of $X$ factoring through $f_{1}$ and $f_{2}$. Use the notation $\rho_{f_{1}} \otimes \rho_{f_{2}}$ for this representation, too. Since $G$ is a finite group, each representation is completely reducible. Any representation of $G_{i}$ induces a representation of $G$ through the canonical projection of $G$ onto $G_{i}$. Write $\rho_{f_{i}}=\oplus_{j=1}^{k_{i}} V_{i, j}, i=1,2$, indicating the irreducible representations of $G$ coming from those of $G_{i}, i=1,2$.

Proposition 8.22. The number of connected components of the fiber product $Y_{1} \times_{X} Y_{2}$ is the same as the number of times the identity appears in $\rho_{f_{1}} \otimes \rho_{f_{2}}$. In turn, this is the same as the number of distinct pairs $\left(j, j^{\prime}\right)$ where $V_{1, j}$ is equivalent to the conjugate of $V_{2, j^{\prime}}$.

If $G=G_{1}=G_{2}$, and $\rho_{1}=\rho_{2}, Y_{1} \times_{X} Y_{2}$ has at least two connected components. In this case it has precisely two if and only if the permutation representation associated with $f_{1}$ (or with $f_{2}$ ) is doubly transitive.

Proof. Apply Lem. 8.11 to conclude there are as many connected components in $Y_{1} \times_{X} Y_{2}$ as the number of orbits in the direct product applied to $G$ of the permutation representations attached to $f_{1}$ and $f_{2}$. This counts the appearances of the identity in the corresponding representation which in turn counts the number of appearances of the identity in $\rho_{f_{1}} \otimes \rho_{f_{2}}$. Use the representation theory reminders in [9.19b] to see this also counts the number of pairs $\left(j, j^{\prime}\right)$ listed in the statement of the proposition. This completes the first part of the proof.

Suppose $\rho_{T}=\oplus_{j=1}^{k} V_{T, j}$ is the decomposition of $\rho_{1} \otimes \rho_{2}$ given in the statement into irreducible representations (over $\mathbb{C}$ ). A permutation representation is the same as its conjugate. So, for each $V_{T, j}$, its conjugate also appears in the summands of $\rho_{T}$. If $\rho_{1}=\rho_{2}$ and $G=G_{1}=G_{2}$, besides the identity in both $\rho_{1}$ and $\rho_{2}$, there must exist at least one other pair indexed by $\left(j, j^{\prime}\right)$ of conjugate representations. From [9.19d], $k=2$ if and only if the permutation representation is doubly transitive. If, however, $k \geq 3$, there will be at least three pairs ( $j, j^{\prime}$ ) indicating corresponding pairs of conjugate represenations. This concludes the proof.

## 9. Exercises

We apply group theory exercises here to geometric applications in Chap. 4. [FH91] contains a hurried encyclopedic account of classical representations. Yet, it doesn't cover our later needs. [Ben91] (very concise) and older relaxed texts like [Ha63] work for Riemann surface applications requiring deeper group theory. We have exercises that prepare some characteristic $p$ representations. These appear in Modular Towers (Chap. 5). Representation theory changes as much as Riemann surface theory. As [Lam98, p. 369] notes, it is about 100 years old. Even such topics as higher characters from its beginnings - unlike linear characters these do determine the group - have still an uncertain place in the theory.
9.1. Constructing manifolds. Call a topological space a pre-manifold if it has coordinate charts, but is not necessarily Hausdorff. We characterize Hausdorff.
(9.1a) Show the space of Ex. 2.4 is not Hausdorff.
(9.1b) Prove Lemma 2.5 using the argument before it.
(9.1c) Let $\left\{\left(X_{\alpha_{i}}, \varphi_{\alpha_{i}}\right)\right\}_{\alpha_{i} \in I_{i}}$ (resp., $\left.\left\{\left(Z_{\alpha}, \varphi_{\alpha}\right)\right\}_{\alpha \in I}\right)$ be topological data for $X_{i}$ (resp., $Z$ ), $i=1,2$. Let $f_{i}: X_{i} \rightarrow Z, i=1,2$ be continuous. Show

$$
\left\{\left(X_{\alpha_{i}} \times X_{\alpha_{j}}\right) \cap\left(X_{1} \times_{Z} X_{2}\right),\left(\varphi_{\alpha_{i}}, \varphi_{\alpha_{j}}\right)\right\}_{\left(\alpha_{i}, \alpha_{j}\right) \in I_{1} \times I_{2}}
$$

gives topologizing data on $X_{1} \times_{Z} X_{2}$ with continuous projections $\mathrm{pr}_{i}$ : $W \stackrel{\text { def }}{=} X_{1} \times_{Z} X_{2} \rightarrow X_{i}, i=1,2$. Further, $W$ is Hausdorff if $X_{1}, X_{2}$ and $Z$ are. Use this to prove Lemma 4.3.
(9.1d) Let $f: X \rightarrow Y$ be continuous, with $X$ and $Y$ pre-manifolds. Let $\gamma$ : $[0,1] \rightarrow Y$ be a path. If a continuous $\gamma_{1}:[0,1) \rightarrow X$ lies over $\gamma_{[0,1)}$ $\left(f \circ \gamma_{1}(t)=\gamma(t)\right.$ for $\left.t \in[0,1)\right)$. Show: For all pairs $\left(\gamma, \gamma_{1}\right)$, there is at most one extension of $\gamma_{1}$ to a path $\gamma_{1}^{*}:[0,1] \rightarrow Y$ if and only if the diagonal in $X \times_{Y} X$ is closed. Call an $f$ satisfying this separated.
(9.1e) With $f$ in d) separated, consider extending $\gamma_{1}$ to $\gamma_{1}^{*}:[0,1] \rightarrow Y$. Show: Such $\gamma_{1}^{*}$ exists (for each $\gamma_{1}$ ) if and only if $f$ is a proper map ( $\S 2.2$ ).
Consider some manifolds (differentiable) from vector calculus.
(9.2a) If $X_{i}$ is $n_{i}$-dimensional, $i=1,2$, show $X_{1} \times X_{2}$ is $n_{1}+n_{2}$-dimensional.
(9.2b) The $n$-sphere is $S^{n}=\left\{\boldsymbol{x} \in \mathbb{R}^{n+1} \mid \sum_{i=1}^{n+1} x_{i}^{2}=1\right\}$. Here is some data for defining a manifold structure on $S^{n}$ :

$$
U^{+}=\left\{\left(x_{1}, \ldots, x_{n+1}\right) \in S^{n} \mid x_{n+1}>0\right\}
$$

and $R_{\boldsymbol{x}}$ is any rotation of the sphere that that takes $\boldsymbol{x}$ to $(0, \ldots, 0,1)$. Let $U_{\boldsymbol{x}}$ be the image of $U^{+}$under $R_{\boldsymbol{x}}^{-1}$, and define $\varphi_{\boldsymbol{x}}$ to be $\mathrm{pr} \circ R_{\boldsymbol{x}}$ where $\operatorname{pr}(\boldsymbol{x})=\left(x_{1}, \ldots, x_{n}\right)$. Show the $\left(U_{\boldsymbol{x}}, \varphi_{\boldsymbol{x}}\right)$ 's are a differentiable atlas on $S^{n}$.
(9.2c) Consider $f \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$ and the set $X_{f}=\left\{\boldsymbol{x} \in \mathbb{R}^{n} \mid f(\boldsymbol{x})=0\right\}$. Let $X_{f}^{0}=\left\{\boldsymbol{x} \in X_{f} \mid \nabla(f)(\boldsymbol{x}) \neq 0\right\}$ (Lemma 3.2). State a differentiable version of the implicit function theorem [Rud76, p. 224] from Chap. 2 §6.2.
(9.2d) Assume $n=3$ in c) and two open sets $U_{1}$ and $U_{2}$ with these properties: $\frac{\partial f}{\partial x_{1}}$ is nonzero in $U_{1}$ and $\frac{\partial f}{\partial x_{3}}$ is nonzero in $U_{2}$. Apply c) to conclude there is a differentiable transition function $\varphi_{2} \circ \varphi_{1}^{-1}$ for the pair $\left(U_{1}, U_{2}\right)$.
(9.2e) If $X_{f}^{0}$ is nonempty, show it is a differentiable $n-1$ dimensional manifold.
(9.2f) State a complex analog of c) for $f \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ using complex partials. How does this show the complex version of $X_{f}^{0}$ is an $n-1$ dimensianal analytic manifold?
(9.2g) Apply the fundamental theorem of algebra [Ah179, p. 122] to show the manifold in f) cannot be compact.
Fiber products and pushouts are categorical constructions. Chap. 4 [11.10] continues this exploration.
(9.3a) The fiber product of two maps $f_{i}: Y_{i} \rightarrow X, i=1,2$, satisfies the following universal property: If $f: Y \rightarrow X$ factors through each of the $f_{i} \mathrm{~s}$, then $f$ factors through $\left(f_{1}, f_{2}\right)$. Further, $\left(f_{1}, f_{2}\right)$ is universal for this property.
(9.3b) The pushout for $f_{i}: Y_{i} \rightarrow X, i=1,2$, satisfies a reverse diagram to the fiber product. It is the maximal object through which both $f_{i}, i=1,2$, factor. For subsets of a set, the pushout would be the union. Show the pushout of pointed covers is exactly as given in Thm. 7.16.
(9.3c) For subgroups of a group, the union is not a group. Show the subgroup generated by the two groups is the pushout.
9.2. Complex structure and torii. Going from $\mathbb{R}$ to $\mathbb{C}$ is partly a linear algebra constraint. Use the identifications $\left\{L_{n}\right\}_{n=1}^{\infty}$ of $\mathbb{R}^{2 n}$ and $\mathbb{C}^{n}$ in §3.1.2. Consider replacing $\left\{L_{n}\right\}_{n=1}^{\infty}$ by the sequence $\left\{L_{n}^{\prime}\right\}_{n=1}^{\infty}$ of linear (invertible) maps (from $\left.\mathbb{R}^{2 n} \rightarrow \mathbb{C}^{n}\right)$. Denote $\left(x_{1}, y_{1}, \ldots, x_{n}, y_{n}\right) \mapsto\left(-y_{1}, x_{1}, \ldots,-y_{n}, x_{n}\right)$ by $J_{n}$.
(9.4a) Show with $L_{n}^{\prime}$ in place of $L_{n}$, though the functions labeled analytic in any neighborhood of an analytic manifold $X$ will change, the set of $n$ dimensional analytic manifolds remains the same.
(9.4b) Show, for analytic manifolds $X$ and $Y$ (possibly of different dimensions), the set of analytic maps $X$ to $Y$ using $\left\{L_{n}\right\}_{n=1}^{\infty}$ map naturally to the corresponding set using $\left\{L_{n}^{\prime}\right\}_{n=1}^{\infty}$.
(9.4c) Show $\left\{L_{n}^{\prime}\right\}_{n=1}^{\infty}$ gives the same analytic functions on each analytic manifold as $\left\{L_{n}\right\}_{n=1}^{\infty}$ if and only if $L_{n}^{\prime}=B_{n} \circ L_{n}$ with $B_{n} \in \mathrm{GL}_{n}(\mathbb{C})$ for all $n$. Further, this is equivalent to $L_{n}^{\prime} \circ J_{n}=i \cdot L_{n}^{\prime}$ for all $n$. Hint: Check on $\mathbb{C}$ linear combinations of $z_{1}, \ldots, z_{n}$ in $\mathbb{C}^{n}$ using $L_{n}^{\prime}$. Also: Invertible $\mathbb{R}$ linear maps $\mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$ are in $\mathrm{GL}_{n}(\mathbb{C})$ if and only if they commute with $i$.
(9.4d) Consider the case $L=L_{n}: \mathbb{R}^{2 n} \rightarrow \mathbb{C}^{n}$ by

$$
\left(x_{1}, y_{1}, \ldots, x_{n}, y_{n}\right) \mapsto\left(x_{1}-i y_{1}, \ldots, x_{n}-i y_{n}\right)=\left(\bar{z}_{1}, \ldots, \bar{z}_{n}\right)
$$

for examples where using $\left\{L_{n}^{\prime}\right\}_{n=1}^{\infty}$ changes a given analytic structure. Hint: See Chap. 4 §??.
Consider the topology of the torus of Fig. 3.
(9.5a) Show the complex torus $\mathbb{C} / L\left(\omega_{1}, \omega_{2}\right)$ of $\S 3.2 .2$ is compact.
(9.5b) Suppose $R>3 r$ with $r, R \in \mathbb{R}$. The torus, $T_{r, R ; \boldsymbol{x}_{0}, \boldsymbol{v}}=$, with radii $(r, R)$ centered at $\boldsymbol{x}_{0}=(0,0,0) \in \mathbb{R}^{3}$ and perpendicular to $\boldsymbol{v}=(0,0,1)$ has this underlying set of points:

$$
\left\{\boldsymbol{x}_{0}+R(\cos (\theta), \sin (\theta), 0)+r(\cos (\theta) \cos (\beta), \sin (\theta) \cos (\beta), \sin (\beta))\right\}_{\theta, \beta \in[0,2 \pi]}
$$

Show $T_{r, R ; \boldsymbol{x}_{0}, \boldsymbol{v}}$ is differentiably isomorphic to $\mathbb{C} / L\left(\omega_{1}, \omega_{2}\right)$.
(9.5c) Consider the two torii in Fig. 2: Assume one is $T=T_{r, R ; \boldsymbol{x}_{0}, \boldsymbol{v}}$, the other $T^{\prime}=T_{r, R ; \boldsymbol{x}_{0}^{\prime}, \boldsymbol{v}^{\prime}}$ for vectors $\boldsymbol{x}_{0}^{\prime}, \boldsymbol{v}^{\prime} \in \mathbb{R}^{3}$ and $T \cap T^{\prime}=\emptyset$. Call $T$ and $T^{\prime}$ unknotted if for any $C>0$ there is a continuous function

$$
F:[0,1] \times \mathbb{R}^{3} \backslash T \rightarrow \mathbb{R}^{3} \backslash T
$$

with $F(0, y)=y$ for $y \in \mathbb{R}^{3} \backslash T$ and $|F(1, y)|>C$ for $y \in T^{\prime}$. Otherwise they are knotted. Show there are two knotted torii in $\mathbb{R}^{3}$.
(9.5d) Regard $\mathbb{R}^{3}$ as in in $\mathbb{R}^{4}$ : It is the set of $\boldsymbol{x} \in \mathbb{R}^{4}$ with $x_{4}=0$. Extend the definitions above to show any pair of torii in $\mathbb{R}^{3}$ is unknotted in $\mathbb{R}^{4}$.
We start discussing the nature of the lattice attached to a complex torus.
(9.6a) Let $\mathbb{C} / L\left(\omega_{1}, \omega_{2}\right)=X$ be a complex torus with lattice $L\left(\omega_{1}, \omega_{2}\right)=L$ as in Ex. 6.18. For $z_{1}, z_{2} \in \mathbb{C}$ define $m\left(z_{1} \bmod L, z_{2} \bmod L\right)$ to be $z_{1}+$ $z_{2} \bmod L$. Define the inverse of $z \bmod L$ to be $-z \bmod L$. Show $X$ is a differentiable group with multiplication $m$.
(9.6b) For $t \in \mathbb{R}$, let $z(t)=\cos (2 \pi t)+\sqrt{-1} \sin (2 \pi t)$. Use $f: X \rightarrow S^{1} \times S^{1}$ by $t_{1} \omega_{1}+t_{2} \omega_{2} \mapsto\left(z_{1}(t), z_{2}(t)\right)$ to conclude that $\pi_{1}(X, 0 \bmod L)$ identifies with $L$ as a group isomorphic to $\mathbb{Z}^{2}$, pairs of integers.
(9.6c) Suppose $x_{1}, x_{2} \in S^{1}$ generate an infinite group $\left\langle x_{1}, x_{2}\right\rangle$. Consider the collection $T_{N}=\left\{x_{1}^{j} x_{2}^{j^{\prime}}\right\}_{-N \leq j, j^{\prime} \leq N}$ for large $N$ to conclude 1 is a limit point for $\left\langle x_{1}, x_{2}\right\rangle$. Conclude: $w_{1}, w_{2} \in \mathbb{C}, \mathbb{C} / L\left(w_{1}, w_{2}\right)$ satisfies the conditions of Lem. 2.3 only if $w_{1}, w_{2}$ lie on different lines through the origin.
Consider comparing two lattices of complex torii. With $L_{i}=L\left(\omega_{1, i}, \omega_{2, i}\right)$, $i=1,2$, continue Ex. 6.18. Assume $\lambda_{i}=\frac{\omega_{1, i}}{\omega_{2, i}} \in \mathbb{C} \backslash \mathbb{R}, i=1,2$.
(9.7a) Assume $\lambda_{2}=\frac{a \lambda_{1}+b}{c \lambda_{1}+d}$ for some $a, b, c, d \in \mathbb{Z}$ for $a d-b c=1$. Show $\mathbb{C} / L_{1}=$ $X_{1}$ and $\mathbb{C} / L_{2}=X_{2}$ are analytically isomorphic. Hint: Map $t_{1} \omega_{1,1}+t_{2} \omega_{2,1}$ to $t_{1}\left(a \omega_{1,1}+b \omega_{2,1}\right) \alpha+t_{2}\left(c \omega_{1,1}+d \omega_{2,1}\right) \alpha$ with $\alpha \in \mathbb{C}$ satisfying

$$
\left(a \omega_{1,1}+b \omega_{2,1}\right) \alpha=\omega_{1,2} \text { and }\left(c \omega_{1,1}+d \omega_{2,1}\right) \alpha=\omega_{2,2}
$$

(9.7b) Why assume $a d-b c=1$ in a)? Why must we have $a, b, c, d$ in $\mathbb{Z}$, rather than just $a, b, c, d \in \mathbb{R}$ ?
(9.7c) Suppose $L_{1} \subset L_{2}$. Consider $f: X_{1} \rightarrow X_{2}$ given in Ex. 6.18. Show there exist $\omega_{1}, \omega_{2} \in L_{2}$ and $n_{1}, n_{2} \in \mathbb{Z}$ with these properties: $L\left(\omega_{1}, \omega_{2}\right)=L_{1}$; and the complex numbers

$$
z\left(k_{1}, k_{2}\right)=\left(\frac{k_{1}}{n_{1}}\right) \omega_{1}+\left(\frac{k_{2}}{n_{2}}\right) \omega_{2}, \quad 0 \leq k_{i} \leq n_{i}, i=1,2
$$

give the $n_{1} n_{2}$ distinct elements $z \bmod L_{1}$ mapping to $0 \bmod L_{2}$. Hint: Apply the Elementary Divisor Theorem Chap. 2 [9.15] to get a basis $\left\{\boldsymbol{u}_{i}\right\}_{i=1}^{2}$ of $L_{2}$ and integers $n_{1}, \ldots, n_{2}$ so that $\left\{n_{i} \boldsymbol{u}_{i}\right\}_{i=1}^{2}$ generates $L_{1}$.
(9.7d) Conclude for $x \in X_{1}$ that $x+z\left(k_{1}, k_{2}\right) \bmod L_{1}$ are the distinct elements of $X_{1}$ mapping $f(x)$ under $f$.
Now we describe holomorphic differentials on a complex torus.
(9.8a) Let $L$ be a lattice in $\mathbb{C}_{z}$. Define $\omega_{\alpha}$ on one of the local coordinate charts $\varphi_{\alpha}\left(U_{\alpha}\right) \subset \mathbb{C}_{z}$ for $\mathbb{C} / L$ to be the differential $d z$ (As in Ex. 6.18). Show this defines a global differential form $\omega_{L}$ on $\mathbb{C} / L$, and the divisor of this form is 0 . Hint: Use that the transition functions, on connected subsets of $\varphi_{\alpha}\left(U_{\alpha} \cap U_{\beta}\right)$ have the form $z \mapsto z+\beta$.
(9.8b) Accept without proof that any meromorphic function has divisor of degree 0 . Conclude: Holomorphic differentials on $\mathbb{C} / L$ have degree 0 divisor; so they are constant multiples of $\omega_{L}$.
(9.8c) A $g$ dimensional complex torus has the form $A=\mathbb{C}^{g} / L$ where $L$ is a $\mathbb{Z}$ module having dimension $2 g$ and such that $\mathbb{R} L=\mathbb{C}^{g}$ (a lattice). Imitate b) to show holomorphic differentials on $A$ form a dimension $g$ vector space.
[9.8c] considers complex torii. Since $\mathbb{C}^{g}$ is contractible, $\pi_{1}(A, 0)$ identifies with $L$. We now see all differentiable groups have an abelian fundamental group.
(9.9a) Suppose that $\gamma_{0, i}$ and $\gamma_{1, i}$ are homotopic paths in a space $X, i=1,2$, and that the end point of $\gamma_{0,1}$ is equal to the initial point of $\gamma_{0,2}$. Show $\gamma_{0,1} \gamma_{0,2}$ is homotopic to $\gamma_{1,1} \gamma_{1,2}$.
(9.9b) Show the associative rule for multiplying paths.
(9.9c) Let $\psi_{1}$ and $\psi_{2}$ be two isomorphisms between $\pi_{1}\left(X, x_{0}\right)$ and $\pi_{1}\left(X, x_{1}\right)$ as in Corollary 1.19. Show $\psi_{2}^{-1} \circ \psi_{1}$ is an inner automorphism of $\pi_{1}\left(X, x_{0}\right)$. That is, it is given by conjugation by an element of $\pi_{1}\left(X, x_{0}\right)$.
(9.9d) A group $G$ is differentiable $G$ if it is a differentiable manifold, and its multiplication and inverse are both differentiable maps. Similarly, there is the notion of analytic group. Show a complex torus $\mathbb{C}^{g} / L$ ( $L$ a lattice) is an analytic group.
(9.9e) Suppose $M$ is a subvariety of $\mathrm{GL}_{n}(\mathbb{C})$ (defined by a finite number of equations in the $n^{2}$ coordinates of the entries), closed under multiplication and inverse. Show $M$ is an analytic group.
(9.9f) For $G$ a differentiable group consider $f_{1}: G \rightarrow(G, 1)$ (resp. $f_{2}: G \rightarrow$ $(1, G))$ by $g \mapsto(g, 1)$ (resp., $g \mapsto(1, g))$. Show for $\left[\gamma_{1}\right],\left[\gamma_{2}\right] \in \pi_{1}(G, 1)$ :

$$
\left.m_{*}\left(\left(f_{1}\right)_{*}\left[\gamma_{1}\right]\right)\left(f_{2}\right)_{*}\left[\gamma_{2}\right]\right)=\left[\gamma_{1}\right]\left[\gamma_{2}\right] .
$$

( 9.9 g ) Continuing b), show $\pi_{1}(G, 1)$ is an abelian group. Conclude: A differentiable manifold $X$ with a nonabelian fundamental group (as often in Chap. 4) has no differentiable group structure.
9.3. $\mathbb{P}^{n}$ compactification. Use the notation of $\S 4.3$.
(9.10a) Consider $h \in \mathbb{C}(w), h=h_{1} / h_{2}$, with $\left(h_{1}, h_{2}\right)=1$. Let $m=h_{2}(w) z-h_{1}(z)$ as in Ex. 4.7. Show the $\mathbb{P}_{z}^{1} \times \mathbb{P}_{w}^{1}$ compactification of $\{(z, w) \mid m(z, w)=$ $0, z \notin z\}$ is a manifold.
(9.10b) Consider $\mathbb{P}^{n}=\mathbb{C}^{n+1} \backslash\{0\} / \mathbb{C}^{*}$. Induct on $n$ to show $\mathbb{P}^{n}=U_{0} \cup \mathbb{P}^{n+1}$. Inductively define a topology: neighborhoods of $\boldsymbol{x} \in \mathbb{P}^{n-1}$ are the image in $\mathbb{P}^{n}$ of neighborhoods of $\left(0, v_{1}, \ldots, v_{n}\right) \in \mathbb{C}^{n+1}$.
(9.10c) Prove directly in $\mathbb{P}^{n}$ : Any infinite sequence has a limit point. Hint: Any infinite sequence has an infinite subsequence in $U_{i}$ for some $i$.
Fiber products help construct new manifolds from old. Use notation of §4.2.3.
(9.11a) Generalize the $\mathbb{P}^{2}$ compactification of $h(w)-g(z)$ from Ex. 4.3.3.
(9.11b) Conclude the proof of Prop. 4.9 by noting $\mathcal{L}_{z^{\prime}}^{h}\left[\left(z-z^{\prime}\right)^{1 / e_{1}},\left(z-z^{\prime}\right)^{1 / e_{2}}\right]$ is a proper subring of $\mathcal{P}_{z^{\prime},\left[e_{1}, e_{2}\right]}^{h}$, though its quotient field equals $\mathcal{P}_{z^{\prime},\left[e_{1}, e_{2}\right]}$.
(9.11c) Finish the hyperelliptic case of Ex. 4.3.3: $\mathbb{P}^{1} \times \mathbb{P}^{1}$-compactification gives a manifold while no $\mathbb{P}^{2}$-compactification ever does.
(9.11d) Apply b) to $f: X \rightarrow \mathbb{P}_{z}^{1}$ of degree at least 3 . Then, $V=X \times_{\mathbb{P}_{z}^{1}} X$ contains the diagonal $\Delta$ and it consists of the union of this and another compact set $V^{\prime}$. Show $V^{\prime}$ has a manifold structure from its embedding in $X \times X$ if and only if there is only one ramified point over each branch point of $f$ and that ramification order is 2 . That is, $f$ is a simple-branched cover.
(9.11e) Show global meromorphic functions on $\mathbb{P}^{n}$ are ratios of (same degree) homogeneous polynomials in the coordinates of $\mathbb{P}^{n}$. Show there is no analytic map $\psi: \mathbb{P}^{2} \rightarrow \mathbb{P}^{1}$. Hint: A ratio of same degree polynomials has a singularity at common zeros.
(9.11f) Assume $\bar{X} \subset \operatorname{pr}_{z, w, u}^{2}$ is a compact manifold, and $\left(z_{0}, w_{0}, u_{0}\right) \in \bar{X}$ is the intersection of $L_{1}$ and $L_{2}$ in Prop. 4.13. Show there is no other value $z_{0}^{\prime} \neq z_{0}$ so $L_{1}-z_{0}^{\prime} L_{2}$ is tangent to $\bar{X}$. Hint: Otherwise, $u^{\prime}=\left(L_{1}-z_{0}^{\prime} L_{2}\right) / z$ and $w^{\prime}=\left(L_{1}-z_{1}^{\prime} L_{2}\right) / z$ give local coordinates for $\bar{X}$ in a neighborhood of $(0,0) \in \in \mathbb{C}_{u^{\prime}} \times \mathbb{C}_{w^{\prime}}$ though both functions ramify at $(0,0)$.
9.4. Paths and vector fields. Let $X$ be a manifold.
(9.12a) Show each (simplicial) path $\gamma:[a, b] \rightarrow X$ is image equivalent to $\gamma_{1}$ : $[0,1] \rightarrow X$. Show each nonconstant path is image equivalent to a path constant on no interval.
(9.12b) Assume $X$ is contractible (Def. 5.8). Suppose $\gamma:[a, b] \rightarrow X$ is a path with initial point $x_{0}$ and endpoint $x_{1}$. Form the function $G:[a, b] \times[0,1] \rightarrow X$ by $G(t, s)=f(\gamma(t), s)$. Use this to show all paths in $X$ with initial point $x_{0}$ and endpoint $x_{1}$ are homotopic.
(9.12c) Let $\gamma:[a, b] \rightarrow \mathbb{R}^{n}$ be a simplicial path. Let $\boldsymbol{f}=\left(f_{1}, \ldots, f_{n}\right): \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be defined and continuous on the image of $[a, b]$. Consider

$$
\sum_{i=1}^{n} \int_{a}^{b} f_{i}(\gamma(t)) \frac{d \gamma_{i}}{d t} d t \stackrel{\text { def }}{=} \int_{\gamma} \boldsymbol{f} \cdot d \boldsymbol{x}
$$

the line integral of $\boldsymbol{f}$ along $\gamma$.
(9.12d) If $\gamma_{1}$ and $\gamma$ are image equivalent paths in $\mathbb{R}^{n}$, show line integrals along them are equal (use change of variables formula from Chap. 2 Lem. 2.3).
(9.12e) Let $F:[a, b] \times[0,1] \rightarrow \mathbb{R}^{n}$ be a homotopy between paths $\gamma_{0}$ and $\gamma_{1}$ (write $\left.F(t, s)=\gamma_{s}(t)\right)$ in $\mathbb{R}^{n}$. Assume $\boldsymbol{f}$ is continuous on the image of $F$. Show the line integral of $\boldsymbol{f}$ along $\gamma_{s}$ is a continous function of $s$.
For a differentiable path $\gamma:[0,1] \rightarrow U$ with $U$ open in $\mathbb{R}^{n}$, there may not exist a vector field $T_{U}$ having $\gamma$ as an integral curve, though locally this is so.
(9.13a) If $T_{U}$ exists explain why $\gamma\left(t_{1}\right)=\gamma\left(t_{2}\right)$ implies $\frac{d \gamma}{d t}\left(t_{1}\right)=\frac{d \gamma}{d t}\left(t_{2}\right)$.
(9.13b) Let $V$ be a neighborhood of the line segment $t \rightarrow(t, 0, \ldots, 0) \in \mathbb{R}_{t}^{n}$, $t \in[0,1]$. Assume there is a one-one differentiable $\Gamma: V \rightarrow U$ with $\Gamma(t, 0, \ldots, 0)=\gamma(t)$. Show $\frac{\partial \gamma}{\partial t_{1}}(\boldsymbol{t})$ (applying $\frac{\partial}{\partial t_{1}}$ to all coordinates of $\Gamma$ ) produces a vector field on $\Gamma(V)$ with $\gamma$ an integral curve.
(9.13c) Assume $\frac{d \gamma}{d t}$ is never 0. Consider $H_{t}=\left\{\boldsymbol{w} \in \mathbb{R}^{n} \left\lvert\, \boldsymbol{w} \cdot \frac{d \gamma}{d t}=0\right.\right\}$. Find differentiable one-one $F: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ with $F(\boldsymbol{t})=\gamma\left(t_{1}\right)+\boldsymbol{w}\left(t_{1}, t_{2}, \ldots, t_{n}\right)$
with $\boldsymbol{w}\left(t_{1}, t_{2}, \ldots, t_{n}\right) \in H_{t_{1}}$ linear in $\left(t_{2}, \ldots, t_{n}\right)$ ( $t_{1}$ fixed). Hint: Apply the chain rule.
(9.13d) How does $F$ give $\Gamma$ in b)?

Returning to (5.3) we relate

$$
\left(f_{\alpha, 1}\left(\boldsymbol{y}_{\alpha}\right), \ldots, f_{\alpha, n}\right)\left(\boldsymbol{y}_{\alpha}\right) \text { to }\left(f_{\beta, i}, \ldots, f_{\beta, n}\right)\left(\psi_{\beta, \alpha}\left(\boldsymbol{y}_{\alpha}\right)\right)
$$

(9.14a) Apply both sides of (5.3) to the coordinate function $y_{\beta, j}$ to get

$$
f_{\beta, j}\left(\psi_{\beta, \alpha}\left(\boldsymbol{y}_{\alpha}\right)\right)=\sum_{i=1}^{n} f_{\alpha, i} \frac{\partial \psi_{\beta, \alpha, j}}{\partial y_{\alpha, i}}\left(\boldsymbol{y}_{\alpha}\right)
$$

where $\psi_{\beta, \alpha, j}$ is the $j$ th coordinate of $\psi_{\beta, \alpha}$. That is, the $f_{\beta}$ s are the result of applying the Jacobian matrix of $\psi_{\beta, \alpha}\left(\boldsymbol{y}_{\alpha}\right)$ to the $f_{\alpha} \mathrm{s}$.
(9.14b) Consider the case $\psi=\psi_{(x, y),(r, \theta)}: \mathbb{R}_{r, \theta}^{2} \rightarrow \mathbb{R}_{x, y}^{2}$ by $(r, \theta) \mapsto(x, y)$. Express $\frac{\partial}{\partial x}$ as $f_{r} \frac{\partial}{\partial r}+f_{\theta} \frac{\partial}{\partial \theta}$ by applying both to $x=r \cos (\theta)$ and $y=r \sin (\theta)$. Do the same for $\frac{\partial}{\partial y}$, expressing it as $f_{r}^{\prime} \frac{\partial}{\partial r}+f_{\theta}^{\prime} \frac{\partial}{\partial \theta}$. Applying $\left(\frac{\partial}{\partial x}, \frac{\partial}{\partial y}\right)$ to $f(x, y)$ and evaluating at $(r \cos (\theta), r \sin (\theta))$ is the same as applying $J\left(\psi_{(x, y),(r, \theta)}\right)\left(\frac{\partial}{\partial r}, \frac{\partial}{\partial \theta}\right)$ to $f(r \cos (\theta), r \sin (\theta))$.
(9.14c) Generalize b) to say (as in (5.4))

$$
J\left(\psi_{\boldsymbol{y}_{\beta}, \boldsymbol{y}_{\alpha}}\right)^{-1}\left(\frac{\partial}{\partial y_{\alpha, 1}}, \ldots, \frac{\partial}{\partial y_{\alpha, n}}\right)=\left(\frac{\partial}{\partial y_{\beta, 1}}, \ldots, \frac{\partial}{\partial y_{\beta, n}}\right) .
$$

9.5. Permutation groups. Suppose $G \leq S_{n}$ is transitive. Def. 7.9 defines primitive subgroup of $S_{n}$.
(9.15a) For $g \in N_{G}(G(1))$, multiplication of $g$ on the left of the distinct right cosets $G(1) \sigma_{1}, \ldots, G(1) \sigma_{n}$ of $G(1)$ permutes these cosets. Conclude: This induces a homomorphism $\psi: N_{G}(G(1)) / G(1) \rightarrow \operatorname{Cen}_{S_{n}}(G)$.
(9.15b) Show $\psi$ is an isomorphism because both groups have order equal

$$
\mid\{i \in\{1,2, \ldots, n\} \mid \sigma(i)=i \text { for each } \sigma \in G(1)\} \mid
$$

(9.15c) Show $N_{G}(G(1)) / G(1)\left(\right.$ or $\left.\operatorname{Cen}_{S_{n}}(G)\right)$ is trivial if $G$ is primitive and $G(1)$ is nontrivial.
(9.15d) Show a nontrivial normal subgroup of a primitive group is transitive.
(9.15e) Show a primitive subgroup of $S_{n}$ containing a 2 -cycle is $S_{n}$. Conclude any transitive group generated by 2 -cycles is $S_{n}$. Hint: Consider the normal subgroup generated by the conjugates of the 2-cycle.
Let $G$ be a centerless group, $\operatorname{Aut}(G)$ its automorphisms and $T: G \rightarrow S_{n}$ faithful transitive permutation representation.
(9.16a) Explain this from [Isa94, p. 43]: In general neither $(g H) A\left(g^{\prime}\right) \stackrel{\text { def }}{=} g H g^{\prime}$ nor $(g H) A\left(g^{\prime}\right) \stackrel{\text { def }}{=}\left(g g^{\prime}\right) H$ define an action on left cosets of $H$ in $G$.
(9.16b) Let $S$ be the collection of conjugates of a subgroup $H$ of the group $G$, with the action by conjugation by elements of $G$ : $S=g^{-1} H g_{g \in G}$ and the right action of $g^{\prime} \in G \mapsto\left(g^{\prime}\right)^{-1} g^{-1} H g g^{\prime}$. What is the coset representation associated with this transitive action, and when is it faithful?
(9.16c) Show (conjugation by) $G$ is normal in $\operatorname{Aut}(G)$. The outer automorphism group $\operatorname{Out}(G)$ of $G$ is the quotient $\operatorname{Aut}(G) / G$. Show the natural map $\psi_{T}$ : $N_{S_{n}}(G) \rightarrow \operatorname{Out}(G)$ has kernel $\operatorname{Cen}_{S_{n}}(G)$ (§7.1.3; compare with [9.15c]).
(9.16d) Denote the image of $\psi_{T}$ in $\operatorname{Out}(G)$ by Out ${ }_{T}(G)$. Show $\operatorname{Out}_{T}(G)=\operatorname{Out}(G)$ if and only if $G(T, 1)$ (§7.1.2) has exactly $n$ images under Aut $(G)$. Hint: Associate to $\alpha \in \operatorname{Aut}(G)$ an element of $S_{n}$ defined up to $\operatorname{Cen}_{S_{n}}(G)$ if it maps among the conjugates of $G(T, 1)$. Show [9.20b] gives examples where $T$ is doubly transitive and $\operatorname{Out}(G) \neq \mathrm{Out}_{T}(G)$.
(9.16e) Case: $G=A_{n}$ (resp. $G=S_{n}$ ), $n \geq 4$, in its standard representation $T$. Show $\operatorname{Out}\left(S_{n}\right)=\{1\}\left(\right.$ resp. $\left.\operatorname{Out}_{T}\left(A_{n}\right)=\operatorname{Out}\left(A_{n}\right)=\mathbb{Z} / 2\right)$ if and only if $S_{n}$ (resp. $A_{n}$ ) has exactly $n$ transitive subgroups of index $n$ under Aut $(G)$. Hint: Intransitive subgroups have small orders. (See [9.17b].)
(9.16f) Set notation in the proof of Prop. 8.19 to change to a left action of $\mathrm{GL}_{n}$.

We will need the following facts later.
(9.17a) For each $i, 2 \leq i \leq n$, consider $L_{i}=\left\{1_{n},(1 i),(2 i), \ldots,(i-1 i)\right\} \subset S_{n}$ ( $1_{n}$ indicates the identity). Show each $x \in S_{n}$ has a unique product representation as $x=x_{1} x_{2} \ldots x_{n}$ with $x_{i} \in L_{i}$. (This gives a technique to generate random elements of $S_{n}$ with uniform distribution.) Hint: For $g \in S_{n}$ if $(n) g=i$, let $h=g(i n)$ and induct on $n$.
(9.17b) [Isa94, p. 79-80] bases $\operatorname{Out}\left(S_{n}\right)=\{1\}$, if $n \neq 6$, on two observations:

- If $\alpha \in \operatorname{Aut}\left(S_{n}\right)$ permutes transpositions, then conjugating by some $g \in S_{n}$ gives $\alpha$. Hint: Elements of $\left(L_{i}\right) \alpha$ in a) then have a unique integer of common support.
- If $n \neq 6$, among elements of order 2 , the conjugacy class of transpositions has a unique cardinality.
(9.17c) Let $T_{H}: G \rightarrow S_{n}$ be a permutation representation. Show all cosets of $H$ have the form $H g^{i}, i=0, \ldots, n-1$, if and only if $g$ is an $n$-cycle in $T_{H}$.
(9.17d) Suppose $F$ has characteristic $p$ which also divides the order of finite group $G$. Show a faithful permutation representation of $G$ cannot be completely reducible. Hint: Reduce to $G=\langle g\rangle$ with $g$ having order $p$.
(9.17e) Suppose $G$ is a free group on $r \geq 2$ generators. Find representations $\varphi: G \rightarrow \mathrm{GL}_{r}(\mathbb{C})$ that are not completely reducible. Hint: Map $G$ into an upper-triangular, not diagonal, matrix group.
9.6. Affine groups as permutation representations. Let $H \leq G L_{k}(F)$ with $F=\mathbb{F}_{q}$. Regard $G=\left\{\left.\left(\begin{array}{cc}A & v \\ 0 & 1\end{array}\right) \right\rvert\, A \in H, v \in V=\mathbb{F}_{q}^{k}\right\}$ as a group $V \times^{s} H$ as in Rem. 7.4. Note: If a nonabelian group replaced $\mathbb{F}_{q}^{k}$, then $A\left(v^{\prime}\right)+v$ should more naturally be written $v+A\left(v^{\prime}\right)$.
(9.18a) Suppose $\{0\}<V_{1}<V$ is an $H$ invariant space. Then, $V_{1} \times^{s} H$ is a subgroup of $G$ properly containing $H$. Show conversely, a group properly between $H$ and $G$ has the form $V_{1} \times{ }^{s} H$ with $H$ invariant $V_{1}$.
(9.18b) Embed $V$ in $G$ by $v \mapsto\left(\begin{array}{cc}1 & v \\ 0 & 1\end{array}\right)$. Have $G$ act on $V$ by $\left(\begin{array}{cc}A^{\prime} & v^{\prime} \\ 0 & 1\end{array}\right)$ maps $v \mapsto A(v)+v^{\prime}=v^{*}$ : equivalent to $\left(\begin{array}{cc}A^{\prime} & v^{\prime} \\ 0 & 1\end{array}\right)$ multiplies $\binom{v}{1}$ to $\binom{v^{*}}{1}$. Show this gives a faithful transitive permutation representation of $G$.
(9.18c) From a) the representation of b) is primitive if and only if $H$ acts irreducibly. Suppose $H=\langle A\rangle$ has a single matrix generator, which we use to makes $V$ into an $F[z]$ module by having $f(z) \in F[z]$ map $v \in V$ to $f(A)(v)$. The elementary divisor theorem (Chap. 2 §9.15) says $V \equiv \oplus_{i=1}^{t} F[z] /\left(f_{i}\right)$
(as an $F[z]$ module). Example: If $v=(a, b) \in F^{2}$, and $A=\left(\begin{array}{ll}1 & 3 \\ 1 & 2\end{array}\right)$,
$f(z)=z^{2}+z+1$, then $f(z)(v)=\left(\begin{array}{cc}6 & 12 \\ 4 & 10\end{array}\right)(v)=(6 a+12 b, 4 a+10 b)$.
We can uniquely choose the $f_{i}$ s monic so that $f_{1}\left|f_{2}\right| \cdots \mid f_{t}$. Show $G$ is primitive if and only if $t=1$ and $f_{1}$ is an irreducible polynomial.
(9.18d) The multiplicative group $\mathbb{F}_{p^{n}}^{*}$ is cyclic. Let $\alpha$ be a generator, and $A$ the matrix of $\alpha$ acting on $\mathbb{F}_{p}^{n}$ by regarding it as $\mathbb{F}_{p^{n}}$. In equation form: $v \in \mathbb{F}_{p^{n}} \mapsto \alpha v \in \mathbb{F}_{p^{n}}$. Show $V \times^{s}\langle A\rangle$ is doubly transitive on $F$.
(9.18e) From Def. 7.9, b) is doubly transitive if and only if $H$ is transitive on $V \backslash\{0\}$. When $H=\langle A\rangle$, show $G$ is doubly transitive if and only if, for some isomorphism of $\mathbb{F}_{p^{n}}$ and $\left(\mathbb{F}_{p}\right)^{n}, A$ acts like multiplication by $\alpha \in \mathbb{F}_{p^{n}}^{*}$.
9.7. Group representations. In this exercise consider representations over any field containing $\mathbb{Q}$.
(9.19a) Show that the direct product of two permutation representations as a group representation is the tensor product of the two group representations. Therefore the trace is the product of the traces.
(9.19b) Finish showing the number of orbits is the same as the number of appearances of the identity.
(9.19c) Let $T_{i}: G \rightarrow S_{n_{i}}, i=1,2$, be permutation representations for which $\mathrm{t}\left(T_{1}(g)\right)=\mathrm{t}\left(T_{2}(g)\right)$ for each $g \in G$ (as in §7.1). Show $n_{1}=n_{2}$ and $T_{1}(g)$ and $T_{2}(g)$ are conjugate in $S_{n_{1}}$ for each $g \in G$. Hint: Induct on the length of the highest disjoint cycles and compare $\mathrm{t}\left(T_{1}(g)\right)$ and $\mathrm{t}\left(T_{1}\left(g^{r}\right)\right)$ for some prime $r$ dividing the order of $g$.
(9.19d) Show $\frac{1}{|G|} \sum_{g \in G} \mathrm{t}(T(g))$ counts the orbits of a permutation representation $T$. Hint: Put the additive operator t on the outside of the sum by regarding $T(g)$ as a permutation matrix. Each orbit $I$ gives a 1-dimensional invariant subspace spanned by $\sum_{i \in I} x_{i}$ (as in §7.1.4).
(9.19e) Show the collection of $L_{\mathrm{C}}=\sum_{u \in C} u$ with C a conjugacy class of $G$, span the $G$ invariant idempotents of $\mathbb{C}[G]$. For $\rho$ any representation, $\frac{1}{|G|} \sum_{g \in G} \mathrm{t}(\rho(g))$ counts appearances of $\mathbf{1}_{G}$ in $\rho$. Hint: $\frac{1}{|G|} \sum_{g \in G} \rho(g)$ is an idempotent, and its trace equals the dimension of its range.
(9.19f) Orthogonality Relations: Let $\rho_{V}$ and $\rho_{W}$ be representations of $G$ on respective spaces $V$ and $W$. Show $\mathrm{t}\left(\rho_{V^{*} \otimes W}(g)\right)=\mathrm{t}\left(\rho_{V}(g)\right) \mathrm{t}\left(\rho_{W}(g)\right)$ gives

$$
\sum_{g \in G} \mathrm{t}\left(\bar{\rho}_{V}(g)\right) \mathrm{t}\left(\rho_{W}(g)\right)=\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{\mathbb{C}[G]}(V, W)
$$

Further, this dimension gives the appearances of $\mathbf{1}_{G}$ in $\operatorname{Hom}_{\mathbb{C}}(V, W)$. Hint: For $\rho^{\prime}$ irreducible, $\mathbf{1}_{G}$ appears exactly once in $\operatorname{Hom}_{\mathbb{C}}\left(V_{\rho^{\prime}}, V_{\rho^{\prime}}\right)$.
(9.19g) Show $V_{T}$ is $\mathbf{1} \oplus V^{\prime}$ with $V^{\prime}$ irreducible if and only if $T$ is doubly transitive. Hint: Apply d) to count appearances of $\mathbf{1}_{G}$ in $V_{T} \otimes V_{T}$; use that $\mathbf{1}_{G}$ appears in $\rho_{1} \otimes \rho_{2}$ with $\rho_{1}, \rho_{2}$ irreducible only if $\rho_{2}=\bar{\rho}_{1}$.
(9.19h) Suppose the representation $T: G \rightarrow S_{n}$ is doubly transitive. Show $G$ does not contain a subgroup $H$ of degree $m<n$ intransitive in $T$. Hint: Count appearances of $\mathbf{1}_{G}$ in $V_{T} \otimes V_{T_{H}}$ using d).
Denote the finite field of $q=p^{r}$ for $p$ a prime by $\mathbb{F}_{q}$. Let $G=\mathrm{GL}_{k}(F)$ be the $k \times k$ invertible matrices with coefficients in the field $F=\mathbb{F}_{q}$. Write $\mathbb{P}^{k-1}(F)$ for
lines through the origin in $\mathbb{F}_{q}^{k}:\{\alpha \boldsymbol{v} \mid \alpha \in F\}$ for some $\boldsymbol{v} \in \mathbb{F}_{q}^{k} \backslash\{\boldsymbol{0}\}$. Then, $G$ has a permutation action $T_{k, F}$ on $\mathbb{P}^{k-1}(F)$ induced from its action on $\mathbb{F}_{q}^{k}$. Let $\psi: \mathbb{F}_{q}^{k} \rightarrow \mathbb{F}_{q}$ be a nonconstant linear map (linear functional). Denote linear functionals up to multiplication by elements of $\mathbb{F}_{q} \backslash\{0\}$ by $\hat{\mathbb{P}}^{k-1}(F)$, with a permutation action $\hat{T}_{k, F}$ : For $\psi \in \mathbb{P}^{k-1}(F)$ and $A \in \mathrm{GL}_{k}(F), \psi^{A}(\boldsymbol{v}) \stackrel{\text { def }}{=} \psi\left((\boldsymbol{v}) A^{-1}\right)$ for $\boldsymbol{v} \in \mathbb{F}_{q}^{k}$.
(9.20a) Show $T_{k, F}$ is doubly transitive of degree $n(q)=\frac{\left(q^{k}-1\right)}{(q-1)}$ for $k>1$.
(9.20b) Show $\hat{T}_{k, F}$ also has degree $n(q)$ and is doubly transitive, though $T_{k, F}$ and $\hat{T}_{k, F}$ are not permutation equivalent. Hint: Show the stabilizer in $G$ of a hyperplane in $\mathbb{P}^{k-1}(F)$ fixes no point.
(9.20c) Show $\mathrm{t}\left(\hat{T}_{k, F}(g)\right)=\mathrm{t}\left(T_{k, F}(g)\right)$, so $\hat{T}_{k, F}$ and $T_{k, F}$ are equivalent as representations [9.19c]. Hint: $\hat{T}_{k, F}(g)$ is induced from the transpose of $g$, and a matrix and its transpose are conjugate.
(9.20d) As in $[9.18 \mathrm{~d}]$, identify $\mathbb{F}_{q}^{k}$ with $\mathbb{F}_{q^{k}}$ as vector spaces over $\mathbb{F}_{q}$ to find $\alpha \in \mathbb{F}_{q^{k}}$ producing $A \in \mathrm{GL}_{k}(F)$ with $T_{k, F}(A)$ and $\hat{T}_{k, F}(A)$ both $n(q)$-cycles.
(9.20e) Assume: $T_{1}, T_{2}$ are inequivalent degree $n$ doubly transitive representations of a group $G$; they are equivalent as group representations; and $T_{1}(g)=$ $T_{2}(g)=(12 \ldots n)$ for some $g \in G$. Let $D$ be the orbit of 1 under $G\left(T_{1}, 1\right)$ in the representation $T_{2}$. Use double transitivity to show $D$ is a difference set: $\left\{d_{i}-d_{j} \mid d_{i} \neq d_{j} \in D\right\}$ contains each nonzero integer $\bmod n$ with the same multiplicity $t[$ Fri73]. Further, $t \cdot(n-1)=|D| \cdot(|D|-1)$. Example: For $k=3, q=2, n=7$ in b), $D=\{1,2,4\}$ and $t=1$.

### 9.8. Easy Galois covers.

(9.21a) Suppose $X$ and $Y_{i}$ are differentiable manifolds, and that $f_{i}: Y_{i} \rightarrow X$ are covering maps, $i=1,2$. Assume $\psi: Y_{1} \rightarrow Y_{2}$ is any continuous map with $f_{2} \circ \psi=f_{1}$. Show $\psi$ is a map of differentiable manifolds. Also: $\psi$ is analytic if $X$ is a complex manifold.
(9.21b) Let $f: Y \rightarrow X$ be a finite cover of degree $n$. Use that $X$ is connected to show $\left|f^{-1}(x)\right|$ is $n$ for each $x \in X$.
(9.21c) Consider $X_{1}=\left\{x+\sqrt{-1} y \in S^{1} \mid y>0\right\}$ and $X_{2}=S^{1}$. Show, for $n>0$, the map of Ex. 6.16 restricted to $X_{1}$ is not a covering map.
(9.21d) Follow the notation of Ex. 6.18 and of [9.7]. Let $L$ and $L_{i}$, with $L_{i} \subseteq L$, $i=1,2$, be lattices. Show that if $f_{i}: X_{i}=\mathbb{C} / L_{i} \rightarrow \mathbb{C} / L$ by $z \bmod L_{i} \mapsto$ $z \bmod L$, then the covers $\left(X_{i}, f_{i}\right)$ are equivalent if and only if $L_{1}=L_{2}$.
(9.21e) Let $X_{i}=X, i=1, \ldots, n$, and let $Y$ be the disjoint union of the $X_{i}$ 's. What is the automorphism group of the cover $Y \rightarrow X$ obtained by mapping each point of $Y$ to its corresponding point in $X$ ?
(9.21f) Let $f: Y \rightarrow X$ be a cover and consider a subgroup $G$ of $\operatorname{Aut}(Y, f)$ of order equal to $\operatorname{deg}(f)$. Assume that, for some point $x_{0} \in X, G$ acts transitively on the set $f^{-1}\left(x_{0}\right)$. Show $f$ restricted to any connected component of $Y$ gives a Galois cover of $X$.
(9.21g) Let $X=Y=\mathbb{C} \backslash\{0\}$. Show $f: Y \rightarrow X$ by $z \mapsto z^{n}$ is a Galois cover. Hint: Consider $\psi_{k}: z \mapsto e^{2 \pi \sqrt{-1} k} z, 0 \leq k \leq n-1$.
(9.21h) Let $X_{i}, i=1,2$, be as in [9.7c] with $L_{1} \subset L_{2}$. Show $f: X_{1} \rightarrow X_{2}$ in Ex. 6.18 is a Galois cover. Hint: Consider $\psi_{k_{1}, k_{2}}: z \bmod L_{1} \mapsto z+$ $z\left(k_{1}, k_{2}\right) \bmod L_{1}$.
(9.21i) Consider a) with $f \in \mathbb{C}[y]$ and $f(y)=y^{n}+c_{n-2} y^{n-2}+\cdots+c_{1} y$. Assume the greatest common divisor of the set $\left\{n\right.$ and $i$ with $\left.c_{i} \neq 0\right\}$ is 1 . Show Aut $(Y, p r)$ is trivial. Hint: Apply Liouville's Theorem [Ahl79, p. 122] to see elements of $\operatorname{Aut}(Y, f)$ have the form $y \mapsto a y+b$ for some $a, b \in \mathbb{C}$.
9.9. Imprimitive and Frattini covers. This discussion on imprimitivity continues in Chap. 4 [11.13]
(9.22a) Let $\pi_{1}\left(X, x_{0}\right)$ be the fundamental group of a connected differentiable manifold $X$. Let $H \sigma_{1}, \ldots, H \sigma_{n}$ be the distinct cosets of a subgroup $H \leq \pi_{1}\left(X, x_{0}\right)$ of index $n$ corresponding to the cover $(Y, f)$ (with fiber $\left\{y_{1}, \ldots, y_{n}\right\}$ over $\left.x_{0}\right)$. Consider the points of $U_{Y, f, n}(\S 8.3 .2)$ over $x_{0}$ that connect by a path to $\left(y_{1}, \ldots, y_{n}\right)$. Show these correspond to distinct $n$ tuples of cosets: $\left\{\left(H \sigma_{1} \sigma, H \sigma_{2} \sigma, \ldots, H \sigma_{n} \sigma\right) \mid \sigma \in \pi_{1}\left(X, x_{0}\right)\right\}$. Why is this the same as $|G|$ ? Conclude $\operatorname{deg}\left(\hat{f}_{i}\right)=|G(Y, f)|$ (as prior to Thm. 8.9).
(9.22b) Show components of $Y \times_{X} Y$ of degree 1 over $Y$ correspond to elements of $\operatorname{Aut}(Y, f)$ (Lem. 8.8). If $f: Y \rightarrow X$ has automorphisms, and $f$ is not a cyclic Galois cover of prime degree, show $G(Y, f)$ is imprimitive. How does [9.21i] give explicit imprimitive covers with no automorphisms?
(9.22c) Show $(Y, f)$ decomposes if and only if $Y \times_{X} Y \rightarrow X$ properly factors through a fiber product of form $Y^{\prime} \times_{X} Y^{\prime}$. If so, show $Y^{\prime} \times_{X} Y^{\prime} \backslash \Delta$ is a nontrivial component of $Y \times_{X} Y$.
Let $K \subset \hat{L} \subset \hat{M}$ be a chain of fields with $\hat{M} / K$ (resp. $\hat{L} / K$ ) Galois with group $G^{*}($ resp. $G)$. This is a Frattini chain if the only subfield $K \leq T \leq \hat{M}$ with $T \cap \hat{L}=K$, is $T=K$. Denote restriction of elements of $G^{*}$ to $\hat{L}$ by rest : $G^{*} \rightarrow G$.
(9.23a) Suppose $T=\hat{M}^{H}$ is the fixed field of a subgroup $H$ of $G^{*}$. Show $T \cap \hat{L}=K$ is equivalent to rest : $H \rightarrow G$. Hint: Use that $T \cap \hat{L}=K$ allows extending any automophism of $\hat{L}$ to $T \cdot \hat{L}$ to be the identity on $T$.
(9.23b) Show a) is equivalent to this group statement: If $H \leq G^{*}$ and $\operatorname{rest}(H)=$ $G$, then $H=G^{*}$ (the map rest : $G^{*} \rightarrow G$ is a Frattini cover). Hint: $\operatorname{rest}(H)=G$ is equivalent to $\hat{M}^{H} \cap \hat{L}=K$.
(9.23c) Suppose $\hat{X} \rightarrow \hat{Y} \rightarrow Z$ is a sequence of covers with $\psi_{X}: \hat{X} \rightarrow Z$ Galois with group $G^{*}$ and $\psi_{Y}: \hat{Y} \rightarrow Z$ Galois with group $G$. Let $\psi: G^{*} \rightarrow G$ be the natural map and assume $\psi$ is a Frattini cover. Show the equivalence with this. For any sequence $\hat{X} \rightarrow W \rightarrow Z$ of covers with $W \neq Z$, there is a proper cover of $Z$ that $W \rightarrow Z$ and $\hat{Y} \rightarrow Z$ factor through.
9.10. Laplacian. The Laplace operator $\nabla^{2}=\frac{\partial}{\partial x} \frac{\partial}{\partial x}+\frac{\partial}{\partial y} \frac{\partial}{\partial y}$ on $\mathbb{R}_{x, y}^{2}$ acts on $C^{\infty}\left(\mathbb{R}^{2}\right)$. It generalize to a Riemann surface $X$ (see Chap. $4 \S 11.11$ for $\wedge$ product). Locally in $z=x+i y$, write a differential 1-form (not necessarily holomorphic) on an open set $U \subset \mathbb{C}$ as $\omega=p(x, y) d x+q(x, y) d y$. Consider $* \omega=-q d x+p d y$. Write $w=u+i v$ for the real and imaginary components of the variable for $\mathbb{C}_{w}$.
(9.24a) With $z=f(w)$, suppose $f: V \subset \mathbb{C}_{w} \rightarrow U \subset \mathbb{C}_{z}$ is analytic, one-one and onto from $V$ to $U$. Write $w=u(x, y)+i v(x, y)$ as the local inverse of $f$. Express $\omega$ as $\left.\Omega_{( } u, v\right)=p(x(u, v), y(u, v)) d x(u, v)+q(x(u, v) y(u, v)) d y(u, v)$. Show $\left.* \Omega_{( } u, v\right)=-Q(u, v) d u+P(u, v) d v$ equals $* \omega$ expressed in $u$ and $v$. Hint: Apply the Cauchy-Riemann equations: $\frac{\partial x}{\partial u}=\frac{\partial y}{\partial v}$ and $\frac{\partial x}{\partial v}=-\frac{\partial y}{\partial u}$.
(9.24b) Conclude from a): On any Riemann surface $X, *$ defines a linear map on differentiable 1-forms.
(9.24c) Show these further properties of $*$ : Its square is multiplication by -1 , $\omega_{\wedge} * \omega_{=}\left(p^{2}+q^{2}\right) d x \wedge d y$, and $* \omega=i \omega$ if $\omega$ is holomorphic. Conclude: $\omega$ is holomorphic if and only if $d \omega=0$ and $* \omega=i \omega$.
(9.24d) Consider $* d=-\frac{\partial}{\partial y} d x+\frac{\partial}{\partial x} d y$ acting on differentiable functions. So, for $f$ differentiable on $X, * d f=* d(f)$ is well-defined, and it extends to 1-forms: $p(x, y) d x+q(x, y) d y \mapsto * d(p(x, y)) \wedge d x+* d(q(x, y)) \wedge d y$. Show $* d(\omega)$ is $-d * \omega$. Define $\nabla^{2}(f)$ by $d * d f=\nabla^{2}(f) d x \wedge d y$. Argue why this defines a (complex) Laplacian $\nabla_{X}^{2}$ on a 1-dimensional complex manifold.
(9.24e) Suppose differentiable $f$ on $X$ has a corresponding $\lambda \in \mathbb{C}$ with $\nabla_{X}^{2}(f)=$ $\lambda f d x \wedge d y$ (everywhere locally). Call $\lambda$ an eigenvalue of $\nabla_{X}^{2}$. If $f_{i}: Y_{i} \rightarrow X$ are inequivalent covers of $X$, with equivalent locally flat bundles (Defn8.6), over $X, i=1,2$, show their Laplacians have the same eigenvalues.

## CHAPTER 4

## RIEMANN'S EXISTENCE THEOREM

This chapter introduces the foundation of the book: The construction of all compact Riemann surfaces through Riemann's classification of the branched covers of the sphere (Thm. 2.6). Still, one cover at a time, won't give us much useful information. We need to know the nature of families of related covers. The Existence Theorem serves well, though it takes additional ideas to find a useful naming scheme for the families. This chapter's nontraditional treatment of modular curves motivates many general ideas in Chap. 5.

## 1. Presentations of fundamental groups of Riemann surfaces

Our command of Riemann's Existence Theorem requires combinatorial ability to list finite quotients of the fundamental group of $U_{\boldsymbol{z}}$. Thm. 1.8 tells us $\pi_{1}\left(U_{\boldsymbol{z}}\right)$ is a free group on $r-1$ generators (with $r=|\boldsymbol{z}|$ ) and more. It is the basis for describing families of covers (Chap. 5) of $\mathbb{P}_{z}^{1}$. Our main computational tools for this are Hurwitz monodromy actions. These are on explicit sets running from types of Nielsen classes (§3.2) to special fundamental group generators of Riemann surfaces defined by Nielsen classes (§9.2).
1.1. Presentations and free products. Most fundamental groups appear as quotients of free groups. Further, we define the kernel of that quotient by listing specific relation elements in the kernel. We recognize the smallest normal subgroup containing these relations as the kernel. A presentation, however, doesn't list all relations from this normal subgroup condition. Presenting groups as quotients of free groups this way is convenient for forming their quotients. To see whether a group $G$ is a quotient of some fundamental group, we need only check if specific generators of $G$ satisfy a tiny list of relations. This suits how we form compact Riemann surfaces from unramified covers of $U_{\boldsymbol{z}}$. Still, this often leaves a tough problem. How to check if an expression from the free group is in that kernel.

For $S$ a set, we first define the group $F(S)$ that $S$ freely generates. The following construction is of a free group with relations. Generalizing this to groups generated freely by subgroups is a categorical rather than quotient construction.

For $s \in S$ and $n \in \mathbb{Z}$, use the symbol $s^{n}$ to denote the pair $(s, n)$. If $t \in S$ and $m \in \mathbb{Z}$ then $s^{n}=t^{m}$ if and only if $s=t, n=m$.

Elements of $F(S)$ are (finite) sequences $\boldsymbol{s}^{\boldsymbol{n}}=\left(s_{1}^{n_{1}}, \cdots, s_{k}^{n_{k}}\right)$ satisfying
$k \in \mathbb{N} ; s_{1}, \cdots, s_{k} \in S ; n_{1}, \cdots, n_{k} \in \mathbb{Z} \backslash\{0\} ;$ and $s_{i} \neq s_{i+1}, i=1, \cdots, k-1$.
Regard the sequence $\emptyset$ with no elements as an element of $F(S)$. Denote $\left(t_{1}^{m_{1}}, \cdots, t_{\ell}^{m_{\ell}}\right) \in F(S)$ by $\boldsymbol{t}^{m}$. Define the product of $\boldsymbol{s}^{\boldsymbol{n}}$ and $\boldsymbol{t}^{m}$ by cancellation to be the elimination of any consecutive terms of the form $t t^{-1}$. Formally, Find
the smallest integer $u$ with this property: $t_{u}^{-m_{u}} \neq s_{k-u+1}^{n_{k-u+1}}$; but $t_{i}^{-m_{i}}=s_{k-i+1}^{n_{k-i+1}}$, $i=1, \cdots, u-1$. Then

$$
\begin{array}{ll}
\boldsymbol{s}^{\boldsymbol{n}} \boldsymbol{t}^{\boldsymbol{m}}= & \left(s_{1}^{n_{1}}, \cdots, s_{k-u}^{n_{k-u}}, \alpha, t_{u+1}^{m_{u+1}}, \cdots, t_{\ell}^{m_{\ell}}\right) \\
\text { where } & \alpha= \begin{cases}\left(s_{k-u+1}^{n_{k-u+1}}, t_{u}^{m_{u}}\right) & \text { for } t_{u} \neq s_{k-u+1} \\
t_{u k-u+1}^{n_{k-u}+m_{u}} & \text { for } t_{u}=s_{k-u+1} .\end{cases} \tag{1.1}
\end{array}
$$

With this multiplication $F(S)$ is a group with $\emptyset$ the identity. For example, an induction on the length of the sequence of the middle term, in a product of 3 terms, suffices to establish the associative law. The inverse of $\boldsymbol{s}^{\boldsymbol{n}}$ is $\left(s_{k}^{-n_{k}}, \cdots, s_{1}^{-n_{1}}\right)$.

For a group $G$ and a subset $S$ of $G$, denote by $\langle S\rangle$ the subgroup of $G$ that $S$ generates. The elements $\boldsymbol{s}^{\boldsymbol{n}} \in F(S)$ for which $s_{1}^{n_{1}} \cdots s_{k}^{n_{k}}$ is the identity in $G$ form a subset $\bar{R}(S)$ called the relations satisfied by $S$. It is a normal subgroup of $F(S)$.

Definition 1.1. Let $S$ be a set of generators of a group $G$. A sequence $\left\{r_{1}, r_{2}, \ldots\right\}$ of $F(S)$ is a presentation of $G$ if $\bar{R}(S)$ is the smallest normal subgroup of $F(S)$ containing $\left\{r_{1}, r_{2}, \ldots\right\}$. We say $\left\{r_{1}, r_{2}, \ldots\right\}$ generates $\bar{R}(S)$. A presentation is finite if both $S$ and $\left\{r_{1}, r_{2}, \ldots\right\}$ are finite sets.

It is standard to denote $\left(s_{1}^{n_{1}}, \ldots s_{k}^{n_{k}}\right)=\boldsymbol{s}^{\boldsymbol{n}} \in F(S)$ by $s_{1}^{n_{1}} \cdots s_{k}^{n_{k}}$ when this symbol could not be confused with the product in another group.

Example 1.2. Let $G=\mathbb{Z}^{2}$, the additive group of integer pairs. Let $s_{1}=(1,0)$ and $s_{2}=(0,1)$. Take for $S$ the set $\left\{s_{1}, s_{2}\right\}$. Then $\left\{s_{1} s_{2} s_{1}^{-1} s_{2}^{-1}\right\}$ is a presentation of $G$. Indeed, $\bar{R}(S)=[F(S), F(S)]$, the commutator subgroup of $F(S)$. [11.7c]

Example 1.3. Take for $S$ the set $\left\{s_{1}, s_{2}, \ldots, s_{r}\right\}$. From now on we denote $F(S)$ by $F_{r}$. There is a natural map from $F_{r}$ to $F_{r-1}$ that maps $s_{i}$ to itself, $i=1, \ldots, r-1$, and $s_{r}$ to $s_{r-1}^{-1} s_{r-2}^{-1} \cdots s_{1}^{-1}$. A nonidentity element of $\bar{R}(S)$ becomes 1 when you make the above substitution for $s_{r}$. Therefore such an element involves $s_{r}$, and $\left\{s_{1} \cdots s_{r}\right\}$ gives a presentation of $F_{r-1}$.

The following treatment on free products of groups, from [Wae48], appears also in [Ma67, p. 97-100]. Let $G_{1}, \ldots, G_{t}$ be groups. We define their free product $G$ by its properties. There are homomorphisms $\alpha_{i}: G_{i} \rightarrow G, i=1, \ldots, t$, satisfying this condition: For any group $H$ and homomorphisms $\beta_{i}: G_{i} \rightarrow H, i=1, \cdots, t$, there exists a unique homomorphism $\beta: G \rightarrow H$ with

$$
\begin{equation*}
\beta \circ \alpha_{i}=\beta_{i}, \quad i=1, \ldots, t \tag{1.2}
\end{equation*}
$$

Modern terminology might suggest the term free sum or pushout; it generalizes for arbitrary groups the direct sum of abelian groups[11.10a]. We now show a free product exists. From the definition it is unique up to isomorphism.

Define $T(\boldsymbol{G})=T\left(G_{1}, \cdots, G_{t}\right)$ as those (finite) sequences $\left(x_{1}, \cdots, x_{n}\right)$ where each $x_{k}$ is a nonidentity element of one of the groups $G_{i}$, and where consecutive terms of the sequence are in different groups. Each $g \in G_{i}$ acts faithfully on the right of $T(\boldsymbol{G})$ as a permutation $\alpha_{i}(g)$ given by the following formula. For $g \in G_{i}$ and $\left(x_{1}, \ldots, x_{n}\right) \in T(\boldsymbol{G}), \alpha_{i}(g)$ maps $\left(x_{1}, \ldots, x_{n}\right)$ to this element:
(1.3a) $\left(x_{1}, \ldots, x_{n} g\right)$ if $x_{n} \in G_{i}$ and $x_{n} g \neq 1_{G_{i}}$;
(1.3b) $\left(x_{1}, \ldots, x_{n-1}\right)$ if $x_{n} \in G_{i}$ and $x_{n}=g^{-1}$;
(1.3c) $\left(x_{1}, \ldots, x_{n}\right)$ if $x_{n} \notin G_{i}$ and $g=1_{G_{i}}$;
(1.3d) $\left(x_{1}, \ldots, x_{n}, g\right)$ if $x_{n} \notin G_{i}$ and $g \notin 1_{G_{i}}$; and
(1.3e) $(g)$ if $\left(x_{1}, \ldots, x_{n}\right)=\emptyset$.

Let $\operatorname{Per}(T(\boldsymbol{G}))$ be the group of (right action) permutations of $T(\boldsymbol{G})$. Then $G$ is the subgroup of $\operatorname{Per}(T(\boldsymbol{G}))$ that the images of the $G_{i}$ under the homomorphisms $\alpha_{i}, i=1, \ldots, t$, generate.

Lemma 1.4. The group $G$ just defined is a free product of $G_{1}, \ldots, G_{t}$.
Proof. Express a given nonidentity element $\gamma$ of $G$ (in reduced form) as

$$
\alpha_{i_{1}}\left(g_{i_{1}}\right) \cdots \alpha_{i_{n}}\left(g_{i_{n}}\right)
$$

where $g_{i_{k}}$ is a nonidentity element of $G_{i_{k}}$ and $i_{k} \neq i_{k+1}, i=1, \ldots, n-1$. This expression is unique. Apply $\gamma$ to $\emptyset$ (as in (1.3e)) to get $\left(g_{i_{1}}, \ldots, g_{i_{n}}\right)$.

Suppose $\beta_{i}: G_{i} \rightarrow H, i=1, \ldots, t$, is any collection of homomorphisms. Define $\beta: G \rightarrow H$ as follows: $\beta\left(\alpha_{i_{1}}\left(g_{i_{1}}\right) \cdots \alpha_{i_{n}}\left(g_{i_{n}}\right)\right)$ is equal to $\beta_{i_{1}}\left(g_{i_{1}}\right) \cdots \beta_{i_{n}}\left(g_{i_{n}}\right)$. Induction on the lengths of the reduced forms of two elements of $G$ shows that $\beta$ is a homomorphism. Clearly $\beta$ is the unique homomorphism satisfying (1.2).
1.2. Fundamental groups of unions of spaces. Let $X$ be a connected union of finitely many differentiable manifolds. Suppose $U$ and $V$ are open subsets of $X$ with $U \cup V=X$, and $U, V$ and $U \cap V$ nonempty and connected. For topological spaces $Y$ and $Z$ with $Y$ a subspace of $Z$ and $y_{0} \in Y$, denote the induced homomorphism $\pi_{1}\left(Y, y_{0}\right) \rightarrow \pi_{1}\left(Z, y_{0}\right)$ by $i(Y, Z)_{*}$.

Theorem 1.5 (Seifert-van Kampen). Let $x_{0} \in U \cap V$. For $H$ a group, let $\beta(U): \pi_{1}\left(U, x_{0}\right) \rightarrow H$ and $\beta(V): \pi_{1}\left(V, x_{0}\right) \rightarrow H$ be two homomorphisms for which

$$
\begin{equation*}
\beta(U) \circ i(U \cap V, U)_{*}=\beta(V) \circ i(U \cap V, V)_{*} . \tag{1.4}
\end{equation*}
$$

Then, there is a unique homomorphism $\beta(X): \pi_{1}\left(X, x_{0}\right) \rightarrow H$ with

$$
\begin{equation*}
\beta(U)=\beta(X) \circ i(U, X)_{*} \text { and } \beta(V)=\beta(X) \circ i(V, X)_{*} . \tag{1.5}
\end{equation*}
$$

In using Thm. 1.5, don't forget $U \cap V$ must be connected. Neglecting this would lead to concluding the torus has trivial fundamental group (Fig. 1).

REMARK 1.6. Commutativity of this diagram characterizes $\pi_{1}\left(X, x_{0}\right)$ :


Figure 1. Two cylinders try to share the fundamental group of a torus, but they connect poorly.

1.3. Proof of Seifert-van Kampen, Thm. 1.5. This is a special case of [Ma67, p.114-22]. We give the proof in four subsections.
1.3.1. $\pi_{1}\left(U, x_{0}\right)$ and $\pi_{1}\left(V, x_{0}\right)$ generate $\pi_{1}\left(X, x_{0}\right)$. Let $\gamma:[a, b] \rightarrow X$ represent an element of $\pi_{1}\left(X, x_{0}\right)$. Find $t_{0}=a<t_{1}<\cdots<t_{n}=b$ so either $U$ or $V$ entirely contains the image of $\gamma_{\left[t_{i}, t_{i+1}\right]} i=1, \ldots, n$. Let $U_{i}$ be either $U$ or $V$, so $U_{i}$ contains the image of $\gamma_{\left[t_{i}, t_{i+1}\right]}$. As $\gamma\left(t_{i}\right)$ lies in both $U_{i-1}$ and $U_{i}$ there is a path $\gamma_{i}$ in $U_{i-1} \cap U_{i}$ joining $x_{0}$ to $\gamma\left(t_{i}\right), i=1, \ldots, n-1$. Then each of the following closed paths is in $U_{i}$ for the corresponding value of $i$ :

$$
\begin{gathered}
\gamma_{0}^{\prime}=\gamma_{\left[t_{0}, t_{1}\right]} \gamma_{1}^{-1}, i=0, \quad \gamma_{i}^{\prime}=\gamma_{i} \gamma_{\left[t_{i}, t_{i+1}\right]} \gamma_{i+1}^{-1}, i=1, \ldots, n-2 \\
\text { and } \gamma_{n-1}^{\prime}=\gamma_{n-1} \gamma_{\left[t_{n-1}, t_{n}\right]}
\end{gathered}
$$

The product $\gamma_{0}^{\prime} \cdots \gamma_{n-1}^{\prime}$ is equivalent to $\gamma$. Write $[\gamma]$ as

$$
\begin{equation*}
i\left(U_{0}, X\right)_{*}\left(\left[\gamma_{0}^{\prime}\right]\right) i\left(U_{1}, X\right)_{*}\left(\left[\gamma_{1}^{\prime}\right]\right) \cdots i\left(U_{n-1}, X\right)_{*}\left(\left[\gamma_{n-1}^{\prime}\right]\right) \tag{1.6}
\end{equation*}
$$

a product of paths, each from $\pi_{1}\left(U, x_{0}\right)$ or $\pi_{1}\left(V, x_{0}\right)$.
1.3.2. Condition for existence of $\beta$. It is natural to define $\beta([\gamma])$ from (1.6):

$$
\begin{equation*}
\beta\left(U_{0}\right)\left(\left[\gamma_{0}^{\prime}\right]\right) \beta\left(U_{1}\right)\left(\left[\gamma_{1}^{\prime}\right]\right) \cdots \beta\left(U_{n-1}\right)\left(\left[\gamma_{n-1}^{\prime}\right]\right) \tag{1.7}
\end{equation*}
$$

We show, if (1.6) is the identity, then so is (1.7); $\beta$ is well-defined.
Let $F:[a, b] \times[0,1] \rightarrow X$ be a homotopy between $\gamma$ and the constant path:

$$
F(t, s)=\gamma_{s}(t), \quad \gamma_{0}(t)=\gamma(t), \text { and } \gamma_{1}(t)=x_{0}
$$

Refine the subdivision $t_{0}=a<t_{1}<\cdots<t_{n}=b$ to find $s_{0}=0<\cdots<s_{m}=1$ so $U_{i, j}$, one of $U$ or $V$, contains the image under $F$ of each rectangle

$$
R_{i, j}=\left\{(t, s) \mid s_{j} \leq s \leq s_{j+1}, t_{i} \leq t \leq t_{i+1}\right\}
$$

Let $V_{i, j}$ be the intersection of $U_{i-1, j}, U_{i-1, j-1}$ and $U_{i, j}$. This refinement doesn't change the value of (1.7). Choose a path $\epsilon_{i, j}:[a, b] \rightarrow V_{i, j}$ with initial point $x_{0}$ and end point $\gamma_{s_{j}}\left(t_{i}\right)=F\left(t_{i}, s_{j}\right)$. When $F\left(t_{i}, s_{j}\right)=x_{0}$, choose $\epsilon_{i, j}$ to be the constant path, and choose $\epsilon_{i, 0}$ to be $\gamma_{i}$ (as in $\S 1.3 .1$ ), $i=1, \ldots, n-1$.

Figure 2. Keeping book along the paths of a grid

1.3.3. Grid following paths. Denote the path $t \in\left[t_{i}, t_{i+1}\right] \mapsto F\left(t, s_{j}\right)$ (resp., $\left.s \in\left[s_{j}, s_{j+1}\right] \mapsto F\left(t_{i}, s\right)\right)$ by $F_{\left[t_{i}, t_{i+1}\right] \times s_{j}}\left(\right.$ resp., $F_{\left.\right|_{t_{i} \times\left[s_{j}, s_{j+1}\right]}}$. Let $\gamma_{i, j}$ be the path $\epsilon_{i, j}\left(F_{\left[t_{i}, t_{i+1}\right] \times s_{j}}\right)\left(\epsilon_{i+1, j}\right)^{-1}$. Let $\delta_{i, j}$ be the path $\epsilon_{i, j}\left(F_{\left.\right|_{t_{i} \times\left[s_{j}, s_{j+1}\right]}}\right)\left(\epsilon_{i, j+1}\right)^{-1}$. Define $g_{i, j}$ to be the image under $\beta\left(U_{i, j}\right)$ of the homotopy class of $\gamma_{i, j}$ in $\pi_{1}\left(U_{i, j}, x_{0}\right)$, $i=0, \ldots, n-1 ; j=0, \ldots, m-1$. Note: (1.4) implies $g_{i, j}$ is also the image under $\beta\left(U_{i, j-1}\right)$ of the class of $\gamma_{i, j}$ in $\pi_{1}\left(U_{i, j-1}, x_{0}\right), i=0, \ldots, n-1 ; j=1, \ldots, m$. So, we consistently define $g_{i, m}$ to be $\beta\left(U_{i, m-1}\right)$, the image of $\gamma_{i, m}$ in $\pi_{1}\left(U_{i, m-1}, x_{0}\right)$, $i=0, \ldots, n-1$. Similarly, $\delta_{i, j}$ gives $h_{i, j} \in H, i=0, \ldots, n ; j=0, \ldots, m-1$.

Since the boundary of $R_{i, j}$ (traversed clockwise) is homotopic to a constant path in $R_{i, j}$, its image under $F$ is homotopic to a constant path in $U_{i, j}$. Therefore

$$
\left(F_{\left.\right|_{t_{i} \times\left[s_{j}, s_{j+1}\right]}}\right)\left(F_{\left.\right|_{\left[t_{i}, t_{i+1}\right] \times s_{j+1}}}\right) \text { is homotopic to }\left(F_{\left.\right|_{\left[t_{i}, t_{i+1}\right] \times s_{j}}}\right)\left(F_{\left.\left.\right|_{t_{i+1} \times\left[s_{j}, s_{j+1}\right]}\right]}\right)
$$

in $U_{i, j}$. Conclude:

$$
\begin{equation*}
\gamma_{i, j} \delta_{i+1, j} \text { is homotopic to } \delta_{i, j} \gamma_{i, j+1} \text { in } U_{i, j} . \tag{1.8}
\end{equation*}
$$

Denote the identity in $H$ by $1_{H}$. An application of $\beta\left(U_{i, j}\right)$ gives
(1.9a) $g_{i, j} h_{i+1, j}=h_{i, j} g_{i, j+1}, i=0, \ldots, n-1 ; j=0, \ldots, m-1$.
(1.9b) As a consequence of $F(t, 1)=F(a, s)=F(b, s)=x_{0}$ :

$$
g_{i, m}=1_{H}, i=0, \ldots, n-1 ; h_{0, j}=h_{n, j}=e_{H}, j=0, \ldots, m-1
$$

Finally, (1.7) is the same as

$$
\begin{equation*}
g_{0,0} g_{1,0} \cdots g_{n-1,0} . \tag{1.10}
\end{equation*}
$$

1.3.4. (1.9a) and (1.9b) imply (1.10) is $1_{H}$. From (1.9b), $g_{0,0} \cdots g_{n-1,0} h_{n, 0}$ equals (1.10). From (1.9a), this is $g_{0,0} \cdots h_{n-1,0} g_{n-1,1}$. Repeat using (1.9a) and (1.9b) to see (1.10) is $g_{0,1} g_{1,1} \cdots g_{n-1,1}$. Inductively: (1.10) is $g_{0, j} g_{1, j} \cdots g_{n-1, j}$ for each $j$. With $j=m$, (1.9b) shows this is $1_{H}$.

Since $\pi_{1}\left(X, x_{0}\right)$ is a pushout of the homomorphisms $\pi_{1}\left(U \cap V, x_{0}\right) \rightarrow \pi_{1}\left(U, x_{0}\right)$ and $\pi_{1}\left(U \cap V, x_{0}\right) \rightarrow \pi_{1}\left(V, x_{0}\right)$, this uniquely defines $\pi_{1}\left(X, x_{0}\right)$ [11.10b] and concludes the proof.
1.4. Classical generators on an $r$-punctured sphere. Let $Y$ be a subspace of a space $X$. Then $Y$ is a retract of $X$ if there is a continuous map $f: X \rightarrow Y$ such that $f(y)=y$ for $y \in Y$. The sequence of maps

$$
Y \xrightarrow{i(Y, X)} X \xrightarrow{f} Y
$$

induces the sequence of homomorphisms of groups

$$
\pi_{1}\left(Y, y_{0}\right) \xrightarrow{i(Y, X)_{*}} \pi_{1}\left(X, y_{0}\right) \xrightarrow{f_{*}} \pi_{1}\left(Y, y_{0}\right)
$$

where $f_{*} \circ i(Y, X)_{*}$ is the identity. This splitting of the sequence of groups means $\pi_{1}\left(X, y_{0}\right)$ is the direct product of $\pi_{1}\left(Y, y_{0}\right)$ and the kernel of $f_{*}$.

Definition 1.7. A retract $Y$ of $X$ is a deformation retract of $X$ if there exists a continuous map $F: X \times[0,1] \rightarrow X$ for which $F(x, 0)=x$ and $F(x, 1)=f(x)$ for $x \in X$, and $F(y, s)=y$ for $y \in Y, s \in[0,1]$.

For each $s \in[0,1]$ the map $F$, restricted to $X \times s$, induces a continuous map $\pi_{1}\left(X, y_{0}\right) \rightarrow \pi_{1}\left(X, y_{0}\right)$. (Regard these fundamental groups as topological spaces with the discrete topology.) Such a map is clearly independent of $s$. For $s=0$ this map is the identity, and for $s=1$ the image of this map identifies with $\pi_{1}\left(Y, y_{0}\right)$. So, $f_{*}$ identifies the fundamental groups $\pi_{1}\left(X, x_{0}\right)$ and $\pi_{1}\left(Y, y_{0}\right)$.
1.4.1. Defining classical generators. Chap. $2 \S 1.1$ introduced the $r$-punctured sphere: $\mathbb{P}_{z}^{1} \backslash\{\boldsymbol{z}\} \stackrel{\text { def }}{=} U_{z}, r$ distinct points $z_{1}, \ldots, z_{r}$ removed from $\mathbb{P}_{z}^{1}$.

Figure 3. Example classical generators based at $z_{0}$


Let $z_{0}$ be a point on $U_{z}$. Let $D_{i}$ be a disc with center $z_{i}, i=1, \ldots, r$. Assume these discs are disjoint and each excludes $z_{0}$. Let $b_{i}$ be a point on the boundary of $D_{i}$. Regard this boundary, oriented clockwise, as a path $\bar{\gamma}_{i}$ with initial and end point $b_{i}$. Finally, let $\delta_{i}$ be a simple simplicial (Chap. 2 Def. 2.1) path with initial point $z_{0}$ and end point $b_{i}$. Assume, also, that $\delta_{i}$ meets none of $\bar{\gamma}_{1}, \ldots, \bar{\gamma}_{i-1}, \bar{\gamma}_{i+1}, \ldots, \bar{\gamma}_{r}$, and it meets $\bar{\gamma}_{i}$ only at its endpoint.

With $D_{0}$ a disc with center $z_{0}$ and disjoint from each of the discs $D_{1}, \ldots, D_{r}$, consider the first point of intersection of $\delta_{i}$ and the boundary $\bar{\gamma}_{0}$ of $D_{0}$. Call this point $a_{i}$. Suppose $\delta_{1}, \ldots, \delta_{r}$ satisfy two further conditions:
(1.11a) they are pairwise nonintersecting, excluding their initial point $z_{0}$; and (1.11b) $a_{1}, \ldots, a_{r}$ appear in order clockwise around $\bar{\gamma}_{0}$.

Since the paths are simplicial this last condition is independent of the choice of $D_{0}$, at least for $D_{0}$ sufficiently small.

With these conditions, the ordered collection of closed paths $\delta_{i} \bar{\gamma}_{i} \delta_{i}^{-1}=\gamma_{i}$, $i=1, \ldots, r$, in Fig. 3 are classical generators (for $\boldsymbol{z}$ ) based at $z_{0}$. We say $\gamma_{i}$ is a classical loop around $z_{i}$. In our case this has a precise meaning.
1.4.2. Main Theorem for classical generators of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Chap. 5 deforms classical generators compatible with deformations of the set $\{\boldsymbol{z}\}=\left\{z_{1}, \ldots, z_{r}\right\}$. Such deformations produce very complicated sets of classical generators. Thus the generality of our next result.

THEOREM 1.8. Let $\left(\gamma_{1}, \ldots, \gamma_{r}\right)$ by any collection of classical generators for $\boldsymbol{z}=\left(z_{1}, \ldots, z_{r}\right)$ based at $z_{0}$ on the $r$-punctured sphere $\mathbb{P}_{z}^{1} \backslash\{\boldsymbol{z}\}$. Then the homotopy classes $\left[\gamma_{1}\right]=s_{1}, \ldots,\left[\gamma_{r}\right]=s_{r}$ generate $\pi_{1}\left(\mathbb{P}_{z}^{1} \backslash\{z\}, z_{0}\right)$ with the one relation $s_{1} \cdots s_{r}$ : The Product-One condition. So, $\pi_{1}\left(\mathbb{P}_{z}^{1} \backslash\{\boldsymbol{z}\}, z_{0}\right)$ is isomorphic to $F_{r-1}$ through the presentation $\left\{s_{1} \cdots s_{r}\right\}$ (Ex. 1.3).

If $\left[\gamma_{1}^{\prime}\right]=s_{1}^{\prime}, \ldots,\left[\gamma_{r}^{\prime}\right]=s_{r}^{\prime}$ is another collection of classical generators, then there is $a \pi \in S_{r}$ so that $s_{i}^{\prime}$ is conjugate to $s_{(i) \pi}, i=1, \ldots, r$.
1.5. Proof of classical generators Thm. 1.8. For the statement on the presentation of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$, induct on $r$. For $r=0$, write $\mathbb{P}^{1}$ as the union of

$$
\mathbb{P}^{1} \backslash\{\infty\}=U_{1} \text { and } \mathbb{P}^{1} \backslash\{0\}=U_{2}
$$

as in Chap. 3 Ex. 3.2.1. Apply Thm. 1.5 (just $§ 1.3 .1$ ). For $r \geq 1$ we show $\pi\left(U_{\boldsymbol{z}}, z_{0}\right)$ :
(1.12a) $\gamma_{1} \cdots \gamma_{r}$ is homotopic (on $U_{\boldsymbol{z}}$ ) to the identity.
(1.12b) $\left[\gamma_{1}\right], \ldots,\left[\gamma_{r-1}\right]$ are free generators of the fundamental group.

These suffice to show the statement gives a correct presentation of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ if we show any relation among $s_{1}, \ldots, s_{r}$ is in the group generated by products of conjugates of the product-one condition. Hints: Do an induction starting with a nontrivial relation containing no subproduct conjugate to the product one relation, and having a minimal number of appearances of $s_{r}$. No appearances of $s_{r}$ is impossible from (1.12b); by conjugating shift the any one appearance of $s_{r}$ to the far right. We divide the proof of (1.12) into 4 parts to separate the conceptual proof from a technical preliminary.
1.5.1. Polygonal paths. We show the set of paths $\gamma_{1}, \ldots, \gamma_{r}$ is (simultaneously) homotopic to a set of simple polygonal paths based at $z_{0}$, intersecting only at $z_{0}$; and that $\gamma_{1} \cdots \gamma_{r}$ is homotopic to a simple polygonal path based at $z_{0}$.

Choose $D_{0}$ so $a_{i}$ is the only intersection of $\delta_{i}$ and $\bar{\gamma}_{0}, i=1, \ldots, r$. This is possible because $\delta_{1}, \ldots, \delta_{r}$ are simplicial. For an integer $n>2$, let $\bar{\gamma}_{i}^{*}$ be the regular $n$-gon inscribed in $\bar{\gamma}_{i}$ as a clockwise path from the vertex $b_{i}$. Chap. 2 Lem. 4.3 allows replacing each $\delta_{i}$ by a polygonal path homotopic to $\delta_{i}$ (with its endpoints fixed), so as to assume our classical generators are polygonal paths.

We explain the formation of the shaded region around the polygonal path $\delta_{i}$ in Fig. 4. The points $b_{i}^{\prime}$ and $b_{i}^{\prime \prime}$ are the vertices of $\bar{\gamma}_{i}^{*}$ next to $b_{i}$. Draw the lines through $b_{i}^{\prime}$ and $b_{i}^{\prime \prime}$ parallel to the last segment of $\delta_{i}$, and let $d=d_{n}$ be the maximum of the distances between these lines and the last segment. Now continue drawing the lines at a distance $d$ parallel to each segment of $\delta_{i}$. For large $n$ : the lines parallel to the last segment meet $\bar{\gamma}_{0}$ at points $a_{i}^{\prime}$ and $a_{i}^{\prime \prime}$; the paths $\delta_{i}^{*}$ and $\delta_{i}^{* *}$ traced by these lines on either side of $\delta_{i}$ are simple and have segments corresponding one-one with the segments of $\delta_{i}$. The shaded region (bounded by $\delta_{i}^{*}, \delta_{i}^{* *}$, the two sides of $\bar{\gamma}_{i}^{*}$ next to $b_{i}$, and the line segments $a_{i}$ to $a_{i}^{\prime}$ and $a_{i}$ to $a_{i}^{\prime \prime}$ ) meets none of the corresponding shaded regions around $\delta_{j}$ for $j \neq i$. In addition, the path going from $a_{i}$ to $a_{i}^{\prime}$, then along $\delta_{i}^{*}$, and then from $b_{i}^{\prime}$ to $b_{i}$ is homotopic (with $a_{i}$ and $b_{i}$ fixed) to $\delta_{i}$ through a homotopy of simple polygonal paths that stay within the shaded region and, until the end, do not meet $\delta_{i}$.

Indeed, with a few choices of lines separating the elbows and ends of the shaded region from the intermediate stretches - this may require a larger value of $n$ - we can make the homotopy canonical. To illustrate, consider the elbow of the last two segments of $\delta_{i}$. The lines $\ell^{\prime}$ and $\ell^{\prime \prime}$ (perpendicular, respectively, to the last and second last segments of $\delta_{i}$ ) that meet at $P$ outline this elbow in Fig. 4. In this region the homotopy takes points along the projection from $P$. In general, the homotopy carries points of $\delta_{i}^{*}$ along the perpendicular to the corresponding segment of $\delta_{i}$.

Let $\lambda_{i}^{*}$ (resp., $\lambda_{i}^{* *}$ ) be a path tracing the ray from $z_{0}$ to $a_{i}^{\prime}$ (resp., $z_{0}$ to $a_{i}^{\prime \prime}$ ). Finally, let $\gamma_{i}^{*}$ be the part of $\bar{\gamma}_{i}^{*}$ with initial point $b_{i}^{\prime}$ and end point $b_{i}^{\prime \prime}$. Then,

$$
\gamma_{i}^{\prime}=\lambda_{i}^{*} \delta_{i}^{*} \gamma_{i}^{*}\left(\delta_{i}^{* *}\right)^{-1}\left(\lambda_{i}^{* *}\right)^{-1}, i=1, \ldots, r,
$$

are simple, polygonal, pairwise nonintersecting (except at $z_{0}$ ) paths that are respectively homotopic to $\gamma_{1}, \ldots, \gamma_{r}$ on $U_{\boldsymbol{z}}$.

Let $\bar{a}_{i}$ be the midpoint of the arc from $a_{i}^{\prime \prime}$ to $a_{i+1}^{\prime}, i=1, \ldots, r-1$. Denote the path along the two straight line segments from $a_{i}^{\prime \prime}$ to $\bar{a}_{i}$, and then from $\bar{a}_{i}$ to $a_{i+1}^{\prime}$ by $\epsilon_{i}^{*}$. Then the following simple polygonal path, $\gamma^{\prime}$, is homotopic on $U_{z}$ to $\gamma_{1}^{\prime} \cdots \gamma_{r}^{\prime}$, and thus to $\gamma_{1} \cdots \gamma_{r}$ :

$$
\begin{equation*}
\lambda_{1}^{*} \delta_{1}^{*} \gamma_{1}^{*}\left(\delta_{1}^{* *}\right)^{-1} \epsilon_{1}^{*} \delta_{2}^{*} \gamma_{2}^{*}\left(\delta_{2}^{* *}\right)^{-1} \epsilon_{2}^{*} \cdots \epsilon_{r-1}^{*} \delta_{r}^{*} \gamma_{r}^{*}\left(\delta_{r}^{* *}\right)^{-1}\left(\lambda_{r}^{* *}\right)^{-1} \tag{1.13}
\end{equation*}
$$

This homotopy shows the interior of a polygonal sector of the disc (marked off clockwise) from $a_{1}^{\prime}$ to $a_{r}^{\prime \prime}$, together with the shaded regions of Fig. 4 and the interiors of $\bar{\gamma}_{1}^{*}, \ldots, \bar{\gamma}_{r}^{*}$ to be part of one connected component of $\mathbb{P}^{1} \backslash \gamma^{\prime}$.
1.5.2. Homotopy of $\gamma_{1} \cdots \gamma_{r}$ to 1 . It suffices to show $\gamma^{\prime}$ is homotopic to the identity. The Jordan curve theorem says the complement of the simple closed path $\gamma^{\prime}$ on $\mathbb{P}^{1}$ consists of two components. For a polygonal path, however, this is fairly easy ([He62, p. 146] or [11.3a]). The Schwartz-Christoffel transformation ([He66, p.351-3] or §6.6) gives a one-one continuous map $\varphi^{\prime}$ from the closed upper hemisphere on $\mathbb{P}^{1}$ to $\mathbb{P}^{1}$, analytic on the open hemisphere, that maps the equator onto $\gamma^{\prime}$. With no loss we assume $\varphi^{\prime}$ maps onto the component excluding $z_{r}$. From the last line of $\S 1.5 .1$, none of $z_{1}, \ldots, z_{r}$ are in the image of $\varphi^{\prime}$. Since the closed upper hemisphere is simply connected, so is the image of the one-one map $\varphi^{\prime}$ on $U_{\boldsymbol{z}}$. Thus, $\gamma^{\prime}$ is homotopic to the identity (see $\S 6.6$ ).
1.5.3. Retraction of $U_{\boldsymbol{z}}$ onto $\gamma_{1}^{\prime} \cup \cdots \cup \gamma_{r-1}^{\prime}$. To simplify our discussion, identity a simple path with its collection of image points. Notice this further use of the Jordan curve theorem (for polygonal paths). The path $\lambda_{r-1}^{* *}$ divides the interior $W$ of $\gamma^{\prime}$ into two parts. The collection of points $\left\{z_{1}, \ldots, z_{r-1}\right\}$ is accessible from one side of $\lambda_{r-1}^{* *}$, and $z_{r}$ from the other. So, $\left\{z_{1}, \ldots z_{r-1}\right\}$ and $\left\{z_{r}\right\}$ lie in distinct components of $W \backslash \lambda_{r-1}^{* *}$ [11.3b]. In the above replace $\gamma^{\prime}$ with following path:

$$
\begin{equation*}
\gamma^{\prime \prime}=\lambda_{1}^{*} \delta_{1}^{*} \gamma_{1}^{*}\left(\delta_{1}^{* *}\right)^{-1} \epsilon_{1}^{*} \cdots \epsilon_{r-2}^{*} \delta_{r-1}^{*} \gamma_{r-1}^{*}\left(\delta_{r-1}^{* *}\right)^{-1}\left(\lambda_{r-1}^{* *}\right)^{-1} \tag{1.14}
\end{equation*}
$$

$\S 1.5 .2$ shows there is a continuous one-one map $\varphi^{\prime \prime}$ from the upper hemisphere mapping the equator onto the path $\gamma^{\prime \prime}$; and mapping onto the component of $\mathbb{P}^{1} \backslash \gamma^{\prime \prime}$ that includes $z_{r}$, but excludes $\left\{z_{1}, \ldots, z_{r-1}\right\}$.

Figure 4. A polygonal thickening of $\delta_{i}$


The upper hemisphere minus $\left(\varphi^{\prime \prime}\right)^{-1}\left(z_{r}\right)$ clearly retracts to the equator. Therefore the closure of the component of $\mathbb{P}^{1} \backslash \gamma^{\prime \prime}$ containing $z_{r}$, with $z_{r}$ removed, retracts to $\gamma^{\prime \prime}$. Denote the closure of the other component by $X^{\prime \prime}$. Similarly, denote the closure of the component of $\mathbb{P}^{1} \backslash \gamma_{i}^{\prime}$ containing $z_{i}$ by $X_{i}, i=1, \ldots, r-1$. Let $Y_{i}$ be the quadrilateral with vertices $a_{i}^{\prime \prime}, \bar{a}_{i}, a_{i+1}^{\prime}$ and $z_{0}, i=1, \ldots, r-2$. Retract $Y_{i}$ onto the union of the two sides defined by $\left\{a_{i}^{\prime}, z_{0}\right\}$ and $\left\{a_{i}^{\prime \prime}, z_{0}\right\}$. Since

$$
X^{\prime \prime}=X_{1} \cup \cdots \cup X_{r-1} \cup Y_{1} \cup \cdots \cup Y_{r-2}
$$

this retracts $X^{\prime \prime}$ onto $X_{1} \cup \cdots \cup X_{r-1}$. Apply the Schwartz-Christoffel transformation to retract $X_{i} \backslash\left\{z_{i}\right\}$ onto $\gamma_{i}^{\prime}, i=1, \ldots, r-1$. This retracts $U_{z}$ onto $\gamma_{1}^{\prime} \cup \cdots \cup \gamma_{r-1}^{\prime}$.
1.5.4. $\left[\gamma_{1}\right], \ldots,\left[\gamma_{r-1}\right]$ generate $\pi_{1}\left(\mathbb{P}^{1} \backslash\{z\}, z_{0}\right)$ freely. The retraction of $\S 1.5 .3$ reduces this to showing $\left[\gamma_{1}^{\prime}\right], \ldots,\left[\gamma_{r-1}^{\prime}\right]$ generate $\pi_{1}\left(\lambda_{1}^{\prime} \cup \cdots \cup \gamma_{r-1}^{\prime}, z_{0}\right)$ freely.

Let $c_{i}$ be a vertex of $\gamma_{i}^{*}$ different from $b_{i}^{\prime}$ or $b_{i}^{\prime \prime}$ (Fig. 4), $i=1, \ldots, r-1$. Take $U$ to be $\gamma_{1}^{\prime} \cup \cdots \cup \gamma_{r-1}^{\prime} \backslash\left\{c_{r-1}\right\}$ and $V$ to be $\gamma_{1}^{\prime} \cup \cdots \cup \gamma_{r-1}^{\prime} \backslash\left\{c_{1}, \ldots, c_{r-2}\right\}$. Then $\gamma_{1}^{\prime} \cup \cdots \cup \gamma_{r-2}^{\prime}$ is a deformation retract of $U ; \gamma_{r-1}^{\prime}$ is a deformation retract of $V$; and $\left\{z_{0}\right\}$ is a deformation retract of $U \cap V$. From Thm. 1.5, $\pi_{1}\left(\gamma_{1}^{\prime} \cup \cdots \cup \gamma_{r-1}^{\prime}, z_{0}\right)$ is a free product of $\pi_{1}\left(U, z_{0}\right)$ and $\pi_{1}\left(V, z_{0}\right)$.

To complete the proof of the theorem, consider another $r$-tuple of classical generators: $\left[\gamma_{1}^{\prime}\right]=s_{1}^{\prime}, \ldots,\left[\gamma_{r}^{\prime}\right]=s_{r}^{\prime}$. Identify the point around which $s_{i}^{\prime}$ loops as the unique point $z^{\prime} \in z$ for which $s_{i}^{\prime} \mapsto 1$ in the natural map $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right) \rightarrow \pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right)$ where $\boldsymbol{z}^{\prime} \dot{\cup}\left\{z^{\prime}\right\}=\boldsymbol{z}$. So, there is a $\pi \in S_{r}$ for which $s_{i}^{\prime}$ loops around $z_{(i) \pi}$. An easy homotopy of both $\gamma_{(i) \pi}$ and $\gamma_{i}^{\prime}$ has these properties.

- It moves only points on these paths within the outermost of $\bar{\gamma}_{(i) \pi}$ and $\bar{\gamma}_{i}^{\prime}$.
- The homotopies end so the respective bounding path to the discs about $z_{(i) \pi}$ are the same.
At time $t$ in the homotopy of $\gamma_{(i) \pi}$ denote the resulting path by $\gamma_{(i) \pi, t}$. In Fig. 5: $\gamma_{i}^{\prime}$ remains constant in the homotopy; $\bar{\gamma}_{(i) \pi, 1}$ is $\bar{\gamma}_{i}^{\prime}$; and only the end portion of $\delta_{(i) \pi, t}$ varies in the homotopy. With $\gamma_{(i) \pi, 1}$ replacing $\gamma_{(i) \pi}=\gamma_{(i) \pi, 0}$ (and the other $r-1$ paths fixed), the equivalence classes in $\pi\left(U_{\boldsymbol{z}}, z_{0}\right)$ give the same elements $s_{1}, \ldots, s_{r}$. With no loss, as in Fig. 5, assume $\gamma_{(i) \pi}$ and $\gamma_{i}^{\prime}$ are respectively $\delta_{(i) \pi} \bar{\gamma}_{(i) \pi}\left(\delta_{(i) \pi}\right)^{-1}$ and $\delta_{i}^{\prime} \bar{\gamma}_{i}^{\prime}\left(\delta_{i}^{\prime}\right)^{-1}$. The homotopy class of $\delta_{(i) \pi, 1}\left(\delta_{i}^{\prime}\right)^{-1}$ conjugates the former to the latter. That completes the proof of the theorem.

Remark 1.9. Massey notes [Ma67, p. 125]:
To actually apply the Seifert-van Kampen Theorem, it is usually necessary to use the properties of deformation retracts.

## 2. Ramified covers from the Existence Theorem

Return to the notation of $\S 2.1$. Let $\psi: Y \rightarrow X$ be a nonconstant analytic map between two connected compact Riemann surfaces. The first part of the Existence Theorem is a combinatorial formula for contructing such ramified covers $\psi$.
2.1. Nonconstant maps of Riemann surfaces. Let $\psi: Y \rightarrow X$ be a nonconstant analytic map of compact connected Riemann surfaces. For any subset $V$ of $X$ denote $\psi^{-1}(V)$ by $Y_{V}$. If $V$ is a point $x \in X$, simplify $Y_{V}$ to be $Y_{x}$, the fiber over $x$. Recall the definition of unramified cover from Chap. 3 Def. 7.12.

Figure 5. Comparing two loops around $z_{(i) \pi}$

2.1.1. The divisor of ramification. We first attach a multiplicity to a point in a fiber. The outcome is that all fibers of $\psi$ will have the same degree.

Lemma 2.1. The map $\psi$ is open and so is surjective. Two analytic functions $\psi_{i}: Y \rightarrow \mathbb{P}_{z}^{1}, i=1,2$, with exactly the same zeros and poles (with multiplicity) on $X$ differ by multiplication by a constant.

For some integer $n,\left|Y_{x}\right|=n$ for all but finitely many $x \in X$. For $x \in X$, $\left|Y_{x}\right| \leq n$. Let $D(\psi)$ be those $x$ with $\left|Y_{x}\right|<n$. Then $Y_{X \backslash D(\psi)} \rightarrow X \backslash D(\psi)$ is an unramified cover.

Representing restriction of $\psi$ around any point $y_{0}$ by an analytic function in a disk allows assigning a multiplicity $e_{y_{0}}$ to $y_{0}$ in $Y_{\psi\left(y_{0}\right)}$. This gives a degree of the fiber $Y_{x}$ by $\operatorname{deg}\left(Y_{x}\right) \stackrel{\text { def }}{=} \sum_{y \in Y_{x}} e_{y}$ and all fibers of $\psi$ have degree $n$.

If $X=\mathbb{P}_{z}^{1}$, then the divisor $(\psi)$ of the meromorphic function $\psi$ has degree 0 . Any meromorphic function on $Y$ comes from an analytic map where $X=\mathbb{P}_{z}^{1}$.

Proof. If $\psi$ maps open sets to open sets, then the range of $\psi$ is open. Since $X$ is compact, the range of $\psi$ is also closed. As $X$ is connected, that means the range is the only possible nontrivial open and closed set, $X$. The statement that $\psi$ is open is local: We have only to show it maps small open sets to small open sets. [Ahl79, p. 131] (as it is used below) shows $\psi$ is locally an open map. Apply this by considering two analytic functions $\psi_{i}: Y \rightarrow \mathbb{P}_{z}^{1}, i=1,2$, with the same divisor of zeros and poles on $Y$. Then, the ratio $\psi_{1} / \psi_{2}$ has no zeros, and no poles. It gives an analytic map to $\mathbb{P}_{z}^{1}$ missing $\infty$ for example. So, it must be constant.

Let $f$ be a nonconstant analytic function on an open connected subset $U$ on $\mathbf{C}$, and let $z_{0} \in U$. There is a neighborhood $V$ of $z_{0}$ on which $f$ is one-one if and only if $\frac{d f}{d z}\left(z_{0}\right) \neq 0\left[\mathbf{A h l 7 9}\right.$, p. 131]. Suppose $\frac{d f}{d z}\left(z_{0}\right) \neq 0$. Then there is a neighborhood $U_{z_{0}}$ of $z_{0}$ for which $\frac{d f}{d z}$ is not 0 and $f$ restricted to $U_{z_{0}}$ is one-one. Let $\left\{\left(U_{\alpha}^{Y}, \varphi_{\alpha}^{Y}\right)\right\}_{\alpha \in I}$ (resp., $\left\{\left(U_{\beta}^{X}, \varphi_{\beta}^{X}\right)\right\}_{\beta \in J}$ ) be an atlas for the manifold $Y$ (resp. $X$ ).

Consider the set $R$ of $y \in Y$ with

$$
\begin{equation*}
\frac{d}{d z}\left(\varphi_{\beta}^{X} \circ \psi \circ\left(\varphi_{\alpha}^{Y}\right)^{-1}\right)\left(\varphi_{\alpha}^{Y}(y)\right)=0 \tag{2.1}
\end{equation*}
$$

for some $\alpha \in I, \beta \in J$ with $y \in U_{\alpha}^{Y} \cap \psi^{-1}\left(U_{\beta}^{X}\right)$. The condition is independent of the choice of $\alpha$ and $\beta$ (as in Chap. 3 Lem. 5.2). If $R$ is infinite, then $R$ has a limit point $y_{0}$. We show this leads to a contradiction.

There exists $\alpha \in I$ and $\beta \in I$ with $y_{0} \in U_{\alpha}^{Y}$ and $\psi\left(y_{0}\right) \in U_{\beta}^{X}$. The zeros of $\frac{d}{d z}\left(\varphi_{\beta}^{X} \circ \psi \circ\left(\varphi_{\alpha}^{Y}\right)^{-1}\right)$ have limit point $\varphi_{\alpha}^{Y}\left(y_{0}\right)$. So $\varphi_{\beta}^{X} \circ \psi \circ\left(\varphi_{\alpha}^{Y}\right)^{-1}$ is constant in a neighborhood of $\varphi_{\alpha}^{Y}\left(y_{0}\right)$ [Ahl79, p. 127], and $\psi$ is constant in a neighborhood of $y_{0}$. The points of $Y$ with a neighborhood on which $\psi$ is constant is an open set contained in $R$. Any accumulation point of it is therefore an accumulation point of $R$. The above argument shows this set is closed. Since $Y$ is connected, the existence of $y_{0}$ shows $\psi$ is constant on all of $Y$, contrary to assumption. So $R$ is finite.

Each $y \in Y \backslash R$ has a connected neighborhood $U_{y}$ of $y$ to which the restriction of $\psi$ is a one-one function. Let $x \in X \backslash \psi(R)$. For each $y \in R$, let $U_{y}$ be a neighborhood of $y$ with $x \notin \psi\left(U_{y}\right)$. As $\psi$ is one-one on $U_{y}, U_{y}$ contains at most one point of $Y_{x}$. The cover $\left\{U_{y}\right\}_{y \in Y}$ of the compact space $Y$ contains a finite subcover. Therefore $Y_{x}$ is finite. Now consider neighborhoods of points of $Y_{x}$.

Let $V_{x}$ be a connected neighborhood of $x$ contained in $\psi\left(U_{y}\right)$ for each $y \in Y_{x}$. Then the connected components of $Y_{V_{x}}$ are $\left\{U_{y} \cap Y_{V_{x}}\right\}_{y \in Y_{x}}$, and the restriction of $\psi$ to each of these is one-one. From Chap. 3 Def. $7.12, \psi$ restricted to $Y_{X \backslash \psi(R)}$ is a cover, and the fibers have constant cardinality (Chap. 3 [9.21b]).

Now consider a fiber $Y_{x}$ with $x \in D(\psi)$. Expression (2.1) generalizes. Any point $y \in Y_{x}$ gives a well-defined integer $e_{y}$ : The minimal $e \geq 1$ with

$$
\frac{d^{e}}{d z^{e}}\left(\varphi_{\beta}^{X} \circ \psi \circ\left(\varphi_{\alpha}^{Y}\right)^{-1}\right)\left(\varphi_{\alpha}^{Y}(y)\right) \neq 0
$$

This is the ramification index of $\psi$ at $y$ (Chap. 2 Def. 7.6). Suppose $\left|Y_{x}\right|=t$. [Ahl79, p. 131] shows $f=\varphi_{\beta}^{X} \circ \psi \circ\left(\varphi_{\alpha}^{Y}\right)^{-1}$ is $e$ to 1 in a neighborhood of $\varphi_{\alpha}^{Y}(y)$ with $y$ removed. So, in some small punctured neighborhood $V_{x}^{0}=V_{x} \backslash\{x\}$ of $x$, the punctured neighborhoods $U_{1}^{0}, \ldots, U_{t}^{0}$ above $V_{x}^{0}$ have this property: $\psi_{U_{i}^{0}}: U_{i}^{0} \rightarrow V^{0}$ is everywhere $e_{i}$ to 1 . Since the degree of each fiber over $x \in V_{x_{0}}$ is $n$, conclude $\sum_{y \in Y_{x}} e_{y}=n$. This is the formula stated in the lemma.

Now assume $X=\mathbb{P}_{z}^{1}$. So, Chap. $4 \S 5.3 .1$ assigns to $\psi$ a well-defined divisor: $Y_{0}-Y_{\infty}$. Its degree is $\operatorname{deg}\left(Y_{0}\right)-\operatorname{deg}\left(Y_{\infty}\right)=n-n=0$. Finally, let $f$ be any global meromorphic function on $Y$. Then, locally $f$ is a ratio of two holomorphic functions on a disk. At each point of the disk this defines a map to $\mathbb{P}_{z}^{1}$ which is analytic, even at the zeros of the denominator (Chap. $2 \S 4.6$ ). So, $f$ is an analytic map to $\mathbb{P}_{z}^{1}$.

We often refer to a cover $\psi: Y \rightarrow X$ by the pair $(Y, \psi)$. With the hypotheses of Lem. 2.1, call $(Y, \psi)$ a ramified cover of $X$ of degree $n: \operatorname{deg}(\psi)=n$. Then $D(\psi)$ consists of the branch points of $\psi$.

Definition 2.2. Let $\psi: Y \rightarrow X$ be an analytic map of 1-dimensional complex manifolds (not necessarily compact or connected). If $(\psi)^{-1}(K)$ is compact for each compact subset $K$ of $X$ and $\left|(\psi)^{-1}(x)\right|=n$ for all but a discrete subset of points $x \in X$, then $(Y, \psi)$ is a finite ramified cover of degree $n$. Denote the set $\left\{x\left|\left|Y_{x}\right| \neq n\right\}\right.$ by $D(\psi)$.
2.1.2. s-equivalence of covers. Let $\psi^{i}: Y^{i} \rightarrow X, i=1,2$, be two finite ramified covers of $X$. Then $\left(Y^{1}, \psi^{1}\right)$ and $\left(Y^{2}, \psi^{2}\right)$ are $s($ trong)-equivalent (as ramified covers of $X$ ) if there is a one-one and onto continuous map $\psi: Y^{1} \rightarrow Y^{2}$ for which $\psi^{2} \circ \psi=\psi^{1}$. Colloquially: There is an isomorphism that commutes with the
projection maps to the base. In $\S 3.2 .2$ this corresponds to the notion of absolute s-equivalence; there is no extra condition on the s-equivalence of these covers.

Then, $\psi$ is automatically an analytic isomorphism [11.2]. Clearly $D\left(\psi^{1}\right)=$ $D\left(\psi^{2}\right)$. Using the phrase s-equivalence differentiates this from other equivalences of covers that appear later. The compactification process for covers of complex manifolds in higher dimensions does not necessarily produce a manifold, as it does in dimension 1 (Thm. 2.6). Still, the notion of s-equivalence makes sense and we extend its use to many situations.

Let $D$ be a finite subset of the connected 1-dimensional compact complex manifold $X$. Cor. 2.9 classifies s-equivalence classes of finite ramified covers $\psi: Y \rightarrow X$ with $D(\psi) \subseteq D$. Restricting $\psi$ to $Y_{X \backslash D(\psi)}$ gives an unramified cover. Therefore explicitly completing such a classification requires explicitly presenting the fundamental group $\pi_{1}\left(X \backslash D, x_{0}\right)$ for $x_{0} \in X \backslash D$.
2.2. Constructing ramified covers. Now take $X$ to be the Riemann sphere, $\mathbb{P}^{1}=\mathbb{P}_{z}^{1}$. Versions of these results work in the general case [11.11].
2.2.1. Product-One Condition. Label points of $D(\psi)$ as $\{\boldsymbol{z}\}=\left\{z_{1}, \ldots, z_{r}\right\}$. Let $z_{0} \in \mathbb{P}^{1} \backslash D(\psi)=U_{\boldsymbol{z}}$. Let $\left(\gamma_{1}, \ldots, \gamma_{r}\right)=\gamma$ be classical generators for $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. A labeling $\boldsymbol{y}=\left(y_{1}, \ldots, y_{n}\right)$ of the points of $Y$ lying over $z_{0}$ determines a transitive permutation representation $T(\boldsymbol{y})$ of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ of degree $n$. This is as in Chap. 3 Thm. 7.16, except we now have additional information. Denote $T(\boldsymbol{y})\left(\left[\gamma_{i}\right]\right)$ by $g_{i} \in$ $S_{n}, i=1, \ldots, r$, and let $G(\boldsymbol{g})$ be the subgroup of $S_{n}$ the $g_{i}$ s generate.

Lemma 2.3. With the hypotheses above, $g_{1} \cdots g_{r}=1$. Conversely, given elements $g_{i} \in S_{n}, i=1, \ldots, r$ satisfying $g_{1} \cdots g_{r}=1$, there exists a unique homomorphism $\psi_{*}: \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right) \rightarrow S_{n}$ mapping $\gamma_{i}$ to $g_{i}, i=1, \ldots, r$. This canonically produces a(n unramified) cover $\psi: Y^{0} \rightarrow U_{\boldsymbol{z}}$ whose components correspond one-one to the orbits of $G(\boldsymbol{g})$ on $\{1, \ldots, n\}$.

Proof. Thm. 1.8 says $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ is a free group on $\gamma$ modulo the product one relation $\left[\gamma_{1} \cdots \gamma_{r}\right]=1$ in the fundamental group. This implies the quotient relation

$$
\left[\gamma_{1} \cdots \gamma_{r}\right]=\left[\gamma_{1}\right] \cdots\left[\gamma_{r}\right]=g_{1} \cdots g_{r}=1
$$

Conversely, the product-one relation on the $g_{i} \mathrm{~s}$ implies there is a homomorphism having the desired properties. The corresponding permutation representations on the orbits of $G(\boldsymbol{g})$ correspond to connected covers of $U_{\boldsymbol{z}}$.

Definition 2.4. We call the r-tuple $\boldsymbol{g}=\left(g_{1}, \ldots, g_{r}\right)$ in Lem. 2.3 a branch cycle description of the cover $\psi: Y \rightarrow \mathbb{P}^{1}$ with respect to $\gamma$.

The group $G(\boldsymbol{g})$ is the monodromy group of the ramified cover $(Y, \psi)$ (with respect to $\boldsymbol{y})$. Refer to an r-tuple $\boldsymbol{g}^{\prime} \in S_{n}^{r}$ for which there is $\beta$ in $S_{n}$ with $\beta^{-1} g_{i} \beta=$ $g_{i}^{\prime}, i=1, \ldots, r$, as absolutely equivalent to $\boldsymbol{g}$.
2.2.2. Compactification of unramified Riemann surface covers. The first part of Riemann's Existence Theorem, the part so technically useful, is that there is a unique compactification of any finite cover $\psi^{0}: Y^{0} \rightarrow U_{\boldsymbol{z}}$ to a cover $\psi: Y \rightarrow \mathbb{P}_{z}^{1}$ of compact Riemann surfaces. We now show this.

Let $D_{i}$ be the disc about $z_{i}$ in Fig. $3, i=1, \ldots, r$. Consider $Y_{D_{i}} \rightarrow D_{i}$, the restriction of $\psi$ over $D_{i}$. Then, $\bar{\gamma}_{i}$ generates $\pi_{1}\left(D_{i} \backslash\left\{z_{i}\right\}, b_{i}\right)$ which maps naturally to $\pi_{1}\left(\mathbb{P}^{1} \backslash D(\psi), b_{i}\right)$. Identify $\pi_{1}\left(\mathbb{P}^{1} \backslash D(\psi), b_{i}\right)$ with $\pi_{1}\left(\mathbb{P}^{1} \backslash D(\psi), z_{0}\right)$ using the path $\delta_{i}$ (of Fig. 3). Apply unique pathlifting along $\delta_{i}$ (Chap. 3 Lem. 7.13). So, the labeling on $\boldsymbol{y}$ uniquely labels points of $Y$ over $b_{i}$.

With this, the permutation from $\bar{\gamma}_{i}$ on the fiber over $b_{i}$ is $g_{i}$. Write $Y_{D_{i} \backslash\left\{z_{i}\right\}}$ as a disjoint union of connected components $\cup_{j=1}^{t_{i}} M_{i, j}$. Up to s-equivalence as a cover of $D_{i} \backslash\left\{z_{i}\right\}$, each $M_{i, j}$ corresponds to an orbit of $\pi_{1}\left(D_{i} \backslash\left\{z_{i}\right\}, b_{i}\right)$ on $\{1,2, \ldots, n\}$. Disjoint cycles in the decomposition of the generator $g_{i}$ determine the orbits (Chap. 2 Prop. 7.4). The degree of $M_{i, j}$ as a cover is the length of its corresponding cycle, $i=1, \ldots, t$. Thus, $g_{i}$ determines the covers $M_{i, 1}, \ldots, M_{i, t_{i}}$ (and their degrees).

Suppose $z_{i}=0$ and $D_{0}$ is a disc about the origin in $\mathbb{C}$. Then, for each integer $e>0$, the s-equivalence class of the connected cover of degree $e$ is represented by

$$
M^{\prime}=\left\{(w, z) \in \mathbb{C} \times \mathbb{C} \mid w^{e}=z\right\}_{D_{0} \backslash\{0\}} \xrightarrow{\text { proj. on } z} D_{0} \backslash\{0\}
$$

For each $z \in D_{0} \backslash\{0\}$, let $D_{z}$ be a disc about $z$ contained in $D_{0} \backslash\{0\}$. The components of $M_{D_{z}}^{\prime}$, with their projections to $D_{0}-\{0\}$, give an atlas on $M^{\prime}$.

Lemma 2.5. The space $M^{\prime} \cup\{(0,0)\}=M$ has a complex manifold structure (extending that of $M^{\prime}$ ) that makes it a ramified cover of $D_{0}$ with exactly one point over 0 . Indeed, $M$ is analytically isomorphic to a disc.

Proof. The mapping $(w, z) \mapsto w$ gives a homeomorphism of $M^{\prime} \cup\{(0,0)\}$ to the subset of $\mathbb{C}$ that lies over $D_{0}$ via the map $w \mapsto w^{e}$. This subset is a disc around the origin, so it is complex analytically isomorphic to $D_{0}$. With this identification of $M$ with $D_{0}$, add it to the atlas to conclude the manifold property. Compactness of the inverse image of a compact subset of $D_{0}$ follows easily (Def. 2.2).
2.2.3. From unramified to ramified covers. Now for Riemann's Existence Theorem: Equivalence classes of ramified covers $\psi: Y \rightarrow X$ with $D(\psi)$ contained in a given set $D^{\prime}$ correspond exactly to classes of permutation representations of $\pi_{1}\left(X \backslash D^{\prime}, z_{0}\right)$ (Chap. $3 \S 7.2 .2$ ). Our next two results give formal restatements.

THEOREM 2.6. Let $\boldsymbol{z}=\left\{z_{1}, \ldots, z_{r}\right\}$ be a collection of $r$ distinct points of $\mathbb{P}_{z}^{1}$. There is a one-one correspondence between connected unramified covers of $U_{\boldsymbol{z}}$ and connected covers of $\mathbb{P}^{1}$ ramified over a subset of $\boldsymbol{z}$.

Proof. From the opening remarks of this subsection we must show that if $\psi^{\prime}: Y^{\prime} \rightarrow \mathbb{P}^{1} \backslash D^{\prime}$ is an unramified cover, then there exists a unique ramified cover $\psi: Y \rightarrow \mathbb{P}^{1}$ such that $Y_{\mathbb{P}^{1} \backslash D^{\prime}}$ is equivalent to $\left(Y^{\prime}, \psi^{\prime}\right)$.

Use the notation prior to Lem. 2.5. For each $i=1, \ldots, r$, it shows how to add just one point $m_{i, j}$ to each component $M_{i, j}, j=1, \ldots, t_{i}$, of $Y_{D_{i} \backslash\left\{z_{i}\right\}}^{\prime}$ to obtain a disjoint union $\cup_{j=1}^{t_{i}} \bar{M}_{i, j}=Y_{i}$ of manifolds with these properties.
(2.2a) There is a ramified covering map $\psi_{i}: Y_{i} \rightarrow D_{i}$.
(2.2b) $\psi_{i}^{-1}\left(D_{i} \backslash\left\{z_{i}\right\}\right)$ is equivalent to $Y_{D_{i} \backslash\left\{z_{i}\right\}}^{\prime}$.
(2.2c) $\bar{M}_{i, j}$ is analytically isomorphic to a disc.

The identification of $\bar{M}_{i, j}$ with a disc in $(2.2 \mathrm{c}), j=1, \ldots, t_{i} ; i=1, \ldots, r$, added to an atlas for $Y^{\prime}$ gives an atlas for $Y=Y^{\prime} \bigcup_{i, j}\left\{m_{i, j}\right\}$. Extend $\psi^{\prime}$ to $\psi: Y \rightarrow \mathbb{P}^{1}$ by mapping $m_{i, j}$ to $z_{i}, j=1, \ldots, t_{i} ; i=1, \ldots, r$. Then $Y_{D_{i}}$ is equivalent to $Y_{i}$, $i=1, \ldots, r$. Now we show $Y$ is a compact manifold.

Since $Y$ has an atlas, it is a manifold if it is Hausdorff. But $\mathbb{P}^{1}$ is Hausdorff. Thus if $y_{1}, y_{2} \in Y$ with $\psi\left(y_{1}\right) \neq \psi\left(y_{2}\right)$, then we get $\psi^{-1}\left(U_{1}\right)$ and $\psi^{-1}\left(U_{2}\right)$, disjoint open sets, respectively, containing $y_{1}$ and $y_{2}$, by taking $U_{1}$ and $U_{2}$ to be disjoint open sets of $\mathbb{P}^{1}$, respectively, containing $\psi\left(y_{1}\right)$ and $\psi\left(y_{2}\right)$. Also, $Y^{\prime}$ is a manifold. Thus we only need consider $y_{1}, y_{2} \in Y$ distinct points with $\psi\left(y_{1}\right)=\psi\left(y_{2}\right)=z_{i}$ for

Figure 6. Virtual neighborhoods awaiting a disc call—see Fig. 7

some $i=1, \ldots, r$. Therefore $y_{1}=m_{i, \ell}$ and $y_{2}=m_{i, k}$ for some $\ell \neq k$ between 1 and $t_{i}$. In particular, $\bar{M}_{i, \ell}$ and $\bar{M}_{i, k}$ are disjoint open sets, respectively, containing $y_{1}$ and $y_{2}$. The Hausdorff property follows.

For $z \in \mathbb{P}^{1}$ let $D_{z}$ be a disc neighborhood of $z$. If $D_{z} \backslash\{z\}$ contains no points of $D^{\prime}$, then each component of $Y_{D_{z}}$ contains a point of $\psi^{-1}(z)$. Thus the open sets $Y_{D_{z}}$ form a neighborhood base for $Y_{z}$. Let $\mathcal{U}=\left\{U_{\alpha}\right\}_{\alpha \in I}$ be an open cover of $Y$. The fiber $Y_{z}$ is contained in a finite union $U_{z}$ of the sets $U_{\alpha}$, so $Y_{D_{z}} \subset U_{z}$ for some choice of $D_{z}$. Since $\mathbb{P}^{1}$ is compact, $\mathbb{P}^{1}=\cup_{i=1}^{t} D_{z_{i}}$ and $Y=\cup_{i=1}^{t} U_{z_{i}}$. Thus $\mathcal{U}$ has a finite subcover, and $Y$ is compact.

The theorem is complete if we show $\psi: Y \rightarrow \mathbb{P}_{z}^{1}$ is unique. Let $\psi_{1}: Y^{1} \rightarrow \mathbb{P}^{1}$ be a ramified cover with $Y_{\mathbb{P}^{1} \backslash D^{\prime}}^{1}$ equivalent to $\left(Y^{\prime}, \psi^{\prime}\right)$, and therefore to $Y_{\mathbb{P}^{1} \backslash D^{\prime}}$. Thus there is an analytic isomorphism $\varphi: Y_{\mathbb{P}^{1} \backslash D^{\prime}}^{1} \rightarrow Y_{\mathbb{P}^{1} \backslash D^{\prime}}$. If $\varphi$ extends to $Y^{1}$ then Lem. 2.1 shows $Y^{1}$ and $Y$ are analytically isomorphic. Let $y \in\left(\psi_{1}\right)^{-1}\left(z_{i}\right)$ for some $i=1, \ldots, r$. Let $U$ be a connected open neighborhood of $y$ contained in some coordinate neighborhood with $\psi_{1}(U)$ contained in $D_{i}$. Since $U$ is connected, $\varphi$ maps $U \backslash \psi_{1}^{-1}\left(z_{i}\right)$ into $\bar{M}_{i, j}$ for some $j$. Riemann's removable singularities theorem extends $\varphi$ to $y$ uniquely [11.2b].

Conspicuous among covers of $U_{\boldsymbol{z}}$ that now compactify to a manifold are those from an algebraic function $f(z) \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$, labeled as $X_{f}^{0}$ in Chap. 3 Prop. 3.12.

Definition 2.7. Call the manifold compactification $X_{f}$ of $X_{f}^{0}$ (or more sloppily, of $f$ ) from Thm. 2.6 its rs-compactification. This theorem says any manifold compactification of $X_{f}^{0}$ will have a unique complex extending structure. Still, this notation differentiates $X_{f}$ from a different compactification that might not have a manifold structure (as in Chap. $3 \S 4.2$ ).
2.3. Combinatorial RET, algebraic and abelian covers. Let $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ be an analytic map of compact Riemann surfaces with $z$ the branch points of $\varphi$. For $z^{\prime} \in \mathbb{P}_{z}^{1}$ consider $D_{\varphi, z^{\prime}}=D_{z^{\prime}}$, the divisor of $\varphi-z^{\prime}$ on $X$ (Chap. 3 §5.3.1). For $z^{\prime}=\infty$, interpret $D_{\varphi, \infty}$, the polar divisor, as counting (with multiplicity) points on $X$ over $\infty$.
2.3.1. An atlas from a compact cover. For $z^{\prime} \notin \boldsymbol{z} \cup\{\infty\}$, and $D_{z^{\prime}}=\sum_{j=1}^{n} x_{j}$, choose a neighborhood $U_{z^{\prime}}$ of $z^{\prime}$ and $U_{x_{i}}$ so $\varphi$ is invertible on $U_{x_{i}}$. As in Chap. 3 Prop. 3.12, use $\left(U_{x_{i}}, \varphi\right)$ as a coordinate chart around $x_{i}$ as $\varphi \stackrel{\text { def }}{=} w_{x_{i}}: U_{x_{i}} \rightarrow U_{z^{\prime}} \subset$
$U_{\boldsymbol{z}} \backslash\{\infty\} \subset \mathbb{C}_{z}$. We extend this around ramified points (when $z^{\prime} \in \boldsymbol{z}$ ) and the possibility $z^{\prime}=\infty$, where $e_{i}$ is the ramification index of $x_{i}$ in the fiber $X_{z^{\prime}}(\S 2.1)$, and $\left\{x_{1}, \ldots, x_{t}\right\}=X_{z^{\prime}}$. First, assume $z^{\prime} \neq \infty$. As in applying (2.2), for some coordinate neighborhood $\left(U_{x_{i}}, \psi_{x_{i}}\right)$ of $x_{i}$, (with $\left.\varphi_{x_{i}}\left(x_{i}\right)=0\right)$ there is a branch of $e_{i}$ th root of the function $\varphi \circ \psi_{x_{i}}^{-1}: \mathbb{C} \rightarrow \mathbb{C}$. So, there is a well defined function —designate it $w_{x_{i}}=\varphi^{1 / e_{i}}$ —one-one in a neighborhood of $x_{i}$ with $w_{x_{i}}$ giving a coordinate chart about $x_{i}$. (Again select $U_{z^{\prime}}$ to avoid $\infty$ and any other points of $z$.) If $z^{\prime}=\infty$, use $w_{x_{i}}=1 / \varphi^{1 / e_{i}}$ instead.

Definition 2.8. Call $\left\{\left(U_{x}, w_{x}\right)\right\}_{x \in X}$ the atlas for $X$ from $\varphi$. In basing a construction on this atlas, we must guarantee the result does not depend on the choice of branches of $e_{i}$ th roots; we have made no canonical choice for these here.
2.3.2. Algebraic and abelian covers of $\mathbb{P}_{z}^{1}$. Combined with Nielsen classes (§3.2), Cor. 2.9 is the statement we use most often in describing types of covers.

Corollary 2.9. Let $\boldsymbol{z}=\left\{z_{1}, \ldots, z_{r}\right\}$ as in Thm. 2.6. Each set of classical generators $\left(\gamma_{1}, \ldots, \gamma_{r}\right)=\gamma$ for $\boldsymbol{z}$ based at $z_{0} \in \mathbb{P}_{z}^{1} \backslash D^{\prime}$ determines a one-one correspondence between equivalence classes of the following sets:
(2.3a) connected covers $\psi: Y \rightarrow \mathbb{P}_{z}^{1}$ with $D(\psi) \subseteq D^{\prime}$ and deg $(\psi)=n$; and
(2.3b) r-tuples $\boldsymbol{g}=\left(g_{1}, \ldots, g_{r}\right) \in S_{n}^{r}$ with $G(\boldsymbol{g})$ transitive, and $g_{1} \cdots g_{r}=1$.

For a representative $\psi: Y \rightarrow \mathbb{P}_{z}^{1}$ of (2.3a) and a labeling $\boldsymbol{y}=\left(y_{1}, \ldots, y_{n}\right)$ of $\psi^{-1}\left(z_{0}\right)$, the correspondence produces a unique representative $\boldsymbol{g}$ of the class of (2.3b); and the disjoint cycles of $g_{i}$ identify with points of $\psi^{-1}\left(z_{i}\right), i=1, \ldots, r$.

Proof. From Thm. 2.6, elements of (2.3a) correspond to equivalence classes of unramified covers of $U_{\boldsymbol{z}}$. Excluding the last line, the corollary follows from the discussion prior to Def. 2.4. Given a representative $\psi: Y \rightarrow \mathbb{P}^{1}$ of a class of (2.3a), and a labeling $\boldsymbol{y}$ of $\psi^{-1}\left(z_{0}\right)$, the discussion following Def. 2.4 shows connected components of $Y_{D_{i} \backslash\left\{z_{i}\right\}}$ correspond uniquely to the disjoint cycles of $g_{i}, i=1, \ldots r$, in the correspondence of (2.3). Then, (2.2) gives a correspondence of the points of $\psi^{-1}\left(z_{i}\right)$ with the components of $Y_{D_{i} \backslash\left\{z_{i}\right\}}, i=1, \ldots, r$. This gives the corollary.

Chap. 2 Thm. 8.8 describes all abelian algebraic functions of $z$. We compare that precise description with Cor. 2.9. An abelian cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is one that is the compactification of a cover of $U_{\boldsymbol{z}}$ with abelian monodromy group. The same terminology is useful in describing nilpotent or solvable covers of any Riemann surface (or of any manifold if there is an appropriate construction of the compactification).

Definition 2.10 (Algebraic cover of $\mathbb{P}_{z}^{1}$ ). Call a cover of compact Riemann surfaces $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ algebraic if there is a second analytic map $f: X \rightarrow \mathbb{P}_{w}^{1}$ so that for some $z^{\prime} \in U_{z}, f$ separates points in the fiber $X_{z^{\prime}}: f\left(x^{\prime}\right) \neq f\left(x^{\prime \prime}\right)$ for distinct points $x^{\prime}, x^{\prime \prime} \in X_{z^{\prime}}$. Then, $\mathbb{C}(z, f) \stackrel{\text { def }}{=} \mathbb{C}(X)$ is the field of functions of $X$.
If $\varphi^{\prime}: X^{\prime} \rightarrow \mathbb{P}_{z}^{1}$ is s-equivalent to $\varphi$ (§2.1.2), then $\varphi$ is algebraic if and only if $\varphi^{\prime}$ is.
Proposition 2.11 (Algebraists' RET). Every algebraic cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is s-equivalent to to an rs-compactification (Def. 2.7) $X_{f}$ of an algebraic function (Chap. 3 Prop. 3.12). The lattice of fields between $\mathbb{C}(z, f(z))$ and $\mathbb{C}(z)$ is dual to the lattice of covers $\varphi_{Y}: Y \rightarrow \mathbb{P}_{z}^{1}$ through which $\varphi$ factors.

Suppose $\hat{L}$ is the Galois closure of $\mathbb{C}(z, f(z))=L$ over $\mathbb{C}(z)$, with branch points $\boldsymbol{z}=\left\{z_{1}, \ldots, z_{r}\right\}$. Then a set of classical generators, $\gamma_{1}, \ldots, \gamma_{r}$, for $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ defines a set of embeddings $\psi_{i}: \hat{L} \rightarrow \mathcal{P}_{z_{i}, e_{i}}$ with $e_{i}$ the ramification index of $f$ over $z_{i}$.

Consider the restrictions $g_{z_{i}, \psi_{i}} \in G_{f}$ of the canonical generator of $G\left(\mathcal{P}_{z_{i}, e_{i}} / \mathcal{L}_{z_{i}}\right)$ to $\hat{L}, i=1, \ldots, r$ (Chap. 2 Lem. 7.9). Then $\left(g_{z_{1}, \psi_{1}}, \ldots, g_{z_{r}, \psi_{r}}\right)=\boldsymbol{g}$ generates $G \hat{L} / \mathbb{C}(z)$ and satifies the product-one condition.

Any abelian cover of $\mathbb{P}_{z}^{1}$ is the rs-compactification of an explicit algebraic function $f$ from branches of log. So, each abelian cover of $\mathbb{P}_{z}^{1}$ is an algebraic cover.

Proof. Consider the function $f: X \rightarrow \mathbb{P}_{w}^{1}$. As in Rem. 2.14, this produces an analytic structure on $X$. The phrase, $f$ is a meromorphic function on $X$, means $f$ and $\varphi$ give same analytic structure on $X$.

As usual form $U_{\boldsymbol{z}} \subset \mathbb{P}_{z}^{1}$. Let $V$ be an open set in $\varphi^{-1}\left(U_{\boldsymbol{z}}\right)$ on which $\varphi$ maps one-one to a disk $D$ in $\mathbb{P}_{z}^{1}$. Use the notation $\varphi_{V}^{-1}$ for the inverse map. Then, $f_{D}=f \circ \varphi_{V}^{-1}: D \rightarrow \mathbb{P}_{w}^{1}$ is meromorphic. Now we show the analytic continuations of $f_{D}$ along paths in $U_{\boldsymbol{z}}$ satisfy Chap. 2 (1.1), properties. Chap. 2 Prop. 6.4 guarantee $f_{D}$ is an algebraic function of $z$.

Let $z_{0} \in D, x_{1} \in V$ over $z_{0}$ and let $\gamma^{*}:[a, b] \rightarrow X$ be the unique lift to $\varphi^{-1}\left(U_{\boldsymbol{z}}\right)$ starting at $x_{1}$. Consider analytic continuation of $f_{D}$ along $\gamma \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right): f_{D, \gamma}(t)$ is the function defined by $f \circ \gamma^{*}(t)$. This gives an analytic continuation according to Chap. 2 Def. 4.1. Further, analytic continuation gives only finitely many possible functions, the functions defined by $f$ at the finite set of points above $z_{0}$. Similarly, test what happens as we approach the points $z^{\prime} \in \boldsymbol{z}$. We evaluate $f$ points with a limit on $X$. So the values remain bounded around a point of the range.

Now consider $\hat{L}$, the Galois closure of $\mathbb{C}(z, f(z))=L$ over $\mathbb{C}(z)$, with branch points $\boldsymbol{z}=\left\{z_{1}, \ldots, z_{r}\right\}$. First note that each element among the $r$ classical generators $\gamma_{1}, \ldots, \gamma_{r}$ defines an embedding of $\hat{L}$ in the corresponding $\mathcal{P}_{z_{i}, e_{i}}$. Write $\gamma_{i}=\delta_{i} \bar{\gamma}_{i} \delta_{i}^{-1}$ (as in Fig. 3), then $\delta_{i}$ gives an analytic contuation of $f$ and all its conjugates to a disk neighborhood about $z_{i}$. Then, Chap. 2 Lem. 7.9, gives the desired embedding $\psi_{i}$. Generation and product-one conditions follow because they hold for the classical generators.

Finally consider when the cover $\varphi$ has abelian monodromy. Chap. 2 (8.8) gives a branch cycle description with values in an abelian group. This was the hypothesis for producing an abelian function through branches of log. So, Chap. 2 Thm. 8.8 says branches of log display this unique cover (up to s-equivalence).
2.3.3. New covers from subfields of algebraic function fields. Def. 3.5 explains normal fiber products of compact Riemann surface covers. This shows Prop. 2.11 directly gives many covers with nonabelian monodromy group as algebraic. §6 explains why any of the competing definitions of algebraic apply to algebraic covers.

Many uses of Riemann's Existence Theorem (including for the Inverse Galois Problem) require knowing covers are algebraic and more. Given $f$ attesting to an algebraic cover, there is a unique $h(w)=w^{n}+\sum_{j=0}^{n-1} u_{j}(z) w^{j} \in \mathbb{C}(z)[w]$ (monic and irreducible in $w$ ) relating $f$ to $z$ in Prop. 2.11. We eventually need the minimal field (of definition) containing all coefficients (in $z$ ) of those $u_{j} \mathrm{~s}, j=0, \ldots, n-1$. We usually want the minimal such field as $f$ varies. It is inefficient (sometimes hopeless), outside special cases, to compute $f$ or $h$ to find this out. There should be a good reason for doing such calculations. For example, theory might show there is a good choice of $f$, yet give reasons for looking more deeply at the algebraic relation. Our examples will show when theory is not yet sufficient to tell everything we want. Then, computing $h$ may give us new clues about theory.

The best situation is that among these fields, as $f$ varies, there is one that is minimal in that any nontrivial isomorphism of that field gives a new cover. This is the situation when the field of moduli is a field of definition ( $\S 6.2$ ); $\S 8.6$ gives the first step in investigating this possibility and variant questions. This is a question that tacitly assumes there is such an $f$ : One reason why Thm. 2.13 is so important.

Given that $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is algebraic, we know that nonconsant elements of $\mathbb{C}(X)$ give all ways that $X$ covers the Riemann sphere.

Corollary 2.12. Each field $L$ properly between $\mathbb{C}$ and $\mathbb{C}(X)$ corresponds to a cover $\psi: X \rightarrow Y$ with $Y$ algebraic and the embedding $f \in \mathbb{C}(Y) \mapsto f \circ \psi \in \mathbb{C}(X)$ identifies $\mathbb{C}(Y)$ with $L$. Conversely, a cover $\psi$ corresponds to subfield $L$.

Proof. For $w \in L$ nonconstant, $x \in X \mapsto w(x)$ gives a cover $\varphi_{w}: X \rightarrow \mathbb{P}_{w}^{1}$. Apply Thm. 2.11 to $L$ between $\mathbb{C}(X)$ and $\mathbb{C}(w)$. Prop. 6.3 shows the converse.

Though we do not complete showing all covers of compact Riemann surfaces are algebraic until Chap. $5 \S$ ??, we record that here. Examples in the remainder of this chapter emphasize aspects of applying Riemann's Existence Theorem. Several concentrate on showing the historical attention given to finding functions displaying covers as algebraic.

Theorem 2.13. Each cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ of compact Riemann surfaces is not only algebraic, it is $\mathbb{P}^{1}$-algebraic.

Remark 2.14 (Warning on constructing $f$ in Def. 2.10). Suppose $X$ is any compact Riemann surface and $f: X \rightarrow \mathbb{P}_{w}^{1}$ is any differentiable map with but finitely many points at which $d f$ is 0 . There are many such maps. Thm. 2.6 says $f$ induces a complex structure on $X$. Chances are, however, that complex structure will differ from that we started with. That is why it is difficult to construct an $f$ that demonstrates a cover of $\mathbb{P}_{z}^{1}$ is algebraic.
2.4. Cuts and impossible pictures. Chap. $3 \S 7.2 .3$ discusses problems with traditional renderings of covers. Even the case when the degree $n$ is 2 , as in $\S 7.1$. Assuming $Y$ has a presentation as a sphere with $g$ handles in $\mathbb{R}^{3}$, presenting the map $\psi$ by a picture in $\mathbb{R}^{3}$ can be confusing. Still, something akin to Fig. 7 appears in many books; for example, [Con78, p. 243].

It includes all the usual elements, especially the cuts. We understand from [Ne81] that Gauss suggested cuts to Riemann (see $\S 10.2$ ). We don’t rely on these cuts. Still, it will be valuable to see what they represent and how we can use symbols from them to draw pictures of the covers. The short and general description in §2.4.3 suffices for an alternate description of the manifold. The slower treatment in §2.4.1 establishes that the idea behind cuts is that covers are a locally constant structure.
2.4.1. The simplest possible cuts. The left of Fig. 7 represents a disc snipped and separated along a radius on the nonpositive real axis $\mathbb{R} \leq 0 \stackrel{\text { def }}{=}\{x \leq 0\} \dot{\cup}\{\infty\}$ from $-\infty$ to 0 . Our perspective is taken from looking along the front edge. So the cut side that is on top has label $T$ and the edge along the bottom has label $B$. The mathematical reality, however, is that (unlike the figure) we shouldn't separate the two sides of the cut (on either disk) by lifting one above the other. Rather, we intend just to remove the cut $\mathbb{R} \leq 0$, including the points $(0$ and $\infty)$ at the ends of the cut. To continue the explanation, call the result of this $U_{z, l}$ and the corresponding figure on the right $U_{z, r}$. At each point $z^{\prime}$ of either of these two figures, the ring of functions we call analytic in a disc about $z^{\prime}$ (entirely within $U_{z, l}$ or $U_{z, r}$ ) identifies
with the ring of analytic functions on $\mathbb{P}_{z}^{1}$ about that same disc regarded as on $\mathbb{P}_{z}^{1}$. Now let $S$ be an open strip on $\mathbb{P}_{z}^{1}$ along the cut. Remove from $S$ the negative real axis $\mathbb{R}^{<0} \stackrel{\text { def }}{=}\{x<0\} \dot{\cup}\{\infty\}$ to leave two open substrips on each side of $S$.

We want to consider the copies $S_{T, l}$ and $S_{B, l}$ (resp. $S_{T, r}$ and $S_{B, r}$ ) on $U_{z, l}$ (resp. $U_{z, r}$ ). These appear in Fig. 7. We also need two copies of $S, S_{l}$ and $S_{r}$. Identify the analytic functions on each with those of $S$, exactly as you would expect from $S$ being on $\mathbb{P}_{z}^{1}$.

The complex space $X^{0}$ we construct to cover $U_{0, \infty}$ consists of four pieces: $U_{z, l}, U_{z, r}$ and $S_{l}$ and $S_{r}$. The map from all four pieces to $\mathbb{P}_{z}^{1}$ is exactly from the identification of each with a subset of $\mathbb{P}_{z}^{1}$. The only item left unsaid is the identification of points of $U_{z, l}, U_{z, r}$ and $S_{l}$ and $S_{r}$ between each of these four pieces. We don't want to identify these with points of $\mathbb{P}_{z}^{1}$ for this purpose. That would just give (two copies of) the manifold $U_{0, \infty}$ back. Here is the right final identification.
(2.4a) Points of $S_{T, l}$ identify with points of the corresponding strip on $S_{l}$, but $S_{B, l}$ identifies with the corresponding strip on $S_{r}$.
(2.4b) Identify points of $S_{T, r}$ with the points of the corresponding strip on $S_{r}$, but identify $S_{B, r}$ with the corresponding strip on $S_{l}$.
(2.4c) Make no further identifications.

Consider the path $\bar{\gamma}:[0,1] \rightarrow U_{0, \infty}$ by $t \in[0,1] \mapsto e^{-2 \pi i t}$ and let $\bar{\gamma}_{1}$ be its lift to $X$ starting at $1 \in U_{z, l}$. We follow it to what happens as it gets to the different pieces. As $t$ increases to $\frac{1}{2}$, within the points of $S_{B, l}$, switch to points we identify with them on $S_{r}$. Now cross $\mathbb{R}^{<0}$ on $S_{r}$, to get to points that identify with points on $S_{T, r}$. Finally, complete $\bar{\gamma}_{1}$ around to 1 on $U_{z, r}$. Total result: Traversing the unique lift of $\bar{\gamma}$ (a clockwise path) starting at $1 \in U_{z, l}$ ends at $1 \in U_{z, R}$. Exercise: Do the same with the lift of $\bar{\gamma}$ starting at $1 \in U_{z, R}$ to see it ends at $1 \in U_{z, l}$.

Figure 7. Connecting two copies of $\mathbb{P}_{z}^{1}$ to double cover $\mathbb{P}_{z}^{1}$

2.4.2. Any cycle, and any one cut. Instead of using the labels l (eft) and r (ight) in $\S 2.4 .1$, we might have used $x_{1}$ and $x_{2}$ by associating everything on the left with the point $1 \in U_{z, l}$ renamed to $x_{1}$, and everything on the right with the point $1 \in U_{z, r}$ renamed to $x_{2}$. There was no loss in our picture of changing left to right to regard it as going from bottom to top. We generalize this in two steps. First: consider any
integer $n$, two distinct points $z_{1}^{\prime}$ and $z_{2}^{\prime}$ on $\mathbb{P}_{z}^{1}$ and any simple, piecewise simplicial path $\delta_{z_{1}^{\prime}, z_{2}^{\prime}}:[a, b] \rightarrow \mathbb{P}_{z}^{1}$ starting at $z_{1}^{\prime}$ and ending at $z_{2}^{\prime}$.

Let $U_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j}$ be a copy of $\mathbb{P}_{z}^{1}$ minus the range of $\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j=1, \ldots, n$. Think of these copies listed from left to right, according to their numbering $\left(U_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j}\right.$ on the far right). Let $S_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}}$ be a thin open strip playing the same role toward $\delta_{z_{1}^{\prime}, z_{2}^{\prime}}$ as did $S$ toward $\mathbb{R}^{\leq 0}$ (starting from $\infty$ and going toward 0 along the negative real axis). Then, consider copies of $S_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}}, S_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j}, j=1, \ldots, n$, with their corresponding substrips $S_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j, T}$ (on the left of $\delta_{z_{1}^{\prime}, z_{2}^{\prime}}$ ) and $S_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j, B}$ (on the right of $\delta_{z_{1}^{\prime}, z_{2}^{\prime}}$ ). Let $\bar{\gamma}:[0,1] \rightarrow U_{z_{1}^{\prime}, z_{2}^{\prime}}$ by $t \mapsto z_{2}^{\prime}+r_{0} e^{-2 \pi i t+t_{0}}$ so $\bar{\gamma}$ meets $\delta_{z_{1}^{\prime}, z_{2}^{\prime}}$ precisely once and not when $t=0\left(\bar{\gamma}(0)=z_{0} \notin \delta_{z_{1}^{\prime}, z_{2}^{\prime}}\right)$. Label the point on $U_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j}$ above $z_{0}$ as $x_{j}$, $j=1, \ldots, n$.

Lemma 2.15. Let $g \in S_{n}$. Then, there is a canonical equivalence on the union of the open sets $U_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j}$ and $S_{\delta_{z_{1}^{\prime}, z_{2}^{\prime}}, j}, j=1, \ldots, n$, so the following holds.
(2.5a) The resulting equivalence classes form a complex manifold $X^{0}$ giving an unramified cover $\varphi^{0}: X^{0} \rightarrow U_{z_{1}^{\prime}, z_{2}^{\prime}}$.
(2.5b) The unique lift of $\bar{\gamma}$ starting at $x_{j}$ ends at $(j) g, j=1, \ldots, n$.

So, $\left(g, \delta_{z_{1}^{\prime}, z_{2}^{\prime}}\right)$ produces a canonical ramified cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ of compact Riemann surfaces, ramified only over $z_{1}^{\prime}$ and $z_{2}^{\prime}$, the completion of $\varphi^{0}$ from Thm. 2.6.

Proof. We do the case $g=(12 \cdots n)$ and leave the adjustments for the general case as an exercise. Most of even this case imitates the case $n=2$. To simplify notation, drop extra reference to the path $\delta_{z_{1}^{\prime}, z_{2}^{\prime}}$. The map of the union of the $S_{j} \mathrm{~s}$ and $U_{j}$ s to $U_{z_{1}^{\prime}, z_{2}^{\prime}}$ is by identifying the points (and the local complex functions) on these sets with those on $\mathbb{P}_{z}^{1}$. The only item left unsaid is the identification of points of the $S_{j} \mathrm{~s}$ with corresponding points of the $S_{T, j} \mathrm{~s}$ and $S_{B, j} \mathrm{~s}$.
(2.6a) Identify points of $S_{T, j}$ with the points of the corresponding strip on $S_{j}$, but identify $S_{B, j}$ with the corresponding strip on $S_{j+1}$.
(2.6b) Make no further identifications, except for $j=n$, we take $j+1$ to be 1 . Do the rest of the lemma as [11.17a] requests.

Remark 2.16 (Locally constant structures). Chap. 3 Ex. 8.18 uses that degree $n$ unramified covers are equivalent to locally constant bundles on $\{1, \ldots, n\}$. Such structures, over $U_{\boldsymbol{z}}$ for example, are equivalent to looking at elements of $\operatorname{Hom}\left(\pi_{1}\left(U_{z}\right), S_{n}\right)$. In Lem. 2.15, the sets $S_{j}$ and $U_{j}$ are simply connected. So above these sets, the cover consists of $n$ connected copies of each of these sets. Using cuts is equivalent to explicitly laying out this locally constant structure.
2.4.3. Any r rooted cuts. Look again at the case of one cut. We may turn this into two rooted cuts by selecting any point $z_{0}$ along the cut. For simplicity assume for now it is not one of the endpoints of the cut. Now follow the procedure below.

Fig. 3 has the notation for the construction of classical generators. We show how the paths $\delta_{1}, \ldots, \delta_{r}$ correspond one-one with $r$ rooted cuts by the following simple device. Extend $\delta_{i}$ to a path $\bar{\delta}_{i}$ by adding the ray from $b_{i}$ to $z_{i}, i=1, \ldots, r$. Thm. 1.8 says the rooted bush formed by the union of $\bar{\delta}_{1}, \ldots, \bar{\delta}_{r}$ has simply connected complement, an essential property for having a collection of cuts on $U_{\boldsymbol{z}}$.

Any sequence of covers $Y \xrightarrow{\varphi_{u}} \mathbb{P}_{u}^{1} \xrightarrow{\varphi_{u, z}} \mathbb{P}_{z}^{1}$ gives three covers for which we would like an algorithm to precisely relate branch cycle descriptions. Especially, we have
applications that allow computing classical generators for $\mathbb{P}_{u}^{1}$ automatically from classical generators for $\mathbb{P}_{z}^{1}$. This would allow constructing a branch cycle description for $\varphi_{u}$ immediately from such a description for $\varphi_{z}=\varphi_{u, z} \circ \varphi_{u}$ (§6.3).

Suppose $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is a ramified cover with branch cycles $\boldsymbol{g}$ from the classical generators that give $\bar{\delta}_{1}, \ldots, \bar{\delta}_{r}$. This assumes a labeling $x_{1}, \ldots, x_{r}$ of $X_{z_{0}}$. Then, we form a cover $\varphi_{c}: X_{c} \rightarrow \mathbb{P}_{z}^{1}$ from the cut construction canonically identifies with $\varphi$. Here are the ingredients.
(2.7a) Label copies of $\mathbb{P}_{z}^{1}$ as $\mathbb{P}_{z, j}^{1}=\mathbb{P}_{j}^{1}, j=1, \ldots, r$. On each remove the points labeled $z_{0}, z_{1}, \ldots, z_{r}$ and call the result $\mathbb{P}_{j}$.
(2.7b) Use each element $g_{i}$ and the cut $\bar{\delta}_{i}$ from $z_{0}$ to $z_{i}$ to attach the $\mathbb{P}_{j}$ s along the lift of the $i$ th cut. When done, compactify what we get.
We use the word triangle on a Riemann surface to mean a(clockwise oriented) boundary of a topological disk with the boundary divided into three oriented simplicial segments (edges) by three points called its vertices (Fig. 8). Call the triangle with its interior (which makes sense as the region to the right of the boundary) a (simplicial) simplex. The proof of Prop. 2.18 consists of describing these attachments and forming from them a natural triangulation of the result.

Definition 2.17. A triangulation of a compact Riemann surface $X$ is a cover of it by simplices satisfying these conditions. The simplex sides meet other simplices in their sides (in opposite orientation), and no two simplices have overlapping interiors. Let $n_{v}$ (resp. $n_{e}, n_{s}$ ) be the number of vertices (resp. edges, simplices). The Euler characteristic of the triangulation is the alternating sum $n_{v}-n_{e}+n_{s}$.

Form a pre-manifold $\mathbb{P}_{j}^{ \pm}$(not Hausdorff) from $\mathbb{P}_{j}$ by replacing each point $z$ along any one of the $\bar{\delta}_{i} \mathrm{~s}$ (minus its endpoints) by two points: $z^{+}$and $z^{-}$. We put a new topology on a quotient relation on the union of $\left\{\mathbb{P}_{j}^{ \pm}\right\}_{j=1}^{n}$. This uses an expected neighborhood basis at all points, except the pairs labeled $z^{+}$and $z^{-}$: Disks not meeting any of the cuts $\bar{\delta}_{1}, \ldots, \bar{\delta}_{r}$. The right neighborhood basis around $z^{+}$and $z^{-}$on a cut use the following. Write $D_{j, z}$, a disk around $z$ (on $\bar{\delta}_{i}$ ), as a union of $D_{j, z}^{+}$and $D_{j, z}^{-}: D_{j, z}^{+}\left(\right.$resp. $\left.D_{j, z}^{-}\right)$is all points on and to the left (resp. right) of $\bar{\delta}_{i}$.

Proposition 2.18. Compactifying $X_{c}^{0}$ gives a cover $\varphi_{x}: X_{c} \rightarrow \mathbb{P}_{z}^{1}$ unramified over $z_{0}$. A map giving the equivalence to $\varphi$ takes $x_{j}$ to the point identified with $z_{0}$ on $\mathbb{P}_{j}^{1}$. Let $t_{i}$ be the number of disjoint cycles in $g_{i}, i=1, \ldots, r$. The cuts from $\bar{\delta}_{1}, \ldots, \bar{\delta}_{r}$ produce a triangulation of $X_{c}$ with $n_{s}=2 n r$ simplices, $3 n r$ sides and $2 n+\sum_{i=1}^{r} t_{i}$ vertices. So the Euler characteristic of $X_{c}$ is $2 n+\sum_{i=1}^{r} t_{i}-n r$.

Precise cut pasting. Form $X_{c}^{0}$ as an equivalence relation on $\cup_{j=1}^{n} \mathbb{P}_{j}^{ \pm}$. Suppose $g_{i}$ maps $k$ to $l$ and $z$ lies on $\bar{\delta}_{i}$. Then, identify $z^{-} \in \mathbb{P}_{k}^{ \pm}$with $z^{+} \in \mathbb{P}_{l}^{ \pm}$. In the resulting set, take a neighborhood of $z^{-}$to be $D_{l, z}^{+} \cup D_{k, z}^{-}$identified along the part of $\bar{\delta}_{i}$ running through $z$.

Interpret the path $\delta_{i} \bar{\gamma}_{i} \delta_{i}^{-1}=\gamma_{i}$ in Fig. 3 as follows.
(2.8a) The lift of $\delta_{i}$ starting at $z_{0}$ on $\mathbb{P}_{k}^{1}$ rides along the right edge of the $g_{i}$-cut on $\mathbb{P}_{k}^{1}$ until it gets to $\bar{\gamma}_{i}$.
(2.8b) The initial point of $\bar{\gamma}_{i}$ is on the ${ }^{-}$-edge of the $g_{i}$-cut on $\mathbb{P}_{k}^{1}$; it ends at the - -edge of the $g_{i}$-cut on $\mathbb{P}_{l}^{1}$.
(2.8c) The lift of $\delta_{i}^{-1}$ starting at $z_{j}$ on $\mathbb{P}_{l}^{1}$ rides along the ${ }^{-}$edge of $g_{i}$-cut on $\mathbb{P}_{l}^{1}$ until it gets to $z_{0}$.

So traversing the lift of $\gamma_{i}$ from $z_{0} \in \mathbb{P}_{k}^{1}$ will end at $z_{0} \in \mathbb{P}_{l}^{1}$. Consider the small clockwise circle about $z_{0}$ denoted $\bar{\gamma}_{0}$ in Fig. 3. Our construction shows that traversing a lift of $\bar{\gamma}_{0}$ has the same effect on the points over $z^{\prime}$ in the range of $\bar{\gamma}_{0}$ as the product $\Pi(\boldsymbol{g})=1$ has on the integers $\{1, \ldots, n\}$. It leaves them fixed. So, a deleted neighborhood of $z_{0}$ has above it $n$ disjoint copies of that neighborhood on $X_{c}^{0}$. According to Lem. 2.5, the compactification does not ramify over $z_{0}$.

Triangulate $\mathbb{P}_{z}^{1}$ using the cuts $\bar{\delta}_{1}, \ldots, \bar{\delta}_{r}$ and the proof of Thm. 1.8. Especially recall $\S 1.5 .2$ showing the outside of the product of the classical generator paths bounds a disk. From this, draw paths $\mu_{i}$ from $z_{i}$ to $z_{i+1}, i=1, \ldots, r-1$, and $\mu_{r}$ from $z_{r}$ to $z_{1}$ with the following properties. The closed path $\mu_{1} \cdot \mu_{2} \cdots \mu_{r}$ bounds a closed (topological) disk $\bar{\Delta}_{\infty}$ that meets the $\bar{\delta}_{i}$ s only at the endpoint $z_{i}$ s. From any point $z_{\infty}$, interior to $\bar{\Delta}_{\infty}$, draw paths $\mu_{1}^{\prime}, \ldots, \mu_{r}^{\prime}$, intersecting only at their beginning point, entirely in $\bar{\Delta}_{\infty}$ from $z_{\infty}$ to the respective $z_{i}$ s.

Triangulate $\mathbb{P}_{j}^{1}$ by listing the three ordered edges of the triangles:

$$
\begin{align*}
\left(\bar{\delta}_{i}, \mu_{i}, \bar{\delta}_{i+1}^{-1}\right), & i=1, \ldots, r-1, & & \left(\bar{\delta}_{r}, \mu_{r}, \bar{\delta}_{1}^{-1}\right), \\
\left(\left(\mu_{i}^{\prime}\right)^{-1}, \mu_{i+1}^{\prime}, \mu_{i}^{-1}\right), & i=1, \ldots, r-1, & & \left.\left(\mu_{r}^{\prime}\right)^{-1}, \mu_{1}^{\prime}, \mu_{r}^{-1}\right) . \tag{2.9}
\end{align*}
$$

Now, triangulate $X_{c}$ using the following simple principle. Each of the $2 r$ triangles in (2.9) bounds a simplex with exactly two endpoints in $\boldsymbol{z}$. Let $S$ be one of these. Remove the two points from $z$ in the boundary; call this $S^{0}$. It is simplyconnected, and $\varphi_{c}: X_{c} \rightarrow \mathbb{P}_{z}^{1}$ is unramified over it. So, $S^{0}$ has $n$ connected components $S_{1}^{0}, \ldots, S_{n}^{0}$ over it. With each take the closure in $X_{c}$ (adding back points of $X_{c}$ over $\boldsymbol{z}$ ). These simplices give the triangulation of $X_{c}$. Just count to get the statement of the proposition.

Figure 8. Cuts for a triangulation of $X_{c}$ when $r=3$


Remark 2.19. The expression for the Euler characteristic in Prop. 2.18 is $2-2 g_{\boldsymbol{g}}=\chi_{X}$, appearing in Prop. 3.10. This shows, all triangulations of a compact Riemann surface $X$ from presenting $\varphi$ using cuts, have the same Euler characteristic. We leave the following observations to the many topology books that treat Euler characteristic in detail and generality. We will do exercises in that direction to illustrate how it works.
(2.10a) $\chi_{X}$ is an invariant of the homeomorphism class of the compact of $X$ (whether from cuts or not).
(2.10b) If the Euler characteristic is 2 then $X$ is topologically a sphere: genus 0 .
(2.10c) If the Euler characteristic is 0 then $X$ is topologically a torus (as in Chap. 3 Fig. 2): genus 1.

To conclude these results from a triangulation of $X$ in either case requires only laying out on the sphere (resp. torus) an equivalent triangulation [11.6].

REmARK 2.20 (Using a branch point as a base point). The beginning literature on Riemann surfaces has figures with cuts. Often the cuts don't have an obvious base point $z_{0}$ attached to them. That early literature is usually about the nature of integrals of meromorphic differentials around closed paths. So the fundamental group action is through the first homology group $H_{1}\left(U_{z}\right)$. As in Lem. 7.1, analytic continuations of the primitive give a complicated analytic continuation action (of course, not through a finite group). Since this is about integration, [11.16b] explores how to use a branch point as a base point for the cuts.
2.5. Residues and traces. Cauchy's Residue Theorem (Chap. 2 §5.4.4) implies the sum of the residues of any meromorphic differential $\omega$ on $\mathbb{P}_{z}^{1}$ is 0 . We prove the same holds on any compact Riemann surface $X$. Then we give Abel's famous necessary condition for a divisor on $X$ to be the divisor of a meromophic function $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$. That it is also sufficient is the cornerstone of the theory of Riemann surfaces ( $\S 7.6$ for surfaces of genus 0 and 1 , and Chap. $5 \S ? ?$ in general).
2.5.1. Sum of the residues is 0 . Let $\omega \in \mathcal{M}^{1}(X)$ be a meromorphic differential on the compact Riemann surface $X$. Chap. $2 \S 4.3$ has the definition of the residue of a meromorphic differential at $z_{0} \in \mathbb{C}_{z}$. Since $X$ is compact, $\omega$ has but finitely many poles (as in the argument for Lem. 2.1). So, it has only finitely many points at which there is a nonzero residue. There are two approaches to showing the sum of the residues of $\omega$ is 0 . We use here Green's Theorem, to have available the exterior calculus for later. Another approach, reducing the sum of the residues to exactly Cauchy's Theorem in the plane comes from uniformization [11.11].
2.5.2. Orientation and Green's Theorem. When we say a path bounds a closed disk $\bar{D}^{\prime}$ in $X$ we mean here that the oriented path has the disc on its left. Suppose $X$ is 2-dimensional differentiable manifold with atlas $\left\{U_{\alpha}, \varphi_{\alpha}\right\}_{\alpha \in I}$. Use $\left(x_{\alpha}, y_{\alpha}\right)$ for the variables of the range of $\varphi_{\alpha}: U_{\alpha} \rightarrow \mathbb{R}^{2}$. For $\alpha, \beta \in I$, use $F_{\beta, \alpha}=\varphi_{\beta} \circ \varphi_{\alpha}^{-1}: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$ for the transition function on $\varphi_{\alpha}\left(U_{\alpha} \cap U_{\beta}\right)$. A differential 2-form on $X$ consists of giving $f_{\alpha}\left(x_{\alpha}, y_{\alpha}\right) d x_{\alpha} \wedge d y_{\alpha}$ for each $\alpha \in I$, satisfying these two conditions.
(2.11a) $f_{\alpha}\left(x_{\alpha}, y_{\alpha}\right): \mathbb{R}^{2} \rightarrow \mathbb{C}$ is differentiable on $U_{\alpha}$.
(2.11b) $f_{\beta}\left(F_{\beta, \alpha}\left(x_{\alpha}, y_{\alpha}\right)\right)=\operatorname{Det}\left(J\left(F_{\beta, \alpha}\left(x_{\alpha}, y_{\beta}\right)\right)\right) f_{\alpha}\left(x_{\alpha}, y_{\alpha}\right)$ where $J(F)$ denotes the Jacobian matrix as in Chap. 3 Lem. 3.2.
[11.4] reminds that 2-forms appear to form integrals over 2-dimensional subsets of $X$. The change of variables $\left(x_{\alpha}, y_{\alpha}\right) \mapsto\left(y_{\alpha}, x_{\alpha}\right)$ would change the sign of this 2-form. In the case $f_{\alpha}\left(x_{\alpha}, y_{\alpha}\right)$ is invariant under this transformation, the new contribution to integrating over $U_{\alpha}$ would subtract from, not add to, the integral. Fortunately, that is not an allowable transformation of coordinates on a 1-dimensional complex manifold. The $\left(x_{\alpha}, y_{\alpha}\right)$ coordinates come from the complex coordinates $x_{\alpha}+i y_{\alpha}$. Any analytic change of $x_{\alpha}+i y_{\alpha}$ leaves the sign of the determinant positive [11.4a].

Definition 2.21 (Orientation). An orientation on a differentiable dimension 2 manifold is a choice of subatlas for which the determinant of the coordinate transformation Jacobian is always positive.

Proposition 2.22 (Green's Theorem). Suppose $\omega$ is meromorphic 1-form in a domain $D \subset X$. The residue of $\omega$ has a well-defined meaning at each $x^{\prime} \in D$. Denote the set of $x^{\prime} \in D$ at which $\omega$ has a nonzero residue by $R_{\omega}(D)$. Let $\gamma$ be a disjoint union $\gamma_{1}, \ldots, \gamma_{t}$ of simple closed paths on $D$, where each $\gamma_{i}$ is the counterclockwise boundary of a closed topological disk $\bar{D}_{i} \subset D$. Assume $\omega$ has no poles on $\gamma$ and all its residues are in $\cup_{i+1}^{t} \bar{D}_{i}$. Then, $\frac{1}{2 \pi i} \int_{\gamma} \omega=\sum_{x^{\prime} \in R_{\omega}(D)} \operatorname{Res}_{x^{\prime}}(\omega)$. In particular, if $D=X$, then $\frac{1}{2 \pi i} \int_{\gamma} \omega=0$.

More generally, let $\omega^{\prime}$ be any differentiable differental 1-form on the domain $D \backslash \cup_{i+1}^{t} \bar{D}_{i}$ as above. Then, there is a differential 2-form d $\omega^{\prime}$ on $D$ so that

$$
\begin{equation*}
\int_{\gamma} \omega=\int_{D \backslash \cup_{i=1}^{t} \bar{D}_{i}} d \omega . \tag{2.12}
\end{equation*}
$$

Proof. We show the last paragraph first. Use the notation from above for a 2-dimensional differentiable manifold. Then, on $\varphi_{\alpha}\left(U_{\alpha}\right)$ from the coordinate chart $\left(U_{\alpha}, \varphi_{\alpha}\right)$, express $\omega^{\prime}$ as $f_{\alpha}\left(x_{\alpha}, y_{\alpha}\right) d x_{\alpha}+g_{\alpha}\left(x_{\alpha}, y_{\alpha}\right) d y_{\alpha}$. The production of the differential 2 -form from $\omega^{\prime}$ comes from the exterior derivative:

$$
\begin{align*}
d\left(\omega^{\prime}\right) & =d f_{\alpha} \wedge d x_{\alpha}+d g_{\alpha} \wedge d y_{\alpha} \\
& =\frac{\partial f_{\alpha}}{\partial y_{\alpha}} d y_{\alpha} \wedge d x_{\alpha}+\frac{\partial g_{\alpha}}{\partial x_{\alpha}} d x_{\alpha} \wedge d y_{\alpha}=\left(\frac{\partial g_{\alpha}}{\partial x_{\alpha}}-\frac{\partial f_{\alpha}}{\partial y_{\alpha}}\right) d x_{\alpha} \wedge d y_{\alpha} \tag{2.13}
\end{align*}
$$

We must establish this is a 2-form: (2.11b) holds [11.23b]. Then, the integration on the right of (2.12) is independent of the coordinate chart. We already know that is true of the integration on the left from Chap. 2 Lem. 2.3. Then, the conclusion is a consequence of Green's Theorem from vector calculus in the plane. While [Rud76, p. 272] has a complete treatment, as our paths are semi-simplicial, the case of bounding by rectangles suffices.

To apply the result to the first paragraph requires only noting that if we are on a 1-dimensional complex manifold, then locally an analytic differential has the form $f_{\alpha}\left(z_{\alpha}\right) d z_{\alpha}$. In that case the Cauchy-Riemann equations immediately imply $d\left(f_{\alpha}\left(z_{\alpha}\right) d z_{\alpha}\right)=0$ [11.4b].

Remark 2.23. Apply Thm. 2.25 to the translate of $\varphi$ by a constant, $\varphi-c$, $c \in \mathbb{C}$. Conclude $\operatorname{deg}\left(D_{z^{\prime}}\right)$ is constant running over all $z^{\prime} \in \mathbb{P}_{z}^{1}$, a case of Lem. 2.1.
2.5.3. Traces of differentials and functions. Let $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ be an analytic map of compact Riemann surfaces. Use notation from the coordinate chart from $\varphi$ (Def. 2.8). Denote meromorphic differentials on $X$ by $\Gamma\left(X, \mathcal{M}^{1}\right)$. Suppose $\omega \in$ $\Gamma\left(X, \mathcal{M}^{1}\right)$, and $\boldsymbol{z}$ is the branch point set of $\varphi$. For $z^{\prime} \notin \boldsymbol{z}$, consider $D_{z^{\prime}}=\sum_{j=1}^{n} x_{j}$. Since $z^{\prime}$ is not a branch point, there is a neighborhood $U_{z^{\prime}}$ of $z^{\prime}$ and $U_{x_{i}}$ so $\varphi$ is invertible on $U_{x_{i}}$. To keep our neighborhoods straight, denote the inverse of $\varphi$ on $U_{x_{i}}$ by $\varphi_{i}^{-1}$. For each $i \in\{1, \ldots, n\}, \varphi_{i}^{-1}: U_{z^{\prime}} \rightarrow U_{x_{i}}$ is a section for $\varphi$. Denote the local variable on $U_{x_{i}}$ by $w_{i}$. On $\varphi_{i}\left(U_{x_{i}}\right)$ write $\omega$ as

$$
h_{i}\left(w_{i} \circ \varphi_{i}^{-1}(z)\right) d\left(w_{i} \circ \varphi_{i}^{-1}(z)\right) .
$$

Define $\mathrm{t}(\omega)$ on $U_{z^{\prime}}$ as a differential in $z$ by $\sum_{i=1}^{n} h_{i}\left(w_{i} \circ \varphi_{i}^{-1}(z)\right) d\left(w_{i} \circ \varphi_{i}^{-1}(z)\right)$.
We extend this around ramified points (when $z^{\prime} \in \boldsymbol{z}$ ) where $e_{i}$ is the ramification index of $x_{i}$ in the fiber $X_{z^{\prime}}$ and $\left\{x_{1}, \ldots, x_{t}\right\}=X_{z^{\prime}}$. Let $\zeta_{e_{i}}=e^{2 \pi i / e_{i}}$, exactly as in Lem. 2.5. To simplify notation designate $\varphi$ as $z$. The only extension of $\mathrm{t}(\omega)$ that
gives the same values over a deleted neighborhood of $U_{z^{\prime}}$ requires the expression $\sum_{i=1}^{t} \sum_{j=0}^{e_{i}-1} h_{i}\left(\zeta_{e_{i}}^{j} z^{1 / e_{i}}\right) d\left(\zeta_{e_{i}}^{j} z^{1 / e_{i}}\right)$ for $\mathrm{t}(\omega)$. Write $d z=e_{i} z^{\frac{e_{i}-1}{e_{i}}} d w_{i}$ to reexpress the contribution around $x_{i}$ as

$$
\begin{equation*}
\sum_{j=0}^{e_{i}-1} \frac{h_{i}\left(\zeta_{e_{i}}^{j} z^{1 / e_{i}}\right)}{e_{i} z^{\frac{e_{i}-1}{e_{i}}}} d z \tag{2.14}
\end{equation*}
$$

So, (2.14) is a Laurent series in $z^{1 / e_{i}}$ times $d z$, symmetric in $\left\{\zeta_{e_{i}}^{j} z^{1 / e_{i}}\right\}_{j=0}^{e_{i}-1}$, the conjugates of $z^{1 / e_{i}}$ over $\mathbb{C}\{\{z\}\}$. Conclude: Each term in $\mathrm{t}(\omega)$, the trace of $\omega$ is a Laurent series in $z$ (times $d z$ ), and $\mathrm{t}(\omega)$ is a differential on $\mathbb{P}_{z}^{1}$.

REmARK 2.24. There is a similar definition of trace for meromorphic functions (elements of $\mathbb{C}(X))$ on $X$. Further, the following extensions are also easy: We may replace $\mathbb{P}_{z}^{1}$ by any Riemann surface $Y$ (not necessarily compact) and $\varphi: X \rightarrow Y$ is a ramified cover. Recall: Regard meromorphic differentials (resp. functions) on $Y$ as meromorphic differentials (resp. functions) on $X$ by pullback (Chap. 3 §5.3.3).

Theorem 2.25. Given a ramified cover $\varphi: X \rightarrow Y$ of Riemann surfaces, the trace $\mathrm{t}=\mathrm{t}_{X / Y}$ from meromorphic differentials on $X$ to those on $Y$ is a $\mathbb{C}$-linear. It maps homolomorphic differentials to holomorphic differentials. In particular, if $Y=\mathbb{P}_{z}^{1}$, then the range of t on holomorphic differentials is 0 .

If $\omega \in \Gamma\left(Y, \mathcal{M}^{1}\right)$, then $\mathrm{t}\left(\varphi^{*}(\omega)\right)=\operatorname{deg}(X / Y) \omega$.
Proof. Consider the statements on holomorphicity. If $\omega$ is holomorphic, each $h_{i}$ above is holomorphic. From $(2.14), \mathrm{t}(\omega)$ has a pole of order no more than $\frac{e_{i}-1}{e_{i}}$ at $z^{\prime}$. The order, however, of the pole must be an integer. That means it has no pole at $z^{\prime}$ and $\omega$ is holomorphic. As there are no holomorphic differentials on the sphere (Chap. 3 Ex. 5.17), t $(\omega)$ vanishes.

More generally, if $\omega$ is any differential, then its trace has the same sum of residues as does $\omega$. This comes back to the case the differential is locally $d x_{i} / x_{i}$ with its trace locally reexpressed as $\sum_{j=0}^{e_{i}-1} h\left(\zeta_{e_{i}}^{j} z^{1 / e_{i}}\right) / e_{i} z^{\frac{e_{i}-1}{e_{i}}} d z$ with $h_{i}=1 / x_{i}$. The final equation is a consequence of the definitions and Rem. 2.24.
2.6. Abel's necessary condition. With $X$ a compact Riemann surface, let $\Gamma(X, \Omega)$ be the vector space of global holomorphic differentials on $X$. We don't know its dimension yet, though Lem. 6.14 shows it is $g_{X}=g_{\boldsymbol{g}}$ (as in Thm. 3.10). Suppose $D^{0}=\sum_{i=1}^{n} x_{i}^{0}$ and $D^{\infty}=\sum_{i=1}^{n} x_{i}^{\infty}$ are two degree $n$ divisors on $X$. We allow some points repeated with multiplicity.

Consider those $n$-tuples of paths $\gamma=\left(\gamma_{1}, \ldots, \gamma_{n}\right)$ for which there is $\sigma \in S_{n}$, with $\gamma_{i}$ having beginning point $x_{i}^{0}$ and end point $x_{(i) \sigma}^{\infty}, i=1, \ldots, n$. Denote these by $\Pi_{1}\left(X, D^{0}, D^{\infty}\right)$. If $\sigma$ is the appropriate permutation, define the endpoint evaluation map by $E_{D^{0}, D^{\infty}}(\gamma)=\sigma$. When the support of $D^{\infty}$ consists of distinct points, this defines $\pi$ uniquely, otherwise it is a coset of the subgroup of permutations stabilizing the ordered set $\left(x_{1}^{\infty}, \ldots, x_{n}^{\infty}\right)$.
2.6.1. Integrating a basis of holomorphic differentials. Abel's necessary condition tests for existence of a meromorphic degree $n$ function on $X$ whose divisor of zeros (resp. poles) is $D^{0}$ (resp. $D^{\infty}$ ). It is tacit that $D^{0}$ and $D^{\infty}$ have no common support and are both positive divisors. Lem. 2.1 says the divisor of zeros and poles determine a function up to multiplication by a constant.

The test on integrals is made efficient by using a basis $\mathcal{B} \stackrel{\text { def }}{=}\left(\omega_{1}, \ldots, \omega_{u}\right)$ for $\Gamma(X, \Omega)$. We integrate the entries of $\mathcal{B}$ along elements of $\Pi_{1}\left(X, D^{0}, D^{\infty}\right)$. Such integrals are equivalent to evaluating analytic continuations of a branch of a primitive (Chap. $2 \S 4.3$ ). So, the monodromy theorem says results will only depend on homotopy classes of such paths (with their endpoints fixed; Chap. 2 Thm. 8.3). Denote these $\pi_{1}\left(X, D^{0}, D^{\infty}\right)$. For these definitions we may allow common support to $D^{0}$ and $D^{\infty}$. When, however, $D^{0}=D^{\infty}$, write $\pi_{1}\left(X, D^{0}\right)$ for the homotopy classes of $n$-tuples of closed paths. In this case, the paths in an ordered $n$-tuples of paths may each have a different end point than beginning point. The case $D^{0}=n x_{0}$ is allowed, to indicate an $n$-tuple of closed paths.

From Thm. 2.6, each meromorphic function on $X$ gives an analytic map $\varphi$ : $X \rightarrow \mathbb{P}_{z}^{1}$. This gives a map from $\gamma \in \pi_{1}\left(X, D^{0}, D^{\infty}\right)$ to the integral of $\mathcal{B}$ over $\gamma$ :

$$
\operatorname{Int}_{D^{0}, D^{\infty}}=\operatorname{Int}_{X, D^{0}, D^{\infty}}(\gamma) \stackrel{\text { def }}{=} \int_{\gamma} \mathcal{B}=\left(\sum_{j=1}^{n} \int_{\gamma_{j}} \varphi_{1}, \ldots, \sum_{j=1}^{n} \int_{\gamma_{j}} \varphi_{u}\right)
$$

THEOREM 2.26. The range of $\operatorname{Int}_{X, m x_{0}}$, for $x_{0} \in X$, is an abelian subgroup $L_{X}$ of $\mathbb{C}^{u}$, independent of either $z_{0}$ or $m \geq 1$. A change of basis for $\Gamma(X, \Omega)$ changes $L_{X}$ by the action (on the left) of some element of $\mathrm{GL}_{n}(\mathbb{C})$.

Suppose there is a nonconstant analytic map $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ with $D^{0}=\varphi^{-1}(0)$ and $D^{\infty}=\varphi^{-1}(\infty)$. Then, $\operatorname{ker}\left(\operatorname{Int}_{D^{0}, D^{\infty}}\right) \neq \emptyset$ and the range of $\operatorname{Int}_{D^{0}, D^{\infty}}$ is $L_{X}$.

Let $\boldsymbol{z}$ be the branch points of $\varphi$, and suppose $0 \notin \boldsymbol{z}$ (resp. $\infty \notin \boldsymbol{z}$ ). Then, $\pi_{1}\left(U_{\boldsymbol{z}}, 0\right)$ (resp. $\pi_{1}\left(U_{\boldsymbol{z}}, \infty\right)$ ) has a faithful left (resp. right) action on $\operatorname{ker}\left(\operatorname{Int}_{D^{0}, D^{\infty}}\right)$. Therefore, $\left\{E_{D^{0}, D^{\infty}}(\boldsymbol{\gamma})\right\}_{\boldsymbol{\gamma} \in \operatorname{ker}\left(\operatorname{Int}_{D^{0}, D^{\infty}}\right)}$ contains the monodromy group $G_{\varphi}$ of $\varphi$. This holds even if $0 \in \boldsymbol{z}$ (resp. $\infty \in \boldsymbol{z}$ ) using a tangential base point at 0 (resp. $\infty$ ).
2.6.2. Proof of Thm. 2.26 and integrations along $\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Consider $\boldsymbol{\gamma}, \boldsymbol{\gamma}^{\prime} \in \pi_{1}\left(X, m x_{0}\right)$. Then, the component wise product $\boldsymbol{\gamma} \cdot \boldsymbol{\gamma}^{\prime}=\left(\gamma_{1} \cdot \gamma_{1}^{\prime}, \ldots, \gamma_{1} \cdot \gamma_{1}^{\prime}\right)$ is in $\pi_{1}\left(X, m x_{0}\right)$. Apply Int to these to see the range is independent of $m$ and is an abelian group. Given another basis $\mathcal{B}^{\prime}$, there exists $A \in \mathrm{GL}_{n}(\mathbb{C})$ so that $A(\mathcal{B})=\mathcal{B}^{\prime}$. Therefore $A\left(\int_{\gamma}(\mathcal{B})\right)=\int_{\gamma} A(\mathcal{B})$ has range in $A\left(L_{X}\right)$.

Now suppose $\varphi$ exists. Start with the case $D^{0}$ and $D^{\infty}$ have $n$ distinct points in their support. Let $\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, 0, \infty\right)$, and define $\gamma=\left(\gamma_{1}, \ldots, \gamma_{n}\right)$ so $\gamma_{i}$ is the unique lift of $\gamma$ starting at $x_{i}^{0}$. Write $\gamma:[0,1] \rightarrow U_{\boldsymbol{z}}$ to define $\left(\gamma_{1}(t), \ldots, \gamma_{n}(t)\right)$ for $t \in[0,1]$, an ordering of $\varphi^{-1}(t)$.

Apply Thm. 2.25 to $\operatorname{Int}_{D^{0}, D^{\infty}}(\gamma)$ by designating the trace from $\varphi$ by $t_{\varphi}$. Then, $\operatorname{Int}_{D^{0}, D^{\infty}}(\gamma)$ is just $\left(\int_{\gamma} \mathrm{t}_{\varphi}\left(\omega_{1}\right), \ldots, \mathrm{t}_{\varphi}\left(\omega_{u}\right)\right)$. Since each of the integrand entries is 0 , this shows that any element of $\pi_{1}\left(U_{\boldsymbol{z}}, 0, \infty\right)$ defines an element of $\operatorname{ker}\left(\operatorname{Int}_{D^{0}, D^{\infty}}\right)$.

If either $D^{0}$ or $D^{\infty}$ has support with multiplicity, connect 0 and $\infty$ by paths $\gamma_{z^{\prime}}$ and $\gamma_{z^{\prime \prime}}$ to respective points $z^{\prime}$ and $z^{\prime \prime}$ that lie (excluding endpoints) entirely in $U_{\boldsymbol{z}}$. Let $\gamma$ denote a path in $U_{\boldsymbol{z}}$ connecting $z^{\prime}$ and $z^{\prime \prime}$. There is still an $n$-tuple of lifts of the path $\gamma_{0} \cdot \gamma \cdot \gamma_{\infty}^{-1}:(0,1) \rightarrow X_{0}$ (avoiding endpoints). Now form the paths to replace $t \mapsto\left(\gamma_{1}(t), \ldots, \gamma_{n}(t)\right)$ by taking the closure of these lifted paths in $X$. The integral is 0 , again from Thm. 2.25.

Similarly, we may compose on the right of $\operatorname{Int}_{D^{0}, D^{\infty}}$ by $\pi_{1}\left(U_{\boldsymbol{z}}, \infty\right)$ so long as 0 and $\infty$ are not in $\boldsymbol{z}$. Suppose, however, $0 \in \boldsymbol{z}$ (the case for $\infty$ is analogous). Then, $\pi_{1}\left(U_{z}, 0\right)$ doesn't make sense.

Chap. $2 \S 8.4$ has the notion of a tangential base point. We need a convenient (nonempty) simply connected open set $D_{\boldsymbol{v}}$ tangent to 0 in $U_{\boldsymbol{z}}$. The choice there was
a disk with 0 on the boundary, defined by a tangent vector $\boldsymbol{v}$ to 0 . Let $\lambda:[0,1] \rightarrow \bar{D}_{\boldsymbol{v}}$ be a path with these properties: $\lambda(0)=0$, restriction to $(0,1]$ has range in $D_{v}$ and $\lambda(1)=z^{\prime} \in D_{\boldsymbol{v}}$. Consider paths (minus beginning and endpoint) given by $\lambda_{(0,1]} \cdots \gamma \lambda_{(0,1]}^{-1}, \gamma$ representing an element of $\pi_{1}\left(U_{\boldsymbol{z}}, D_{\boldsymbol{v}}, z^{\prime}\right)$. This has $n$ distinct lifts to $X$. Their closures have their beginning and end points in $D^{0}$.

Up to homotopy, these paths don't depend on $z^{\prime}$ or $\lambda$ (though it does on $D_{\boldsymbol{v}}$ ). So, up to homotopy, composition of these paths defines a group $\pi_{1}\left(U_{\boldsymbol{z}}, D_{\boldsymbol{v}}\right)$ with an action on the left of $\operatorname{ker}\left(\operatorname{Int}_{D^{0}, D^{\infty}}\right)$. The isomorphism class of the group is the same no matter the choice of $D_{\boldsymbol{v}}$. There is, however, no canonical isomorphism between the groups if you change $D_{\boldsymbol{v}}$ to another tangential disk [11.16a].

## 3. Nielsen classes and Hurwitz monodromy

This section introduces combinatorial group theory that helps display the myriad covers from Cor. 2.9. $\S 4$ uses this to illustrate Riemann's Existence Theorem. We suggest the reader go between the two sections on a first reading; we put many concepts together in this section. That includes interpretation of the genus of a compact surface, and the related fiber product and Galois closure of compact covers topics. Braid and Hurwitz monodromy representations are critical to this book. [Ar25], [Ar47], [Bi75], [Boh47], [Ch47], [KMS66], [Ma34], [Mar45] hint at early literature on the Braid group. None, however, of these sources apply these to the families of Riemann surface covers. Further, the Hurwitz monodromy group is a modest player in them though some of their combinatorics, especially [Boh47] and [KMS66], appears in our picture.
3.1. Artin Braids and Hurwitz monodromy. Let $F_{r}$ be the free group on the elements of $S=\left\{s_{1}, \ldots, s_{r}\right\}$. Since $F_{r}$ is a free group, any $r$ words $w_{1}, \ldots, w_{r}$ in $S$ determine a homomorphism of $F_{r}$ into itself by mapping the ordered $r$-tuple $\left(s_{1}, \ldots, s_{r}\right)=\boldsymbol{s}$ respectively to $\left(w_{1}, \ldots, w_{r}\right)$. So, given $\boldsymbol{s}$, any other $r$-tuple, $\left(s_{1}^{\prime}, \ldots, s_{r}^{\prime}\right)$, of generators of $F_{r}$ determines an element of the automorphism group $\operatorname{Aut}\left(F_{r}\right)$ of $F_{r}$. Denote the set of (ordered) $r$-tuples of generators of $F_{r}$ by $\mathcal{G}_{F_{r}}$.
3.1.1. Automorphisms of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ permuting classical generators. Certain automorphisms of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ play a big role from here on. Chap. 5 describes the geometry that produces them. Here they are a combinatorial tool.

Let $Q_{i}$ be the permutation of $\mathcal{G}_{F_{r}}$ that sends entries of $\left(s_{1}, \ldots, s_{r}\right)=\boldsymbol{s}$ (in order) to the new r-tuple of generators

$$
\begin{equation*}
\left(s_{1}, \ldots, s_{i-1}, s_{i} s_{i+1} s_{i}^{-1}, s_{i}, s_{i+2}, \ldots, s_{r}\right), i=1, \ldots, r-1 \tag{3.1}
\end{equation*}
$$

The Artin braid group (of degree $r$ ), is the subgroup of permutations of $\mathcal{G}_{F_{r}}$ that $Q_{1}, \ldots, Q_{r-1}$ generate. We denote it $B_{r}$.

Lemma 3.1. Any $\boldsymbol{s} \in \mathcal{G}_{F_{r}}$ gives a faithful map $\psi_{\boldsymbol{s}}: B_{r} \rightarrow \operatorname{Aut}\left(F_{r}\right): Q \in B_{r}$ maps to the automorphism that takes $\boldsymbol{s}$ to $(\boldsymbol{s}) Q$. Suppose $Q \in B_{r}$ and $\alpha \in \operatorname{Aut}\left(F_{r}\right)$. Then, $Q$ acts on $\alpha$ by this formula: $(\boldsymbol{s}) \alpha^{Q} \stackrel{\text { def }}{=}(\boldsymbol{s}) Q^{-1} \alpha Q$. The action of $B_{r}$ on inner automorphisms of $F_{r}$ is trivial. Also, $\psi_{\boldsymbol{s}}$ is a 1-cocycle on the group $B_{r}$ : $\left(Q Q^{\prime}\right) \psi_{\boldsymbol{s}}=\left(\left(Q^{\prime}\right) \psi_{\boldsymbol{s}}\right)(Q) \psi_{\boldsymbol{s}}^{Q^{\prime}}$.

Proof. The effect of $Q \in B_{r}$ on any one $\boldsymbol{s} \in \mathcal{G}_{F_{r}}$ determines it. So, $\psi_{\boldsymbol{s}}$ is faithful. Notice that conjugation by $w$ commutes with the action of $Q_{i}$ on $\boldsymbol{s}$. As
these are generators of $B_{r}$, this implies $\alpha^{Q}=\alpha$ for $Q \in B_{r}$ if $\alpha$ is conjugation by $w$. Check how both sides of the cocycle condition act on $s$ :

$$
(\boldsymbol{s}) Q Q^{\prime}=(\boldsymbol{s})\left(Q Q^{\prime}\right) \psi_{\boldsymbol{s}}=((\boldsymbol{s}) Q) Q^{\prime}=\left(\left((\boldsymbol{s}) Q^{\prime} Q^{\prime-1}\right)(Q) \psi_{\boldsymbol{s}}\right) Q^{\prime}=\left((\boldsymbol{s})\left(Q^{\prime}\right) \psi_{\boldsymbol{s}}\right)(Q) \psi_{\boldsymbol{s}}^{Q^{\prime}}
$$

This concludes the lemma.
3.1.2. Hurwitz monodromy quotient of the braids. The word cocycle in Lem. 3.1 has a more complicated meaning than in Chap. 3 §5.4.1 where it was a condition on transition functions. This is a group cocycle, for a group acting on a nonabelian group (rather than on a module). Our emphasis is that $\psi_{\boldsymbol{s}}$ is a cocycle, not a homomorphism. The Hurwitz monodromy group (of degree $r$ ) is the quotient of $B_{r}$ by the normal subgroup generated by

$$
\begin{equation*}
Q(r)=Q_{1} Q_{2} \cdots Q_{r-1} Q_{r-1} \cdots Q_{2} Q_{1} \tag{3.2}
\end{equation*}
$$

Denote this quotient group by $H_{r}$.
Observations from the following proposition will appear in examples of $\S 4$. It simplifies reading Chap. 5 to be already acquainted with these. Let $\bar{R}$ be the normal subgroup of $F_{r}$ that $s_{1} \cdots s_{r}=u_{\boldsymbol{s}}$ generates (Ex. 1.3). Denote $F_{r} / \bar{R}$ by $G_{r}$.

Proposition 3.2. The following properties hold for $B_{r}$ (acting on $\mathcal{G}_{r}$ ).
(3.3a) Each $Q \in B_{r}$ maps $s_{1} \cdots s_{r}$ to itself and $s_{i}$ to a conjugate of $s_{j}$ for some $j$ (dependent on $i$ ). This induces a homomorphism $\Psi_{r, *}: B_{r} \rightarrow S_{r}$ (the Noether representation) mapping $Q_{i}$ to $(i i+1) \in S_{r}, i=1, \ldots, r$.
(3.3b) The $Q_{i}$ s have these relations: $Q_{i} Q_{j}=Q_{j} Q_{i}, 1 \leq i \leq j \leq r-1 ; j \neq i-1$ or $i+1$, and $Q_{i} Q_{i+1} Q_{i}=Q_{i+1} Q_{i} Q_{i+1}, i=1, \ldots, r-2$.
(3.3c) Elements of $\operatorname{ker}\left(B_{r} \rightarrow H_{r}\right)$ induce inner automorphisms of $G_{r}$.

Proof. Each formula is a simple computation on the effect of sides of the equation on elements of $S$. For example, since $S$ is a set of generators, to see (3.3a) note that the result of applying any $Q_{i}$ to $S$ is another generating set. Then, induct on the length of a word in the $Q_{i}$ s to conclude the result from the application of $Q_{i}$ which maps $s_{i}$ to a conjugate of $s_{i+1}$ and $s_{i+1}$ to $s_{i}$.

The first relation of (3.3b) is obvious, for $Q_{i}$ and $Q_{j}$ with $i$ and $j$ separated, move indices with no common support. The other formula follows from a renaming of the indices and showing that $Q_{1} Q_{2} Q_{1}=Q_{2} Q_{1} Q_{2}$ in its application to $\left(s_{1}, s_{2}, s_{3}\right)$ :

$$
\begin{aligned}
& \left(s_{1}, s_{2}, s_{3}\right) Q_{1} Q_{2} Q_{1}=\left(s_{1} s_{2} s_{1}^{-1}, s_{1} s_{3} s_{1}^{-1}, s_{1}\right) Q_{1}=\left(s_{1} s_{2} s_{3} s_{2}^{-1} s_{1}^{-1}, s_{1} s_{2} s_{1}^{-1}, s_{1}\right) \\
& \left(s_{1}, s_{2}, s_{3}\right) Q_{2} Q_{1} Q_{2}=\left(s_{1} s_{2} s_{3} s_{2}^{-1} s_{1}^{-1}, s_{1}, s_{2}\right) Q_{2}=\left(s_{1} s_{2} s_{3} s_{2}^{-1} s_{1}^{-1}, s_{1} s_{2} s_{1}^{-1}, s_{1}\right)
\end{aligned}
$$

Finally, consider an extension of this computation.

$$
\begin{aligned}
\left(s_{1}, \ldots, s_{r}\right) Q(r) & =\left(s_{1} s_{2} s_{1}^{-1}, s_{1} s_{3} s_{1}^{-1}, \ldots, s_{1} s_{r} s_{1}^{-1}, s_{1}\right) Q_{r-1} \cdots Q_{1} \\
& =\left(u_{s} s_{1} u_{\boldsymbol{s}}^{-1}, s_{1} s_{2} s_{1}^{-1}, s_{1} s_{3} s_{1}^{-1}, \ldots, s_{1} s_{r} s_{1}^{-1}\right)
\end{aligned}
$$

As $u_{\boldsymbol{s}}$ has image the identity in the group $G_{r}, Q(r)$ induces conjugation by $s_{1}^{-1}$ in $G_{r}$. If $Q \in B_{r}$ maps $\left(s_{1}, \ldots, s_{r}\right)$ to $\left(s_{1}^{\prime}, \ldots, s_{r}^{\prime}\right)$, then $Q Q(r) Q^{-1}$ gives this chain of mappings: $\boldsymbol{s} \mapsto \boldsymbol{s}^{\prime} \mapsto\left(s_{1}^{\prime} \boldsymbol{s}^{\prime}\left(s_{1}^{\prime}\right)^{-1}\right) Q^{-1}=s_{1}^{\prime} \boldsymbol{s}\left(s_{1}^{\prime}\right)^{-1}$. Everything in $\operatorname{ker}\left(B_{r} \rightarrow H_{r}\right)$ is a product of powers of elements of form $Q Q(r) Q^{-1}$. So, this shows (3.3c).
3.2. s-equivalences on Nielsen classes. The original definition of Nielsen class is from [Fri77]. Special cases appearing in [Fri73], and many illustrating examples related to elliptic curves in [Fri78]. They loom large in the books of Matzat-Malle and Voelklein. The former calls them generating s-systems [MM95, p. 26] (our $r$ is their $s$ ) and the latter uses the name ramification type [V̈̈96, p. 37] for the most closely related definition.

An old literature on simple branched covers influenced classical geometers ([Cl1872], [Hu1891]). This continued through papers of Lefschetz, Segre and Zariski. Simple branched covers apply to the moduli space of genus $g$ curves, knot types and Lefschetz pencils (of surfaces). Our interest came through complex multiplication and modular curves. We found every finite group produces a modular curve-like setup (Chap. 5 §??). S(trong)-equivalences and r(educed)-equivalences classes on elements of Nielsen classes give geometric meanings to some valuable group properties. These showed the Inverse Galois Problem fit very generally with many classical problems. A reader will require time to acclimate to these.
3.2.1. Setup for Nielsen classes. Consider any cover of compact connected Riemann surfaces $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ with $r$ branch points $\boldsymbol{z}$. Denote the degree of the cover by $n$. Thm. 2.6 shows one way to picture how that cover arises. Choose an ordered $r$-tuple of classical generators $\boldsymbol{s}$ for $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Then $\varphi$ and an ordering of the points of $X$ over $z_{0}$ determines the image of the entries of $\boldsymbol{s}$ in the monodromy group $G$ of the cover: Each $s_{i}$ in $\boldsymbol{s}$ maps to some $g_{i} \in G$.

Conversely, given $\boldsymbol{s}$ and $\boldsymbol{g}=\left(g_{1}, \ldots, g_{r}\right)$ generators of $G$ satisfying the productone condition $g_{1} \cdots g_{r}=1$, interpreting $\boldsymbol{s}$ as cuts (§2.4.3) attached according to the branch cycle description $\boldsymbol{g}$ produces $\varphi$ (Def. 2.4).

As $\boldsymbol{s}$ runs over all classical generators, Thm. 1.8 gives this data attached to $\varphi$ :
(3.4a) an associated group $G=G(\boldsymbol{g})$;
(3.4b) a permutation representation $T: G \rightarrow S_{n}$; and
(3.4c) conjugacy classes $\mathbf{C}=\left(\mathrm{C}_{1}, \ldots, \mathrm{C}_{r}\right)$ of $G$ into which entries of $\boldsymbol{g}$ fall in some order (denoted $\boldsymbol{g} \in \mathbf{C}$ ).
Further, running over all possible classical generators $\boldsymbol{s}$, the collection of images of $\boldsymbol{s}$ (branch cycle descriptions $\boldsymbol{g}$ ) that correspond to $\varphi$ all fall in this set:

$$
\begin{equation*}
\mathrm{Ni}(G, \mathbf{C}, T)=\left\{\left(g_{1}, \ldots, g_{r}\right) \mid \prod_{i=1}^{r} g_{i}=1, G(\boldsymbol{g})=G \leq S_{n}, \boldsymbol{g} \in \mathbf{C}\right\} \tag{3.5}
\end{equation*}
$$

We often use $\Pi(\boldsymbol{g})$ in place of $\prod_{i=1}^{r} g_{i}$. Then, (3.5) is the Nielsen class of $(r$ tuples in $G$ ) corresponding to $(G, \mathbf{C}, T)$. Elements in this set are Nielsen class representatives.
3.2.2. The s(trong)-equivalences on $\mathrm{Ni}(G, \mathbf{C}, T)$. Consider the subgroup of $S_{n}$ that normalizes $G$ and permutes entries of $\mathbf{C}$. Denote this $N_{S_{n}}(G, \mathbf{C})=N_{T}(G, \mathbf{C})$. For convenience we list some equivalences on a Nielsen class that will appear later. For $N$ any group between $G$ and $N_{T}(G, \mathbf{C})$, let $n \in N$ act on $\boldsymbol{g} \in \mathrm{Ni}(G, \mathbf{C}, T)$ by

$$
\boldsymbol{g} \mapsto n \boldsymbol{g} n^{-1} \stackrel{\text { def }}{=}\left(n g_{1} n^{-1}, \ldots, n g_{r} n^{-1}\right)
$$

Denote the orbits for this action by $\operatorname{Ni}(G, \mathbf{C}, T) / N$.
We reserve a special notation, for two cases:
(3.6a) $\mathrm{Ni}(G, \mathbf{C}, T)^{\text {abs }}$ when $N=N_{T}(G, \mathbf{C})$ for absolute s-equivalence classes (of Nielsen class representatives); and
(3.6b) $\mathrm{Ni}(G, \mathbf{C}, T)^{\mathrm{in}}$ when $N=G$ and $T$ is the regular representation (acting on cosets of the identity subgroup), for inner s-equivalence classes.
In applyings Prop. 3.2, for an element $Q \in B_{r}$, when possible use the notation $q$ for its image in $H_{r}$. For all s-equivalences, Prop. 3.2 gives an action of $H_{r}$ that preserves these equivalence classes. Here is how the generator $q_{i} \in H_{r}$ acts on $\boldsymbol{g} \in \mathrm{Ni}(G, \mathbf{C}, T) / N$, corresponding to (3.1):

$$
\begin{equation*}
(\boldsymbol{g}) q_{i} \stackrel{\text { def }}{=}\left(g_{1}, \ldots, g_{i-1}, g_{i} g_{i+1} g_{i}^{-1}, g_{i}, g_{i+2}, \ldots, g_{r}\right), i=1, \ldots, r-1 \tag{3.7}
\end{equation*}
$$

As in Chap. $3 \S 7.1 .2$ denote the elements $g \in G$ with $(1) T(g)=1$ by $G(T, 1)$.
Suppose $G$ is abelian. Then, the action of $q \in H_{r}$ permutes the entries of $\boldsymbol{g}$ according to $\Psi_{r, *}(q) \in S_{r}$. This holds for inner classes. We give some standard situations that generalize this using the commutator notation $\left(g_{1}, g_{2}\right) \stackrel{\text { def }}{=} g_{1} g_{2} g_{1}^{-1} g_{2}^{-1}$. Let $G^{*}=(G, G)=\left\langle\left(g_{1}, g_{2}\right) \mid g_{1}, g_{2} \in G\right\rangle$ be the commutator subgroup of $G$.

For $G$ any finite group and $H \triangleleft G$, suppose $T$ is a transitive permutation representation of $G$ and $T^{G / H}$ is the induced representation of $G / H$ from the cosets of $G(T, 1) /(G(T, 1) \cap H \equiv G(T, 1) \cdot H / H$. The next lemma follows from the definitions. We will see this situation come up often. We do not assume $\mathbf{C}$ is a set of conjugacy classes whose elements lie outside $H$. So it is possible some entries of $\mathbf{C}$ will become trivial mod $H$.

Lemma 3.3. Mapping $\boldsymbol{g} \in \mathrm{Ni}\left(G, \mathbf{C}, T_{G}\right)$ to the $r$-tuple with entries reduced modulo $H$ produces a natural map $\psi_{G, \mathbf{C}, T_{G} ; H}: \mathrm{Ni}\left(G, \mathbf{C}, T_{G}\right) \rightarrow \mathrm{Ni}\left(G / H, \mathbf{C} / H, T_{G / H}\right)$. This commutes with the action of $B_{r}: \psi_{G, \mathbf{C}, T_{G} ; H}$ is $B_{r}$ equivariant (Chap. 3 §7.1.3).

Any $N$ between $G$ and $N_{S_{n}}(G, \mathbf{C})$ that also normalizes $H$ produces an $H_{r}$ equivariant map $\mathrm{Ni}\left(G, \mathbf{C}, T_{G}\right) / N \rightarrow \mathrm{Ni}\left(G / H, \mathbf{C}, T_{G / H}\right) /(N / H)$.
3.3. Normal fiber products and Galois closure. We inspect the fiber product of two compact Riemann surfaces $\varphi_{i}: X_{i} \rightarrow \mathbb{P}_{z}^{1}$ by comparing two natural choices. According to Prop. 4.9 the naive fiber product $X_{1} \times_{\mathbb{P}_{z}^{1}} X_{2}$ will produce an analytic manifold at a point $\left(x_{1}^{\prime}, x_{2}^{\prime}\right)$ lying over $z^{\prime} \in \mathbb{P}_{z}^{1}$ if and only if at least one of the corresponding pairs of ramification orders $e_{x_{i}^{\prime} / z^{\prime}}$ is $1, i=1,2$. It also showed there really should be be $d=\left(e_{x_{1}^{\prime} / z^{\prime}}, e_{x_{2}^{\prime} / z^{\prime}}\right)$ distinct points (with ramification orders $\left[e_{1}^{\prime}, e_{2}^{\prime}\right]$ over $\left.z^{\prime}\right)$ in this fiber product corresponding to the pair ( $x_{1}^{\prime}, x_{2}^{\prime}$ ). Riemann's Existence Theorem combinatorially gives that by forming a fiber product in the category of compact Riemann surfaces (Prop. 3.4).
3.3.1. Fiber products of compact Riemann surfaces. For a given compact Riemann surface $Y$ let $\mathcal{C}_{X}^{c}$ be the category of finite covers $\varphi: X \rightarrow Y$ of compact Riemann surfaces where a map between two $\varphi_{i}: X_{i} \rightarrow Y, i=1,2$, is a map of Riemann surfaces $\psi: X_{1} \rightarrow X_{2}$ that commutes with the maps to $Y: \varphi_{2} \circ \psi=\varphi_{1}$. Let $\boldsymbol{y}$ be the union of the branch points for $\varphi_{1}$ and $\varphi_{2}$, and denote $Y \backslash\{\boldsymbol{y}\}$ by $U_{\boldsymbol{y}}$. It is often useful to indicate lengths of disjoint cycles of an element $g \in S_{n}$ by symbols like $\left(s_{i, 1}\right) \cdots\left(s_{i, t_{i}}\right)$ (Chap. 3 §7.1.4).

Let $\varphi_{i}^{0}: X_{i}^{0} \rightarrow U_{\boldsymbol{y}}$ be the restriction of $\varphi_{i}$ over $U_{\boldsymbol{y}}$. Compatible with Def. 1.3, form the unramified fiber product map $\varphi_{1}^{0} \times_{U_{\boldsymbol{y}}} \varphi_{2}^{0}: X_{1}^{0} \times_{U_{\boldsymbol{y}}} X_{2}^{0} \rightarrow U_{\boldsymbol{y}}$. This may have several components, even if each of the $X_{i}^{0}$ are connected (see and Chap. 3 $\S 8.6 .1$ and $\S 5.1$ ). Thm. 7.16 uses an ordering of points above some base point $y_{0}$. With this it corresponds to components of the fiber product a pair of subgroups $H_{1}$ and $H_{2}$ of $\pi_{1}\left(U_{\boldsymbol{y}}, y_{0}\right)$. The component of the fiber product corresponds to the subgroup $H_{1} \cap H_{2}$. The maximal pointed cover of $U_{\boldsymbol{y}}$ through which both $\mathrm{pr}_{1}$ and
$\mathrm{pr}_{2}$ factor comes from the subgroup $\left\langle H_{1}, H_{2}\right\rangle=H$ generated by $H_{1}$ and $H_{2}$. Then, the monodromy group of the fiber product component defined by $\left(H_{1}, H_{2}\right)$ is the fiber product $G_{H_{1}} \times{ }_{G_{H}} G_{H_{2}}$.

Proposition 3.4. Let $\varphi_{1} \times{ }^{c} \varphi_{2}: X_{1} \times{ }_{Y}^{c} X_{2} \rightarrow Y$ be the extension of $\varphi_{1}^{0} \times_{U_{\boldsymbol{y}}} \varphi_{2}^{0}$ to the unique manifold completion of $X_{1}^{0} \times_{U_{y}} X_{2}^{0}$ given by Cor. 2.9. This satisfies the categorical fiber product in the category $\mathcal{C}_{Y}^{c}$.

Suppose $Y=\mathbb{P}_{z}^{1}$ (write $\boldsymbol{z}$ for $\boldsymbol{y}$ ), and ${ }_{1} \boldsymbol{g}$ and ${ }_{2} \boldsymbol{g}$ are respective branch cycles relative to a classical set of generating homotopy classes for $\pi_{( }\left(Y_{z}, z_{0}\right)$ and orderings of the points $X_{i, z_{0}}$ of $X_{i}$ above $z_{0}, i=1,2$. Branch cycles for $\varphi_{1} \times{ }^{c} \varphi_{2}$ are then

$$
\left(\left({ }_{1} g_{1},{ }_{2} g_{1}\right), \ldots,\left({ }_{1} g_{r},{ }_{2} g_{r}\right)\right) \in G_{H_{1}} \times{ }_{H} G_{H_{2}}
$$

given by their action on the orbit of points on the component over $z_{0}$.
Let $z_{i} \in \boldsymbol{z}$ and let $x_{1}^{\prime}$ (resp. $x_{2}^{\prime}$ ) be a point of $X_{i}$ above $z_{i}$. Assume $x_{k}^{\prime}$ corresponds to the orbit of ${ }_{k} g_{i}$ labeled by its disjoint cycle ${ }_{k} g_{i}^{\prime}$ (of length ${ }_{k} s_{i}^{\prime}$ ) in the disjoint cycle decomposition of $k g_{i}, k=1,2$. Then, points of $\varphi_{1} \times{ }^{c} \varphi_{2}: X_{1} \times_{\mathbb{P}_{z}^{1}}^{c} X_{2}$ over both $x_{1}^{\prime}$ and $x_{2}^{\prime}$ correspond one-one with orbits of $\left({ }_{1} g_{i}^{\prime},{ }_{2} g_{i}^{\prime}\right)$ on pairs of letters in the respective orbits of the cycles $1 g_{i}^{\prime}$ and ${ }_{2} g_{i}^{\prime}$.

Proof. Since $\varphi_{1} \times{ }^{c} \varphi_{2}$ is a map of compact Riemann surfaces, it is in the right category. To show it is a fiber product consider what happens if we have maps of compact Riemann surfaces $\varphi: W \rightarrow Y$, and $\psi_{i}: W \rightarrow X_{i}, i=1,2$, so that $\varphi_{i} \circ \psi_{i}=\varphi, i=1,2$. We only need show there is a unique map $\alpha: W \rightarrow X_{1} \times{ }_{Y}^{c} X_{2}$ that suits the other maps. Restrict all the existing maps and Riemann surface covers over $U_{\boldsymbol{y}}$, and use ${ }^{0}$ superscripts to indicate that. Our previous understanding of fiber product produces the corresponding $\alpha^{0}: W^{0} \rightarrow\left(X_{1} \times_{Y}^{c} X_{2}\right)^{0}$. Now apply the unique completion property of Cor. 2.9 to get $\alpha$ which then automatically has all desired properties.

Almost everything else is a restatement of previous propositions, though we comment further on the last paragraph of the statement. By relabeling the points in the fibers of $X_{i}$ over $z_{0}$, assume with no loss that ${ }_{1} g_{i}^{\prime}$ acts as $\left(a_{1} \ldots a_{e_{1}}\right)$ and ${ }_{2} g_{i}^{\prime}$ acts as $\left(b_{1} \ldots b_{e_{1}}\right)$. The final statement says that $\left({ }_{1} g_{i}^{\prime},{ }_{2} g_{i}^{\prime}\right)$ has $d=\left(e_{1}, e_{2}\right)$ orbits of length $\left[e_{1}, e_{2}\right]$ on the pairs $\left\{\left(a_{u}, b_{v}\right)\right\}_{1 \leq u \leq e_{1}, 1 \leq v \leq e_{2}}$ [11.12a].

Definition 3.5. In Prop. 3.4, $\varphi_{1} \times{ }^{c} \varphi_{2}: X_{1} \times{ }_{Y}^{c} X_{2} \rightarrow Y$ is the normal fiber product of $\varphi_{1}$ and $\varphi_{2}$.

REMARK 3.6 (Use of the word normal). In many problems the fiber product appears as an auxiliary construction. Whether the naive or normal is a better choice depends on circumstances. Usually, however, the normal is best. In our category $\mathcal{C}_{Y}^{c}$ it would appear we are stuck with considering only manifolds. For higher dimensional manifolds this result does not work, because it is possible that two manifold (ramified) covers $\varphi_{i}: X_{i} \rightarrow \mathbb{P}^{n}$, with $n \geq 2, i=1$, 2 , have no manifold fiber product. That is, there is no manifold completion of the fiber product with these properties:
(3.8a) It is the expected fiber product restricted over the unramified locus.
(3.8b) It is a finite cover of $\mathbb{P}^{n}$.

The correct extension of the Prop. 3.4 uses normal analytic sets (§8.5).
3.3.2. Geometry of the Galois closure. Consider a cover $f: Y \rightarrow X$ of degree $n=\operatorname{deg}(f)$ with an attached permutation representation $T_{f}=T: G \rightarrow S_{n}$. When $f$ is an unramified cover, Chap. $3 \S 8.3 .2$ constructs the Galois closure of this cover. We want to do the same when the cover ramifies. While the construction goes through using either the naive or normal fiber product (§3.3), we emphasize the latter. So, from this point, when we say fiber product of two covers, we are referring to the normal fiber product.

When $f$ is unramified, we took the fiber product $Y_{f}^{n} \stackrel{\text { def }}{=} Y_{X}^{n}$ of $\varphi, n$ times. Now take the normal fiber product, so the resulting set is a manifold. Then, $Y_{X}^{n}$ has components where each point has at least two of the coordinates identical. These form the fat diagonal. Remove components of this fat diagonal to give $Y^{*}$, which (exactly as in Chap. 3 Thm. 8.9) has as many components as $\left(S_{n}: G\right)$. List one of these components as $\hat{Y}$. Points in $\hat{Y}$ over the branch points no longer have the form of an $n$-tuple of points in $Y$. The stabilizer in $S_{n}$ of $\hat{Y}$ is a conjugate of $G$. Normalize by choosing $\hat{Y}$ so the stabilizer is actually $G$.

Lemma 3.7. Then, $\hat{\varphi}: \hat{Y} \rightarrow X$ is Galois with group $G$.
If $X=\mathbb{P}_{z}^{1}$ and the cover was in the Nielsen class $\mathrm{Ni}(G, \mathbf{C}, T)$, with $T: G \rightarrow S_{n}$ a faithful permutation representation, the cover $\hat{\varphi}$ has the same conjugacy classes $\mathbf{C}$, but the representation is the regular representation. The Galois cover $\hat{Y} \rightarrow Y$ has group $G(1)=G(T, 1)$ where $T$ acts on $G(1)$ cosets. The next lemma (from [Fri77, Lem. 2.1]) is just the compactified version of Chap. 3 Lem. 8.8.

Lemma 3.8. The centralizer of $G$ in $N_{S_{n}}(G, \mathbf{C})$ induces the automorphisms of $X$ that commute with $\varphi_{T}$.

Consider any permutation representation $T^{\prime}: G \rightarrow S_{n^{\prime}}$. This provides $\varphi_{T^{\prime}}:$ $X_{T^{\prime}} \rightarrow \mathbb{P}_{z}^{1} ; X_{T^{\prime}}$ is the quotient $\hat{X} / G\left(T^{\prime}, 1\right)$ (with $G\left(T^{\prime}, 1\right)$ as in $\S 3.2 .1$ ).

From Thm. 2.6 the next observations follow from the analogous statements for unramified covers in Chap. $3 \S 8.3$. A cover $(Y, \psi)$ is Galois if the order of $\operatorname{Aut}(Y, \psi)$ is $n$, as big as it can be. The construction above gives a unique minimal Galois cover $\hat{Y} \xrightarrow{\psi} X$ fitting in a commutative diagram, the Galois closure diagram


Suppose $X=\mathbb{P}_{z}^{1}$, and $\boldsymbol{g}$ is a branch cycle description of the cover with respect to canonical generators of $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. The group $\operatorname{Aut}(\hat{Y}, \hat{\psi})$, isomorphic to $G(\boldsymbol{g})$, canonically identifies with elements of $S_{\hat{n}}$ that centralize the image of $G(\boldsymbol{g})$ in its right regular representation where $\hat{n}=\operatorname{deg}(\hat{\psi})$.

For any subgroup $H$ of $\operatorname{Aut}(\hat{Y}, \hat{\psi})$ let $\bar{H}$ be the subgroup of $\pi_{1}\left(\mathbb{P}^{1} \backslash D(\psi), z_{0}\right)$ that maps onto $H$. From $\bar{H}$ we obtain a cover $\psi_{H}: Y_{H} \rightarrow \mathbb{P}^{1}$ (Chap. 3 Thm. 8.9) that fits in a commutative diagram

$$
\begin{array}{ccc}
\hat{Y} & \xrightarrow{\hat{\psi}_{H}} & Y_{H}  \tag{3.10}\\
\hat{\psi} \searrow & \varliminf_{\mathbb{P}_{z}^{1}} \psi_{H}
\end{array}
$$

where $Y \rightarrow Y_{H}$ is Galois with group isomorphic to $H$. This is a version of the classical Galois correspondence.

Corollary 3.9. Let $T_{H}$ be the coset representation of the group $G(\boldsymbol{g})$ corresponding to a subgroup $H$. Then $T_{H}(\boldsymbol{g})=\left(T_{H}\left(g_{1}\right), \ldots, T_{H}\left(g_{r}\right)\right)$ is a description of the branch cycles for the cover $\psi_{H}: Y_{H} \rightarrow \mathbb{P}^{1}$.

Proof. Let $T_{\bar{H}}$ be the coset representation of $\pi_{1}\left(\mathbb{P}^{1} \backslash D(\psi), z_{0}\right)$ corresponding to the subgroup $\bar{H}$, and let $\hat{H}$ be the kernel of the map from $\pi_{1}\left(\mathbb{P}^{1} \backslash D(\psi), z_{0}\right)$ given by $\left[\gamma_{i}\right] \rightarrow \sigma_{i}, i=1, \ldots, r$, as in Cor. 2.9. Recall that $\hat{H}$ is the maximal normal subgroup of $\pi_{1}\left(\mathbb{P}^{1} \backslash D(\psi), z_{0}\right)$ contained in $\bar{H}$, and the quotient $\bar{H} / \hat{H}$ is isomorphic to $H$. Then $\left(T_{\bar{H}}\left(\left[\gamma_{1}\right]\right), \ldots, T_{\bar{H}}\left(\left[\gamma_{r}\right]\right)\right)=T_{H}(\boldsymbol{g})$. Since the left side consists of a branch cycle description for $\left(Y_{H}, \psi_{H}\right)$, this concludes the corollary.
3.4. Riemann-Hurwitz and the genus of a cover of $\mathbb{P}_{z}^{1}$. Let $\boldsymbol{g}$ correspond to $\psi: Y \rightarrow \mathbb{P}_{z}{ }^{1}$ as in Cor. 2.9. Indicate lengths of disjoint cycles of $g_{i}$ by the symbol $\left(s_{i, 1}\right) \cdots\left(s_{i, t_{i}}\right)$ (Chap. 3 §7.1.4). Points of $Y$ corresponding to cycles of length greater than 1 are ramified points of $\psi$. The index of $g_{i}, \operatorname{ind}\left(g_{i}\right)$, is the $\operatorname{integer} \sum_{j=1}^{t_{i}}\left(s_{i, j}-1\right)=n-t_{i}$.
3.4.1. The appearance of $g_{\boldsymbol{g}}$. Consider the quantity $g_{\boldsymbol{g}}$ defined by the RiemannHurwitz formula:

$$
\begin{equation*}
2\left(n+g_{\boldsymbol{g}}-1\right)=\sum_{z_{i} \in D(\psi)} \operatorname{ind}\left(g_{i}\right) . \tag{3.11}
\end{equation*}
$$

Note!: The following lemma requires $Y$ to be connected. Chap. 3 Ex. 5.12 defines the differential $d \psi$ of the function $\psi$.

Proposition 3.10. The expression $t_{\psi}=\sum_{z_{i} \in D(\psi)} \operatorname{ind}\left(g_{i}\right)-2 n$ is even. So, $g_{\boldsymbol{g}}$ in (3.11) is an integer. Further, $t_{\psi}$ is the degree of the divisor $(d \psi)$. Finally, $t_{\psi}=t_{\boldsymbol{g}}$ depends only on $Y$, and not on $\psi$ or $n$, and $g_{\boldsymbol{g}}=\left(t_{\psi}+2\right) / 2$ is nonnegative.

Proof. The determinant of (the matrix for) $g_{i}$ is $(-1)^{\operatorname{ind}\left(g_{i}\right)}$ (Chap. 3 §7.1.4); check for each disjoint cycle. The product-one condition implies an even number of $g_{i} \mathrm{~s}$ have determinant -1 . So, for an even number of $g_{i} \mathrm{~s}$, $\operatorname{ind}\left(g_{i}\right)$ is odd. In particular, $\sum_{i=1}^{r} \operatorname{ind}\left(g_{i}\right)$ is even, and $g_{\boldsymbol{g}}$ is an integer.

Suppose $\left\{\varphi_{\alpha}, U_{\alpha}\right\}_{\alpha \in I}$ is the coordinate chart for $Y$ from $\psi$ (Def. 2.8). We may assume the local expression for $\psi$ at $y \in Y$ is $\psi \circ \varphi_{\alpha}^{-1}\left(z_{\alpha}\right)$ with $\varphi_{\alpha}(y)=0$. Then, the leading term is $a_{u} z_{\alpha}^{u}\left(a_{u} \neq 0\right)$ and the divisor of $d \psi$ at $y$ is $y^{u-1}$. For $y$ over $z$, if $z \in \mathbb{C}$, then $u=e_{y}$. If, however, $z=\infty$, then $u=-e_{y}$, and the divisor of $d \psi$ at $y$ is $-e_{y}-1$. The expression $-e_{y}-1$ summed over $y \in Y_{\infty}$ is the same as the sum over $e_{y}-1-2 e_{y}$. Since $\sum_{y \in Y_{\infty}} e_{y}=n$ (Lem. 2.1), this gives the formula.

Now use that $Y$ is connected so that $G(\boldsymbol{g})$ is transitive. We show

$$
\sum_{i=1}^{r} \operatorname{ind}\left(g_{i}\right)-2(n-1)
$$

is nonnegative. When all the $\sigma$ 's are 2-cycles the result follows if $r \geq 2(n-1)$. That is immediate from the first part of Lem. 3.11. To reduce to that case, write each $g_{i}$ as $\prod_{u=1}^{\operatorname{ind}\left(g_{i}\right)} h_{u, i}$ with each $h_{u, i}$ a 2 -cycle. With $\boldsymbol{h}_{i}=\left(h_{1, i}, \ldots, h_{\operatorname{ind}\left(g_{i}\right), i}\right)$, replace $\boldsymbol{g}$ by the juxtaposition of these $\boldsymbol{h}_{i} \mathrm{~s}$ : $\boldsymbol{h}=\left(\boldsymbol{h}_{1}, \ldots, \boldsymbol{h}_{r}\right)$. Then, $\boldsymbol{h}$ satisfies the product-one condition and $\langle\boldsymbol{h}\rangle$ is transitive. (It is $S_{n}$ : Chap. 3 [9.15e].) Further, $g_{\boldsymbol{h}}=g_{\boldsymbol{g}}$. So, the general formula for the genus of $\boldsymbol{g}$ follows from the case for 2 -cycles.

We have only to show $g_{\boldsymbol{g}}$ does not depend on $\psi$. If $\psi^{*}$, however, is another function, then $t_{\psi}$ and $t_{\psi^{*}}$ are the respective degrees of the two differentials $d \psi$ and
$d \psi^{*}$ on the compact Riemann surface $Y$. The result follows from the statement in §5.3.1 that these degrees are equal.
3.4.2. Non-negativity of $g_{\boldsymbol{g}}$. Let $\operatorname{Ni}(G, \mathbf{C}, T)$ be the Nielsen class for the group $G$ and $r$ of its conjugacy classes $\mathbf{C}$, with $T: G \rightarrow S_{n}$ faithful and transitive.

Lemma 3.11 (2-cycle Braids). For $\boldsymbol{g} \in \mathrm{Ni}(G, \mathbf{C}, T), t_{\boldsymbol{g}}=\sum_{i=1}^{r} \operatorname{ind}\left(g_{i}\right)-2 n$ is independent of the choice of $\boldsymbol{g}$. When $\boldsymbol{g}$ consists of 2-cycles in $S_{n}$ generating a transitive subgroup, $\left(t_{\boldsymbol{g}}+2\right) / 2=g_{\boldsymbol{g}} \geq 0$.

Proof. The index of an element in $S_{n}$ is independent of its conjugacy class. Since the conjugacy classes of entries of any $\boldsymbol{g} \in \mathrm{Ni}(G, \mathbf{C}, T)$ differs only by permutation from any other, the expression $t_{\boldsymbol{g}}$ is independent of the choise of $\boldsymbol{g}$.

Now apply transitivity of $G(\boldsymbol{g})$ and assume $\boldsymbol{g}$ has entries consisting of 2-cycles. There must be a series of $n-1$ entries of $\boldsymbol{g}$ so that, after the first, each consists of $\left(i_{1}, i_{2}\right)$ with $i_{1}$ in the support of the previous 2-cycles, and $i_{2}$ is not. Apply an element $Q \in B_{r}$ to $\boldsymbol{g}$ to braid these so the $n-1$ entries just chosen come together as the first $n-1$ of the 2 -cycles (for help, see [11.8]). Then, the product of the first $n-12$-cycles is an $n$-cycle.

Now we use the product-one condition: $\prod_{i=1}^{n-1} g_{i} \prod_{i=n}^{r} g_{i}=1$. Since $\prod_{i=1}^{n-1} g_{i}$ is an $n$-cycle, that implies $\prod_{i=n}^{r} g_{i}$ is also. Therefore $\left\langle g_{i}, i \geq n-1\right\rangle$ is also transitive. Now apply the previous argument to $\left(g_{n}, \ldots, g_{r}\right)$ to conclude there are at least $n-1$ of them, giving a total of at least $2(n-1)$. This concludes the proof.

Definition 3.12 (The genus). Prop. 3.10 defines the genus $g_{\psi}$ of a compact Riemann surface $Y$ presented as a cover $\psi: Y \rightarrow \mathbb{P}_{z}^{1}$.

Other books on Riemann surfaces give examples of computing $g_{Y}$ from (3.11). Rarely, however, do they discuss a branch cycle description of $\psi$ and such examples are usually abelian covers from branches of logs (Thm. 8.8 as in Prop. 2.11).

A topologist might say they have an easier proof of the Riemann-Hurwitz formula. That suggested proof is likely dependent on having a triangulation of $Y$. The formula then interprets as expressing the Euler characteristic of $Y$ (see Rem. 2.19). There are many ways to prove this formula. No matter what the proof, interpreting the integer $g_{\boldsymbol{g}}$, the genus of $Y$, is the key point.
3.5. Hurwitz spaces; inner s-equivalence and conjugacy classes. For each s-equivalence we must consider sets of corresponding covers.
3.5.1. Notation for Hurwitz spaces. Suppose $\operatorname{Ni}(G, \mathbf{C})$ is a Nielsen class with $r$ conjugacy classes. Then, any cover in the Nielsen class has an attached set $\boldsymbol{z}$ of $r$ distinct branch points. Label the space of these unordered branch points as $U_{r}$. $\S 4.2 .1$ identifies $U_{r}$ with $\left(\mathbb{P}_{z}^{1}\right)^{r} \backslash \Delta_{r} / S_{r}$. For each s-equivalence, the classes of covers with a given $\boldsymbol{z}$ as branch point set is the same as the number of s-equivalence classes in the Nielsen class.

Label the collection of equivalence classes of covers in a given s-equivalence class by using the notation $\mathcal{H}$, denoting a Hurwitz space, usually with extra decoration to indicate the type of s-equivalence classes.

There are $\mid \mathrm{Ni}(G, \mathbf{C}, T)$ abs $\mid$ absolute s-equivalence classes of covers with branch points $\boldsymbol{z} \in U_{r}$ with the data $(G, \mathbf{C}, T)$ attached to them. Prop. 2.18 shows it requires a choice of classical generators (or cuts) canonically corresond these two sets. Denote the set of classes of $\operatorname{Ni}(G, \mathbf{C}, T)$ abs covers by $\mathcal{H}(G, \mathbf{C}, T)^{\text {abs }}$. A point $\boldsymbol{p} \in \mathcal{H}(G, \mathbf{C}, T)^{\text {abs }}$ has as a representative a cover $\varphi_{\boldsymbol{p}}: X_{\boldsymbol{p}} \rightarrow \mathbb{P}_{z}^{1}$.

Inner s-equivalence of covers, corresponds exactly to (3.6b). The following pairs correspond to a point $\boldsymbol{p} \in \mathcal{H}_{G} \stackrel{\text { def }}{=} \mathcal{H}(G, \mathbf{C})^{\text {in }}$ :

$$
\begin{equation*}
\left(\hat{\varphi}: \hat{X} \rightarrow \mathbb{P}_{z}^{1}, G\left(\hat{X} / \mathbb{P}_{z}^{1}\right) \xrightarrow{\alpha} G\right) \tag{3.12}
\end{equation*}
$$

A given such pair is quivalent to $\left(\hat{\varphi}^{\prime}: \hat{X}^{\prime} \rightarrow \mathbb{P}_{z}^{1}, G\left(\hat{X}^{\prime} / \mathbb{P}_{z}^{1}\right) \xrightarrow{\alpha^{\prime}} G\right)$ if
(3.13) $\hat{\psi}: \hat{X} \rightarrow \hat{X}^{\prime}$ with $\hat{\varphi}^{\prime} \circ \hat{\psi}=\hat{\varphi}$ induces $\alpha^{\prime}$.

For example: Suppose $g \in G$ maps $\hat{X} \rightarrow \hat{X}$, changing $\alpha$ by conjugation by $g$. Then, composing $\alpha$ with $g$ gives a cover inner equivalent to (3.12). On the other hand, composing $\alpha$ with an outer automorphism of $G$ gives a new equivalence class.

Proposition 3.13. Given a (faithful) permutation representation $T: G \rightarrow S_{n}$, there is a natural map $\Psi_{\text {in,abs }}: \mathcal{H}(G, \mathbf{C})^{\text {in }} \rightarrow \mathcal{H}(G, \mathbf{C}, T)^{\text {abs }}$ by

$$
\left(\hat{\varphi}: \hat{X} \rightarrow \mathbb{P}_{z}^{1}, G\left(\hat{X} / \mathbb{P}_{z}^{1}\right) \xrightarrow{\alpha} G\right) \mapsto \varphi: \hat{X} / \alpha^{-1}(G(T, 1)) \rightarrow \mathbb{P}_{z}^{1}
$$

This map is $\left.\mid N_{S_{n}}(G, \mathbf{C}) / G\right) \mid$ to 1 over every point of $\mathcal{H}(G, \mathbf{C}, T)^{\text {abs }}$.
Definition 3.14. An element $\boldsymbol{g} \in \mathrm{Ni}(G, \mathbf{C})$ is a H-M (Harbater-Mumford) representative if $r$ is even and $\boldsymbol{g}=\left(g_{1}, g_{1}^{-1}, \ldots, g_{r / 2}, g_{r / 2}^{-1}\right)$.

Example 3.15 (Comparing $H_{r}$ inner and absolute orbits). There is a general problem that arises when applying prop. 3.13. Suppose $\boldsymbol{g}_{1}$ and $\boldsymbol{g}_{2}$ represent two distinct elements of $\mathcal{H}(G, \mathbf{C})^{\text {in }}$ that lie over the same element of $\mathcal{H}(G, \mathbf{C}, T)^{\text {abs }}$. When is there a $q \in H_{r}$ that takes $\boldsymbol{g}_{1}$ to $\boldsymbol{g}_{2}$ ?

With $T_{n}$ the standard representation of $S_{n}$, each element of $\mathcal{H}\left(A_{n}, \mathbf{C}, T_{n}\right)^{\text {abs }}$ has exactly two from $\mathcal{H}\left(A_{n}, \mathbf{C}\right)^{\text {in }}$ above it. Suppose $\boldsymbol{g}$ is an H-M rep. from $g_{1}, \ldots, g_{r / 2}$ where there is $\alpha \in S_{n} \backslash A_{n}$ such that $\alpha g_{i} \alpha^{-1}=G_{i}^{-1}, i=1, \ldots, r / 2$. Then, $(\boldsymbol{g}) q=\alpha \boldsymbol{g} \alpha^{-1}$ with $q=q_{1} q_{3} \cdots q_{r-1} \in H_{r}$. That solves relating the inner and absolute $H_{r}$ orbits in this case.
3.5.2. Conjugacy classes and multiplier groups. Many times one conjugacy class will appear several times in $\mathbf{C}$. It is easy to label conjugacy classes in $S_{n}$. One tricky event is when several entries of $\mathbf{C}$ are distinct conjugacy classes in $G(\boldsymbol{g}) \leq S_{n}$, but generate the same conjugacy class in $S_{n}$. We give easy examples here.

Suppose C is a conjugacy class in a group $G$ consisting of elements having order $m$. Then, for $k \in(\mathbb{Z} / m)^{*}$ denote the $k$ th powers of elements of C by $\mathrm{C}^{k}$. For a collection $\mathbf{C}$ of conjugacy classes use the notation $\mathbf{C}^{k}, k \in \hat{\mathbb{Z}}^{*}$ (integers relatively prime to the order of elements in $\mathbf{C}$ ).

Definition 3.16. Call $\mathbf{C}$, conjugacy classes in $G$, a rational union if $\mathbf{C}^{k}=\mathbf{C}$ (both sides counted with multiplicity) for all $k \in \hat{\mathbb{Z}}^{*}$. There is always a natural rationalization $\mathbf{C}^{\prime}$ of $\mathbf{C}$ : The minimal rational collection of conjugacy classes containing $\mathbf{C}$.

Let $T_{n}$ be the standard representation of $S_{n}, n \geq 3$. As in Chap. $3 \S 7.1 .4$, indicate conjugacy classes in $S_{n}$ with a simple notation. Give C $C_{i}$ by its cycle type: $\left(s_{i, 1}\right) \cdots\left(s_{i, t_{i}}\right), i=1, \ldots, r$. As $\sum_{j+1}^{t_{i}} s_{i, j}=n$, it is often (not always) convenient to order the $s_{i, j}$ by size: $s_{i, j} \leq s_{i, j+1}$. Recall: This class is in $A_{n}$ if and only if $n-t_{i}$ (its index, §3.4) is even.

For any homomorphism $\psi: H \rightarrow G$ (containment of $H$ in $G$ is the standard case) a conjugacy class C in $H$ generates a conjugacy class $\mathrm{C}_{G}$ in $G$ : For $g \in \mathrm{C}$, $\mathrm{C}_{G}$ is the collection of conjugates of $g$ in $G$.

Definition 3.17 (Multiplier group). Let C be a conjugacy class C in $G$ whose elements have order $m$. The multipier group of C is $M_{\mathrm{C}} \stackrel{\text { def }}{=}\left\{k \in(\mathbb{Z} / m)^{*} \mid \mathrm{C}^{k}=\mathrm{C}\right\}$. The multiplier field $K_{\mathrm{C}}$ is the fixed field in $\mathbb{Q}\left(e^{2 \pi i / k}\right)$ of $M_{\mathrm{C}}$.
3.5.3. Multiplier groups and fields in $A_{n}$. Each conjugacy class in $S_{n}$ is rational. It is more complicated for $A_{n}$. The following results give valuable examples.

Lemma 3.18. For a conjugacy class C in $A_{n}$, there are two possibilities for $\mathrm{C}_{S_{n}}=\left(s_{1}\right) \cdots\left(s_{t}\right): \mathrm{C}_{G}=\mathrm{C}$, or $\mathrm{C}_{G}=\mathrm{C} \dot{\mathrm{U}} \mathrm{h} \mathrm{C} h$ with $h=(12)$. The former happens if and only if there is an even length cycle or a product of an odd number of disjoint 2-cycles that centralizes any $g \in \mathrm{C}$. The latter happens if and only if
(3.14) all the $s_{j} s$ are odd, $j=1, \ldots, t$, and distinct.

Proof. Suppose $h$ is either an $m$-cycle with $m$ even or it is product of $m$ disjoint 2-cycles with $m$ odd. Then $S_{n}=A_{n} \dot{\cup} h A_{n}$. If $h$ centralizes $g \in \mathrm{C}$, then the orbit of $h A_{n}$ on $g$ is the same as that of $A_{n}$ and $\mathrm{C}_{S_{n}}=\mathrm{C}$.

Conversely, by the class equation if $\mathrm{C}_{S_{n}}$ is larger than C , some nontrivial element of $S_{n} \backslash A_{n}$ centralizes $g$. Suppose $m$ is the length of a disjoint cycle in $g$ and there are $t_{m}$ of these. Denote by $g_{m}$ the product of all these disjoint $m$-cycles in $g$. Write $g$ as the product of these $g_{m} \mathrm{~s}$ running over all distinct integers $m$. Denote the centralizer of $\left(1 m+1 \ldots\left(t_{m}-1\right) m+1\right) \ldots\left(m 2 m \ldots t_{m} m\right)$ by $C_{m}$. Then, the centralizer of $g$ is isomorphic to the direct product of the $C_{m} \mathrm{~s}$.

Now we check that the group $C_{m}$ is the wreath product

$$
\mathbb{Z} / m \imath S_{t_{m}}=(\mathbb{Z} / m)^{t_{m}} \times{ }^{s} S_{t_{m}}(\text { Chap. } 3 \S 8.4)
$$

regarded as a subgroup of $S_{m t_{m}}$. The copy of $(\mathbb{Z} / m)^{t_{m}}$ identifies with products of powers of the disjoint cycles in $g_{m}$. A $\pi \in S_{t_{m}}$ maps $\left(i_{1}, \ldots, i_{t_{m}}\right) \in(\mathbb{Z} / n)^{t_{m}}$ to $\left(i_{(1) \pi}, \ldots, i_{\left(t_{m}\right) \pi}\right)$. Example: $\pi=(12)$ acts in $S_{m t_{m}}$ as $(1 m+1)(2 m+2) \cdots(m 2 m)$, a product of $m$ disjoint 2 -cycles. If $m$ is even then $C_{m}$ contains an $m$-cycle, that is not in $A_{m t_{m}}$. If $m$ is odd, but larger than 1, a 2-cycle $\pi \in S_{t_{m}}$ acts as a product of $m$ disjoint 2-cycles in $A_{m t_{m}}$. So, $C_{m}$ is in $A_{m t_{m}}$ if and only $t_{m}$ is 1 and $m$ is odd. That concludes the proof.

Assume $g \in \mathrm{C}$ with $\mathrm{C}_{S_{n}}=\left(s_{1}\right) \cdots\left(s_{t}\right)$ satisfies (3.14), the only possible nonrational conjugacy classes in $A_{n}$. The next proposition checks which of those are rational when $\mathrm{C}=(n)$ ( $n$ is odd); [11.18b] outlines the general case [Fri95b, p. 332].

Recall: $p^{u}$ exactly divides $n$ (written $p^{u} \| n$ ) if $p^{u}$ divides $n$, but $p^{u+1}$ does not. Also, use Euler's Theorem that if $p$ is an odd prime, the invertible integers $\left(\mathbb{Z} / p^{u}\right)^{*}$ (of $\mathbb{Z} / p^{u}$ ) is a cyclic group.

Proposition 3.19 (Irrational Cycles). Consider the case $n>4$ is odd and $g \in \mathrm{C}$ with $C_{S_{n}}=(n)$. Suppose $n$ is not a square. Let $J$ be those primes $p$ that exactly divide $n$ to an odd power $p^{u(p)}$. For any $p \in J$, let $k \in(\mathbb{Z} / n)^{*}$ have these properties: its image in $\left(\mathbb{Z} / p^{u(p)}\right)^{*}$ generates this cyclic group; and its image in $\left(\mathbb{Z} / p^{\prime u^{\prime}}\right)^{*}$ is 1 for primes $p^{\prime} \neq p$ that divide $n$. Then, $g^{k}$ and $g$ are not conjugate in $A_{n}: \mathrm{C}$ is not a rational conjugacy class.

Denote $\sqrt{\prod_{p \in J}(-1)^{(p-1) / 2} p}$ by $\alpha_{n}$. For all odd $n, K_{\mathrm{C}}=\mathbb{Q}\left(\alpha_{n}\right)$.
Conversely, if $n$ is an odd square, $g^{k}$ is conjugate to $g$ in $A_{n}$ for all $k \in(\mathbb{Z} / n)^{*}$ : C is a rational conjugacy class.

Proof. Suppose $n$ is not a square. With $k$ (and $p \in J$ ) as in the statement, we show $g^{k}$ and $g$ aren't conjugate in $A_{n}$. With no loss, $g=(1 \ldots n)$. So $g^{k}$ maps
$i \mapsto i+k \bmod n, i=1, \ldots, n$. Multiplication by $k$ gives a permutation $\tau_{k}$ of the integers modulo $n$. Then, $\tau_{k}^{-1} g \tau_{k}$ equals $g^{k}$ :

$$
(k i) \tau_{k}^{-1} g \tau_{k}=(i) g \tau_{k}=(i+1) \tau_{k}=k i+k
$$

We characterize those $k$ with $\tau_{k}$ not in $A_{n}$. Apply the Chinese remainder theorem to write $(\mathbb{Z} / n)^{*}=\prod_{i=1}^{t}\left(\mathbb{Z} / p_{i}^{u_{i}}\right)^{*}$ with $p_{1}, \ldots, p_{t}$ distinct (odd) primes. So, it suffices to check if $\tau_{k} \in A_{n}$ for $k=\boldsymbol{k}_{i}=\left(1, \ldots, 1, k_{i}, 1, \ldots, 1\right)$; the only nonidentity entry is $k_{i}$, a generator of the cyclic group $\left(\mathbb{Z} / p_{i}^{u_{i}}\right)^{*}$, in the $i$-th position. Consider what happens with $k$ equal $\left(k_{1}, 1, \ldots, 1\right)$.

First, assume $t=1, u_{1}=u$ and $k_{1}=k$. Consider the cycle structure of $\tau_{k}$ $\mathbb{Z} / p^{u}$. Multiplication by $k$ on integers of $\mathbb{Z} / p^{u}$ exactly divisible by $p^{i}, i<u$, gives one orbit of length $p^{u-i}-p^{u-i-1}$. For each $i$ between 0 and $u-1$, this cycle has even length—not in $A_{n}$. (The orbit for $i=u$ has length 1.) Thus, the permutation is a product of $u$ elements not in $A_{n}$ (and it fixes exactly one integer). The total permutation from multiplication by $k$ is in $A_{n}$ if and only if $u$ is even.

For the general case, write $\mathbb{Z} / n$ as $\mathbb{Z} / p_{1}^{u_{1}} \times \mathbb{Z} / n^{\prime}$. Multiplication by $k$ is the identity on the second coordinate. Thus, it stabilizes each coset $\mathbb{Z} / p_{1}^{u_{1}} \times k^{\prime}$ with $k^{\prime} \in \mathbb{Z} / n^{\prime}$. In particular, $\tau_{k}$ is the product of $n^{\prime}$ elements coming from the first case above. Thus, $\tau_{k} \in A_{n}$ if and only if $u_{1}$ is even. The converse comes by noting it suffices to check the elements $\boldsymbol{k}_{i}$ above.

Finally, we identify the field $\hat{\mathbb{Q}}_{n}$. Identify the kernel of $\mu:(\mathbb{Z} / n)^{*} \rightarrow \mathbb{Z} / 2$ by $k \in(\mathbb{Z} / n)^{*}$ maps to $\tau_{k} \bmod A_{n}$. In the above notation, $\boldsymbol{k}_{i}$ goes to 1 if and only if $i \in J$. The unique quadratic extension of $\mathbb{Q}$ inside $\mathbb{Q}\left(\zeta_{p_{j}}\right)$ is $\mathbb{Q}\left(\sqrt{(-1)^{\left(p_{j}-1\right) / 2} p_{j}}\right)$. Conclude by noting the kernel of $\mu$ is of index 2 in $(\mathbb{Z} / n)^{*}$ and it fixes $\alpha_{n}$.

Example 3.20. Suppose $C_{1}, C_{2}$ and $C_{3}$ are respectively the conjugacy classes of the 5 -cycles in $A_{5}$ given by $g_{1}=(12345), g_{2}=(13524)$ and $g_{1}$ again. Then, $C_{1}, C_{2}, C_{3}$ is not a rational union because the conjugacy class of $g_{1}$ appears with multiplicity 2 , while its square appears only with multiplicity 1 . The collection $\mathbf{C}^{\prime}=\left(C_{1}, C_{2}, C_{1}, C_{2}\right)$ is its rationalization.

Example 3.21 (Rational conjugacy classes in $A_{9}$ ). The conjugacy classes of $A_{9}$ that don't remain the same in $S_{9}$ are those that become (1)(3)(5) of (9) in $S_{9}$. In general, counting the partitions of $n$ into distinct odd integers is a nontrivial combinatorial business (see [11.18d]). [A199] says the number of partitions of $n$ by odd distinct integers equals partitions of $n$ with all parts $\neq 2$, at least 6 apart and at least seven apart if both parts are even. For $n=25$ this count is

$$
12=|(25),\{(i, 25-i), 1 \leq i \leq 9, i \neq 2,(1, k, 25-k-1), 7 \leq k \leq 9\}|
$$

According to Prop. 3.19, there are two rational conjugacy classes $A_{9}$ that become (9) in $S_{9}$. From [11.18b] the two conjugacy classes C for which $\mathrm{C}_{S_{n}}=(1)(3)(5)$ are not rational and $M_{\mathrm{C}}=\mathbb{Q}(\sqrt{-3 \cdot 5})$.

## 4. Applications of the Existence Theorem

This section should surprise the reader at how simple group theory, starting with dihedral groups, reveals serious classical topics. We develop two skills.

- Creating notation for calculating collections of covers.
- Finding algebraic functions to give coordinates on such collections.

For any group $G$ denote by $\operatorname{Aut}(G)$ the full set of automorphisms of $G$, and by $\operatorname{Inn}(G)$ the autmorphisms induced by conjugation by $G$. The first nonabelian group that comes up in Galois theory is the dihedral group. Prop. 2.11 shows all abelian covers are algebraic. Covers $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ with dihedral monodromy group, even when $X$ has genus 0 , are not obviously algebraic. Part of Abel's Theorem is equivalent to displaying functions that show this. There is more to such covers than one would expect from its group theory alone.

We start slowly with dihedral covers, because there is so much history in them, especially about coordinates. $\S 4.1$ is a case that function theoretically is almost trivial, though its applications require careful coordinates.
4.1. Dihedral - a ka Tchebychev - polynomials. Suppose a degree $n$ cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ has genus $0\left(g_{X}=0\right)$ and branch cycles $\boldsymbol{g}=\left(g_{1}, \ldots, g_{r}\right)$ (relative to some choice of classical generators) with at least one totally ramified place. That means some $g_{i}$, say $g_{r}$, is an $n$-cycle in $G(\boldsymbol{g}) \leq S_{n}$. At first examples use the standard representation $T_{n}$ of $S_{n}$ restricted to $G(\boldsymbol{g})$. Apply Riemann-Hurwitz to conclude $\sum_{i=1}^{r-1} \operatorname{ind}\left(g_{i}\right)=n-1$.
4.1.1. Cyclic covers and Redei functions. An element of $S_{n}$ has index $n-1$ if and only if it is an $n$-cycle. We draw conclusions from this and the productone condition, $\Pi(\boldsymbol{g})=1$. If there is another $n$-cycle among the branch cycles, then $r=2$. By conjugating by an element of $S_{n}$ we may take $g_{1}=(1 \ldots n)$ and $g_{2}=g_{1}^{-1}$. There is unique absolute Nielsen class of genus 0 covers with at least two $n$-cycles: $\mathrm{Ni}\left(\mathbb{Z} / n, \mathbf{C}_{n, n}, T_{n}\right)^{\text {abs }}$. Further, in that class there is exactly one absolute s-equivalence class representing the Nielsen class: $\mathbf{C}$ consists of C and $\mathrm{C}^{-1}$, a conjugacy class in $\mathbb{Z} / n$ and its inverse. The case $n=2$ is trivial.

For $n \geq 3$, there are $\varphi(n) / 2$ inner Nielsen classes of such covers,

$$
\mathrm{Ni}\left(\mathbb{Z} / n,\left(\mathrm{C}^{j}, \mathrm{C}^{-j}\right)\right)^{\text {in }}, \text { with }(j, n)=1, j \leq n / 2
$$

As C contains one element, there are two inner s-equivalence class representing each Nielsen class: One with $g \in \mathrm{C}$ (resp. $g^{-1} \in \mathrm{C}^{-1}$ ) the branch cycle for $z_{1}$ (resp. $z_{2}$ ); another with the branch cycles switched.

These abelian covers we can produce by hand. Cases like this where $G$ has a nontrivial center present special problems, as we'll see later. Just consider $\mathbb{P}_{w}^{1} \rightarrow \mathbb{P}_{z}^{1}$ by $w \mapsto w^{n}$ : 0 and $\infty$ map respectively to 0 and $\infty$. Put the branch points anywhere using $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$ (say $\alpha=\frac{z-z_{1}}{z-z_{2}}$ ) that maps $z_{1}, z_{2}$ to $0, \infty$. Then, $w \mapsto \alpha^{-1}\left((\alpha(w))^{n}\right)$ gives a representing cover $\varphi_{\mathrm{C}, \mathrm{C}^{-1}, \boldsymbol{z}}: \mathbb{P}_{w}^{1} \rightarrow \mathbb{P}_{z}^{1}$ in the absolute s-equivalence class with branch points $z=\left\{z_{1}, z_{2}\right\}$. Further, it is $z_{i}$ that maps to $z_{i}$, $i=1,2$, by $\varphi_{\mathrm{C}, \mathrm{C}^{-1}, z}$. We've explicitly written a representative of very s-equivalence class of covers in the Nielsen class.
$\S 4.2 .1$ discusses r-equivalence classes. In this equivalence, all the covers $\varphi_{\mathrm{C}, \mathrm{C}^{-1}, \boldsymbol{z}}$ are equivalent. There is just one element in any Nielsen class, for we can put the branch points where we want, and switch the branch points, too. Recall: $\mathbb{P}_{z}^{1}\left(\mathbb{F}_{q}\right)$ the values on the Riemann sphere in the finite field $\mathbb{F}_{q}$ (Chap. 2 [9.19]).

Example 4.1 (Redei functions). The problem solved by Redei functions is to consider the collection of covers $\varphi_{\mathrm{C}^{-\mathrm{C}^{-1}, \boldsymbol{z}}}$ up to changing $\varphi$ to $\alpha^{-1} \varphi \circ \alpha$ with $\alpha \in \mathrm{PGL}_{2}(\mathbb{Q})$. Assume $n \geq 3$ is odd. That is we trying to describe rational requivalence repesentatives in this Nielsen class. If $\varphi$ has coefficients in $\mathbb{Q}$, then the set $\left\{z_{1}, z_{2}\right\}$ is a $\mathbb{Q}$-set (see Lem. 6.4). [LN83] discusses Redei functions in detail. They give the easiest examples of exceptional functions $f: \mathbb{P}_{w}^{1} \rightarrow \mathbb{P}_{z}^{1}$ that map
one-one when restricted to $\mathbb{P}_{w}^{1}\left(\mathbb{F}_{q}\right)$, for infinitely many prime powers $q$. They are perfect for standard cryptography applications, as are Dickson polynomials and other dihedral cover examples.

The branch points $\{0, \infty\}$ and $\left\{z_{1}, z_{2} \mid z_{1}=\sqrt{m}, z_{2}=-\sqrt{m}, m\right.$ a square-free integer represent the $\mathbb{Q}$ absolute r-equivalence classes [11.15a].
4.1.2. Twisted Chebychev - a ka Dickson - polynomials. Here are the conditions for absolute Nielsen classes of Chebychev covers $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ :
(4.1a) $X$ has genus $g_{X}=0$ and $\operatorname{deg}(\varphi)$ is an odd prime $p$;
(4.1b) $G \leq S_{p}$ is a subgroup of $\mathbb{Z} / p \times^{s}(\mathbb{Z} / p)^{*} \stackrel{\text { def }}{=} \mathbb{A}_{p}$ (acting on $\left.\mathbb{Z} / p\right)$;
(4.1c) $\mathbf{C}$ has an entry, say $\mathrm{C}_{r}$, that is a $p$-cycle; and
(4.1d) $\varphi$ is not a cyclic cover.

Tacitly the permutation representation throughout is the degree $p$ representation $T_{p}$ on $\mathbb{Z} / p$. Recall, we represent elements of $\mathbb{A}_{p}$ by $2 \times 2$ matrices $\left(\begin{array}{ll}a & 0 \\ b & 1\end{array}\right)$ with multiplication from matrix multiplication (Chap. 3 Rem. 7.4). Using §4.1.1, we have just one $p$-cycle of conjugacy classes. Elements of order $p$ are conjugate in $\mathbb{A}_{p}$ to $\left(\begin{array}{ll}1 & 0 \\ 1 & 0\end{array}\right)$. Denote this conjugacy class $\mathrm{C}_{p}$. The other conjugacy classes in $\mathbb{A}_{p}$ correspond one-one with non-identity elements of $(\mathbb{Z} / p)^{*}$. Denote the corresponding conjugacy class to $a \in(\mathbb{Z} / p)^{*}$ by $C_{a}$. For $A$ a subgroup of $(\mathbb{Z} / p)^{*}, \mathbb{Z} / p \times^{s} A$ is the corresponding subgroup of $\mathbb{A}_{p}$. Prop. 4.2 and Cor. 4.3 is from [Fri70].

Proposition 4.2. There is only one absolute Nielsen class satisfying (4.1). It is $\mathrm{Ni}(G, \mathbf{C})^{\text {abs }}$ with $\mathbf{C}=\left(\mathrm{C}_{-1}, \mathrm{C}_{-1}, \mathrm{C}_{p}\right) \stackrel{\text { def }}{=} \mathbf{C}_{(-1)^{2} \cdot p}$ and $G=\mathbb{Z} / p \times^{s}\langle-1\rangle$. Further, there is one element in this Nielsen class. More generally, for any odd $n>0$, there is a unique absolute representative in the absolute Nielsen class of $\mathrm{Ni}\left(D_{n}, \mathbf{C}_{(-1)^{2} \cdot n}\right)^{\text {abs }}$.

Proof. With no loss in an absolute Nielsen class take branch cycles so that $\boldsymbol{g}$ has $g_{r}=\left(\begin{array}{ll}1 & 0 \\ 1 & 0\end{array}\right)$. The other $g_{i} \mathrm{~s}$ are in $C_{a}, a \in(\mathbb{Z} / p)^{*} \backslash\{1\}$, which acts as multiplication by $a$ on $\mathbb{Z} / p$. If $m_{a}$ is the order of $a$, then this action has $\frac{p-1}{m_{a}}$ orbits of length $m_{a}$, and one orbit of length 1 . The index of such a $g_{i}$ is thus $\frac{p-1}{m_{a}}\left(m_{a}-1\right)$. Now apply Riemann-Hurwitz (3.11) using that $g_{X}=0: p-1=\sum_{i=1}^{r-1} \frac{p-1}{m_{a_{i}}}\left(m_{a_{i}}-1\right)$. The expression $\frac{m-1}{m}(m \geq 2)$ is at least $\frac{1}{2}$, with equality if and only if $m=2$ ( $m_{a}=-1$. Since $r-1 \geq 2$, the result is $r-1=2$, and $g_{i} \in \mathrm{C}_{-1}, i=1,2$.

Now we see there is only one element in this Nielsen class. Fix $g_{3}$, and take $g_{1}=$ $\left(\begin{array}{cc}-1 & 0 \\ b & 1\end{array}\right)$, with $g_{2}$ determined by the product-one relation: $g_{1} g_{2} g_{3}=1$. Normalize further by conjugating the 3 -tuple by a power of $g_{3}$ :

$$
\left(\begin{array}{cc}
1 & 0 \\
-k & 1
\end{array}\right)\left(\begin{array}{cc}
-1 & 0 \\
b & 1
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
k & 1
\end{array}\right)=\left(\begin{array}{cc}
-1 & 0 \\
b+2 k & 1
\end{array}\right) .
$$

So, by choosing $k$ so $b+2 k=0 \bmod p$ gives a unique representative of $\operatorname{Ni}\left(\mathbb{Z} / p \times^{s}\right.$ $\left.\langle-1\rangle, \mathbf{C}_{(-1)^{2} \cdot p}\right)^{\text {abs }}$.

Up to reduced equivalence, we may place the three branch points $\boldsymbol{z}=\left\{z_{1}, z_{2}, z_{3}\right\}$ by whatever three points we want. Given $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ (unique up to equivalence)
in this equivalence class, we find a polynomial $T_{p}(w)$ with branch points $-2,+2, \infty$ reduced equivalent to it. Further, from the branch cycle description there is exactly one unramified point of $X$ over each of -2 and +2 (use the corresponding between points over branch points and disjoint cycles of the branch cycles). So, by ordinary equivalence, put these at $-1,+1, \infty$ respectively. This is a less trivial case than previously for producing a function on the covers to show they are algebraic.

## Corollary 4.3.

Proof. There is one element in the absolute s-equivalence classes of polynomials with dihedral group cover. Suppose $f$ is a monic degree $n$ polynomial over $\bar{F}$ that gives a branched cover $\mathbb{P}_{T}^{1} \rightarrow \mathbb{P}_{z}^{1}$ with two finite branch points $z_{1}, z_{2} \in \bar{F}$, both ramified of order 2. The following observations occur in [Fri70]. The geometric Galois group of the Galois closure is a dihedral group. If $n$ is odd, then the Nielsen class of the cover is $\mathrm{Ni}\left(D_{n}, \mathbf{C}_{n \cdot 2^{\frac{n-1}{2}} \cdot 2^{\frac{n-1}{2}}}\right)$. Further, since the normalizer of $D_{n}$ in $S_{n}$ has no center, any cover with branch points $\left\{\infty, z_{1}, z_{2}\right\}$ in this Nielsen class is determined up to a unique isomorphism. So, if the unordered branch points are defined over $F$, then the cover is represented by a unique polynomial over $F$. As $z_{1}+z_{2}$ are defined over $F$, changing $z$ to $z-\left(\frac{z_{1}+z_{2}}{2}\right)$ normalizes further to assume the branch points sum to 0 . Call these normalized Chebychev polynomials. From these observations the following are clear. For any $d \in F^{*}$, and odd positive integer $n$ define the Dickson Polynomial $D_{n}(a, w)$ to be $a^{n / 2} T_{n}\left(a^{-1 / 2} w\right)$. As $a$ varies we get all the normalized Chebychev polynomials. Clearly two such polynomials are isomorphic over $F$ if and only multiplication by some $b \in F$ maps the branch points of one to the other.

If, however, $n$ is even, the conjugacy classes defining the Nielsen class are distinct and the branch points all are defined over $F$. A compensating fact is that $N_{S_{n}}\left(D_{n}\right)$ has a nontrivial centralizer $Z_{S_{n}}\left(D_{n}\right)=\left\langle(1 \ldots n)^{n / 2}\right\rangle$ in $S_{n}$ (multiplication by -1 leaves $1+n / 2$ invariant modulo $n$. [Tu95] [Wel69] [LN73] [Mu80-02]
4.2. $\mathrm{PGL}_{2}(\mathbb{C})$ action, r-equivalence and hyperelliptic covers. $\S 3.5$ explains the sets of covers in $\mathcal{H}(G, \mathbf{C}, T)^{\text {abs }}$ and other s-equivlaence classes. The group $\mathrm{PGL}_{2}(\mathbb{C})$, as one-one analytic maps of $\mathbb{P}_{z}^{1}$ enters immediately to give from each s-equivalence class, a new equivalence ( r (educed)-equivalence) from it.
4.2.1. $\mathbb{P}^{r}$ as $\left(\mathbb{P}_{z}^{1}\right)^{r} / S_{r}$ and $r$-equivalence. Identify the elements of $\mathbb{P}^{r}$ (Chap. 3 $\S 4.3)$ as nonzero monic polynomials in a variable $z$ of degree at most $r$. For example, if $\left(a_{0}, a_{1}, \ldots, a_{r}\right)$ represents a point of $\mathbb{P}^{r}$, and $z_{0} \neq 0$, by scaling it by $\frac{1}{z_{0}}$ assume with no $\operatorname{loss} z_{0}=1$. Then, take the polynomial associated to this point as $z^{r}+$ $\sum_{i=0}^{r-1}(-i)^{r-i} a_{r-i} z^{i}$. There is a natural permutation action of $\pi \in S_{r}$ on the entries of $\left(\mathbb{P}_{z}^{1}\right)^{r}: \pi:\left(z_{1}, \ldots, z_{r}\right) \mapsto\left(z_{((1) \pi}, \ldots, z_{(r) \pi}\right)$. Denote the set of distinct $r$-tuples of elements of $\left(\mathbb{P}_{z}^{1}\right)^{r}$ by $U^{r}=\left(\mathbb{P}_{z}^{1}\right)^{r} \backslash \Delta_{r}$. Call $\Delta_{r}$ the fat diagonal: The locus were two or more equal entries.

Proposition 4.4. Represent the natural quotient map

$$
\Psi_{r}:\left(z_{1}, \ldots, z_{r}\right) \in\left(\mathbb{P}_{z}^{1}\right)^{r} \mapsto\{z\} \in\left(\mathbb{P}_{z}^{1}\right)^{r} / S_{r}
$$

by sending $\left(z_{1}, \ldots, z_{r}\right)$ to the polynomial $\prod_{i=1}^{r}\left(z-z_{i}\right)$ in $z$ : If $z_{i}=\infty$, replace $\left(z-z_{i}\right)$ by 1. This canonically identifies $\Psi_{r}$ with degree $n$ ! analytic map of complex manifolds $\left(\mathbb{P}_{z}^{1}\right)^{r} \rightarrow \mathbb{P}^{r}$ [11.14a]. Identify unordered sets of $r$ branch distinct points
as an affine subspace $U_{r}$ of $\mathbb{P}^{r}$; the complement of the classical discriminant locus $D_{r}$ identified with the image of $\Delta_{r}$ [11.14b].

If $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ represents an s-equivalence class of covers in a given Nielsen class Ni , then the collection $\left\{\alpha \circ \varphi: X \rightarrow \mathbb{P}_{z}^{1}\right\}_{\alpha \in \mathrm{PGL}_{2}(\mathbb{C})}$ gives the set of covers r-equivalent to $\varphi$. To any cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$, associate the unordered collection of its branch points $z \in \mathbb{P}^{r}$. This branch point map produces a map we will never lose sight of in the rest of this book.

Suppose we have a group $G$, conjugacy classes $\mathbf{C}$ in $G$, a permutation representation $T: G \rightarrow S_{n}$ and $G \leq N \leq N_{S_{n}}(\mathbf{C})$. For covers in the set of r-equivalence classes use the notation $\mathcal{H}(G, \mathbf{C}, T) / N^{\text {rd }}$. We have a special notation $\mathcal{H}^{\text {abs,rd }}$ and $\mathcal{H}^{\text {in,rd }}$ for the associated reduced absolute and inner equivalence classes of covers.

The action of $\alpha \in \mathrm{PGL}_{2}(\mathbb{C})$ on $\left(\mathbb{P}_{z}^{1}\right)^{r}$ by $\left(z_{1}, \ldots, z_{r}\right) \mapsto\left(\alpha\left(z_{1}\right), \ldots, \alpha\left(z_{r}\right)\right)$ maps $U^{r}$ into itself.

Proposition 4.5. The actions of $\mathrm{PGL}_{2}(\mathbb{C})$ and of $S_{r}$ on $U^{r}$ commute. This gives a complex analytic map $\Psi_{r}^{\text {rd }}:\left(\mathbb{P}_{r}^{1}\right)^{r} \rightarrow \mathrm{PGL}_{2}(\mathbb{C}) \backslash\left(\mathbb{P}_{r}^{1}\right)^{r} / S_{r} \stackrel{\text { def }}{=} J_{r}$ factoring through the space $\mathrm{PGL}_{2}(\mathbb{C}) \backslash\left(\mathbb{P}_{r}^{1}\right)^{r} \stackrel{\text { def }}{=} \Lambda_{r}$. For all $r$ the spaces $\Lambda_{r}$ and $J_{r}$ are normal affine varieties, though for $r \geq 5$, neither is a manifold.

Then, $\Psi_{r}^{\mathrm{rd}}$ induces a natural map of any Hurwitz space $\mathcal{H}(G, \mathbf{C}, T) / N^{\mathrm{rd}}$ to $J_{r}$.
Proof. [11.14c]
Refer to the induced map $\Lambda_{r} \rightarrow J_{r}$ also as $\Psi_{r}^{\text {rd }}$ when that causes no confusion.
Corollary 4.6. The space $J_{4}$ (resp. $\Lambda_{4}$ ) naturally identifies with $\mathbb{P}_{j}^{1} \backslash\{\infty\}$ (resp. $\mathbb{P}_{\lambda}^{1} \backslash\{0,1 \infty\}$ ) and $\Lambda_{4} \rightarrow J_{4}$ compactifies to a Galois covering map with group $S_{3}$, ramified over $j=0, j=1$ and $j=\infty$ with branch cycles identified with ( $(135)(246),(12)(34)(56), R E T U R N$.
4.2.2. Hyperelliptic covers. For $G \leq S_{n}$ denote its intersection with $A_{n}$ by ${ }^{+} G$ : Indicating the elements of positive sign in this representation. For any degree $n$ cover $\varphi: X \rightarrow Y$, its monodromy group $G_{\varphi}$ is a subgroup of $S_{n}$. Similarly, for any finite group $G$ consider the collection of faithful transitive permutation representations (up to permutation equivalence, ${ }^{+} P_{G}$ that give an embedding of $G$ in an alternating group. In that case

Lemma 4.7. If $G$ has no normal subgroup of index 2, then ${ }^{+} P_{G}$ consists of all faithful permutation representations of $G$. This holds if $G$ is generated by elements of odd order, or if $G$ is 2-perfect.

Conversely, suppose $H \triangleleft G$ and has index 2. Then, RETURN
Proof. ${ }_{T}^{+} G=G$ for each $T \in{ }^{+} P_{G}$ if and only if consists of faithful permutation representations.

Example 4.8 ( $\mathrm{H}-\mathrm{M}$ reps. and r-equivalence). Recall the definition of $\mathrm{H}-\mathrm{M}$ reps. from Def. 3.14. In the Nielsen class $\operatorname{Ni}(G, \mathbf{C}, T)^{\text {abs }}$, consider an H-M rep. Suppose $r=4$. Reduced equivalence when $\mathrm{r}=4$ comes about because of the linear fractional transformations that flip any two pairs of branch points. In particular, the setup for $q_{1} q_{3}^{-1}$ action on an H-M rep. takes the representative $\boldsymbol{g}=\left(g_{1}, g_{1}^{-1}, g_{2}, g_{2}^{-1}\right)$ to $\left(g_{1}^{-1}, g_{1}, g_{2}^{-1}, g_{2}\right)$. The action of this element fixes the absolute class $\boldsymbol{g}$ if there is some element $g \in N_{S_{n}}(\mathbf{C})$ that conjugates the first 4 -tuple to the second. (If we were doing inner r-equivalence, $g$ would be in $G$.)

In all cases you get the following conclusion about any cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ (in an H-M rep. orbit) if this happens. There is another map $\varphi^{\prime}: X \rightarrow \mathbb{P}_{z}^{1}$ and $\alpha: \mathbb{P}_{z}^{1} \rightarrow \mathbb{P}_{z}^{1}$ having order 2 , so $\alpha \circ \varphi^{\prime}=\varphi$. Take $\mathbb{P}_{w}^{1}$ tobethequotientof $\mathbb{P}_{z}^{1}$ by $\alpha$ : Giving a map $\mu: \mathbb{P}_{z}^{1} \rightarrow \mathbb{P}_{w}^{1}$ branched at two points. Example case: the cover is cyclic of odd prime degree $p$. The degree 2 cover, however, is not a Galois cover with group $D_{p}$.
4.3. Involution dihedral covers. Last week we produced modular curves from dihedral involution-covers. Today (Wed. Mar. 19, at 3PM) I'll generalize this to a much larger class of j-line covers that we may compare with modular curves. This retains observations from Riemann's generalization of Abel's Theorem.

We do this by reflecting on application valuable properties of ALL finite groups. The analog: Dihedral groups are to modular curves as general p-perfect groups are to generalizing j-line covers. Modular curves come in series related to a prime p. The analog says each finite group comes with a series related to any prime $p$ dividing its order.

This approach to the regular version of the Inverse Galois problem has a valuable structure.

1. It includes the most famous theorems in diophantine geometry as special cases.
2. It exposes difficult modular representation problems beyond what group theorists classically study. (John Thompson noted some are analogs to such topics as the Golod-Shafarevich class field tower.)
3. Modular representation theory interprets simple properties of the Modular Tower levels.

Once we get by the initial definitions (like the mapping class group), the relation to group theory comes clear. We saw last time that modular curves are algebraic precisely because of the relation between j-invariants that comes from the dihedral involution realizations. So, parameter spaces for dihedral involutions covers are algebraic using specific coordinates from Abel's Theorem. The next step is developing analogs of this to Modular Towers through theta nulls. I'll conclude today with the first topic in that direction - half-canonical classes.
2. Finding all abelian covers of a compact surface is equivalent to finding all functions on the surfaces. How is it tacit in this description to know the surface is algebraic? Example: From the fundamental group, you know about abelian covers. Yet, ... The distinction between hyperelliptic and general: You know a hyperelliptic surface is algebraic.

## 5. Braid orbits

Take $r=4$ and $G=S_{5}$. Let $\mathrm{C}_{1}$ and $\mathrm{C}_{3}$ be the conjugacy classes of 2-cycles in $S_{5}, \mathrm{C}_{2}$ the conjugacy class of a 3-cycle and $\mathrm{C}_{4}$ the conjugacy class of a 5-cycle. Consider the Nielsen class $\operatorname{Ni}\left(S_{5}, \mathbf{C}\right) / S_{5}=\mathrm{Ni}^{+}$:

$$
\left\{\boldsymbol{g}=\left(g_{1}, \ldots, g_{4}\right) \mid g_{1} g_{2} g_{3} g_{4}=1,\langle\boldsymbol{g}\rangle=S_{5} \text { and } \boldsymbol{g} \in \mathbf{C}\right\} / S_{5}
$$

(5.1a) How many elements are in $\mathrm{Ni}^{+}$?
(5.1b) Let $\psi: \pi_{1}\left(U_{\boldsymbol{z}}\right) \rightarrow S_{5}$ map a fixed set $\bar{g}_{1}, \ldots, \bar{g}_{4}$ into some element of $\mathrm{Ni}^{+}$. Why is the cover corresponding to such a homomorphism a genus 0 compact Riemamn surface minus a finite set of points?
(5.1c) Represent $S_{5}$ on the 10 unordered distinct pairs of integers from $\{1, \ldots, 5\}$ : $T: S_{5} \rightarrow S_{10}$. Example: (12345) has two orbits on these 10 pairs. What are the lengths of the disjoint cycles of $T$ applied to an element of the conjugacy class of a 3 -cycle in $S_{5}$ ?
(5.1d) Compose $\psi$ with $T$ to get $T \circ \psi=\psi^{\prime}: \pi_{1}\left(U_{\boldsymbol{z}}\right)$. What is the genus of the curve at the top of the corresponding cover $X=X_{\psi} \rightarrow \mathbb{P}_{z}^{1}$ ?
(5.1e) Does the isomorphism class of $X_{\psi}$ depend on $\psi$ (assuming $\psi$ is in the Nielsen class $\left.\mathrm{Ni}^{+}\right)$?
5.0.1. Genus of the corresponding degree 10 covers. Let $\boldsymbol{g}$ be a branch cycle description of the cover from $\mathrm{Ni}^{+}$in [11.20]. Compute the genus $g$ of ${ }_{+} \mathcal{T}_{p}^{(2)}$ from Riemann-Hurwitz:

$$
\begin{equation*}
2(10+g-1)=\sum_{i=1}^{4} \operatorname{ind}\left(R\left(g_{i}\right)\right) \tag{5.2}
\end{equation*}
$$

Suppose $g_{1}$ and $g_{3}$ are 2 -cycles from $S_{5}$. Then, $R\left(g_{i}\right)$ has shape $(2)(2)(2)$ in the representation $R, i=1,3$. Similarly, if $g_{2}$ is a 3 -cycle, $R\left(g_{3}\right)$ has shape (3)(3)(3). Finally, $R\left(g_{4}\right)$ has shape (5)(5). Thus, the total contribution to the right side of (11.21) is $2 \cdot 3+6+2 \cdot 4=20$ and $g=1$.

Next: Compute $\mathrm{Ni}^{+}$modulo conjugation by $S_{5}$. Choose $S_{5}$ representatives with $g_{4}$ equal $g_{\infty}=(12345)^{-1}$. Divide $\mathrm{Ni}^{+}$into two sets $T_{1}$ and $T_{2}: \boldsymbol{g} \in T_{1}$ has $g_{1}$ and $g_{2}$ with no integers of common support, and $\boldsymbol{g} \in T_{2}$ has $g_{1}$ and $g_{2}$ with one integer of common support. Conjugate by a power of $\boldsymbol{g}_{\infty}$ to assure elements of $T_{1}$ have $g_{1}=(1 j)$ with $j=2$ or 3 . Similarly, elements in $T_{2}$ have 1 as common support of $g_{1}$ and $g_{2}$. From this, list $\mathrm{Ni}^{+, \text {abs }}$.
5.0.2. Covers with group $A_{5}$. (3)(3)(3)(5): Suppose $g_{3}=(123)$.
(5.3a) Ramification: $g_{1} g_{2}$ is $(2)(2)$, assume missing integer is 1 , so to get product a 5 -cycle: may assume $g_{1} g_{2}$ is $(25)(34)$. Now everthing is fixed and need only count number of ways to write $g_{1} g_{2}$ is a product of two three cycles. Hint: Products of two 3 -cycles giving (25)(34): You get one element from $(425)(234)$. Now conjugate the pair $((425),(234))$ by the centralizer of $(25)(34)$, the group $\langle(25)(34),(24)(35)\rangle$.
(5.3b) If $g_{1} g_{2}$ is (3), then conjugate by $\left\langle g_{3}\right\rangle$ to assume common integer is 1 , and $g_{1} g_{2}=(145)$. Hint: Take $\left(g_{1}, g_{2}\right)=((143),(135)$, and then conjugate by $\langle(23),(145)\rangle$.
(5.3c) If $g_{1} g_{2}$ is (5). Then, product can't be of type (2)(3) (Riemann-Hurwitz), and have only to assure the (5) times $g_{3}$ doesn't fix anything. That means can't have $2 \mapsto 1,3 \mapsto 2$ or $1 \mapsto 3$. Also, since by conjugation by $\langle(45),(123)\rangle$ can assume (15???) resulting in (15243) or (15324). Hint: For each of (15243) or (15324), we need to count all the ways to write this 5 -cycle as a product of two 3 -cycles. For (12345), assume the integer 1 is the common integer to the 3 -cycles. So, $\left(g_{1}, g_{2}\right)=((123),(145))$. Then, by conjugating by $\langle(12345)\rangle$, gives the five cases where $g_{1}$ and $g_{2}$ have any desired integer in common.
(5.3d) Up to equivalence, there are exactly 4 covers from a), 6 covers from b) and 10 covers from c), or 20 total covers. Also, by applying powers of $q_{1}$ to case c) you get 10 total in two orbits of length five. Same for b), two orbits of length 3 , and for a), two orbits of length two.
5.0.3. Non-rigid $A_{n}$ covers. Consider $n \geq 5$, odd and squarefree. Let $\mathbf{C}$ be conjugacy classes of $\left(g_{1}, g_{2}, g_{3}\right) \in A_{n}^{3}$ with $g_{1}=(12)(34), g_{2}=(13567 \ldots n)$ and $g_{3}=(12 \ldots n)^{-1}$. Check: Geometric monodromy is $A_{n}$. Representatives for conjugation of $S_{n}$ on $\mathrm{Ni}\left(S_{n}, \mathbf{C}\right) / S_{n}, i=3, \ldots,(n+1) / 2$ :

$$
\boldsymbol{g}_{j}^{\prime}=\left((12)(j j+1),(13 \ldots j j+2 j+3 \ldots n), g_{3}^{\prime}\right)
$$

## RETURN

Question 5.1. Exists $f: \mathbb{P}_{y}^{1} \rightarrow \mathbb{P}_{z}^{1}$ in $\mathbb{Q}[y]$ ?
If yes, derivative is $g(y)=(y-a)(y-b) y^{n-3} \in \mathbb{Q}[y]$. Conclude: $a, b$ either in $\mathbb{Q}$ or conjugate over $\mathbb{Q}$. Further $f(x)=$ :

$$
y^{n} / n-(a+b) y^{n-1} /(n-1)+a b y^{n-2} /(n-2)+d
$$

and $f(a)=f(b)$. With $d \in \mathbb{Q}, b / a=\alpha$, simplify: $(n-2)\left(\alpha^{n}-1\right)=n\left(\alpha^{n-1}-\alpha\right)$. Divide by $\alpha-1$ :

$$
h_{n}(\alpha)=(n-2) \alpha^{n-1}-2\left(\alpha+\cdots+\alpha^{n-2}\right)+(n-2),
$$

divisible by $(\alpha-1)^{2}$. Then, $f$ over $\mathbb{Q}$ exists when $\in \mathbb{Q}[\alpha]$ of degree 2 divides $h_{n}(\alpha) /(\alpha-1)^{2}$. Mathematica: $h_{n}(\alpha) /(\alpha-1)^{2}$ irreducible over $\mathbb{Q}$ for odd $n \leq 31$.

### 5.1. Nontrivial components of fiber products.

5.2. Reduced Nielsen classes and mapping class orbits. Automorphisms of $\hat{\varphi}$ identify with the centralizer of $G$ in $N_{R}(G, \mathbf{C})$. Point over $z_{0}$ gives $G\left(X / \mathbb{P}_{z}^{1}\right) \xrightarrow{\alpha} G$. A $\hat{\psi}$ (in (3.12)) is unique if $G$ has no center: There exists a unique total family

$$
\begin{equation*}
\mathcal{T}_{G, \mathbf{C}}^{\mathrm{in}}=\mathcal{T}^{\mathrm{in}} \rightarrow \mathcal{H}^{\mathrm{in}} \times \mathbb{P}_{z}^{1} \tag{5.4}
\end{equation*}
$$

$\mathcal{H}^{\text {in }}$ is a fine moduli space. The (minimal) field of definition of $\mathcal{T}_{\boldsymbol{p}}^{\text {in }} \rightarrow \boldsymbol{p} \times \mathbb{P}_{z}^{1}$ is $\mathbb{Q}(\boldsymbol{p})$ [FV91].

Proposition 5.2. Get $(G, \mathbf{C})$ regular realizations over $\mathbb{Q}$ from $\boldsymbol{p} \in \mathcal{H}^{\text {in }}(\mathbb{Q})$. Necessary: $\mathcal{H}^{\text {in }}$ has a $\mathbb{Q}$ component ( $\mathbf{C}$ is a rational union).
5.2.1. Reduced Nielsen classes. Notation for $M_{4}$ generators: $\gamma_{0}=q_{1} q_{2}, \gamma_{1}=$ $q_{1} q_{2} q_{1}, \gamma_{\infty}=q_{2}$. Product one:

$$
1=q_{1} q_{2} q_{1} q_{1} q_{2} q_{1}=\gamma_{1}^{2}=q_{1} q_{2} q_{1} q_{2} q_{1} q_{2}=\gamma_{0}^{3}
$$

Compute $Q_{i} \mathrm{~s}[\mathbf{F r i 9 0}]$ :

$$
Q_{1}=(25364)(798), Q_{2}=(14985)(367), Q_{3}=(25364)(798)
$$

Consider w-equivalence classes $\mathrm{Ni}\left(A_{5}, \mathbf{C}_{3^{4}}\right)^{\text {abs }} / \mathcal{Q}$ for $\mathrm{Ni}\left(A_{5}, \mathbf{C}_{3^{4}}\right)^{\text {abs }}$. Action of $M_{4}=H_{4} / \mathcal{Q}$ produces a (ramified) cover $\mathcal{H}^{\text {abs,rd }} \rightarrow \mathbb{P}_{j}^{1} \backslash\{\infty\}$. Compactify to $\bar{\varphi}^{\text {abs,rd }}: \overline{\mathcal{H}}^{\text {abs,rd }} \rightarrow \mathbb{P}_{j}^{1}$ with branch cycles $\left(\gamma_{0}^{\prime}, \gamma_{1}^{\prime}, \gamma_{\infty}\right)[$ Fri99, §7.4]:

$$
\gamma_{0}^{\prime}=(214)(378)(569), \gamma_{\infty}^{\prime}=(14985)(367)
$$

Note: The monodromy group is $A_{9}$. This is a cover with $(0,1, \infty)$ as branch points; the cusp widths are 1,3 and 5.
5.2.2. Organizing braid orbits with the sh-incidence matrix. First take $r=4$. The sh-incidence matrix summarizes a pairing using sh on $\gamma_{\infty}$ orbits.

For a general reduced Nielsen class, list $\gamma_{\infty}$ orbits as $O_{1}, \ldots, O_{n}$. The shincidence matrix $A(G, \mathbf{C})$ has $(i, j)$ term $\left|\left(O_{i}\right) \mathbf{s h} \cap O_{j}\right|$. Since sh has order two on reduced Nielsen classes, this is a symmetric matrix. Equivalence $n \times n$ matrices $A$ and $T A^{\mathrm{t}} T$ running over permutation matrices $T$ ( ${ }^{\mathrm{t}} T$ is its transpose) associated to elements of $S_{n}$. List $\gamma_{\infty}$ orbits as

$$
O_{1,1}, \ldots, O_{1, t_{1}}, O_{2}, 1, \ldots, O_{2}, t_{2}, \ldots, O_{u, 1}, \ldots, O_{u, t_{u}}
$$

corresponding to $\bar{M}_{4}$ orbits. Choose $T$ to assume $A(G, \mathbf{C})$ is arranged in blocks along the diagonal.

Lemma 5.3. If $A_{j}$ is the $j$ th block of $A(G, \mathbf{C})$, then $A_{j}$ doesn't break into smaller blocks. So, $\bar{M}_{4}$ orbits form irreducible blocks in the sh-incidence matrix.

Proof. With no loss assume one $\bar{M}_{4}$ orbit and two blocks, with orbit listings as $O_{1}, \ldots, O_{k}, O_{k+1}, \ldots, O_{t}$. As, however, there is one orbit, for some $j \leq k$, $\left|\left(O_{i}\right) \mathbf{s h} \cap O_{j}\right| \neq 0$ for some $i>k$. This contradicts there being two blocks.

In practice it is difficult to list the $\gamma_{\infty}$ orbits. So, we start with the H-M reps., apply sh, then complete the $\gamma_{\infty}$ orbits and check $\left|\left(O_{i}\right) \mathbf{s h} \cap O_{j}\right|$. Sometimes we'll then be done. The case $\left(A_{5}, \mathbf{C}_{3^{4}}\right)$ illustrates this. Denote (as above) the $\gamma_{\infty}$ orbits of $\boldsymbol{g}_{1}$ and $\boldsymbol{g}_{2}$ by $O(5,5 ; 1)$ and $O(5,5 ; 2) ; \gamma_{\infty}$ orbits of

$$
((513),(245),(154),(123)) \text { and }((324),(513),(154),(123))
$$

by $O(3,3 ; 1)$ and $O(3,3 ; 2)$; and of $\left(\boldsymbol{g}_{1}\right)$ sh by $O(1,2)$.
Table 1. sh-Incidence Matrix for $\mathrm{Ni}_{0}$

| Orbit | $O(5,5 ; 1)$ | $O(5,5 ; 2)$ | $O(3,3 ; 1)$ | $O(3,3 ; 2)$ | $O(1,2)$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $O(5,5 ; 1)$ | 0 | 2 | 1 | 1 | 1 |
| $O(5,5 ; 2)$ | 2 | 0 | 1 | 1 | 1 |
| $O(3,3 ; 1)$ | 1 | 1 | 0 | 1 | 0 |
| $O(3,3 ; 2)$ | 1 | 1 | 1 | 0 | 0 |
| $O(1,2)$ | 1 | 1 | 0 | 0 | 0 |

5.2.3. The $\mathbf{s h}$-incidence matrix for general $r$. For general $r$, denote $q_{1} \cdots q_{r-1}$ at the Hurwitz monodromy level to be the shift $\mathbf{s h}_{r}$, so $\mathbf{s h}_{4}$ is what we call the shift above. Ideas for $r=4$ generalize to indicate cusp geometry for general $r$.

The element $\mathbf{s h}_{r}$ plays the role of a shift in two ways. Consider an $r$-tuple $\overline{\boldsymbol{\sigma}}=\left(\bar{\sigma}_{1}, \ldots, \bar{\sigma}_{r}\right)$ of free generators of $F_{r}$. The effect of $\mathbf{s h}_{r}$ on $\overline{\boldsymbol{\sigma}}$ is to give

$$
(\overline{\boldsymbol{\sigma}}) q_{1} \cdots q_{r-1}=\left(\bar{\sigma}_{1} \bar{\sigma}_{2} \bar{\sigma}_{1}^{-1}, \ldots, \bar{\sigma}_{1} \bar{\sigma}_{r} \bar{\sigma}_{1}^{-1}, \bar{\sigma}_{1} \bar{\sigma}_{1} \bar{\sigma}_{1}^{-1}\right)
$$

In specializing to a Nielsen class the effect of $\mathbf{s h}_{r}$ is to the shift the Nielsen class representative entries by 1 . Iterate this $r$ times to see the effect of $\mathbf{s h}_{r}^{r}$ is conjugation on $\boldsymbol{\sigma}$ by the product $\bar{\sigma}_{1} \cdots \bar{\sigma}_{r}$ of these generators.

Such a conjugation commutes with the action of the braid group. So we have an interesting interpretation for the action of conjugating by $\mathbf{s h}_{r}$ on the generators
$q_{1}, \ldots, q_{r}$. Define $q_{0}$ to be $\mathbf{s h}_{r}^{-1} q_{1} \mathbf{s h}_{r}$. Then, conjugation by $\mathbf{s h}_{r}$ on the left of the array $\left(q_{0}, q_{1}, \ldots, q_{r-2}, q_{r-1}\right)$ maps its entries to

$$
\begin{aligned}
\mathbf{s h}_{r}\left(q_{0}, \ldots, q_{r-1}\right) \mathbf{s h}_{r}^{-1} & =\left(q_{1}, q_{1} q_{2} q_{1} q_{2}^{-1} q_{1}^{-1}, q_{1} q_{2} q_{3} q_{2} q_{3}^{-1} q_{2}^{-1} q_{1}^{-1}, \ldots\right) \\
& =\left(q_{1}, q_{2}, \ldots, q_{r-1}, q_{0}\right)
\end{aligned}
$$

To see the effect of conjugation of $\mathbf{s h}_{r}$ on $q_{r-1}$ use that $\mathbf{s h}_{r}^{r}$ is in the center of $H_{r}$ (or of $B_{r}$ ). Then, $q_{0}=\mathbf{s h}_{r}^{r}\left(q_{0}\right) \mathbf{s h}_{r}^{-r}=\mathbf{s h}_{r} q_{r-1} \mathbf{s h}_{r}^{-1}$.

Denote $q_{r-1} q_{r-2} \cdots q_{1}$ by $\mathbf{s h}_{r}^{\prime}$. Notice $\left(\mathbf{s h}_{r}^{\prime}\right)^{r}$ has exactly the same effect on $\boldsymbol{\sigma}$ as does $\mathbf{s h}_{r}^{r}$. In $H_{r}$ use that $q_{1} \cdots q_{r} q_{r} \cdots q_{1}=1$ to see $\mathbf{s h}_{r}^{r}\left(\mathbf{s h}_{r}^{\prime}\right)^{r}=1$, so $\mathbf{s h}_{r}^{r}=z$ has its square equal to 1 . When $r=4$ the group $M_{4}$ is exactly $H_{r} /\left\langle\mathbf{s h}_{4}^{4}\right\rangle=H_{r} /\langle z\rangle$. An especially handy description of $z$ in this case is $q_{1} q_{3}^{-1}$. In general there is a $\mathbf{s h}_{r}$-incidence matrix. As in the case $r=4$, it suffices to choose the image of $q_{v}$ in $\bar{M}_{r}$ for some value of $v$. It doesn't make any difference which $v$, though for $r=4$ it was convenient to take $v=2$. Call the resulting element $\gamma_{\infty}$. List the $\gamma_{\infty}$ reduced orbits as $O_{1}, \ldots, O_{t}$ and define $A(G, \mathbf{C})$ to be the matrix with $(i, j)$ term $\left|\left(O_{i}\right) \mathbf{s h}_{r} \cap O_{j}\right|$. For general $r$ it won't be symmetric.

## 6. Coordinates and covers

We will use covers of $\mathbb{P}_{z}^{1}$ as a record of an algebraic relation in one variable. Finding convenient ways to label such covers is necessary for applications. The test for the effectiveness of the labeling is how well it answers old questions and helps formulate new approaches to old topics. The first natural label to attach to a cover is its unordered set $\boldsymbol{z}$ of branch points. We know a lot about an $r$ branch point cover if we know its Nielsen class and its branch points. Yet in practise that isn't enough information to answer questions that have guided 200 years of intensive work on genus 1 curves. The following points related to coordinates will occur in the remainder of this chapter.
(6.1a) Relation and implication of the definition of algebraic cover of $\mathbb{P}_{z}^{1}$ in Chap. 2 to that of cover in $\S 2$.
(6.1b) Genus 0 dihedral involution covers ( of $\mathbb{P}_{z}^{1}$ ) correspond to rational functions $\mathbb{P}_{w}^{1} \rightarrow \mathbb{P}_{z}^{1}$ and implications of this for explicitly presenting genus 0 and genus 1 covers of $\mathbb{P}_{z}^{1}$.
(6.1d) Isomorphism classes of 1-dimensional complex tori correspond to $\mathrm{PGL}_{2}(\mathbb{C})$ equivaleces classes of points of $U_{r}$, and its analog for all Nielsen classes.
(6.1e) The universal covering space of $U_{\boldsymbol{z}}$ is the upper half plane and related examples of uniformization.
6.1. Algebraic covers and projective space. Chap. 2 (1.1) and (1.2) gave two definitions of analytic function $f(z)$ being algebraic function and related by an equation $m(w, z)=0$ to $z$. From this Chap. 3 Prop. 3.12 preduced an unramified cover $\varphi^{0}: X^{0} \rightarrow U_{\boldsymbol{z}}$. Its key point is that there are $n=\operatorname{deg}_{w}(m)$ distinct values $w^{\prime}$ for which $m\left(w^{\prime}, z^{\prime}\right)=0$ for $z^{\prime} \in U_{\boldsymbol{z}}$.

Example 6.1. From Chap. $2 \S 8.2$ (see Chap. 34.3 .3 ), a genus 1 degree 2 cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is algebraic. The unique cover of $U_{z}$ up to s-equivalence ( $\S 2.1 .2$ ) is the algebraic set $X_{\boldsymbol{z}}^{0}=\{(z, w) \mid m(z, w)=0\}$ with $m(z, w)=w^{2}-\prod_{i=1}^{4}\left(z-z_{i}\right)$. If one of the $z_{i}$ s is $\infty$ replace $z-z_{i}$ by 1 . Since the completion of $X^{0}$ to a ramified cover is unique, this also gives $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ up to s-equivalence. Abel's Theorem $\S 7.6$
shows they are also 1-dimensional complex tori. It further shows all 1-dimensional tori are algebraic. RETURN How can we use this to see all genus 1 surfaces are complex tori, and that all genus 0 curves are analytically isomorphic to $\mathbb{P}_{z}^{1}$ ?

Definition 6.2 (General algebraic cover). Continue Def. 2.10. Let $\varphi: Y \rightarrow X$ be an analytic map of compact Riemann surfaces with $\boldsymbol{x}$ the branch points of $\varphi$. We say $\varphi$ is algebraic if $Y$ is an algebraic cover (with some unspecified map to $\mathbb{P}^{1}$ ) and there exists an analytic map $\psi: Y \rightarrow \mathbb{P}_{w}^{1}$ so that for some $x^{\prime} \in X_{\boldsymbol{x}}, \psi$ separates the points in the fiber $X_{z^{\prime}}$.

Proposition 6.3. A cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is algebraic if and only if $\mathbb{C}(X)$ has sufficient algebraic functions to separate all points on $X$. This implies that if $X$ is algebraic and $X \rightarrow Y$ is a cover, then $Y$ is also algebraic.

A tentative definition of algebraic, called $\mathbb{P}^{1}$-algebraic, appears in Chap. 3 Def. 3.3. We start by explaining why it is typical to use the phrase algebraic on a manifold to mean it has an embedding in $\mathbb{P}^{N}$. Then, we consider the classical Luroth Theorem as a use of coordinates that leads us to note the complication in practical checking for decompositions of a cover.
6.1.1. Invariants and automorphisms of $\left(\mathbb{P}^{1}\right)^{N}$. Recall the definition of $\mathbb{P}^{1}$ algebraic from Chap. 3 (3.3). Ideas: Every algebraic manifold has a Galois cover by a $\mathbb{P}^{1}$-algebraic manifold, and every $\mathbb{P}^{1}$-algebraic manifold is algebraic (Segre embedding). Both $\mathbb{P}^{N}$ and $\left(\mathbb{P}^{1}\right)^{N}$ are simply connected [11.9d]. Yet, the spaces $\mathbb{P}^{N}$ are not $\mathbb{P}^{1}$-algebraic, for they have no analytic maps to $\mathbb{P}_{z}^{1}$ Chap. 3 [9.11e]. What do we get from algebraic that is better than uniformization?

### 6.2. Fields of definition, fields of moduli and Branch Cycle Lemma.

Lemma 6.4. If $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ has definition field $K$, then its branch points $\boldsymbol{z}$ form a $K$-set.

What you need from an algebraic structure to define the field of moduli. What you need from a cover to define these two quantities. One thing we can define is the definition field of the branch points $\boldsymbol{z}$ over a cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$.

Lemma 6.5 (Branch point control). Suppose $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is algebraic, and the branch points $\boldsymbol{z}$ have $K$ as a field of definition. Then, there is a cover $\varphi^{\prime}: X^{\prime} \rightarrow \mathbb{P}_{z}^{1}$ over $\bar{K}$ with these properties.
(6.2a) $\varphi^{\prime}$ is algebraic and has $\bar{K}$ as its definition field.
(6.2b) There exists analytic $\psi: X \rightarrow X^{\prime}$ with $\varphi^{\prime} \circ \psi=\varphi$.

Proof.
Give a cover that has field of moduli $\mathbb{R}$, but not field of definition $\mathbb{R}$.
6.2.1. Labeling Riemann surfaces and counting representatives of Nielsen classes. Prop. 2.18 gives a precise way to label a particular Riemann surface: Show the branch points and cuts on $\mathbb{P}_{z}^{1}$, and give $\boldsymbol{g} \in \mathrm{Ni}\left(G, \mathbf{C}, T_{G}\right)$. As we see in Chap. 5, if there is more than one s-equivalence class, equations for the cover won't have coefficients in functions of the branch points alone. That is very significant, and likely counter intuitive to the reader at this point. Many applications require finding a cover that has special coefficients (like over $\mathbb{Q}$ ). Inspecting actual equations for such a matter is rarely helpful if $r \geq 4$. So, what is the irrelevant information, and what to display?

Notation that memorably labels the conjugacy classes can be helpful in displaying expectations from the Branch Cycle Lemma about the position of the branch
points for a cover over $\mathbb{R}$, or the $p$-adics or over $\mathbb{Q}$. This is especially significant when two distinct conjugacy classes $\mathrm{C}_{i}, \mathrm{C}_{j}$ happen to be the same when extended to $S_{n}$ through $T_{G}: G \rightarrow S_{n}$. A graphic using a particular $\boldsymbol{g}$ representative of the Nielsen class can be revealing if it displays the real points lying over the real line on $\mathbb{P}_{z}^{1}$. This is pretty much the game when $r=3$, and it will be especially fruitful in applying to covers of the $\mathbb{P}_{j}^{1}$ when they are reduced Hurwitz spaces of some Nielsen class. Our most crucial cases have several $H_{r}$ orbits. When we know how to do so, we might present select representatives of those orbits.

Often our conjugacy classes have special shapes that allow computing the number of s-equivalence classes in a Nielsen class directly. Sometimes, however, it is good to know there is a pure computation for this count coming directly from the structure constant formula for the group ring $\mathbb{Z}[G]$ (for example, $[\mathbf{S e 9 2}, \S 7.2]$ or [Vö96, p. 54]). Let $m$ be constant on conjugacy classes $\mathrm{C}_{1}, \ldots, \mathrm{C}_{r}$ and $g \in G$. Value of $m$ on $\mathrm{C}_{i}$ is $m\left(\mathrm{C}_{i}\right)$. Denote $\sum_{i=1}^{r} \sum_{u_{i} \in C_{i}} m\left(u_{1} \cdots u_{r} g\right)$ by $I(m ; \mathbf{C}, g)$. If $m=\chi$ is an irreducible character of $G, I(\chi ; \mathbf{C}, g)=$

$$
\sum_{i=1}^{r} \sum_{u_{i} \in C_{i}} \chi\left(u_{1} \cdots u_{r} g\right)=\chi(y) \prod_{i=1}^{r} \chi\left(u_{i}\right) / \chi(1)^{r}
$$

Take $\chi_{1}, \ldots, \chi_{s}$ the irreducible complex characters of $G$. Then, $I(m ; \mathbf{C}, g)=$ $\sum_{i} m_{i} I\left(\chi_{i} ; \mathbf{C}, g\right)$. Write $m=m_{i} \chi_{i}$. Consider $\psi_{G}=\frac{1}{|G|} \sum_{i=1}^{s} \chi_{i}(1) \chi_{i}: 1$ at $1_{G}$ and 0 otherwise. So, $I\left(\psi_{G} ; \mathbf{C}, g\right)$ counts solutions of $u_{1}, \ldots, u_{r} g=1$ with $u_{i} \in \mathrm{C}_{i}$ :

$$
N\left(\mathrm{C}_{1}, \ldots, \mathrm{C}_{r}, g\right)=|G|^{r-1} \sum_{i=1}^{s} \prod_{j=1}^{r} \chi_{i}\left(\mathrm{C}_{j}\right) \chi_{i}(g)
$$

6.3. Branch cycles for sequences of genus $\mathbf{0}$ covers. $\S 2.4 .3$ describes the association of $r$ rooted cuts with classical generators for $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. We acknowledge [MP93] and [CG95] for their approach to the following computational problem.

Problem 6.6. Starting from $\varphi: \mathbb{P}_{u}^{1} \rightarrow \mathbb{P}_{z}^{1}$ and given $\bar{g}_{1}, \ldots, \bar{g}_{r} \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ classical generators, find an algorithm computing classical generators for $\mathbb{P}_{u}^{1}$ from among the following collection of elements.

## RETURN

Genus 0 and how to handle the appearance of such without knowing it is analytically OR topologically isomorphic to the sphere. The key point is to have that any oriented triangle bounds a disk, and then it is possible to replace a general tree with vertices by a collection of simple rooted cuts that then give a system of classical generators. Triangulate classical generators, and then consider lifts of the triangles to a cover using the cut version of Fig. 3 (§2.4.3).
6.3.1. Luroth's Theorem. Consider a rational function $f \in \mathbb{C}(w)$ of degree $n$. From Chap. 2 [9.4a],

$$
\boldsymbol{z}=\left\{z^{\prime} \in \mathbb{C} \left\lvert\,\left(f(x)-z^{\prime}, \frac{d f}{d x}\right)=1\right.\right\} \cup\{\infty\}
$$

Let $S=\left\{f^{-1}(\boldsymbol{z})\right\} \cup\{\infty\}$. From Ex. 3.14 (and Chap. 2 Thm. 6.4), $f: \mathbb{P}_{w}^{1} \backslash S \rightarrow U_{z}$ is a cover that extends to a map of compact manifolds $\bar{f}: \mathbb{P}_{w}^{1} \rightarrow \mathbb{P}_{z}^{1}$. Suppose $f$ is indecomposable: Not a composition $f_{1}\left(f_{2}(z)\right)$ with $\operatorname{deg}\left(f_{i}\right)>1, i=1,2$. Being indecomposable is equivalent to the Galois closure group $G_{f}$ being a primitive
subgroup of $S_{n}$. show exactly what permutations of points in the fiber of this cover extend to an automorphism of the cover (for most $f$, none) [9.21i].
6.4. Belyi's covers of $\mathbb{P}_{z}^{1}$. Suppose $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is a map of compact Riemann surfaces with $r$ branch points $\boldsymbol{z}$. Define the $J_{r}$-invariant to be the corresponding point $J(\varphi)$ of $U_{r} / \mathrm{PGL}_{2}(\mathbb{C})=J_{r}$. Suppose $J(\varphi)$ is in $K$. The algebraic half of the Existence Theorem says $X$ is w-equivalent to a cover over $\bar{K}$. In particular, two independent functions $\varphi_{1}$ and $\varphi_{2}$ on a compact Riemann surface $X$ give a canonical embedding of $X$ as a projective variety over $\bar{K}$. If $J\left(\varphi_{1}\right) \in \bar{K}$, then so is $J\left(\varphi_{2}\right)$. RETURN
6.4.1. Covers defined over $\overline{\mathbb{Q}}$. Special case of Lem. 6.5: 3 branch point covers strongly equivalent to covers over $\overline{\mathbb{Q}}$. From Lem. 6.5 we know that if any cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ has branch points over $\overline{\mathbb{Q}}$, then it is equivalent to a cover with definition field in $\overline{\mathbb{Q}}$. So, to construct a cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ equivalent to the given one over $\overline{\mathbb{Q}}$, with only three branch points, it suffices to find a map $h: \mathbb{P}_{z}^{1} \rightarrow \mathbb{P}_{w}^{1}$ that map the set $z$ into the set $\{0,1, \infty\}$.

Proposition 6.7. Each $\boldsymbol{z} \in \mathbb{P}_{z}^{1}(\overline{\mathbb{Q}})$ has

$$
\mathbb{P}_{z}^{1} \xrightarrow{f_{1}} \mathbb{P}_{w_{1}}^{1} \rightarrow \cdots \xrightarrow{f_{u}} \mathbb{P}_{w_{u}}^{1}
$$

with $f_{u} \circ \cdots \circ f_{1}$ having branch points $0,1, \infty$ and $\boldsymbol{z} \subset f^{-1}(0,1, \infty)$. More generally, if $\boldsymbol{z} \in \mathbb{P}_{z}^{1}(K)$ with $K$ of transendence dimension $t$ over $\mathbb{Q}$, then there is a similarly result, with $f$ having $t+3$ (or fewer) branch points.

Main use: $G_{\mathbb{Q}}$ faithful on projective systems of 3 branch point covers. Use $g=g(X)$ for genus of Riemann surface $X$. Subtopics:

- List of 3 branch point $S_{n, g}$ covers
- Observations giving Belyi's result
- Compare $S_{n, g}$ covers with Belyi maps with Guralnick's genus $g$ problem
- Some nonrigid 3 branch point covers
(6.3a) Find $g_{1}: P_{z}^{1} \rightarrow P_{u_{1}}^{1}$ with $f_{1}(\boldsymbol{z})$ in branch point locus, branch points of $g_{1}$ in $\mathbb{Q}$.
(6.3b) With $g_{1}(\boldsymbol{z})$ in $\mathbb{Q}$, compose with $g_{2}: P_{u_{1}}^{1} \rightarrow P_{u_{2}}^{1}$ so $g_{2} \circ g_{1}$ has 3 branch points.
Induct on maximal degree of $\boldsymbol{z}$ support and $r$. Always assume $\left\{z_{1}, z_{2}, z_{3}\right\}=$ $\{0,1, \infty\}$.

Step a: Map by $f$, irreducible polynomial for branch point of maximal degree over $Q$. Adds branch points of $f$ to list, but new branch points are zeros of $\frac{d f}{d z}$ : have lower degree.

Step b: Write $z_{4}=a / b, a, b \in \mathbb{Z}$. Choose $\psi(z)=z^{u}(z-1) v$ : logarithmic derivative is $u / z+v /(1-z)$. Assure $z_{4}$ is a branch point by choosing $u, v \in \mathbb{Z}$ so $u / z_{4}+v=\left(1-z_{4}\right)$. This reduces the branch points by one.
6.4.2. Three branch point $S_{n, g}$ covers. Call a cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ an $S_{n, g}$ cover if $g(X)=g$ and $\operatorname{deg}) \psi)=n$.

Lemma 6.8. Fix $g$. There are infinitely many $n \geq 1$ for which there are three branch point covers $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$, with $g(X)=g$ and having monodromy $S_{n}$.

Proof. Take $n=m_{1}+\cdots+m_{s}$. Modify $S_{n, 0}$ covers:

$$
g_{1}=\left(1 \ldots m_{1}\right) \cdots\left(m_{1}+\cdots+m_{s-1}+1 \ldots n\right)
$$

Genus 0: Take

$$
\begin{aligned}
g_{2}= & \left(m_{1}-1 \ldots 1\right)\left(m_{1}+m_{2}-1 \ldots m_{1}+1\right) \ldots \\
& \left(n-1 \ldots n-m_{s}+1\right)\left(m_{1} m_{1}+m_{2} \ldots n\right)
\end{aligned}
$$

Then, $\operatorname{ind}\left(g_{1}\right)=n-s$ and $\operatorname{ind}\left(g_{2}\right)=n-s-1$. Compute $g_{1} g_{2}=g_{3}$ :

$$
\left(m_{1} m_{1}-1 m_{1}+m_{2} m_{1}+m_{2}-1 \ldots n n-1\right)
$$

So, $\operatorname{ind}\left(g_{3}\right)=2 s-1$. RET gives genus 0 covers. For $S_{n}$, select $m_{1}, \ldots, m_{s}$ accordingly.

For $g=1$ covers, switch 1 and 2 in $g_{1}$, but not in $g_{2}$. This changes nothing from conclusions, except adding 2 to index of $g_{3}$ :

$$
\left(m_{1} 21 m_{1}-1 m_{1}+m_{2} m_{1}+m_{2}-1 \ldots n n-1\right)
$$

That concludes the proof.
6.4.3. Comparison with the genus $g$ problem. Belyi produces three branch point $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$, usually composed of many maps between $\mathbb{P}^{1} \mathrm{~s}$. The construction rarely provides covers that are primitive.

Question 6.9. Does $X$, over $\overline{\mathbb{Q}}$ have a three branch point cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ with primitive monodromy.

Help from the literature. Fix $g$ with $g>6$. By [GN95], only $<\infty$ ly many $X$ of genus $g$ have three branch point covers with solvable monodromy ([Fri99, $\S 5$ - $\S 6]$ for genus 0 problem).

CONJECTURE 6.10 (Guralnick). Genus $g$ three branch point primitive covers of $\mathbb{P}_{z}^{1}$ with monodromy neither $S_{n}, \mathbb{Z} / 2 \imath S_{n}, A_{n}$ or $\mathbb{Z} / 2 \imath A_{n}$ (for some $n$ ) are finite.

The hardest case will be the 3 branch point case, though the likely result is more precisely the following. There is a function $N(g)$, quadratic in $g$, and either $n \leq N(g)$ or one of the following holding.
(6.4a) The Galois closure of $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ has genus at most 1 and $G$ is cyclic, dihedral with $n$ a prime or $n$ is a prime or prime squared and the Galois closure is an elliptic curve (affine case).
(6.4b) $g=0$ or 1 and the monodromy group $G$ satisfies $G=A_{n}, S_{n}, A_{m} 乙 S_{2} \leq$ $G \leq S_{m} \imath S_{2}\left(n=m^{2}\right)$
(6.4c) $g=0$ or 1 and $G=A_{m}$ or $S_{m}$ with $n=m(m-1) / 2.4 . g>1$ and $G=A_{n}$ or $S_{n}$.
Of course, if we have $f: X \rightarrow Y$ primitive and $Y$ has genus $h>0$, then for $h>1, n \leq g$, so only finitely many and if $h=1$, then either $n$ is bounded in terms of $g$ or $G=A_{n}$ or $S_{n}$ (these last two statements are in Guralnick-Neubauer from the first Seattle meeting and these hold in positive characteristic as well which is noted in my MSRI paper [the other volume].

The work left is in dealing with symmetric and alternating groups acting on $k$-sets with 3 or 4 branch points and in dealing with some product actions where the group is essentially $S_{m}$ 乙 $S_{2}$ acting with $n=(m(m-1) / 2)^{2}$. Kay, Dan and I should have shortly a complete list of $g=0$ aside from the open cases listed above (i.e. ones not involving symmetric groups) and John Shareshian and I are working on symmetric groups with possibly all of us working on the product case.

Question 6.11 (Easier Question). Compatible with Lem. 6.8, are the $[X] \in$ $\mathcal{M}_{g}$ having an $S_{n, g}$ cover of $\mathbb{P}_{z}^{1}$ dense?
6.4.4. Non-rigid $A_{n}$ covers. Consider $n \geq 5$, odd and squarefree. Let $\mathbf{C}$ be conjugacy classes of $\left(g_{1}, g_{2}, g_{3}\right) \in A_{n}^{3}$ with $g_{1}=(12)(34), g_{2}=(13567 \ldots n)$ and $g_{3}=(12 \ldots n)^{-1}$. Check: Geometric monodromy is $A_{n}$. Representatives for conjugation of $S_{n}$ on $\mathrm{Ni}\left(S_{n}, \mathbf{C}\right) / S_{n}, i=3, \ldots,(n+1) / 2$ :

$$
\boldsymbol{g}_{j}^{\prime}=\left((12)(j j+1),(13 \ldots j j+2 j+3 \ldots n), g_{3}^{\prime}\right)
$$

Question 6.12. Exists $f: \mathbb{P}_{y}^{1} \rightarrow \mathbb{P}_{z}^{1}$ in $\mathbb{Q}[y]$ ?
If yes, derivative is $g(y)=(y-a)(y-b) y^{n-3} \in \mathbb{Q}[y]$. Conclude: $a, b$ either in $\mathbb{Q}$ or conjugate over $\mathbb{Q}$. Further $f(x)=$ :

$$
y^{n} / n-(a+b) y^{n-1} /(n-1)+a b y^{n-2} /(n-2)+d
$$

and $f(a)=f(b)$. With $d \in \mathbb{Q}, b / a=\alpha$, simplify: $(n-2)\left(\alpha^{n}-1\right)=n\left(\alpha^{n-1}-\alpha\right)$. Divide by $\alpha-1$ :

$$
h_{n}(\alpha)=(n-2) \alpha^{n-1}-2\left(\alpha+\cdots+\alpha^{n-2}\right)+(n-2)
$$

divisible by $(\alpha-1)^{2}$. Then, $f$ over $\mathbb{Q}$ exists when $\in \mathbb{Q}[\alpha]$ of degree 2 divides $h_{n}(\alpha) /(\alpha-1)^{2}$. Mathematica: $h_{n}(\alpha) /(\alpha-1)^{2}$ irreducible over $\mathbb{Q}$ for odd $n \leq 31$.
6.4.5. Belyi's Theorem in positive characteristic.

Proposition 6.13. Let $K$ be an algebraically closed field of characteristic $p>$ 0. A projective curve $X$ over $K$ has definition field $K$ if there is a finite map $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ over $K$ with only tame ramification and at most three branch points. Further, if a projective curve $X$ has field of definition $\overline{\mathbb{F}}_{p}$, then it admits a finite map $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ with only tame ramification and at most three branch points if there is at least one finite map $\varphi^{\prime}: X \rightarrow \mathbb{P}_{z}^{1}$ over $\overline{\mathbb{F}}_{p}$ with with only tame ramification.

Proof. The first statement follows from knowing that a deformation of tame covers that leaves the branch points fixed doesn't change the equivalence class of the cover. This is one half of Grothendieck's main theorem on the fundamental group of a cover. The second statement is very simple. The map $C_{p^{n}-1}: \mathbb{P}_{z}^{1} \rightarrow \mathbb{P}_{w}^{1}$ by $z \mapsto z^{p^{n}-1}-1$ maps all elements of $\mathbb{F}_{p^{n}}^{*}$ to 1 . The cover $C_{p^{n}-1}$ is tamely ramified (and ramified only over 0 and $\infty$. So, by choosing $n$ so that $\mathbb{F}_{p^{n}}$ contains all branch points, $C_{n}$ maps these down to $0,1, \infty$.

The hypothesis of existence of a tame cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ is Saïdi's; he noted more than the simple argument in the proof of Prop. 6.13. In characteristic $p \geq 3$, such tame covers exist by an analog of an argument of Lefschetz. This say you may project the curve $X$ from a nonsingular projective embedding in $\mathbb{P}^{3}$ to get a simple branched cover over the algebraic closure of the finite field. (The arithmetic form of the argument of [FJ78, Lem. 2.1] applies since we may take the finite field cardinality to be arbitrarily large.) [Schr02, §6] raises all these points, then takes it one step further to deal with removing the hypothesis of existence of a tame cover when $p=2$.
[Schr02, §5] allows going from showing the general curve of genus $g$ over $\mathbb{C}$ has a presentation in the Nielsen class $\mathrm{Ni}\left(A_{n}, \mathbf{C}_{3^{r}}\right)$ to finding an open algebraic set of curves of genus $g=r+1$ over $\overline{\mathbb{F}}_{2}$ in this Nielsen class. [Schr02, Cor. 5.3] proves the space of isomorphism classes of curves over $\overline{\mathbb{F}}_{2}$ in this Nielsen class has dimension at least $2 g-3$. In the language of Chap. $5 \S ? ?$ and $[\mathbf{F r i 8 9}]$, the moduli dimension of the Nielsen class $\operatorname{Ni}\left(A_{n}, \mathbf{C}_{3^{r}}\right)$ is $2 g-3$ for $r \geq n$. The expectation is that the moduli dimension is actually $3 g-3(r \geq n+1)$ : full moduli. Schroer's argument, though it quotes [FKK01], does not identify the two components - separated by a
spin lifting invariant ([Fri95a, Ex. III.12], [Fri96] and Chap. 5 §??) - that occur in these two families. We expect full moduli for both components. Even if the Nielsen class $\mathrm{Ni}\left(A_{n}, \mathbf{C}_{3^{r}}\right)$ has full moduli dimension, unlike the case where $p \geq 3$, it is unclear (even over $\mathbb{C}$ ) if that implies each genus $g$ curve has a representing cover with only odd order ramification.
6.5. Higher genus versions of Thm. 1.8 and uniformization. We now generalize Thm. 1.8 for any compact Riemann surface $X$. A proof along the lines of that theorem works, but with some technical difficulties (see [11.11]).
6.5.1. Homology of a manifold and triangulations. We tacitly assume a compact Riemann surface has a triangulation, say it is given by an analytic map to $\mathbb{P}_{z}^{1}$.

Lemma 6.14. Let $X$ be a compact Riemann Surface and $\Gamma\left(X, \Omega_{X}\right)$ its space of global holomorphic differentials. With $\chi_{X}$ the Euler characteristic of $X$, the dimension $u \leq\left(2-\chi_{X}\right) / 2=g_{X}$.

Theorem 6.15. Let $X$ be a compact Riemann surface. Let $\left\{x_{1}, \ldots, x_{r}\right\}=\{\boldsymbol{x}\}$ be $r$ distinct points on $X$. There is a number $g=g(X)$ such that for $x_{0} \in X \backslash\{\boldsymbol{x}\}=$ $X^{0}$, there are closed paths

$$
\begin{equation*}
\alpha_{1}, \ldots, \alpha_{g} ; \beta_{1}, \ldots, \beta_{g} ; \gamma_{1}, \ldots, \gamma_{r} \tag{6.5}
\end{equation*}
$$

based at $x_{0}$ so their homotopy classes generate $\pi_{1}\left(X^{0}, x_{0}\right)$ with the one relation

$$
\begin{equation*}
\left[\alpha_{1}\right]\left[\beta_{1}\right]\left[\alpha_{1}\right]^{-1}\left[\beta_{1}\right]^{-1} \cdots\left[\alpha_{g}\right]\left[\beta_{g}\right]\left[\alpha_{g}\right]^{-1}\left[\beta_{g}\right]^{-1}\left[\gamma_{1}\right] \cdots\left[\gamma_{r}\right] \tag{6.6}
\end{equation*}
$$

That is, Thm 6.15 says $\pi_{1}\left(X, x_{0}\right)$ is isomorphic to the quotient of the free group on (symbols given by) homotopy classes of the (6.5) paths, by the smallest normal subgroup containing expression (6.6). In addition, let $(U, \varphi)$ be any coordinate neighborhood (with trivial fundamental group) in an atlas for $X$ (Chap. 3 Def. 1.5) so that $U$ contains $\left\{x_{0}, x_{1}, \ldots, x_{r}\right\}$. Then we may take $\gamma_{1}, \ldots, \gamma_{r}$ any set of paths that $\varphi$ maps to a collection of classical generators relative to $\left(\varphi\left(x_{1}\right), \ldots, \varphi\left(x_{r}\right)\right)$ based at $\varphi\left(x_{0}\right)$ in $\varphi(U)$.

Example 6.16 (A sphere with $g$ handles-the case $r=0$ ). Cut out 2 g disjoint discs from the sphere; then join the boundaries of these discs in pairs by cylinders that slightly flare at the ends. We may embed any compact Riemann surface (having genus $g$ ) in $\mathbb{R}^{3}$ as a sphere with $g$ handles. This is easy to prove for the complex torus of Chap. 3 Ex. 3.2.2 - it is homeomorphic to a sphere with one handle. The general fact, however, is more difficult [Spr57, Chap. 5]. If follows, however, quite easily from the uniformization of such a Riemann surface $(g \geq 2)$ by the disk. PUT THE ARGUMENT OF A POLYGONAL DOMAIN HERE.

Nevertheless, assuming this, it is easy to draw paths representing generators of the fundamental group of a sphere with $g$ handles (Fig. 9). Take $g$ nonintersecting paths $\bar{\alpha}_{1}, \ldots, \bar{\alpha}_{g}$ each of which goes like a bracelet around the handles; and $g$ further nonintersecting paths $\bar{\beta}_{1}, \ldots, \bar{\beta}_{g}$ where $\bar{\beta}_{i}$ travels along the ith handle from the unique point of intersection, $m_{i}$, of $\bar{\alpha}_{i}$ and $\bar{\beta}_{i}$ to the edge of the first hole cut in the sphere, then along the sphere to the other hole defining the handle, and, finally, back along the handle to $m_{i}$. In addition, these paths are chosen so that $\bar{\alpha}_{i}$ and $\bar{\beta}_{j}$ do not intersect for $i \neq j$, and the crossproduct $\boldsymbol{t}\left(\alpha_{i}\right) \times \boldsymbol{t}\left(\beta_{i}\right)$ of the tangent vectors to $\bar{\alpha}_{i}$ and $\bar{\beta}_{i}$ at $m_{i}$ (§1.b) points outward.

Figure 9. A sphere with one handle $(g=1)$ pretties itself


To obtain paths based at some specific point $x_{0}$ draw paths $\delta_{i}$ from $x_{0}$ to $m_{i}$ and let $\alpha_{i}=\delta_{i} \bar{\alpha}_{i}\left(\delta_{i}\right)^{-1}$ (resp., $\beta_{i}=\delta_{i} \bar{\beta}_{i}\left(\delta_{i}\right)^{-1}$ ), $i=1, \ldots, g$. It is convenient to choose $x_{0}$ off the handles.
6.6. The Schwarz-Christoffel Transformation. Given a polygon $P$ with vertices $A_{1}, \ldots, A_{n}$ and clockwise interior angles $\pi \alpha_{1}, \ldots, \pi \alpha_{n}$, we want to map the interior $\Pi$ of this polygon on the Upper-half plane $\mathbb{H}$ in a one-one conformal fashion onto $\mathbb{H}$. We actually do the opposite, $f: \mathbb{H} \rightarrow \Pi$. Let $a_{1}<a_{2}<\cdots<a_{n}$ map respectively to $a_{1}, \ldots, a_{n}$ by $f$. We can assign three of these points arbitrarily, and $f$ is one-one on $\mathbb{H}$. [Hil62, Thm. 17.5.3] says: Let $D$ be the interior of a simple closed curves $C$ and $f: K \rightarrow D$ onto conformally. Then, $f$ is continuous on $\bar{K}$ and the correspondence between $\bar{K}$ and $\bar{D}$ is one-one and bi-continuous.

Since such an $f$ is locally one-one, $f^{-1}$ is locally defined on $D$ which is simplyconnected. So, by the Monodromy Theorem (Chap. 3 Thm. 6.11) any branch of $f^{-1}$ extends to $D$, and it is one-one. The next two subsections (based on [?, p. 372-374]) show $f$ is the same as the function

$$
\begin{equation*}
F(z)=C_{1} \int_{a_{1}}^{z}\left(t-a_{1}\right)^{\alpha_{1}-1} \cdots\left(t-a_{n}\right)^{\alpha_{n}-1} d t+C_{2} \tag{6.7}
\end{equation*}
$$

To define $F$ locally around any convenient point in $\mathbb{H}$ (say $z_{0}=i$ ) as in Chap. 2 $\S 3.4$, use a branch $F_{i}(z)$ of $\log \left(z-a_{i}\right)$ around $z_{0}$. Then, interpret the integral for $F$ to be a primitive (antiderivative) of $\prod_{i=1}^{n} e^{\left(\alpha_{i}-1\right) F_{i}(z)}$. Since $\mathbb{H}$ is analytically isomorphic to a disk, this choice of $F$ extends analytically to all of $\mathbb{H}$ (Chap. 2 Prop. 3.6).
6.6.1. Differential properties of $F$. There are two components to the complement of a simple closed polygonal path [11.3a]. The interior is the component $U_{P}$ that consists of points with nonzero winding number with respect to $P$. So, $U_{P}$ is simply connected according to Chap. 2 §8.3. Apply the Riemann mapping theorem to show there exists an analytic one-one $f: \mathbb{H} \rightarrow U_{P}$. We want to identify $f$.

Consider the differential equation

$$
\begin{equation*}
\frac{d\left(\log \left(h^{\prime}(z)\right)\right.}{d z}=\frac{h^{\prime \prime}}{h^{\prime}}=\sum_{j=1}^{n} \frac{\alpha_{j}-1}{z-a_{j}} \stackrel{\text { def }}{=} g(z) . \tag{6.8}
\end{equation*}
$$

Note that $g$ is a linear fractional transformation with simple poles at $a_{j}$ and corresponding residue $\alpha_{j}-1$. Also, $g$ vanishes at $\infty$ because the sum of the interior angles adds to $2 \pi$. Then, $F=h$ is a solution of (6.8). Then, To show that $f$ is also, we show the logarithmic derivative $\frac{f^{\prime \prime}}{f^{\prime}}$ of $f^{\prime}(z)$ has the residues and poles of $g$ and it vanishes at $\infty$. As the properties of $g$ come from its being meromorphic on the whole plane, that first requires showing we can analytically continue $f$ to anywhere in the whole plane minus $a_{1}, \ldots, a_{n}$.

Remark 6.17 (Point of the logarithmic derivative). The expression $\frac{h^{\prime \prime}}{h^{\prime}}$ is invariant under composition by elements of $\mathrm{PGL}_{2}(\mathbb{C})$ having the form $z \mapsto a h(z)+b$ for $a \in \mathbb{C}^{*}, b \in \mathbb{C}$. Conversely, if $\mathcal{D}_{\mathcal{A}}(h)=\frac{h^{\prime \prime}}{h^{\prime}}=\frac{h_{1}^{\prime \prime}}{h_{1}^{\prime}}$, then $h_{1}(z)=a h(z)+b$ : This differential equation $\mathcal{D}_{\mathcal{A}}(h)=\mathcal{D} \mathcal{A}(a h(z)+b)$ characterizes invariance of functions under the action of the affine group $\mathcal{A}=\mathbb{C} \times^{s} \mathbb{C}^{*}$ as in [11.25c].
6.6.2. Schwarz's famous reflection principle. Define $f$ in the lower half plane by crossing $\left(a_{j}, a_{j+1}\right)$ and reflecting its values across that line exactly as Schwartz did it [Sc1890, A paper from 1866]. Suppose $h$ is any function on a domain $D$ in the upper half plane that extends continuously to the line segment $(a, b) \subset \mathbb{R}$ (on the boundary of $D$ ). Denote the set of points of $D$ reflected in the $x$-axis by $\bar{D}$. Further, assume $h$ maps $(a, b)$ to $L_{h(a), h(b)} \stackrel{\text { def }}{=}\{h(a)+t(h(b)-h(a)) \mid t \in(0,1)\}$. Let $-L_{h(a), h(b)}$ denote reflection of points $z \in \mathbb{C}$ in the line through $L_{h(a), h(b)}$. Then, $-L_{h(a), h(b)}: z \mapsto A\left(\overline{A^{-1}(z)}\right)$ with $A(z)=\left(h(b)-h(a)\left(\frac{(z-a)}{(b-a)}\right)+h(a)\right.$.

LEMMA 6.18. The formula $h(w) \stackrel{\text { def }}{=} h(\overline{\bar{w}})^{L_{h(a), h(b)}}$ defines a function analytic on $D \cup \bar{D} \cup(a, b)$ and equal to $h(z)$ for $z \in D$.

Proof. If $A^{-1}(h(z))$ extends analytically to $D \cup \bar{D} \cup(a, b)$, then so does $h(z)$. This reverts us to the case $h$ takes $(a, b)$ to the line segment $(0,1)$. [Ahl79, p. 172173] emphasizes in the proof of this case, that it comes to extending the real and imaginary part of such an $h$ to be harmonic. This he does by the formulas $\Re(h)(\bar{z})=$ $u(z)$ and $-\Im(h)(\bar{z})=v(z)$ for $z \in D$. The tricky part is showing the points of $(a, b)$ are also in the domain of harmonicity of the extending function $V(z): V(z)=v(z)$ for $z \in D,-v(\bar{z})$ for $z \in \bar{D}$, and 0 on $(a, b)$. For $t \in(a, b)$, this is an application of the Poisson integral [Ahl79, p. 168] $P_{V}$ defined on any small disk $D_{t}$ about $t$ by the boundary values of $V$ on that disk. The function $V-P_{V}$ vanishes on the intersection of $D_{t}$ and the real line, and on the boundary of the disk in the upper half-plane. By the maximum principle for harmonic functions, $V-P_{V}$ is identically zero on the upper half of the disk, and similarly on the lower half of the disk, concluding the proof.

Apply this to $D=\mathbb{H}$ : We analytically continue $f$ to the lower plane by extending $f=f^{(0)}$ to a function $f_{i}^{(0)}$ on $D \cup \bar{D} \cup\left(a_{i}, a_{i+1}\right)=U_{i}$ for some $i$. The effect is that we analytically continued along a path crossing $\left(a_{i}, a_{i+1}\right)$ from the upper half plane to the lower half plane, and then took the unique function defined on the simply connected (even contractible) set $U_{i}$. From the reflection principle, the analytic function $f_{i}^{(0)}$ maps the lower half plane to the original (open) polygon reflected in the line through $f\left(a_{i}\right)$ and $f\left(a_{i+1}\right)$. Now work with $f_{i}^{(0)}$ in the lower half plane, and apply the same principles. We can extend it back up into the upper
half plane using any of the line segments $\left(a_{j}, a_{j+1}\right)$. Denote $f_{i}^{(0)}$ by $f^{(1)}$. This analytically continues to the whole upper half plane (along a path) through $\left(a_{j}, a_{j+1}\right)$. The result is a function $f_{j}^{(1)}$ on $\mathbb{H}$ that maps $\mathbb{H}$ to the interior of a new polygon.

By the reflection principle, the new polygon is just the result of reflecting the original polygon in two different lines. Note: Composition of two such reflections is an affine transformation: $M: \mathbb{C} \rightarrow \mathbb{C}$ by $z \mapsto a^{\prime} z+b^{\prime}$ for some $a^{\prime} \in \mathbb{C}^{*}$ on the unit circle and $b^{\prime} \in \mathbb{C}$. So the effect of the analytic continuations that end in the upper half plane is to give polygons congruent to the original polygon. On $\mathbb{H}$, since $f^{(1)}$ is composition with an affine transformation, we have $\mathcal{D}_{\mathcal{A}}(f)=\mathcal{D}_{\mathcal{A}}\left(f^{(1)}\right.$ (see Rem. 6.17). The differential operator provides a function $\mathcal{D}_{\mathcal{A}}(f)$ invariant under the different analytic continuations of $f$ in the domain $\mathbb{C} \backslash\left\{a_{1}, \ldots, a_{n}\right\}$.
6.6.3. Singularities of analytic continuations of $f$. The Scharz-Christoffel transformation that maps $\mathbb{H}$ onto the sector at angle $0<\alpha \pi<\pi$ is $f(z)=z^{\alpha}$. We know $f$ at $a_{j}$ maps an interior angle of $\pi$ radians to one of $\alpha_{j} \pi$ radians, and it analytically continues around $a_{j}$. The residue of $\frac{f^{\prime \prime}}{f^{\prime}}$ for such a function $f$ is well-defined up to a change of variables in $f$ in a neighborhood of a punctured disk around $a_{j}$. With no loss, for the computation of the residue, $f$ to be a branch of $\left(z-a_{j}\right)^{\beta}=e^{\beta \log (w)}$ with $w=z-a_{j}$ (Chap. $2 \S 8.2$ ). Knowing that analytic continuation around the counterclockwise upper disk, gives a change of angle of $\alpha_{j} \pi$, shows $\beta$ is $\alpha_{j}$. So, the residue of $\frac{f^{\prime \prime}}{f^{\prime}}$ at $a_{j}$ is given by the residue of $\alpha_{j}\left(\alpha_{j}-1\right) z^{\alpha_{j}-1} / \alpha_{j} z^{\alpha_{j}}=\left(\alpha_{j}-1\right) 1 / z$, which is $\alpha_{j}-1$. Conclude that $\frac{f^{\prime \prime}}{f^{\prime}}$ has the right residues and is meromorphic everywhere. So, $f=F$ for some choice of $C_{1}$ and $C_{2}$.
6.7. Monodromy and hypergeometric functions. [Ahl79, p. 315-321] Starts with a discussion of homogeneous linear ordinary differential equations with meromorphic function coefficients of degree $n$ :

$$
\begin{equation*}
\sum_{k=0}^{n} a_{i}(z) w^{(k)}=0 \tag{6.9}
\end{equation*}
$$

A meromorphic solution $w(z)$ in a neighborhood of $z_{0}$ is ordinary if $a_{0}\left(z_{0}\right) \neq 0$, and around such a point there are $n$ linearly independent solutions. At an ordinary point if you also specify the values of $w^{(k)}\left(z_{0}\right), k=0, \ldots, n-1$, then there will be a unique solution $w(z)$. The standard proof of this is given by writing a power series solution, solving for the coefficients inductively while establishing the series converges [Ahl79, p.310]. The conditions on the values of the first $n$ terms in the powers series are linear, establishing that the space of solutions around $z_{0}$ is $n$ dimenional. While general (not necessarily linear) differential equations generalize our basic study of algebraic functions, linear such do not. For example, in studying algebraic equations attached to genus 1 curves, we treat the differential equation $\left(w^{\prime}\right)^{2}=a w^{3}+b w+c$ with $a, b, c \in \mathbb{C}$. Linear equations arise here in a different context.
6.7.1. Monodromy from a differential equation. Let $\boldsymbol{z}$ be the finite set of points at which at least one of the $a_{i}$ s is not analytic or at which $a_{0}(z)$ has a zero. Then, as with any other analytic continuation situation we may analytically continue a basis $f_{1}, \ldots, f_{n}$ of solutions at $z_{0}$ around any element of $\Pi\left(U_{z}, z_{0}, z_{1}\right)$. At least for such linear differential equations, Riemann suggested establishing similar properties as for analytic continuation of meromorphic algebraic functions. Under certain conditions we can expect the monodromy action to determine the differential equation.

This, however, requires inspecting the nature of the differential equation near the support of $\boldsymbol{z}$ in a new way. By multiplying through by any denominators, we may with no loss assume the coefficients are analytic in $\mathbb{C}_{z}$, and that they have no common zero. In a neighborhood of $z_{0}$ we can analytically continue the solution, and since the space has dimension $n$, we are getting a representation of $\mathbb{Z}$ through a matrix $M_{z_{0}}$. Suppose $n=2$ and the monodromy matrix has distinct eigenvalues $e^{2 \pi i \alpha_{1}}$ and $e^{2 \pi i \alpha_{2}}$. Then, the solutions would have local expressions as $z^{\alpha_{1}} h_{1}(z)$ and $z^{\alpha_{2}} h_{2}(z)$, with $h_{1}$ and $h_{2}$ analytic and nonzero at $z_{0}$. As usual, write $\ln \left(z-z_{0}\right)=m(z)$, a a branch of $\log$ in a disc about a point $z^{\prime}$ near $z_{0}$ to see the effect of a clockwise analytic continuation about $z^{\prime}$ about $z_{0}: e^{m(z) \alpha} \mapsto$ $e^{(m(z)+2 \pi i) \alpha}$. Looking at the term of highest order, gives the indicial equation for such $\alpha$, which is the characteristic polynomial for the monodromy action matrix. A regular singular point is one for which the characteristic polynomial is computable from $\alpha$ satisfying the highest order solution from undetermined coefficients. State this condition in terms of the orders at $z_{0}$ of the coefficients of the differential equation. There are complications, however, if the solutions of that indicial equation have $e^{2 \pi i) \alpha_{1}}=e^{2 \pi i) \alpha_{2}}$. One possibility is that there are two values of $\alpha$ and they do give independent solutions. Another is that if you take the larger value (difference is positive integer) you get the solution $w_{1}(z)$ above, and you get another one by taking a solution of form $w_{2}(z)=C w_{1}(z) \ln \left(z-z_{0}\right)+\left(z-z_{0}\right)^{\alpha_{2}} h_{2}(z)$, and finally if $\alpha_{1}=\alpha_{2}$, this last situation does occur for certain. The general case is given by a similar Jordan canonical form.

We can also change variables and consider the possibility of having a regular singular point at $\infty$. If there are regular singular points everywhere (and only finitely many of them), then all coefficients are rational functions in $z$. When $n=2$ this says $w^{\prime \prime}-p w^{\prime}-q w$ with $p$ and $q$ rational functions with $q$ having at most a double pole, and $p$ having at most a single pole anywhere in the finite plane. Then, at $\infty$ with $z=1 / u$, write $W(u)=w(1 / u)$ and $\frac{d w}{d z}$ after as $-u^{2} \frac{d W(u)( }{d u}$ and $\frac{d^{2} w}{d z^{2}}=$ after substitution $z \mapsto 1 / u$ as $2 u^{3} \frac{d W(u)( }{d u}+u^{4} \frac{d^{2} W(u)}{d^{2} u}$. So, now we check if $u=0$ is a regular singular point for this equation. A regular singular point is equivalent to $2 / u+u^{-2} p(1 / u)$ has a pole of order at most 1 at 0 , and $u^{-4} q(1 / u)$ has a pole of order at most 2 at 0 . Bessel's equation: $z w^{\prime \prime}+w^{\prime}+z w$ therefore has regular singular points everywhere except at $\infty$.
6.7.2. Regular singular points everywhere. If there are but two regular singular points, put them at 0 and $\infty$. Then, $p(z)=A / z$ and $q=B / z^{2}$. Of course, in this case the global monodromy group is given by the action on $\left(z^{\alpha_{1}}, z^{\alpha_{2}}\right.$ if the solutions of the indicial equation are distinct, and by $\left(z^{\alpha}, \ln (z) z^{\alpha}\right)$ if they are not.

Now suppose there are exactly three regular singular points, and place them at $0,1, \infty$.
6.7.3. Application to integrals. Let $a_{1}, \ldots, a_{r}$ be any complex numbers and assume we have chosen $g_{i}(z)$ to be a branch of $\left(z-z_{i}\right)^{-a_{i}}$ in a neighborhood of $z_{0}$. Then, $\prod_{i=1}^{r} g_{i} d z \stackrel{\text { def }}{=} \omega_{\boldsymbol{z}}$ is a differential 1-form in a neighborhood of $z_{0}$ that analytically continues along any $\gamma \in \Pi_{1}\left(U_{\boldsymbol{z}}, z_{0}, z^{\prime}\right)$ to give a differential 1-form in the neighborhood of $z^{\prime}$. Then, $\int_{\gamma} \stackrel{\text { def }}{=} I_{\boldsymbol{a}}(\gamma)$ makes sense. The monodromy theorem says this depends only on the homotopy class of the path. Further, if $\gamma_{i, j}$ is a piecewise differentiable path on $\mathbb{P}_{z}^{1}$ from $z_{i}$ to $z_{j}$, then even $I_{\boldsymbol{a}}\left(\gamma_{i, j}\right)$ makes sense by taking
the integral to be $\lim \epsilon \mapsto 0 I_{\boldsymbol{a}}\left(I_{\boldsymbol{a}}\left(\delta_{e, \epsilon}\right)-\delta_{b, \epsilon}\right)$ with with $\delta_{b, \epsilon}$ going from $z_{0}$ to $\gamma(\epsilon)$ in $U_{\boldsymbol{z}}$ and $\delta_{e, \epsilon}$ being the composite of $\delta_{b, \epsilon}$ and $\gamma_{[\text {epsilon,1- }] \text {. }}$.

Notice, unless all the $a_{i} \mathrm{~s}$ are the same, this is placing an ordering on them to form $\omega_{\boldsymbol{a}}$. Suppose $\varphi: U_{\boldsymbol{z}} \rightarrow U_{\boldsymbol{z}}$ is a diffeomorphism that fixes each of the points in $\boldsymbol{z}$. To see the effect of the diffeomorphisms it suffices to take the case $\left(z_{1}, z_{2}, z_{2}\right)=(0,1, \infty)$, and let $z_{4}=\lambda$.

## 7. Abel's contributions and modular curves

This culminated in Abel's beautiful characterization of analytic functions, the precise form that the Riemann-Roch theorem takes, on a complex torus. Fine though it is in its classical form, we now take a nontraditional view that soon will reveal modular curves and some of their far-reaching generalizations. Our approach motivates the whole theory of moduli spaces of Riemann surfaces. We start with finding the essential parameters that characterize these integrals up to algebraic transformations. Functions that naturally live on a (1-dimensional) complex torus are algebraically related to inverses of special cases of functions given by integrals in (6.7). Compatible with other lessons about coordinates, the functions that live on one complex torus cannot be elementarily related to those on an analytically nonisomorphic torus.

RETURN If you have a function that separates points at one fiber, why does it separate at all but finitely many fibers? Answer: Suppose $g: X \rightarrow \mathbb{P}_{w}^{1}$, and for infinitely many $z^{\prime} \in \mathbb{P}_{z}^{1}$, there are $w_{1}\left(z^{\prime}\right), w_{2}\left(z^{\prime}\right) \in \varphi^{-1}\left(z^{\prime}\right)$ such that $g\left(w_{1}\left(z^{\prime}\right)\right)=$ $g\left(w_{2}\left(z^{\prime}\right)\right)$.
7.1. Integrals of primitives. In $\S 6.7$ we have the special case appearing in Chap. 2 (6.6): $\int_{\gamma} \frac{d z}{\left(z^{3}+c z+d\right)^{\frac{1}{2}}}$. with $c, d \in \mathbb{C}$. Use the notation $\mathcal{A}_{h}\left(U_{\boldsymbol{z}}\right)$ for the analytic continuations of $h$ based at $z_{0}$ (Chap. $2 \S 4.5$ ).

Lemma 7.1. Suppose $h(z) \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)$ is an algebraic function of degree $n$. Form $H(z)_{\lambda}=\operatorname{Int}(h(z))_{\gamma}$ to mean the analytic continuation of a primitive $H(z)$ for $h(z)$ along $\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ (Chap. 2 §4.3). Define $G_{H}$ to be the monodromy group of this action: The collection of permutations on $\mathcal{A}_{H}\left(U_{z}\right)$. Then, there is a natural map $G_{H} \rightarrow G_{h}$ given by taking the derivative: $\frac{d H_{\gamma}}{d z} \mapsto h_{\gamma}$. Denote by $G_{h}(1)$ the stabilizer of $h$ in $G_{h}$ and let $L_{h}$ be the pullback of $G_{h}(1)$ in $G_{H}$. Then, $L_{h}$ is an abelian group and, $L_{h}=L_{h_{\gamma}}$ for $\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. This identifies $G_{H}$ as a subgroup of the wreath product $L_{h} \prec G_{f}=\left(L_{h}\right)^{n} \times{ }^{s} G_{f}$ (§8.4).

Proof. The effect of operating by $\gamma^{*} \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ on $\left\{\operatorname{Int}(h(z))_{\gamma}\right\}_{\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)}$ maps to the action on $\left\{h(z)_{\gamma}\right\}_{\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)}$ by taking the derivative: $h_{\gamma}=\frac{d H_{\gamma}}{d z}$. Further, suppose $h_{\gamma}=h: \gamma$ is in the pullback of $G_{h}(1)$. Then, $H_{\gamma}$ is a primitive of $h_{\gamma}=h$, and so $H_{\gamma}(z)=H(z)+c_{\gamma}$ with $c_{\gamma} \in \mathbb{C}$. Since $c_{\gamma}$ is $\int_{\gamma} h(z) d z$, for $\gamma_{1}, \gamma_{2} \in \operatorname{ker}(\Gamma), c_{\gamma_{1}}+c_{\gamma_{2}}=c_{g a m m a_{1}+\gamma_{2}}$. Use the notation $L_{h}$ for this abelian group of integration constants.

Consider a conjugate $h_{\gamma}$ of $h$ given by $\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$. Then, $\int_{\gamma^{*}} h_{\gamma} d z$ is the same as $\int_{\gamma \cdot \gamma^{*}} h d z$. As $[\gamma] \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)=\pi_{1}\left(U \boldsymbol{z}, z_{0}\right)$, the final set of integrals around closed paths is the same. Finally, let us embed $G_{H}$ in $L_{h} \prec G_{f}=\left(L_{h}\right)^{n} \times{ }^{s} G_{f}$. Let $\left\{h_{i}\right\}_{i=1}^{n}$ be represent the distinct elements of $\mathcal{A}_{h}\left(U_{\boldsymbol{z}}\right)$. It has a map to $G_{h}$ and the kernel of that map fixes each of $L_{h_{\gamma}}$, each of which we can identify. The argument is now the same as in Chap. 38.14 that this gives a wreath product.

Example 7.2 (Branch cycles in $S_{n} \backslash A_{n}$ ). Assume each class in $\mathbf{C}$ is in $S_{n} \backslash A_{n}$. Let $\boldsymbol{p} \in \mathcal{H}(G, \mathbf{C})^{\mathrm{in}, \mathrm{rd}}$ lie over $\boldsymbol{z}$. Take the regular representation of $G$ as giving a map $G \rightarrow S_{|G|}$. The cover $\varphi_{\boldsymbol{p}}: X_{\boldsymbol{p}} \rightarrow \mathbb{P}_{z}^{1}$ naturally factors through $E_{\boldsymbol{z}} \rightarrow \mathbb{P}_{z}^{1}$ : Quotient $X_{p}$ by $G \cap A_{|G|}$. (This works for any even $r ; E_{z}$ is then hyperelliptic.)

Suppose $r=4$ and $\boldsymbol{g} \in \operatorname{Ni}(G, \mathbf{C})$ with $G \leq A_{n}$. Choose $h_{1}, h_{2} \in S_{n} \backslash A_{n}$. Then, $\left(h_{1} g_{1}, g_{2} h_{2}, h_{2}^{-1} g_{3}, g_{4} h_{1}^{-1}\right)$ satisfies the product-one condition. It produces a Nielsen class (for some new group) with moduli problem directly recognizing the $j$-line as parameterizing elliptic curves.

Write out the covering property for $w \mapsto h(w)$ and $y \mapsto y^{2}$. Then, take the fiber product over $\mathbb{P}_{z}^{1}$. Use the path lifting property for each separately, then join them and ask how you can understand this to be like an exponential map. Problem: The set is more complicated, and not looking like an open subset of the plane. Then, consider the integral $F(w)=\int_{\gamma} \frac{d w}{\sqrt{h(w)}}$ in analogy with the discussion right after Prop. 3.5.

Let $H(w)$ be the inverse function of $F(w): \quad G(F(w))=w$. Conclude this uniformization of

$$
\begin{equation*}
E_{g}=\left\{(w, y) \mid y^{2}-h(w)\right\}: z \mapsto\left(G(z), G^{\prime}(z)\right) \tag{7.1}
\end{equation*}
$$

Lemma 7.3. Periods of $H(w)$ and why they have rank two. FINISH THIS
7.2. Starting Abel's Theorem. We start with an interpretation of Abel's the basic problem which was about the nature of antiderivatives. Analyze elementary antiderivatives, like the watershed example $\int \frac{d x}{\sqrt{x^{3}+a x+b}}$. Specifically, what is the dependence of these antiderivatives on the parameters $a$ and $b$ ?

Here $m(z, w)=w^{2}-\left(z^{3}+a z+b\right)$. Let $\boldsymbol{z}$ be $\infty$ together with the zeros of $z^{3}+a z+$ $b$. Write $G(z)=\frac{1}{\sqrt{z^{3}+a z+b}}$, a branch of this square root defined in a neighborhood of $z_{0} \notin \boldsymbol{z}$ (Chap. $2 \S 8.2$ ). Consider $F(z)=F_{a, b}(z)$, an antiderivative of $G(z)$, locally. As in Chap. $2 \S 4.3$, it has analytic continuations along elements of $\Pi\left(U_{\boldsymbol{z}}, z_{0}\right)$ and it depends only on the homotopy class of the path. These continuations produce an abelian group of periods (Chap. 2). Chap. 4 shows the group is $\mathbb{Z} \times \mathbb{Z}$. Further, its fit with the analytic continuations of $G(z)$ appears in the semidirect product $\mathbb{Z} \times \mathbb{Z} \times{ }^{s}\{ \pm 1\}(\S 8)$. Let $D_{n}$ be the dihedral group of order $2 n$.
7.2.1. The antiderivatives: when are they equivalent by substitutions? This was Abel's version of the problem. He already had experience with showing functions might be new: The new functions that produced solutions to the general quintic. His formulation of old included taking compositions of rational functions, multiplication by constants, various other functions regarded as known, and the conceptual addition of
(7.2) functional inverse.

We have learned in Chap. 2 that there really are but two elementary functions if you use (7.2): as functions of a complex variable they are $z$ and $\log (z)$. Also, you may profitably consider the integral locally as a function of z by regarding it as as an antiderivative in $z$. By allowing this you agree that you will stay within elementary functions. Still, we sceptically scrutinize the antiderivative operation, for we are also asking when that applied to elementary functions takes us out of them.

Example 1: . Example: $z^{1 / n}=e^{z / n}$.
7.3. Substitutions by elementary functions. We may take a branch of inverse of a locally one-one analytic function Chap. $2 \S 6.1$. We did that in $\S 6.6$ and found that Substitutions in the antiderivative variable $w=F_{a, b}(z)$ become compositions of the inverse function $F_{a, b}^{-1}(w)=z$. So Abel took a more restricted version of the problem by equivalencing by rational function substitutions and derivatives in w.

Problem 7.4. When is $L_{a, b, c} \stackrel{\text { def }}{=} \mathbb{C}\left(F_{a, b, c}^{-1}(w), \frac{d F_{a, b, c}^{-1}}{d w}, \ldots\right)$ isomorphic to $L_{a^{\prime}, b^{\prime}, c^{\prime}}$ ? Only one derivative needed by Abel's use of the chain rule. The fields $L_{a, b, c}$ are called function fields. We can describe every element of $L_{a, b, c}$.
7.4. Explicit functions for dihedral covers. We have applied the Existence Theorem in $\S 4.3$ to see how many different involution dihedral covers of $U_{\boldsymbol{z}}$ we get when $r=4$. Now we consider another approach to dihedral covers. The two approaches are complementary. In $\S 4.3$ we have a degree $p^{k+1}$ function to $\mathbb{P}_{u}^{1}$ (to retain the classical $z$ variable for Abel's approach, we switch the variable in Riemann's Existence Theorem to $u$ ). In Abel's approach, the connection to the $u$ variable is nonobvious.

Let $L\left(\omega_{1}, \omega_{2}\right) \stackrel{\text { def }}{=} L_{\boldsymbol{\omega}}=\left\{m_{1} \omega_{1}+m_{2} \omega_{2} \mid m_{1}, m_{2} \in \mathbb{Z}\right\}$ be a lattice in $\mathbb{C}$. First observation: With no loss we may assume $\tau=\omega_{2} / \omega_{1}$ has positive imaginary part. We say $\tau$ lies in the upper half plane. Elliptic functions $f(z)$ on $L_{\omega}=L$ have all elements of $L$ as periods. In the next lemma we consider $f(z)$ in a fundamental domain for $L$.

Theorem 7.5. Residues of an elliptic function sum to 0 . Conclude: A nonconstant elliptic function has as many zeros as it has poles. Finally, if $a_{1}, \ldots, a_{n}$ are the zeros of $f$, and $b_{1}, \ldots, b_{n}$ are its poles, then $\sum_{j=1}^{n} a_{j}-\sum_{j=1}^{n} b_{j} \in L$.

Proof. Take $B$ as the clockwise boundary of the fundamental domain. Let $f(z)$ be an elliptic function for $L$. The sum of the residues is therefore $\frac{1}{2 \pi i} \int_{B} f(z) d z$. Integrals in opposite directions on opposite sides of the parallelogram cancel. Therefore, this is 0 . Apply this to $f^{\prime}(z) / f(z)$, also an elliptic function. For the last sentence, compute $\frac{1}{2 \pi i} \int_{B} z f^{\prime}(z) / f(z) d z$ carefully.

Here is the Weierstrass $\wp$-function:

$$
\begin{equation*}
\wp\left(z ; \omega_{1}, \omega_{2}\right)=\frac{1}{z^{2}}+\sum_{\omega \neq 0}\left(\frac{1}{(z-\omega)^{2}}-\frac{1}{(\omega)^{2}}\right) \tag{7.3}
\end{equation*}
$$

This actually converges, except at $z=0$. The observation is there exists $u>0$ with $\left|m_{1} \omega_{1}+m_{2} \omega_{2}\right| \geq u\left(\left|m_{1}\right|+\left|m_{2}\right|\right)$. Clearly, $\wp\left(z ; \omega_{1}, \omega_{2}\right)$ is an even function. It's derivative $\frac{d \varsigma}{d z}\left(z ; \omega_{1}, \omega_{2}\right)=\wp^{\prime}(z)=-2 \sum_{\omega} \frac{1}{(z-\omega)^{3}}$ is an odd doubly periodic function. Conclude that $\wp$ is also doubly periodic.
7.5. The function $\zeta(z)$. Since $\wp$ has residues equal to 0 , there exists $\zeta(z)$ whose derivative is $\wp$. We see easily:

$$
\zeta_{( }(z)=\frac{1}{z}+\sum_{\omega \neq 0}\left(\frac{1}{z-\omega}+\frac{1}{\omega}+\frac{z}{\omega^{2}}\right)
$$

Thus, $\zeta\left(z+\omega_{1}\right)=\zeta(z)+\eta_{1}$ and $\zeta\left(z+\omega_{2}\right)=\zeta(z)+\eta_{2}$ for some constants $\eta_{1}$ and $\eta_{2}$. Integration of $\zeta$ around $B$ shows

$$
\begin{equation*}
\eta_{1} \omega_{2}-\eta_{2} \omega_{1}=2 \pi i \tag{7.4}
\end{equation*}
$$

Further, take $\sigma(z)$ to be a function whose logarithmic derivative is $\zeta(z)$ :

$$
\begin{equation*}
\sigma_{\tau}(z)=\sigma(z)=z \prod_{\omega \neq 0}\left(1-\frac{z}{\omega}\right) e^{z / \omega+\frac{1}{2}(z / \omega)^{2}} \tag{7.5}
\end{equation*}
$$

Clearly, $\sigma(z)$ is an odd function. Also, $\sigma\left(z+\omega_{1}\right)=-\sigma(z) e^{\eta_{1}\left(z+\omega_{1} / 2\right)}$, etc. The values $z=\omega_{i} / 2$ already play a special role.

### 7.6. Abel's construction of functions when $g=1$.

Theorem 7.6 (Abel, 1837). An elliptic function with periods $\omega_{1}$ and $\omega_{2}$, zeros $a_{1}, \ldots, a_{n}$ and poles $b_{1}, \ldots, b_{n}$ is a constant times $\prod_{k=1}^{n} \frac{\sigma\left(z-a_{k}\right)}{\sigma\left(z-b_{k}\right)}$ [Ah179, p. 267].

The functions $\wp(z)$ and $\wp^{\prime}(z)$ satisfy this simple equation:

$$
\begin{equation*}
\wp^{\prime}(z)^{2}=4 \wp(z)^{3}-60 G_{2} \wp(z)-140 G_{3} \tag{7.6}
\end{equation*}
$$

with $G_{k}=\sum_{\omega \neq 0} \frac{1}{\omega^{k}}$. The formula is important, but the principle easy: create a doubly periodic function with no poles.

The point: $\left(\wp(z), \wp^{\prime}(z)\right)$ maps complex analytically from the complex torus to the points of the cubic equation (7.6). Thus, the complex torus and the points on the cubic equation represent the same Riemann surface. A most important step is to interpret ideas on the complex torus through (7.6).
7.7. The unique $\theta$ function with odd characteristic. On the complex torus $\mathbb{C} / L_{\boldsymbol{\omega}}$, we may characterize $\wp(z)$ among elliptic functions with a pole of order 2 at the origin, and no other pole up to a linear change of variable. Let $W(z)$ be another, and normalize by linear change so $\wp(z)-W(z)$ has no poles at all and expansion of $W(z)$ around $\infty$ looks like $\frac{1}{z^{2}}+a_{1} z+$ higher terms. Then, $\wp(z)-W(z)$ is identically zero. The same argument shows $\wp(z)$ is an even function: $\wp(z)-\wp(-z) \equiv 0$.

Then, $\zeta(z)$ is the unique antiderivative that is odd: $\zeta(z)=-\zeta(-z)$ for $z \in \mathbb{C}$. Finally, $\sigma_{\tau}(z)$ has the property its logarithmic derivative equals $\zeta(z)$. Thus, it is defined up to multiplicative constant, and it must be an odd function. Conclusion:

THEOREM 7.7. Up to a multiplicative constant, there is a unique odd function $\sigma_{\tau}(z)$ that plays the role of a $\theta$ function on $X_{\tau}=\mathbb{C} / L_{\tau}$. Suppose $D_{\boldsymbol{a}, \boldsymbol{b}}=\sum_{j=1}^{n} a_{j}-$ $\sum_{j=1}^{n} b_{j}$ represents any degree 0 divisor on $X_{\tau}$. Then,

$$
\omega_{\boldsymbol{a}, \boldsymbol{b}}=\sum_{j=1}^{n} \zeta\left(z-a_{i}\right)-\sum_{j=1}^{n} \zeta\left(z-b_{j}\right) d z
$$

is a differential form on $X_{\tau}$ with these properties.
(7.7a) It is a logarithmic differential with polar divisor $D_{a, b}$.
(7.7b) Its periods are pure imaginary.

Proof. We've shown everything except the properties of $\omega_{\boldsymbol{a}, \boldsymbol{b}}$. We have $\zeta(z+$ $\left.\omega_{i}\right)=\zeta(z)+\eta_{i}$ for some complex number $\eta_{i}, i=1,2$. Therefore, $\sum_{j=1}^{n} \zeta\left(z-a_{i}\right)-$ $\sum_{j=1}^{n} \zeta\left(z-b_{j}\right)$ is already a well-defined elliptic function on $X_{\tau}$.

DO WE KNOW ITS PERIODS ARE PURE IMAGINARY?
7.8. Uniformization and the $j$-line. Addition on $\mathbb{C} / L$ is obvious.
7.8.1. The cubic equation for the torus. How do we express this in coordinates from the cubic equation? There is this formula [Ahl79, Ex. 2-7, p. 269]:

$$
\operatorname{det}\left|\begin{array}{ccc}
\wp(z) & \wp^{\prime}(z) & 1  \tag{7.8}\\
\wp(u) & \wp^{\prime}(u) & 1 \\
\wp(u+z) & -\wp^{\prime}(u+z) & 1
\end{array}\right|=0
$$

Consider the zeros of the cubic on the right side of (7.6). Call these $e_{1}, e_{2}, e_{3}$. Since $\wp^{\prime}(z)$ is an odd function, we find these $e_{i}$ 's are $\wp\left(\frac{\omega_{1}}{2}\right), \wp\left(\frac{\omega_{2}}{2}\right)$ and $\wp\left(\frac{\omega_{1}+\omega_{2}}{2}\right)$. Since $\wp(z)$ assumes each value with multiplicity 2 (on a fundamental domain), the $e_{i}$ 's are distinct. Thus, $\lambda(\tau)=\frac{e_{3}-e_{2}}{e_{1}-e_{2}}$ is a well-defined function of $\tau$.

We want a well-defined function of $L$ : equivalently, a function preserved under unimodular transformations of $\tau$. Notice these properties of $\lambda(\tau)$.
(7.9a) Unimodular transformations of $\left(\omega_{1}, \omega_{2}\right)$, represented by

$$
\left(\begin{array}{ll}
a & b  \tag{7.10}\\
c & d
\end{array}\right) \equiv\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right) \bmod 2
$$

leave $\lambda(\tau)$ invariant.
(7.10b) Unimodular transformations don't satisfying (7.9a) don't leave $\lambda(\tau)$ invariant.
(7.10c) $\lambda(\tau)$ takes all values except 0 or 1 for $\tau$ in the upper half plane.
(7.10d) Each value $\lambda(\tau)$ takes on, it assume locally with multiplicity 1.

Define $j(\tau)$ to be $\frac{4}{27}\left(\frac{e_{1} e_{2}+e_{2} e_{3}+e_{3} e_{1}}{\left(e_{1}-e_{2}\right)^{2}\left(e_{2}-e_{3}\right)^{2}\left(e_{3}-e_{1}\right)^{2}}\right)$. Consider the subfield of $\mathbb{C}(\lambda(\tau))$ invariant under the full unimodular group (as in (7.9b)). Then, $j(\tau)$ generates this field [Ahl79, Ex. on p. 274].
7.8.2. The function $\lambda(\tau)$. Finally, from (7.9c) and (7.9d) we have Picard's big theorem.

THEOREM 7.8. Suppose $f(z)$ is entire and it omits at least two values. Then, $f(z)$ is constant.

Proof. Suppose $f(z)$ omits at least two values. With no loss, take these to be 0,1 . For any other value of $f$, say, $f\left(z_{0}\right)$, (3) shows there is $\tau_{0}$ in the upper half plane with $\lambda\left(\tau_{0}\right)=f\left(z_{0}\right)$. Apply (4). The complex version of the implicit function theorem gives an analytic function $h(z)$ defined in a neighborhood of $z_{0}$. It has these properties: $\lambda(h(z))=f(z)$ for $z$ close to $z_{0}$; and $h\left(z_{0}\right)=\tau_{0}$. Since the complex plane is simply connected, the monodromy theorem says $h$, defined locally for each $z_{0}$, is the restriction of one entire function $H(z)$. Note, however, $H(z)$ takes its values in the upper half plane. This violates the maximum modulus principle: $e^{i H(z)}$ has absolute value at most 1 .

### 7.8.3. Uniformization from the Existence Theorem.

Corollary 7.9. Uniformization of $U_{\boldsymbol{z}}, r \geq 3$, by a disk following the deformation talk at MSRI.
7.8.4. Involution dihedral covers again. Let $f(z)=(z-1) \cdots(z-r)$, let $D$ be an open connected set in $\mathbb{C}-\{1, \ldots, r\}$ and let $A_{1}, \ldots, A_{k}$ be the connected components of $\mathbb{P}_{z}^{1} \backslash D$. Suppose $A_{k}$ is the component that contains $\infty$. Now take $r=3$. Assume that $\gamma$ is a path defined on $[0,1]$. We can define $g_{\gamma}(t)$, a branch of $\log ((z-1)(z-2)(z-3))$ along $\gamma$, by the following formula:

$$
\frac{1}{2 \pi i}\left(\int_{\gamma} \frac{d z}{z-1}+\int_{\gamma} \frac{d z}{z-2}+\int_{\gamma} \frac{d z}{z-3}\right)
$$

Let $a_{i}(\gamma)$ be the index of $\gamma$ with respect to $i=1,2,3$. Thus $e^{g_{\gamma}(t) / n}$ is independent of $\gamma$ if and only if $a_{1}+a_{2}+a_{3} \equiv 0 \bmod n$ for every closed path in $D$. For example, if $n=3$, then $g(z)$ exists if $D=\mathbb{C} \backslash[1,3]$, but it doesn't exist if $D=\mathbb{C} \backslash\{1\} \cup[2,3]$.

Consider closed paths $\gamma$ based at 0 in the case that $n=2$ and $D=\mathbb{C} \backslash\{1,2,3\}$. If a closed path has index 1 with respect to 1 and index 0 with respect to 2 and 3 , we must have $\int_{\gamma} f(z)^{\frac{1}{2}} d z \neq 0$. Otherwise there would be an analytic function $F(z)$ defined on $D=\mathbb{C} \backslash[1,3]$ such that $\frac{d F}{d z}$ would be a branch of $(f(z))^{\frac{1}{2}}$, contrary to the above argument.

## 8. Algebraic coordinates

8.1. Points about algebraic varieties. 2. Segre embedding of two projective varieties. P. 66

For prevariety, I need to talk about the patching maps.
3. Variety: p. 68: A variety is a prevariety $X$ (finite cover of affine varieties irreducible and have a coordinate ring that is an integral domain) together with the property that if $f: Y \rightarrow X$ and $g: Y \rightarrow X, Y \times_{X} Y \cap \Delta$ is closed in $Y \times Y$. The case when $f$ and $g$ are the projections is just the case $\Delta_{X}$ is closed in $X \times X$. The general case follows from this case because it is also factors through $Y \times:(f, g) \rightarrow X \times X$ and the set in question is just the pullback of $\Delta_{X}$, which being closed has its inverse closed.

Prop. 5 (p. 71): If $X$ is a prevariety and any two points are in an affine open piece, then X is a variety. Reason: Get Hausdorff in the affine pieces, and so this applies to projective varieties.

Prop. 6: X a variety, $U=\operatorname{Spec}(R)$ and $V=\operatorname{Spec}(S)$, then $U \cap V$ has coordinate ring $R \cdot S$ with the composite in the function field. $K(X)$,
p. 86: Local set-theoretic complete intersection: If $Z$ has codimension $r$ there is an open set $U$ of any given $z \in Z$ so that there are $r$ functions in the coordinate ring that set-theoretically describes the set $Z$ around $z$.
p. 89: Suppose a morphism $\mathbb{P}^{n} \xrightarrow{f} \mathbb{P}^{m}$ with closed image ( $\S 9$ says this is always true), then the image is a point or has dimension $n$. Reason: If the image has dimension less than $n$, then we can find $n$ hypersurfaces $H_{1}, \ldots, H_{n}$ in $\mathbb{P}^{m}$ defined by homogeneous polynomials so that the intersection of these with the image is trivial. Pull these back to the $\mathbb{P}^{n}$ to conclude you have $n$ hyperplanes in $\mathbb{P}^{n}$ that intersect trivially.

Equations may not give us what we expect, so it is necessary to have simple criteria to assure that what is essential from some particular type of equation compactification does not depend on the particular choice of compactification.

### 8.2. Completion of the fundamental group.

8.3. Functions on the universal covering space. We follow the treatment of [BL92, Appendix].

Let $X$ be a complex manifold, $\tilde{X}$ its universal covering space. Then, $\pi_{1}(X)$ acts on $\tilde{X}$ on the left and on $\mathbb{C}(\tilde{X})$ on the right. Identify $H^{1}\left(\pi_{1}(X), \mathbb{C}(\tilde{X})^{*}\right)$ with factors of automorphy. For $f \in H^{1}\left(\pi_{1}(X), \mathbb{C}(\tilde{X})^{*}\right), f\left(g_{1} g_{2}\right)=\left(f\left(g_{1}\right)\right) g_{2} f\left(g_{2}\right)$.

Lemma 8.1. Each element of $H^{1}\left(\pi_{1}(X), \mathbb{C}(\tilde{X})\right)$ defines a line bundle on $X$.
Proof. Note: $\left.f \in H^{1}\left(\pi_{1}(X)\right), \mathbb{C}(\tilde{X})^{*}\right)$ determines a function $h$ from $\pi_{1}(X) \times$ $\tilde{X} \rightarrow \mathbb{C}^{*}$ by $h(g, x)=f(g)(x)$, and $h\left(g_{1} g_{2}, x\right)=h\left(g_{1}, g_{2}(x)\right) h\left(g_{2}, x\right)$. Define an
action of $\pi_{1}(X)$ on $\tilde{X} \times \mathbb{C} b y g(x, t)=(g(x), f(g(x)) t)$. The action is free because $f(g(x))$ is nonvanishing. Also, the quotient of this action is a bundle $\mathcal{L}_{f}$.

To obtain a trivialization of $\mathcal{L}_{f}$, choose an open covering $\mathcal{U}_{\alpha \in I}$ such that on each open $U$ of the set, the universal cover has the form $U \times \pi_{1}(X)$. Choose a connected $W_{\alpha}$ in $\pi^{-1}\left(U_{\alpha}\right)$ respecting the naturalprojection $\pi: \tilde{X} \rightarrow X$. For $W_{1}$ covering $U_{1}$ and $W_{2}$ covering $U_{2}$, there exists $g_{1,2} \in \pi_{1}(X)$ such that $g_{1,2}$ takes the point $w(x)$ of $W_{1}$ over $x$ to a point of $W_{2}$ over $x \in U_{1} \cap U_{2}$. We take the transition functions as $s_{1,2}(x)=f\left(g_{1,2}\right) w_{2}$. Check the cocycle condition holds.
8.4. Some comparisons with [Har77] and [Mu66]. [Spr57] has no standard notation for transition functions on a complex manifold. Apparantly such a standard notation, such as that for transition functions $\psi_{\beta, \alpha}=\varphi_{\beta} \circ \varphi_{\alpha}^{-1}$ (as in Lem. 2.2) seems to have awaited the sheaf formulation of manifolds. Early places where US students could see this notation include [Gun66, p. 15]. As in Def. 3.6 we tend to record our analytic or meromorphic functions on $U_{\alpha}$ as those on $\varphi_{\alpha}\left(U_{\alpha}\right)$. What, however, is the precise ring of analytic functions on the overlap $U_{\alpha} \cap U_{\beta}$, of two coordinate charts $\left(U_{\alpha}, \varphi_{\alpha}\right)$. Using the ring $\mathcal{H}_{\mathcal{U}}\left(U_{\alpha}\right)=\left\{f \circ \varphi_{\alpha} \mid f \in \mathcal{H}\left(\varphi_{\alpha}\left(U_{\alpha}\right)\right)\right\}$ exactly solves this problem of having a ring directly defined on the open set $U_{\alpha}$. For each open set $U \subset X$ we define the analytic functions on $U$ by sying they are functions on $U$ whose restriction to each $U \cap U_{\alpha}$ is the pullback by $\varphi_{\alpha}$ of a function analytic on $\varphi_{\alpha}\left(U \cap U_{\alpha}\right)$. This gives us the essential properties of the sheaf of holomorphic functions, $\mathcal{H}_{X}$.

For example, in [Mu66, p. 33] the maps defining a presheaf for a chain of open sets $U_{1} \subset U_{2} \subset U_{3}$ are called restriction and the cocycle condition $\psi_{\gamma, \beta} \circ \psi_{\beta, \alpha}=\psi_{\gamma, \alpha}$ for transition functions takes the form $\operatorname{rest} U_{2}, U_{1} \circ \operatorname{rest} U_{3}, U_{2}=\operatorname{rest} U_{3}, U_{1}$. That is, restriction is from the bigger set to the smaller, and the direction is said this way in the subsecriptions. This is at the ring level. At the point level it goes the other way.
8.5. The lemmas of Noether and Chow. In response to your questions about $U_{r}=\mathbb{P}^{r} \backslash D_{r}$, here are the relevant answers. Include the definition of normal variety.

PART A. $D_{r}$ is an algebraic variety: There are two ways to see that. First way: There is an equation that describes it. The start of that equation comes from its basic description.

Write $\mathbb{P}^{r}$ as the collection of nonzero polynomials (up to multiplication by a nonzero constant) of degree at most $r$. If a polynomial has degree less than $r$, regard $\infty$ as being a zero of it. Then, $D_{r}$ is the locus of polynomials with two or more zeros that are equal. If you write the condition for a polynomial $f(z)$ to have repeated roots, it is that the gcd of $f$ and $f^{\prime}$ have a common factor. The Euclidean algorithm allows you to write the gcd of $f$ and $f^{\prime}$ as $h f+g f^{\prime}$ for some relatively prime polynomials $h$ and $g$. Finally, the condition that there are such relatively prime $h$ and $g$ giving a nontrivial linear combination can be written as a condition on a matrix (formed from the coefficients of $f$ and $f^{\prime}$ ) having a nonzero determinant. So, $D_{r}$ is a polynomial equation in the coefficients of $f$ formed from setting a determinant expression to 0 .

Here is how that relates to Chap. 3 and Chap. 4. In Chap. 3 I talk about the concept of $P^{1}$-algebraic (Def. 3.10). I've always wanted to put that discussion in a book, for $\mathbb{P}^{1}$-algebraic is a naive definition of algebraic that is close to being
algebraic but not quite. The definition - which you would never have seen before - is equivalent to a compact manifold having an embedding in $\left(\mathbb{P}^{1}\right)^{N}$ for some $N$. The definition of algebraic is that there is an embedding in $\mathbb{P}^{N}$ and $\S 4.1 .2$ discusses the difference between $\mathbb{P}^{1}$-algebraic and being algebraic. I've left the full motivation - as I say in $\S 4.1 .2$ - for the definition of algebraic for Chap. 4. See in Part C why I wait until Chap. 4 to do this. Also, I mention a second more abstract reason why $D_{r}$ is algebraic there.

PART B. The complex structure on $U_{r}$ is from the complex structure on $\mathbb{P}^{r}$ (§4.2.2). An open subset of a complex manifold is a complex manifold. PART C. Your question "A big question we have is how Hurwitz spaces are algebraic varieties."

ANSWER: How about an easier question? Suppose you have a cover $X$ of $\mathbb{P}^{1} \backslash \boldsymbol{z}=U_{\boldsymbol{z}}$ (Riemann sphere minus a finite set of points). Why is $X$ algebraic? Here is one answer: Because you can compactify $X$ in a unique way to a compact Riemann surface (this is at the beginning of Chap. 4), and all compact Riemann surfaces are algebraic. Better yet, you can see its algebraic structure by relating it to the algebraic structure of $\mathbb{P}^{1}$. All of this is Riemann's Existence Theorem, the topic of Chap. 4. This is not an easy theorem. Understanding the significance and how to use this are the main topics of the book. The question you are asking is a generalization of this: Any cover of $\mathbb{P}^{r}$ minus an algebraic subset is an open subset of an algebraic variety. The result is due to Grauert and Remmert and it plays a big role in my first big paper on relating the Inverse Galois Problem to Hurwitz spaces in 1977. So, you must wait to Chap. 4 for a full discussion.

The simple relation between $\left(\mathbb{P}^{1}\right)^{r}$ and $\mathbb{P}^{r}$ is that the latter is the quotient by an $S_{r}$ action of the former. At first I intended that to be in Chap. 3. Late in the game I saw it made a more coherent discussion to be in Chap. 4. Once, however, you know that, then you have that $D_{r}$ is the quotient by of the fat diagonal on $\left(\mathbb{P}^{1}\right)^{r}$. A valuable lemma is this. Let $X$ be an algebraic variety and let $G$ be a finite group acting algebraically on $X$. Then $X / G$ is also an algebraic variety. Indeed, showing the value of having this fact is what is my motivation to someone reading these chapters for expanding beyond $\mathbb{P}^{1}$-algebraic to the definition involving $\mathbb{P}^{N}$.

Additional Point: Chow's Lemma - which I'm sure Mirroslav will do in his class - says that any complex analytic subset of an algebraic variety is also algebraic. This is NOT deep, though conceptually very valuable.
8.6. The Branch Cycle Lemma. If a cover $\varphi: X \rightarrow \mathbb{P}_{z}^{1}$ in this family has definition field $K$, then the Galois closure of the cover has Galois (arithmetic monodromy) group a subgroup of $N_{S_{n}}(G, \mathbf{C})$.
8.6.1. BRANCH CYCLE SETUP.

Question $8.2\left((G, \mathbf{C})\right.$-cover?). Quest. A: Does $\boldsymbol{z}$ and $\varphi_{\boldsymbol{g}}: X_{\boldsymbol{g}} \rightarrow \mathbb{P}_{z}^{1}$ exist with branch points $\boldsymbol{z},\langle\boldsymbol{g}\rangle=G$ and $\boldsymbol{g} \in \mathbf{C}(\boldsymbol{g} \in \mathrm{Ni}(G, \mathbf{C}))$ over $\mathbb{Q}$ ?

Quest. B: A $(G, \mathbf{C})$-Galois cover over $\mathbb{Q}$ ?
Q. A or B requires $\boldsymbol{z}$ is a $\mathbb{Q}$ set. For $z_{0} \in \mathbb{P}_{z}^{1}(\mathbb{Q}): \sigma \in G_{\mathbb{Q}}$ acts on $\gamma \in \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)$ through what the result does to $f \in \mathcal{E}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {alg }}$ :

$$
f \mapsto f_{\sigma^{-1} \circ \gamma \circ \sigma}=f_{\gamma^{\sigma}} .
$$

Extend $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right) \rightarrow G$ to $\psi_{\boldsymbol{z}^{\prime}, z_{0}}: \pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right)^{\text {alg }} \rightarrow G$. As a profinite group, $\pi_{1}\left(U_{\boldsymbol{z}^{\prime}}, z_{0}\right)^{\text {alg }}$ is also free on $r$ (topological) generators modulo a product-one relation.

Notation: $\sigma \in G_{K}$ maps to $n_{\sigma} \in \hat{\mathbb{Z}}^{*}=G\left(\mathbb{Q}^{\text {cyc }} / \mathbb{Q}\right)$. For any $K \leq \mathbb{C}$ (like $\mathbb{R}$ or $\left.\mathbb{Q}_{p}\right)$ consider

$$
\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right)^{\text {ar }} \stackrel{\text { def }}{=} \pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right) \times{ }^{s} G_{K}
$$

BRANCH CYCLE ARGUMENT: [Fri73] or [Fri77]
Lemma 8.3 (Branch Cycle Argument). For $Q$. B: $\psi$ must extend to a homomorphism $\psi: \pi_{1}\left(U_{\boldsymbol{z}, z_{0}}\right)^{\text {ar }} \rightarrow G$.

Find $\pi \in S_{r}$ to satisfy $z_{i}^{\sigma}=z_{(i) \pi}$. Then, affirmative to $Q$. $B$ requires

$$
\begin{equation*}
C_{(i) \pi}^{n_{\sigma}}=C_{i}, i=1, \ldots, r \tag{8.1}
\end{equation*}
$$

Affirmative for $Q$. A requires a choice of Galois closure group $\hat{G}$ between $G$ and $N_{S_{n}}(G, \mathbf{C})$. For each such $\hat{G}$, replace (8.1) by

$$
g_{\sigma} C_{(i) \pi}^{n_{\sigma}} g_{\sigma}^{-1}=C_{i}, i=1, \ldots, r, \text { for some } g_{\sigma} \in \hat{G}
$$

Let $G=A_{5}, \mathrm{C}_{5}^{+}$the class of (12345), $\mathrm{C}_{5}^{-}$the class of (13524), $\mathrm{C}_{3}$ the class of 3 -cycles. Example four branch point covers:
(8.1a) C $_{5_{+}^{2} 3^{2}}$ : No for Q. B, yes for Q. A.
(8.1b) $\mathbf{C}_{5_{+5-3}}$ : Yes for Q. A and B.
(8.1c) $\mathbf{C}_{5_{+}^{2} 5_{-}^{2}}$ : Yes for Q. A and B.
8.6.2. Field of moduli. Ingredients for a field of moduli (over $\mathbb{Q}$ ):

- A collection of algebraic objects $\mathcal{P}$ over $\overline{\mathbb{Q}}$ invariant under $G_{\mathbb{Q}}$.
- A $G_{\mathbb{Q}}$ equivalence relation $\mathcal{E}$ on $\mathcal{P}$.

For $\boldsymbol{p} \in \mathcal{P}$ let $\mathcal{E}_{\boldsymbol{p}}$ be its equivalence class: $H_{\boldsymbol{p}}^{\mathcal{E}}$ is the subgroup of $G_{\mathbb{Q}}$ stabilizing $\mathcal{E}_{\boldsymbol{p}}$; $K_{\boldsymbol{p}}^{\mathcal{E}}$ is the fixed field of $H_{\boldsymbol{p}}^{\mathcal{E}}$.

FACT 8.4. For each $\boldsymbol{p}^{\prime} \in \mathcal{E}_{\boldsymbol{p}}$, any field of definition of $\boldsymbol{p}^{\prime}$ contains $K_{\boldsymbol{p}}^{\mathcal{E}}$.
From a Nielsen class $\mathrm{Ni}(G, \mathbf{C}, T)^{\text {abs }}$ : $\mathbf{C}$ a rational union of classes and $\boldsymbol{z} \in$ $\mathbb{P}^{r} \backslash D_{r}(\mathbb{Q})$.
(8.1a) $\mathcal{E}_{\boldsymbol{z}}^{\text {abs }}(G, \mathbf{C}, T)$ : s-classes in $\mathrm{Ni}(G, \mathbf{C}, T)$ with branch points $\boldsymbol{z}$.
(8.1b) Same as (a) except $\mathcal{E}_{\boldsymbol{z}}^{\text {abs,rd }}(G, \mathbf{C}, T)$ is w-equivalence classes.
(8.1c) As in (b) except drop $T$; replace ${ }^{\text {abs }}$ by ${ }^{\text {in }}$.
(8.1d) As in the previous, except drop $\boldsymbol{z}$.

## COURSE MODULI:

Question 8.5. For an equivalence class $\mathcal{E}$, and $\boldsymbol{p} \in \mathcal{P}$ how to compute $K_{\boldsymbol{p}}^{\mathcal{E}}$ ? Does $\mathcal{E}_{\boldsymbol{p}}$ contain something over $K_{\boldsymbol{p}}^{\mathcal{E}}$ ? What is the lattice $\mathcal{L}_{\boldsymbol{p}} / K_{\boldsymbol{p}}^{\mathcal{E}}$ of definition fields for elements of $\mathcal{E}_{\boldsymbol{p}}$ ?

Most algebraically defined equivalence classes, including those defined by Nielsen classes have a reasonable (course) moduli space $\mathcal{H}_{\mathcal{E}}$.

FACT 8.6. In the covering space equivalences, $\mathcal{H}_{\mathcal{E}}$ is an affine algebraic variety with a well-defined field of definition - as a moduli space. If $\mathbf{C}$ is a rational union it is $\mathbb{Q}$.

Proposition 8.7. The Branch Cycle Lemma gives the correct field of definition for moduli spaces defined by Nielsen classes.

Meaning: For any family of covers in this equivalence class, you can define a natural map to the moduli space so the space is locally - for the étale topology the pullback of a family over the given space.
8.6.3. How Hurwitz spaces arise. Given $\varphi: X \rightarrow \mathbb{P}^{1}$, branch points $\boldsymbol{z}_{0}$, from RET: What if you wanted equations for it?

- Would you regard $z$ as variable?
- Express equations in $\boldsymbol{z} \in U_{r}=\mathcal{P}^{r} \backslash D_{r}$ ?

Why concentrate on $\boldsymbol{z}_{0}$ ? Want $\varphi_{\boldsymbol{z}}$ for all $\boldsymbol{z}$. If possible, then analytically continuing $\varphi_{\boldsymbol{z}}$ around $\mathcal{P} \in \pi_{1}\left(U_{r}, \boldsymbol{z}_{0}\right)$ returns to $\varphi_{\boldsymbol{z}_{0}}$.

Homotopy class of $\mathcal{P}$ is $Q_{\mathcal{P}} \in H_{r}$ : Hurwitz monodromy group. At end of $\mathcal{P}$ the cover has branch cycle description $(\boldsymbol{g}) Q_{\mathcal{P}}$. Compute with starting classical generators of $\pi_{1}\left(U_{\boldsymbol{z}_{0}}\right)$. So, $\varphi_{\boldsymbol{z}}$ valid for all $\boldsymbol{z}$ requires $(\boldsymbol{g}) Q_{\mathcal{P}}$ be $\boldsymbol{g}$ (modulo conjugation by $G$ or closely related). Check: Is $(\boldsymbol{g}) Q$ essentially $\boldsymbol{g}$ for all $Q \in H_{r}$.

Example 8.8. Consider $\boldsymbol{g}=$

$$
((123),(321),(145),(154)) \in \operatorname{Ni}\left(A_{5}, \mathbf{C}_{3^{4}}\right)
$$

My next talk computes $(\boldsymbol{g}) Q_{\mathcal{P}}$.
Find: $\varphi_{\boldsymbol{z}}$ coefficients have coordinates for nontrivial $U_{r}$ cover. What cover?
Explain why nonsingular conics in $\mathbb{P}^{n}$ are isomorphic to $\mathbb{P}^{n}$ over $\mathbb{C}$ (use Lem. 4.13), but not over $\mathbb{Q}$.

## 9. Using algebraic coordinates and higher monodromy

The complements are topics where we are incomplete with corroborating proofs, giving, however, enough to use them in Chap. 5.

### 9.1. Complements on algebraic coordinates.

9.2. Fundamental groups from branch cycles and higher monodromy.
9.2.1. Computing the fundamental group from branch cycles. Suppose $\varphi: X \rightarrow$ $\mathbb{P}_{z}^{1}$ has $X$ of genus 0 and $z$ as branch points. Then $\S ? ?$ has a procedure for computing classical generators for $X \backslash\left\{\varphi^{-1}(\boldsymbol{z})\right\}$ from those for $U_{\boldsymbol{z}}$. In particular, given $Y \rightarrow$ $X \rightarrow \mathbb{P}_{z}^{1}$, this gives a uniform procedure for computing branch cycles for the cover $Y \rightarrow X$, and thus expresses these branch cycles from those for $\varphi$. A fairly simple procedure allows computing
9.2.2. Action on the fundamental group and homology of a fiber.
9.2.3. Action on periods of integrals.
9.3. Flat bundles and complete reducibility. Check out of [Gr70, $\S 3$ and §4] on complete reducibility in the action on flat bundles.
9.3.1. The genus 0 problem. Suppose $f(w)$ is rational function in $w$. It maps points on the $w$-sphere to the $z$-sphere. The Galois closure group of the splitting field of $f(w)-z$ over $\mathbb{C}(z)$ (monodromy group of $f$ ) is special. That is the gist of the genus 0 problem over $\mathbb{C}$. (The same qualitative statement holds for any fixed genus.) Guralnick and Thompson's original version is this. With finitely many exceptions, the simple composition factors of the monodromy group of such a map must be alternating or cyclic groups. The solution of this left three big problems.
(9.1a) What are the precise monodromy groups, with finitely many exceptions, of indecomposable rational functions? Guralnick's 0-Conjecture: We only get alternating groups, symmetric groups, cyclic and dihedral groups. These should come only with special degree representations.
(9.1b) What groups must one add for rational functions over fields of positive characteristic? Guralnick's $p$-Conjecture: In characteristic $p$ add Chevalley groups over extensions of $\mathbb{F}_{p}$ to alternating and cyclic groups.
(9.1c) Mumford's Question: What function fields in one variable over $\mathbb{C}$ have uniformizations by the Galois closure field of a rational map?
[Fri99] discusses all of these and the history from the Davenport Problem motivations to the complete resolution of the genus 0 problem. Further, Davenport's problem in positive characteristic corroborates Guralnick's inspired p-conjecture. So does the voluminous work of Abhyankar toward his exponent mantra for producing Chevalley groups over $\mathbb{F}_{p}$ from genus 0 covers. A gem from 1995 is Müller's listing of the monodromy groups of polynomials [Mül95].

Like the genus 0 problem, Mumford's question has several forms. For example: Any curve defined by a separated variables equation $f(w)-g(u)$ would have its function field in the composite of splitting fields over $\mathbb{C}(z)$ of functions $f(w)-z$ and $g(u)-z$. That includes all hyperelliptic curves. Directly, the description of modular curves as moduli of genus 0 curves [Fri78] produces elliptic curves from systems of rational function Galois closures, no composite required. Mumford's question has no representation in this volume. It remains untouched in that no function field has been excluded from the genus 0 closure field.

### 9.4. Unramified Frattini covers.

9.4.1. Projectives in the category of profinite group covers. Projectives in the category of profinite group covers of a given group, and here projectives that are Frattini exist. Considered covers with kernel pro-p and kernel elementary p. $\mathcal{C}_{p \infty}(G), \mathcal{C}_{p}(G)$. Other category $K[G]$ modules $\mathcal{C}_{k[G]}(M)$ covering $M$. Porjective profinite objects exist and are unique up to isomorphism. The name for the last is $\mathbb{P}(M)$. Frattini subgroup corresponds to radical of $M$.
$\hat{F}_{2}(2) \times{ }^{s} \mathbb{Z} / 3$ has an embedding of $F_{2} \times{ }^{s} \mathbb{Z} / 3$ in it, though this will not always be the case. Actually the character of the abelian quotient of ${ }_{2} \tilde{A}_{5}$ is not $\mathbb{Q}$ - rational, so this doesn't have a corresponding lattice. Holt and PLeskin:

Gruenberg 1970s: $\mathcal{C}(G) \equiv \mathcal{C}_{\mathbb{F}_{p}[G]}(\omega)$ with $\omega$ Aug. ideal. $\operatorname{Ext}^{1}(\omega, M)=\operatorname{Ext}^{2}(1, M)=$ $H^{2}(G, M) . M \rightarrow N \rightarrow \omega$ is equivalent to $0 \rightarrow M \rightarrow H \rightarrow G \rightarrow 1$. It works with longer sequences too.

Corollary 9.1 (Gaschutz). $M_{0}=\Omega^{2}\left(\mathbb{F}_{p}\right)$.
Draw the conclusion that if $\operatorname{dim}\left(M_{n}\right)>1$, then the exact trivial action part of $G_{n}$ (in $G_{0}$ ) is the $O_{p^{\prime}}(G)$. To see the nonobvious direction, restriction to a $p$-element $g$ in this action, and conclude that $g$ would be acting trivially on a projective.

If $\operatorname{dim}\left(M_{n}\right)>1, Z\left(G_{n}\right) \cap O_{p^{\prime}}(G)=1, G_{n}$ is $p$-perfect $\Longrightarrow Z\left(G_{n+1}\right)=1$, and $G_{n+1}$ is $p$-perfect and $\operatorname{dim}\left(M_{n+1}\right)>1$.

Proposition 9.2 (Griess-Schmid). $\operatorname{dim}\left(M_{n}\right)=1$ if and only if $G$ is $p$-super solvablle with cyclic p-Sylow.
9.5. Equations, coordinates and cryptography. Sometimes the intense studies that engage mathematicians hide the forest behind the trees.

Much of physical science started with the study of relations between theoretically measurable quantities. The applications came because some of these quantities were practically measurable or had in-hand control (like the force and angle of a rocket launch) while the other was more theoretical (like the final destination of
the rocket). Further, there are so many physical systems and so many possible trajectories of different characteristic, that once the idea of analytic coordinates was available, it was inescapable to ask when one system could be transformed to another and how one would test for that. Any one of the following topics has characteristics that resonate into an area of investigation.
(9.2a) Deciding if there are substitutions of variables that change one integral into another.
(9.2b) Finding normal forms for collections of trajectories of a dynamical system.
(9.2c) Creating large isolated systems of interacting particles that emulate logical computer functions.
(9.2d) Finding uniformization situations that allow effectively encoding data in special ways in the set of solutions of equations.

## 10. A piece of the historical record

We start with some historical comments, first personal then some gleaned from [ Ne 81 ] and other sources.
10.1. The career view. When I was a student 1964-1967, I roamed the library at University of Michigan, gleaning from a great store of mathematics, the topics that interested me greatly. These, of course, would have inundated me in my quest to get done, and get done quickly, despite my background as an electrical engineer (who worked for three years in aerospace companies). While my primary interest at the time was in algebraic number theory, there was a quickening surge when I saw the papers of Abel and Galois in the beginning volumes of Crelle's journal. What I immediately saw was that Galois' historical legend was of a different nature than that supported by those of his works Crelle published.

Books on the algebraic theory of curves didn't appeal to me much, for they tended to handle one curve at a time, in intricate detail. There wasn't much at the time on Riemann surfaces in English, except [Spr57]. Further, I saw that the only way to escape the excessive reliance on special forms was through moduli. Moduli formulations of problems were instituted in the early 1960's by Grothendieck and his school. In few 21st century libraries can a student roam through the beginning volumes of Crelle's journal. The difference between being able to order them upon demand and having them there for the curious roaming student (includes curious faculty) is akin to the difference between having snow in California's mountains when you want it, and having God bring it to your home.

In looking back I have the memory that few other topics interested me. Partly that was to prevent the domination of the library volumes calling "Read me, read me!" Those books were real to me, and the urgency of all that mathematics was an insistent plea, that called for my resistence. I let my intuition guide me. So, on certain topics very important to the last fifteen years of my work, I learned little at all, sometimes thinking - as with profinite groups - the subjects weren't that tough. Other topics, related to Lie algebras, differential geometry and partial differential equations seemed to have hordes already committed to them. Even if I avoided their study, others assured they weren't neglected.

These personal comments clarify - though I was commonly praised for my openness to all kinds of mathematics - that my education went through a many year process before my desire to see, and add to, the connectedness between topics
grew. As, however, that happened, the resources for feeing that desire were being removed from the libraries in the late seventies on. Further, the tremendous growth of mathematics was leading mathematicians - and all other academicians -into a personal contest whose criterion for success was dominated by a count of numbers of papers.

Suppose for example, one mathematician in her/his career writes 70 papers, each of roughly 40 pages. Contrast that with another who writes 560 papers, each on average five pages. Without even looking at the papers, my vision of the first writer is that $\mathrm{s} / \mathrm{he}$ built upon projects that defined underappreciated portions of extant areas, or even created new additions to these. I would expect the ratio of theorems to definitions and supporting examples to be slightly smaller in the former, though I would also expect much less starting and stopping from paper to paper to present a common setting. Two other points I would expect, without seeing the papers. The journals, departments and peer reviews of the former career would have many more hesitancies about the evaluation of that career; while the mathematical community would have more focused comments on the nature of that career. It seems significant, however, as to whether these two alternative careers are conscious decisions or merely manifest of the talent and organizational contingencies faced by the mathematician. Are they what you can do, or what you would do?
10.2. Influences on Riemann. We list some significant events in the theory of complex variables.

- A.-L. Cauchy 1789-1857: By 1825 had command of the definite integral between complex limits and presented the Cauchy Integral Theorem.
- P.A. Laurent (1813-1854): In 1843 discovered the Laurent expansion of an analytic function in the deleted neighborhood of an isolated singularity.
- J. Liouville (1809-1882): Formulated many theorems in the theory of elliptic functions.
- V. Puiseux (1820-1883): Investigated analytic continuation in studies of the behaviour of algebraic functions in the neighborhood of one of their branch points.
- C.A. Briot (1817-1882) and J.C. Bouquet (1819-1885): Assembled the previous topics in a series of articles in 1856, bringing them together in the influential [BB1856].
[ Ne 81 ] says it took from 1814 to 1846 to expand a special case of Cauchy's Theorem to integration over a general closed path. Cauchy didn't recognize the significance of the Cauchy-Riemann equation until 1851. Like Weiertrass, Ch. Méray (1835-1911, a Briot and Bouquet disciple), emphasized that continuity is insufficient. As expected, he emphasized the need for a theory based on Taylor series.
B. Riemann (1826-1866) from his thesis 1851 and his 1857 articles on abelian functions, used the Cauchy-Riemann equations exclusively. He basing many of his proofs on potential theory. [Ne81, p. 89]: It was Gauss' (1777-1855) writings that the young Riemann studied wth special zeal. From these he drew signficiant inspirations for his [doctoral] thesis. He wrote his father how he found these papers. What he especially appreciated was Gauss' contributions to conformal mapping using essentially a Dirichlet principle.

According to Betti, Riemann said he got the idea of cuts from conversations with Gauss (1777-1855; see §2.4) [Ne81, p. 90]. Letters of Klein and Schering attest
to Gauss' influence on Riemann's theory of hypergeometric series. This influence came partly from Gauss' papers. Still, it is striking to consider the over 70 year old Gauss sketching plans for such an etherial construction to the very young Riemann.

During his time in Berlin (1847-1849) P.G. Dirichlet (1805-1859), G. Eisenstein (1823-1852) and C.G.J. Jacobi (1804-1851) especially influenced Riemann. He attended Dirichlet's lectures on partial differential equations, and Eisenstein and Jacobi lectures on elliptic integrals. Riemann read Cauchy and Legendre on elliptic functions. [Ne81, p. 91]:

Riemann was suitable, as no other German mathematician then was, to effect the first synthesis of the "French" and "German" approaches in general complex function theory.
His introductory lectures started with these topics: the Cauchy integral formulae; operations on infinite series; the Laurent series; and analytic continuation by power series. [Ne81, p. 92] includes a photocopy of a famous picture on analytic continuation from Riemann's own hand. Picard and Lefschetz both used this picture (from Riemann's collected works) in autobiographies of what influenced critical theorems of theirs. Riemann also lectured on the argument principle, the product represention of an entire fuction with arbitrarily prescribed zeros and the evaluation of definite integrals by residues. His most advanced lectures were from his published papers solving the Jacobi inversion problem (§??).
10.2.1. Competition between Riemann and Weierstrass. [Ne81, p. 93]: K. Weierstrass (1815-1897) himself stressed above all the great influence of N.H. Abel (18021829) on him. At first Weierstrass was an unknown. Only after his 1856 paper on abelian functions did he get his position in Berlin. It was in 1856 that the competition between Riemann and Weierstrass became intense, around the solution of the Jacobi Inversion problem.
[Ne81, p. 93]: May 18 and July 2, 1857, Riemann submitted his two part solution to Jacobi's general inversion problem with these carefully measured words:

Jacobi's inversion problem, which is settled here, has already been solved for the hyperelliptic integrals in several ways through the persistent and regally successful work of Weierstrass, of which a survey has been communicated in Vol. 47 of the Journ für Math. (p. 289). Until now, however, only a part of these investigations has been fully worked out and published (vol. 52, p. 285), namely the part that was outlined in $\S 1$ and $\S 2$ of the earlier paper and in the first half of $\S 3$, pertaining to elliptic functions. Only after the full publication of the promised paper shall we be able to tell to what extent the later parts of the presentation agree with my article not only in results but also in the methods leading to them.
Weierstrass consequently withdrew the 3rd installment of his investigations, which he had in the meantime finished and submitted to the Berlin Academy. He explained this (much later) in his collected works as follows.

I withdrew [the 1857 manuscript] for, a few weeks later, Riemann published an article on the same topic, [...] on entirely different foundations from mine and did not make immediately clear that it agreed completely with mine in its results. The proof
for it entailed some investigations of chiefly an algebraic nature, whose execution was not altotether easy for me ... But after this difficulty was overcome it seemed to me that a thorough going overhaul of my paper was necessary. ... I could only toward the end of 1869 give to the solution of the general inversion problem that form in which I have treated it from then on in my lectures.
10.2.2. Soon after Riemann died. Publicly they seemed to have gotten along [Ne81, p. 95]. Professionally the mutual influence was unquestionably great. It would be entirely conceivable that the general systematic construction of the Weierstrassian function theory, achieved around 1860, could have been inspired by the works of Riemann perstaining to the same set of ideas.
[Ne81, p. 96]: After Riemann's death, Weierstrass attacked his methods quite often and even openly. July 14, 1870 was when he read his now famous critique on the Dirichlet Principle before the Royal Academy in Berlin. Weierstrass showed there did not always exist a function among those admitted [in variation problems] whose expression in question attained the lower bound, as Riemann had assumed. A letter to H. A. Schwarz on Oct. 3, 1875 says:

The more I think about the principles of function theory, the firmer becomes my conviction this must be based on the foundation of algebraic truths, and that, consequently, it is not the right way if instead of building on simple and fundamental algebraic theorems, one appeals to the "transcendental" [by which Riemann has discovered so many of the most important properties of aglebraic functions].
During its heydey (1870-1890), the Weierstrassian school took over nearly every position in Germany. For instance, Schwarz was at Göttingen.
[Ne81, p. 98] asserts it was the Goursat part of Cauchy's theorem that renovated Riemann's approach, starting around 1900. [Ahl79, p. 111] with no precise citation, refers to Goursat's contribution as,

This beautiful proof, which could hardly be simpler is due to É. Goursat, who discovered that the classical hypothesis of a continuous $f^{\prime}(z)$ is redundant.
Curiously, there is precisely one reference [Ahl79, p. 121] in all of [Ahl79]. This is a footnote:

Without use of integration R. L. Plunkett proved the continuity of the derivative (BAMS 65, 1959). E. H. Connell and P. Porcelli proved the existence of all derivatives (BAMS 67, 1961). Both proofs lean on a topological theorem due to G. T. Whyburn.
That unique quote suggests Ahlfors supports the significance of Goursat to Riemann's renovation. Yet, there is a complication in analyzing Neuenschwanden's thesis. Wow would one document that this event turned mathematicians to the geometric/analytic view of Riemann? Historically it seems sensible to investigate the span from [AG1895] to [Wey55] as a shift from genus 1 to higher genus. Yet, that period is clearly insufficient to deal with an aspect of the true shift, from moduli of genus 1 curves (including modular curves) to general moduli. Theories toward the latter include Teichmüller theory (analytic) and geometric invariant theory (algebraic) or expedient precursors of the Hurwitz space approach like the

Schiffer-Spencer deformation theories of varying the complex structure around a single point of a Riemann surface.

I suspect Goursat's theorem is a simple explanation that first year graduate students can follow. Likely, however, serious applications and resonant questions required understanding the variation of structures on a Riemann surfaces with the variation of the surface itself. My experiences are that not only do these issues confound graduate students, often specialists in complex variables struggle with these. Both technically and conceptually handling the hidden monodromy considerations (see Chap. $1 \S 5.4 .3$ ) is a tough topic. One practical approach to this topic, hidden in combinatorial actions of the braid group (et. al.) in this chapter, appears undiluted in Chap. 5. The only tool flexible enough to handle the complexity of the structure was that of Riemann. If that is right, then it is the documentation of these applications and questions that would illuminate on the story of the resurrection of Riemann's work. This makes it all look like slow continual progress. When, however, we come to Galois, the story has a different nature. We see it through modular curves which still to this day herald those works that accrue the most prestige.
10.3. The place of Galois. One thing is certain: Mathematicians often use Galois' name. By contrast, the most often told stories of the circumstances of his death appear unlike the essence of Galois.

I give the gist of what [Rig96] says about Galois' suicide. Galois, despondent from the suicide of his father, and the rejection of his papers by the Academy of Sciences, primarily from the negligence of Cauchy, committed a heroic suicide. She says: "offering his body against the politics of the Bourbon restoration." His rejection by Stéphanie Poterin-Dumotel exacerbated his despondency. She was the daughter of a doctor, Jean-Louis Poterin-Dumotel, who lived on the same street where Galois was transferred during a parole from prison for his major political escapade. She wasn't, in anyone's eyes, a "tart."
10.3.1. Removing the ethereal from what happened. [Rig96, p. 112] has the description of Galois' sacrifice, the morning of May 30, 1832. For shear detail, it takes your breath away. It's so solid about the climate of his sacrifice by comparison with the legend. I've never quite seen how most mathematician's credit the dual story as a romance. It lacks much of the drama of Rigatelli's analysis. Next, my own words try to capture the essense of about 50 pages from [Rig96].

While still on parole, Galois could have joined Auguste and Michel Chevalier in the Saint-Simonian community at Ménilmontant. Still, the rules imposed by Bezard and Enfantin, the leaders of the movement, would have taken away his independence. The larger picture requires some familiarity with court politics of the time. Marie-Corline, duchesse de Berry, had recently returned to France. (She didn't know Galois personally, yet we see she plays a real role in what Galois was about.) As the widow of Charles X's son, at the time of his assassination, she was expecting an heir. The boy, now 12, was living in exile in Prague, under the guidance of Cauchy. Yes, that is our Cauchy from $\S 10.2$. He was was demonstrating his devotion to the Bourbon dynasty. Galois offered himself as a hero sacrifice for the necessity of taking up dramatic arms for the republic (not the monarchy). Rigatelli poses that he arranged a dual with his friend L. D. (is that all we know of him?). A complicated piece in the tragedy, was that only his opponant's pistol would be loaded. He left several letters plausibly corroborating the dual [Rig96,
p. 109]. To accomplish certain aims his group, Friends of the People, needed only to spread the news the duel was actually a police ambush.

He did not tell Chevalier of his plot. Rather, he said only that the rejection of Stéphanie was devastating. Rigatelli emphasizes the skillful writing in these letters disguising the true situation. These letters gave rise to all the legends. They spoke with certainty of his death. Rigatteli suggests this was a sure sign of contrivance. The newspaper Le Précurseur actually told the story as it happened:

At point blank range, each [of these friends] was given a pistol and fired. Only one of the pistols was loaded [Rig96, p. 113].

The following day, at midday, 3,000 people were present at the cemetary of Montparnasse. The plan was to attack the police, when the coffin lowered into the grave. While Plaignol and Pinel, leaders of the Friends of the People, were delivering the eulogy in honour of Galois, word passed that General Lamarque had died. They decided this second funeral would attract a much larger, more emotionally involved, crowd. A swift decision brought Galois' funeral to a hasty, silent end.

The National Academy rejected Galois' famous memoire, on his solvability criterion for construction by radicals had been rejected. It was only published 14 years later. A few mathematicians conceive it, beyond doubt, as the foundation of modern algebra. Many, however, do not. Reason: It is common to view it as sketching some general abstract idea of group and permutation representations. Reality, again, is emphatically more mathematically precise and problem oriented. Reality takes account of someone having to understand its contents.
10.3.2. Group theory highlights in Galois' works. [Rig96, p. 133] One of Galois' results was that primitive representations of solvable groups must have prime power degree. He saw that degree $p$ (prime) equations have group of type $\mathbb{Z} / p \times^{s}(\mathbb{Z} / p)^{*}$ (the elementary semi-direct product). He thought to also do this for general primitive representations. He did this by considering the (Galois) field of order $p^{v}$ and he looked at the roots as listed by the congruences $\bmod p$ in this field two subscripts equal if they are given by $k^{p^{v}} \equiv k \bmod p$. So, he hoped to show the roots were permuted to take $x_{k}$ to $x_{(a k+b)^{p^{r}}}$ with $a^{p^{v}-1} \equiv 1, b^{p^{v}} \equiv b \bmod p$ and $r$ an integer. He was trying to say that if you take the roots to be $x_{i_{1}, i_{2}, \ldots, i_{v}}$ (vector space designation over the field $\mathbb{Z} / p$ of dimension $v$ ), then the group is in affine transformations augmented by the Frobenius. Rigatelli says if and only if, though certainly these are not usually solvable. She says Galois realized this later.
10.3.3. Solvability criterion and Modular curves. [Rig96][p. 137]: The third memoire is in the letter to Chevalier. In this he starts by considering integrals of the three kind. He understood (from having read Legendre and Abel - as he says in his papers) that if $g$ is the number of integrals of 1st kind then the number of periods is $2 g$ [the number of global holomorphic differentials is half the rank of the first homology: §6.5]. He was considering the monodromy on the periods, though Rigatelli does not note this. The equation giving the division of periods into $p$ equal parts has degree $p^{2 n}-1$, and its group is GL $_{p^{2 n}}$. [BAD62, p. 162-165] starts with this quote from Galois' letter to Chevalier:

La condition que j'ai indiquée dans le bulletin de Ferussac pour que l'equation soit soluble par radicaux ist trop restreinte. Il y a peu d'exceptions, mais il y en a.

My English translation: The condition that I have indicated in the bulletin of Ferussac for the equation to be solvable by radicals is too restrictive. There are few exceptions, but there are some.

After explaining the idea of primitive equations he says the following (to Chevalier). We may, however, thank historians for struggling for us with the meaning of difficult language translation combined with archaic mathematical formulations.
(10.1a) For an equation of prime degree to be solvable, it is necessary and sufficient that from any two known roots, the others are rational functions of them.
(10.1b) If an equation of degree $m$ is solvable by radicals, then $m$ is a prime power.
(10.1c) Further, in the case of prime-power degree, the equation is solvable if any two roots rationally give the others.
Rule (10.1c) overlooks the particular cases $m=9$ and $25, m=4$ and generally that where $u^{t}$ is a divisor of $\frac{p^{v}-1}{p-1}$ with $u$ prime and $\frac{\left(p^{v}-1\right) v}{u^{t}(p-1)} \equiv p \bmod u^{t}$. Galois asserts in his letter that this case, nevertheless, deviates very little from the general rule. It must always be that with two of the roots known, the others are deduced from them, by means of a number of radicals of degree $p$, equal to the number of divisors of the type $u^{t}$ satisfying the equations above. Galois says all these results come from his theory of permutations.

Finally, he says, let $k$ by the modulus of an elliptic function, $p>3$ a prime. In order that the equation of degree $p+1$ that gives the diverse modules of functions transformed relatively to the prime number $p$, be solvable by radicals, it is necessary from two choices, either one of the roots is rationally known, or each is a rational functions from any other. It does not matter here, of course, what are the particular values of the modulus $k$. It is evident that this does not hold in general. He also says it is remarkable that the general modular equation of degree 6 , corresponding to the prime 5, may be reduced to a fifth degree equation. This does not hold for any higher degree modular equations.

## 11. Exercises

11.1. Topology of covers. Let $\varphi: X \rightarrow Y$ be an unramified cover.
(11.1a) Suppose $Y$ has a countable basis for its topology. Show the same holds for $X: X$ is second countable.
(11.1b) Now assume $Y$ is a second countable Riemann surface and $X$ is also a Riemann surface where $\varphi$ a ramified cover. Show $X$ is second countable. Show $X$ is compact if and only it has the limit point property: An infinite sequence has a convergent subsequence.
The following shows a strong equivalence $\psi: Y^{1} \rightarrow Y^{2}$ between two finite ramified covers $\left(Y^{i}, \psi^{i}\right)$ of $X$ (as in $\S 3.2 .2$ ), $i=1,2$, is automatically an analytic isomorphism.
(11.2a) Restrict both covers to $X \backslash D\left(\psi^{\prime}\right)$ : show that $\psi: Y_{X \backslash D\left(\psi^{1}\right)}^{1} \rightarrow Y_{X \backslash D\left(\psi^{2}\right)}^{2}$ is analytic.
(11.2b) Use Riemann's removable singularities theorem [Ahl79, p. 129] to extend the map of a) to an analytic map including the discrete set $Y_{D\left(\psi^{1}\right)}^{1}$.
We now investigate the Jordan curve theorem.
(11.3a) Let $\gamma^{\prime}$ be a simple closed polygonal path on $\mathbb{P}^{1}=\mathbb{P}_{z}^{1}$ and let $U$ be a connected component of $\mathbb{P}^{1} \backslash \gamma^{\prime}$. Show the points of $\gamma^{\prime}$ in the boundary of $U$ are both open and closed. Let $z_{0}$ be in the range of $\gamma^{\prime}$. Thus show that
$\mathbb{P}^{1} \backslash \gamma^{\prime}$ has at most two connected components, told apart as the points that connect in a neighborhood on the left of $z_{0}$ versus those that connect to a neighborhood on the right of $z_{0}$. The component connected to $\infty$ consists of points that relative to $\gamma$ have winding number 0 . So, it suffices to show there is a point of $\mathbb{P}_{z}^{1} \backslash \gamma^{\prime}$ with nonzero winding number relative to $\gamma^{\prime}$. See [Ahl79, p. 118].
(11.3b) Let $\gamma^{\prime}$ be a simple simplicial closed path in $\mathbb{P}^{1}$ and let $W$ be the interior of $\gamma^{\prime}$. Let $\gamma^{\prime \prime}$ be a simple simplicial path that meets $\gamma^{\prime}$ only at the end points, $x_{0}$ and $x_{1}$, of $\gamma^{\prime}$. Note that $\gamma^{\prime} \backslash\left\{x_{0}, x_{1}\right\}$ consists of two connected components. Conclude from the Jordan curve theorem that each component together with $\gamma^{\prime \prime}$ defines a simple closed curve whose interior consists of one of the two components of $W \backslash \gamma^{\prime \prime}$. Each component consists of the points path connected to any given point of the component.
Relate back to Chap. 2 [9.17] to complete a discussion of orientability of complex manifolds. Also Def. 2.21.
(11.4a) Use the complex structure to get the orientation from the sign of the expression $d x \wedge d y$. That works by substituting variables along the paths. Show that if $Y \rightarrow X$ is a finite ramified cover, an orientation of $X$ gives an orientation of $Y$, and if $X$ is a complex manifold, so is $Y$.
(11.4b) Show $d(f(z) d z)=0$.

The Brouwer separation theorem [?, p. 11-21] states that a compact differentiable 2-dimensional manifold $M$ in $\mathbb{R}^{3}$ separates. That is, $\mathbb{R}^{3} \backslash M$ consists of two components, an inside that is bounded, and an outside that is unbounded.
(11.5a) Use the separation theorem to show that there is a continuously varying vector $\boldsymbol{w}_{m}$ of length 1 for each $m \in M$ such that the dot product $\boldsymbol{w}_{m} \cdot \boldsymbol{v}_{m}=$ 0 for every vector $\boldsymbol{v}_{m}$ tangent to some path on $M$ through $m$ (i.e., $\boldsymbol{w}_{m}$ is normal to $M$ at $m ; \S 1 . c)$.
(11.5b) Use the notation of Chap. $3 \S 2.21$. Show that a compact differentiable manifold in $\mathbb{R}^{3}$ is orientable. Hint: Restrict to coordinate charts that at each point have the RETURN
(11.5c) Orientation for higher dimensional manifolds.
(11.5d) Show that there is an unramified cover $\mathbb{P}^{1} \rightarrow X$ of degree 2 for which $X$ is not orientable. Conclude that such an $X$ cannot be embedded in $\mathbb{R}^{3}$. Hint: $\mathbb{P}^{1}$ is homeomorphic to the 2 -sphere, $S^{2}$, in $\mathbb{R}^{3}$. Make the set whose points consist of pairs of endpoints of diameters of $S^{2}$ into a manifold.
The next exercise series shows how a little combinatorics of triangles reveals that the Euler characteristic of a Riemann surface being genus 0 or 1 shows it is topologically a sphere or a torus. This continues the discussion of Rem. 2.19. Suppose two Riemann surfaces $X_{i}, i=1,2$, have triangulations have respective triangulations $T_{i}, i=1,2$. Suppose there is a numbering of the simplices, edges and vertices in both so that the numbering for one is exactly the same as the numbering for the other. Call the triangulations equivalent.
(11.6a) Suppose for the triangulation of $X_{c}$ in Prop. 2.18 the Euler characteristic is $2\left(g_{\boldsymbol{g}}=0\right)$. Lay out a triangulation on the sphere equivalent to this triangulation. Conclude an Euler characteristic 2 implies the surface is homeomorphic to the sphere.
(11.6b) Do the same for concluding about $X_{c}$ when its genus is 1 that it is homeomorphic to a torus.
11.2. Artin braids and Hurwitz monodromy. Notation is from Def. 1.1.
(11.7a) Show any group requiring at most $|S|$ generators is a quotient of $F(S)$.
(11.7b) Let $[G, G]$ denote the commutator subgroup of a group $G$ : Elements $g_{1} g_{2} g_{1}^{-1} g_{2}^{-1}$ for all $g_{1}, g_{2} \in G$ generate $[G, G]$. Let $A(S)=F(S) /[F(S), F(S)]$ (or $A_{n}$ if $|S|=n$ ) be the free abelian group on $S$. Show any abelian group requiring at most $|S|$ generators is a quotient of $A(S)$.
(11.7c) Show $A(S)=A_{n} \cong \mathbb{Z}^{n}$ if $S=\left\{s_{1}, \ldots, s_{n}\right\}$.
(11.7d) Let $F(S)=F_{2 r}$ and let $\bar{R}(S)$ be the normal subgroup generated by

$$
s_{1} s_{r+1} s_{1}^{-1} s_{r+1}^{-1} s_{2} s_{r+2} s_{2}^{-1} s_{r+2}^{-1} \cdots s_{r} s_{2 r} s_{r}^{-1} s_{2 r}^{-1}
$$

Show that $G=F_{2 r} / \bar{R}(S)$ is not a free group. Hint: $G /[G, G]=A_{2 r}$.
Refer to the properties in Prop. 3.3.
(11.8a) other automorphisms of $F_{r}$ not included in the braid group.
(11.8c) Define the straight (or pure) braids to be the elements of $B_{r}$ in the kernel of $\Psi_{r, *}$.
(11.8d) Let $\operatorname{Inn}\left(F_{r}\right)$ be the normal subgroup of $\operatorname{Aut}\left(F_{r}\right)$ generated by conjugations of elements of $F_{r}$ on itself. The mapping class group (of degree $r$ ) is the image in $\operatorname{Aut}\left(F_{r}\right) / \operatorname{Inn}\left(F_{r}\right)$ of automorphisms of $F_{r} / \bar{R}$ induced by the action of $H_{r}$ (or $B_{r}$ ) on $F_{r}$. Denote this group by $M_{r}$. Show

$$
\begin{aligned}
\tau_{1}=\left(Q_{2} Q_{3} \cdots Q_{r-1}\right)^{1-r}, \ldots, \tau_{\ell+1} & = \\
\left(Q_{1} \cdots Q_{\ell}\right)^{\ell+1}\left(Q_{\ell+2} \cdots Q_{r-1}\right)^{\ell+1-r}, \ldots, \tau_{r-1} & =\left(Q_{1} \cdots Q_{r-2}\right)^{r-1}
\end{aligned}
$$

and $\tau=\left(Q_{1} \cdots Q_{r-1}\right)^{r}$ are in the kernel of the natural map $H_{r} \rightarrow M_{r}$.
(11.8e) Show there is a unique group (the dihedral group of degree $n$ ) of order $2 n$ and generated by two elements $\sigma_{1}, \sigma_{2}$ of order 2 for which $\sigma_{1} \sigma_{2}$ is of order $n$. Similarly, show there is a unique group (the dicyclic group of degree $2 n$ ) of order $4 n$ and generated by $\sigma_{1}, \sigma_{2}$ of respective orders $2 n$ and 4 , and for which $\sigma_{2}^{-1} \sigma_{1} \sigma_{2}=\sigma_{1}^{-1}$ and $\sigma_{2}^{2}$ is in the group generated by $\sigma_{1}$.
(11.8f) Show $H_{3}$ is isomorphic to the dicyclic group of degree 6, and that $M_{3}$ is isomorphic to $S_{3}$. Hint: $Q_{1} Q_{2}$ and $Q_{1} Q_{2} Q_{1}$ are also generators of $H_{3}$.

### 11.3. Seifert-van Kampen theorem and fiber products.

(11.9a) Give an example to show why $U \cap V$ must be connected for Thm. refthm7 to hold. Hint: Look at Fig 2.1; but it's not the easiest example.
(11.9b) Show that if $\pi_{1}\left(U \cap V, x_{0}\right)$ is trivial in Thm. 1.5, then $\pi_{1}\left(X, x_{0}\right)$ is the free product of $\pi_{1}\left(U, x_{0}\right)$ and $\pi_{1}\left(V, x_{0}\right)$. Conclude in this case that if $\pi_{1}\left(U, x_{0}\right)$ and $\pi_{1}\left(V, x_{0}\right)$ are, respectively, isomorphic to $F_{r}$ and $F_{s}$ (Ex. 2.5) then $\pi_{1}\left(X, x_{0}\right)$ is isomorphic to $F_{r+s}$.
(11.9c) If $\pi_{1}\left(V, x_{0}\right)$ is trivial, show that $i(U, X)_{*}: \pi_{1}\left(U, x_{0}\right) \rightarrow \pi_{1}\left(X, x_{0}\right)$ is surjective with kernel the smallest normal subgroup $N$ of $\pi_{1}\left(U, x_{0}\right)$ containing the image of $\pi_{1}\left(U \cap V, x_{0}\right)$ by the map $i(U \cap V, V)_{*}$. Hint: Take $H$ to be $\pi_{1}\left(U, x_{0}\right) / N$. Let $\beta(U): \pi_{1}\left(U, x_{0}\right) \rightarrow H$ be the natural map $(\beta(V)$ is the trivial homomorphism). Conclude that the kernel of $\beta(X)$ is $N$.
(11.9d) Show $\pi_{1}\left(\mathbb{P}^{n}\right)$ is trivial for $n \geq 0$. Hint: It is a union of pieces with trivial fundamental group. There is another approach. If $V$ is a projective manifold, and $V_{1}$ is a codimension 1 subvariety, then the natural map $\pi_{1}\left(V \backslash V_{1}\right) \rightarrow \pi_{1}(V)$ is surjective. Now use that $\mathbb{C}^{n}$ is contractible.
Consider fiber products and pushouts in the category of finite groups covering a given group. This is entirely different than the similar categorical notions that appear in Chap. 3 [9.3].
(11.10a) Show the free product of groups $G_{1}, \ldots, G_{t}$ in Lem. 1.4 is a pushout in the category of groups and homomorphisms by taking the map of $\{1\}$ into each of these groups. Note: (1.2) is therefore a sum, rather than a product. Product would be a group $G$ with maps to all the $G_{i}$ s so that any $H$ that maps to all the $G_{i}$ s would map to $G$.
(11.10b) Show uniqueness of pushouts is general in categories of groups. Apply this to Thm. 1.5 to see why this defines $\pi_{1}\left(X, x_{0}\right)$ uniquely from the maps $\pi_{1}\left(U \cap V, x_{0}\right) \rightarrow \pi_{1}\left(U, x_{0}\right)$ and $\pi_{1}\left(U \cap V, x_{0}\right) \rightarrow \pi_{1}\left(V, x_{0}\right)$.
11.4. Residues and uniformization for covers of curves of genus 1. Thm. 6.15 gives the fundamental group of an $r$ punctured Riemann surface $X$ of genus $g$. Thm. 2.6 started with a statement about the fundamental group of $U_{\boldsymbol{z}}$.
(11.11a) Genus 1 Curve.
(11.11b) This also applies to any Riemann surface uniformized by a disc or by the complex plane.
(11.11c) State and prove a generalization of Thm. 2.6 and Cor. 2.6 characterizing the ramified covers of $X$ ramified over a finite subset $\boldsymbol{x}=\left(x_{1}, \ldots, x_{r}\right)$ of $X$, where $X$ is any compact Riemann surface. Use $X_{\boldsymbol{x}}^{0}$ for $X \backslash \boldsymbol{x}$.
(11.11d) Consider the case $X$ has genus 1. Suppose $\varphi: Y \rightarrow X$ is a ramified cover, with $\boldsymbol{x}$ its branch points and $\left\langle a, b, \bar{g}_{1}, \ldots, \bar{g}_{r}\right\rangle$ classical generators of $\pi_{1}\left(X_{\boldsymbol{x}}^{0}, x_{0}\right)$.
(11.11e) HELP Assume $N$ is odd, and label a conjugacy class in $S_{N}$ as of type $2^{k}$ if it is the conjugacy class of $k$ disjoint 2 -cycles. Let $\mathbf{C}_{3 \cdot n, n-1,1}$ be the conjugacy classes of type $\left(2^{n}, 2^{n}, 2^{n}, 2^{n-1}, 2\right)$ in $S_{N}$ with $n=(N-1) / 2$. Recall also the Nielsen class notation $\boldsymbol{g} \in \mathbf{C}_{3 \cdot n, n-1,1}$ to mean in some order the entries of $\boldsymbol{g}=\left(g_{1}, \ldots, g_{5}\right)$ are in the respective conjugacy classes.
11.5. Reducible fiber products. Assume $\varphi: Y_{i} \rightarrow X$ are ramified covers of compact Riemann surfaces, and let $W=Y_{1} \times_{X} Y_{2}$ be the fiber product with its topology coming from it being a subspace of $Y_{1} \times Y_{2}$. Chap. 3 [9.11c] gives simple situations where fiber products $W$ are reducible. We develop more substantial examples by applying RET to the groups of Chap. 3 [9.20].
(11.12a) Finish the proof of Prop. 3.4 by considering $g_{k}$, an $e_{k}$ cycle in $S_{e_{k}}, k=1,2$. Show $\left(g_{1}, g_{2}\right) \in S_{e_{1}} \times S_{e_{2}}$, acting on $\{(i, j)\}_{1 \leq i \leq e_{1}, 1 \leq j \leq e_{2}}$ is a product of $\left(e_{1}, e_{2}\right)$ disjoint $\left[e_{1}, e_{2}\right]$ cycles.
(11.12b) How does a) describe the irreducible components of the (normal) fiber product of $\mathbb{P}_{w_{i}}^{1} \rightarrow \mathbb{P}_{z}^{1}$ be $w_{i} \mapsto w_{i}^{e_{i}}, i=1,2$ ?

Galois correspondence and primitivity II: Consider the components of $Y \times_{X}$ $Y \backslash \Delta$ of form $Y^{\prime} \times_{X} Y^{\prime} \backslash \Delta$. Each of these attests to a decomposition of $f: Y \rightarrow X$
according to Chap. 3 [9.22c]. We point out how the use of coordinates gives a more practical test.
(11.13a) Reference p. 37, R. Rosario, Thesis: Unirational Fields, Univ. of Cantabria: $(x-\alpha) \prod_{j=1}^{k} P_{i_{j}}$ with these being the factors of $P(x)$ over $K(\alpha)$. If the coefficients generate a proper subfield then imprimitive. Looking for zeros forming a set of imprimitivity.

This exercise considers the branch point and reduced branch point maps $\Psi_{r}$ : $\left(\mathbb{P}_{z}^{1}\right)^{r} \rightarrow \mathbb{P}^{r}$ and $\Psi_{r}^{\text {rd }}:\left(\mathbb{P}_{z}^{1}\right)^{r} \rightarrow J_{r}$.
(11.14b) Do the discriminant locus $D_{r}$.

Look back at Ex. 4.1, and do Schur conjecture.
(11.15a) The $\mathbb{Q}$ absolute r-equivalence classes are represented by their branch points $\{0, \infty\}$ and the collection $\left\{z_{1}, z_{2} \mid z_{1}=\sqrt{m}, z_{2}=-\sqrt{m}, m\right.$ a square-free integer.
(11.15b) Do the setup for the Schur conjecture.
(11.15c) Finish the Schur conjecture.
11.6. Cuts, tangential base points and symbolic pictures. Consider the use of the fundamental group $\pi_{1}\left(U_{\boldsymbol{z}}, D_{\boldsymbol{v}}\right)$ in the proof of Thm. 2.26.
(11.16a) Show that if $D_{\boldsymbol{v}^{\prime}}$ is another tangential disk to 0 , giving an isomorphism between $\pi_{1}\left(U_{\boldsymbol{z}}, D_{\boldsymbol{v}}\right)$ and $\pi_{1}\left(U_{\boldsymbol{z}}, D_{\boldsymbol{v}^{\prime}}\right)$ depends on how you regard $D_{\boldsymbol{v}}$ connected to $D_{\boldsymbol{v}^{\prime}}$.
(11.16b) Use as a base point for cuts one of the branch points.

This exercise builds from Chap. 3 [7.2.3]. The point is that we often need notation to differentiate between more subtle appearance of conjugacy classes in $S_{n}$. Use notation of $\S 2.4$ for discussing cuts.
(11.17a) For the situation of one cut, complete the proof of Lem. 2.15.
11.7. Alternating group conjugacy classes. We first finish considering the rationality of conjugacy classes in alternating groups.
(11.18a) Assume $\boldsymbol{g} \in S_{n}^{r}, n \geq 3$ and $\langle\boldsymbol{g}\rangle$ is transitive. Show $\langle\boldsymbol{g}\rangle \geq A_{n}$ if $\boldsymbol{g}$ contains a 3-cycle. Hint: Show $\langle\boldsymbol{g}\rangle$ must be primitive and imitate Chap. 3 [9.15e].
(11.18b) Let C be an $A_{n}$ conjugacy class with $\mathrm{C}_{S_{n}}=\left(m_{1}\right) \cdots\left(m_{t}\right)$ and all $m_{i} \mathrm{~s}$ distinct and odd. Write $m=\left[m_{1}, \ldots, m_{t}\right]$ as $\prod_{i=1}^{s} q_{i}^{v_{i}} ; q_{i}$ s distinct primes and the $v_{i}$ s positive. Suppose $q_{i}^{v_{i, j}}$ exactly divides $m_{j}$ by $v_{i, j}$. Denote $\left(v_{i, 1}, \ldots, v_{i, t}\right)$ by $\boldsymbol{v}_{i}$. Define $\mu: \mathbb{Z}^{t} \rightarrow \mathbb{Z} / 2$ by $\left(a_{1}, \ldots, a_{t}\right) \mapsto$ $\sum_{j=1}^{t} a_{j} \bmod 2$. Follow Prop. 3.19 using $k$, generating $\left(\mathbb{Z} / q_{i}^{v_{i}}\right)^{*}$, acting on $\oplus_{j=1}^{t} \mathbb{Z} / q_{i}^{v_{i, j}}$. Identify $k$ with a permutation $\tau_{k} \in S_{V_{i}}$, with $V_{i}=$ $\oplus_{j=1}^{t} q_{i}^{v_{i, j}}$. Show $\tau_{k} \in A_{n}$ if and only if $\mu\left(\boldsymbol{v}_{i}\right)=0$. Hint: $\tau_{k}$ comes from the product of the actions of $k$ on $\mathbb{Z} / q_{i}^{v_{i, j}}$.
(11.18c) With C as in b), show $M_{\mathrm{C}}$ is nontrivial if and only if $\mu\left(\boldsymbol{v}_{i}\right)$ is nonzero for some $i$ between 1 and $s$. Let $J$ be those $i$ with $\mu\left(\boldsymbol{v}_{i}\right) \neq 0$. Denote $\sqrt{\prod_{i \in J}(-1)^{\left.q_{i}-1\right) / 2} q_{i}}$ by $\alpha_{\mathrm{C}}$. Show $M_{\mathrm{C}}$ is $\mathbb{Q}\left(\alpha_{\mathrm{C}}\right)$.
(11.18d) Consider the 12 pairs C of conjugacy classes C of $A_{25}$ for which $\mathrm{C}_{S_{25}} \neq \mathrm{C}$. Imitate Ex. 3.21 to show that of these the only pair consisting of rational conjugacy classes is that with $\mathrm{C}_{S_{25}}=(25)$. Show that for C with $\mathrm{C}_{S_{25}}=$ $(1)(9)(15)($ resp. $(1)(3)(21)), M_{\mathrm{C}}=\mathbb{Q}(\sqrt{-3 \cdot 5})\left(\right.$ resp. $\left.M_{\mathrm{C}}=\mathbb{Q}(\sqrt{-7})\right)$.
Subtler issues about conjugacy classes in $A_{n}$.
(11.19a) Consider $\boldsymbol{g} \in S_{n}^{r}, n \geq 4$, that consist of products of two disjoint 2-cycles. If $n$ is even show that there are examples with $\langle\boldsymbol{g}\rangle$ transitive, but not primitive. If $n=7$, show $\langle\boldsymbol{g}\rangle$ could be $\mathrm{PSL}_{3}(\mathbb{Z} / 2)$ instead of $A_{7}$.
(11.19b) Let $g_{1}=(1 n), g_{2}=(2 \cdots n), g_{3}=(12 \cdots n)^{-1}, r=3$ and $H=A_{n}$ in Corollary 2.17. Show that $\psi_{H}: Y_{H} \rightarrow \mathbf{P}^{1}$ has a description of its branch cycles given by $((12),(12))$. Find $f(y) \in \mathbf{Q}[y]$ such that $\psi: \mathbb{P}_{y}^{1} \rightarrow \mathbb{P}_{z}^{1}$ by $y \mapsto f(y)=z$ has $\left(\sigma_{1}, \sigma_{2}, \sigma_{3}\right)$ as a description of its branch cycles.
Take $r=4$ and $G=S_{5}$. Let $\mathrm{C}_{1}$ and $\mathrm{C}_{3}$ be the conjugacy classes of 2-cycles in $S_{5}, \mathrm{C}_{2}$ the conjugacy class of a 3 -cycle and $\mathrm{C}_{4}$ the conjugacy class of a 5 -cycle. Consider the Nielsen class $\mathrm{Ni}\left(S_{5}, \mathbf{C}\right) / S_{5}=\mathrm{Ni}^{+}$:

$$
\left\{\boldsymbol{g}=\left(g_{1}, \ldots, g_{4}\right) \mid g_{1} g_{2} g_{3} g_{4}=1,\langle\boldsymbol{g}\rangle=S_{5} \text { and } \boldsymbol{g} \in \mathbf{C}\right\} / S_{5}
$$

(11.20a) How many elements are in $\mathrm{Ni}^{+}$?
(11.20b) Let $\psi: \pi_{1}\left(U_{\boldsymbol{z}}\right) \rightarrow S_{5}$ map a fixed set $\bar{g}_{1}, \ldots, \bar{g}_{4}$ into some element of $\mathrm{Ni}^{+}$. Why is the cover corresponding to such a homomorphism a genus 0 compact Riemamn surface minus a finite set of points?
(11.20c) Represent $S_{5}$ on the 10 unordered distinct pairs of integers from $\{1, \ldots, 5\}$ : $T: S_{5} \rightarrow S_{1} 0$. Example: (12345) has two orbits on these 10 pairs. What are the lengths of the disjoint cycles of $T$ applied to an element of the conjugacy class of a 3 -cycle in $S_{5}$ ?
(11.20d) Compose $\psi$ with $T$ to get $T \circ \psi=\psi^{\prime}: \pi_{1}\left(U_{\boldsymbol{z}}\right)$. What is the genus of the curve at the top of the corresponding cover $X=X_{\psi} \rightarrow \mathbb{P}_{z}^{1}$ ?
(11.20e) Does the isomorphism class of $X_{\psi}$ depend on $\psi$ (assuming $\psi$ is in the Nielsen class $\left.\mathrm{Ni}^{+}\right)$?
Now we discuss the genus of the corresponding degree 10 covers. Let $\boldsymbol{g}$ be a branch cycle description of the cover from $\mathrm{Ni}^{+}$in [11.20]. Compute the genus $g$ of $+\mathcal{T}_{\boldsymbol{p}}^{(2)}$ from Riemann-Hurwitz:

$$
\begin{equation*}
2(10+g-1)=\sum_{i=1}^{4} \operatorname{ind}\left(R\left(g_{i}\right)\right) \tag{11.21}
\end{equation*}
$$

Suppose $g_{1}$ and $g_{3}$ are 2 -cycles from $S_{5}$. Then, $R\left(g_{i}\right)$ has shape $(2)(2)(2)$ in the representation $R, i=1,3$. Similarly, if $g_{2}$ is a 3 -cycle, $R\left(g_{3}\right)$ has shape (3)(3)(3). Finally, $R\left(g_{4}\right)$ has shape $(5)(5)$. Thus, the total contribution to the right side of (11.21) is $2 \cdot 3+6+2 \cdot 4=20$ and $g=1$.

Next: Compute $\mathrm{Ni}^{+}$modulo conjugation by $S_{5}$. Choose $S_{5}$ representatives with $g_{4}$ equal $g_{\infty}=(12345)^{-1}$. Divide $\mathrm{Ni}^{+}$into two sets $T_{1}$ and $T_{2}: \boldsymbol{g} \in T_{1}$ has $g_{1}$ and $g_{2}$ with no integers of common support, and $\boldsymbol{g} \in T_{2}$ has $g_{1}$ and $g_{2}$ with one integer of common support. Conjugate by a power of $\boldsymbol{g}_{\infty}$ to assure elements of $T_{1}$ have $g_{1}=(1 j)$ with $j=2$ or 3 . Similarly, elements in $T_{2}$ have 1 as common support of $g_{1}$ and $g_{2}$. From this, list $\mathrm{Ni}^{+, \text {abs }}$.

Now we consider some genus covers with group $A_{5}$ and branch cycles having the following type. $(3)(3)(3)(5)$ : Suppose $g_{3}=(123)$.
(11.22a) Ramification: $g_{1} g_{2}$ is $(2)(2)$, assume missing integer is 1 , so to get product a 5 -cycle: may assume $g_{1} g_{2}$ is $(25)(34)$. Now everthing is fixed and need only count number of ways to write $g_{1} g_{2}$ is a product of two three cycles. Hint: Products of two 3-cycles giving (25)(34): You get one element from $(425)(234)$. Now conjugate the pair $((425),(234))$ by the centralizer of $(25)(34)$, the group $\langle(25)(34),(24)(35)\rangle$.
(11.22b) If $g_{1} g_{2}$ is (3), then conjugate by $\left\langle g_{3}\right\rangle$ to assume common integer is 1 , and $g_{1} g_{2}=(145)$. Hint: Take $\left(g_{1}, g_{2}\right)=((143),(135)$, and then conjugate by $\langle(23),(145)\rangle$.
(11.22c) If $g_{1} g_{2}$ is (5). Then, product can't be of type (2)(3) (Riemann-Hurwitz), and have only to assure the (5) times $g_{3}$ doesn't fix anything. That means can't have $2 \mapsto 1,3 \mapsto 2$ or $1 \mapsto 3$. Also, since by conjugation by $\langle(45),(123)\rangle$ can assume (15???) resulting in (15243) or (15324). Hint: For each of (15243) or (15324), we need to count all the ways to write this 5 -cycle as a product of two 3 -cycles. For (12345), assume the integer 1 is the common integer to the 3 -cycles. So, $\left(g_{1}, g_{2}\right)=((123),(145))$. Then, by conjugating by $\langle(12345)\rangle$, gives the five cases where $g_{1}$ and $g_{2}$ have any desired integer in common.
(11.22d) Up to equivalence, there are exactly 4 covers from a), 6 covers from b) and 10 covers from c), or 20 total covers. Also, by applying powers of $q_{1}$ to case c) you get 10 total in two orbits of length five. Same for b), two orbits of length 3, and for a), two orbits of length two.
11.8. Differentials and differential equations. The space of holomorphic differentials has dimension bounded by $g$.
(11.23a) Finish the pairing with homology classes.
(11.23b) Show $d f$ in (2.13) is a 2 -form.

Let $\mathcal{H}(D)$ be functions analytic on a domain $D$. A differential equation on $D$ comes from $m \in \mathcal{H}(D)\left[w_{0}, w_{1}, \ldots, w_{k}\right]$ : a polynomial with coefficients in $\mathcal{H}(D)$. Solutions of the equation are functions $f(z)$, analytic in some disk in $D$ with $m\left(f(z), \frac{d f}{d z}, \ldots, \frac{d^{k} f}{d z^{k}}\right) \equiv 0$ on this disk. Especially interesting are equations defined by $m$ linear in the variables $w_{0}, w_{1}, \ldots, w_{k}$, with coefficients in $\mathbb{C}[z]$ : linear, algebraic differential equations. Examples: $m_{1}\left(w_{0}, w_{1}\right)=w_{0}-w_{1}, m_{2}\left(w_{0}, w_{1}\right)=w_{0}-z w_{1}$ and $m_{3}\left(w_{0}, w_{1}\right)=z w_{0}-w_{1}$.
(11.24a) Suppose $m \in \mathbb{C}[z]\left[w_{0}, \ldots, w_{k}\right]$ defines a linear algebraic differential equation. Let $f(z)$, analytic on $D$, solve the equation. Show there exists a finite set $\boldsymbol{z} \subset \mathbb{P}_{z}^{1}$ satisfying (1.1a). Hint: Let $\frac{d}{d z}$ act on functions analytic on $D$. Produce a matrix operator $\mathcal{D}$ on $\mathcal{H}(D)^{k+1}$ :

$$
\left(f_{0}, f_{1}, \ldots, f_{k}\right) \mapsto\left(m\left(f_{0}, f_{1}, \ldots, f_{k}\right), f_{1}-\frac{d f_{0}}{d z}, \ldots, f_{k}-\frac{d f_{k-1}}{d z}\right)
$$

Find $\boldsymbol{z}$ from the determinant of $\mathcal{D}$.
(11.24b) Show the vector space of solutions (analytic in a disk centered at $z_{0} \in$ $\left.\mathbb{P}_{z}^{1} \backslash \boldsymbol{z}\right)$ of the differential equation $m$ has dimension $k$.
(11.24c) Consider the case $m=P_{1}(z) w_{0}-P_{2}(z) w_{2}$. where $g=P_{1} / P_{2}$ is a rational function satisfying some well know conditions that permit its solutions to be continued over every path on the $z$-sphere. When the denominator of $g$ has degree $2 p+2$ then there are $2 p-1$ unknown conditions that prevent the unique specification of the numerator - the goal being that
the ratio of a pair of independent solutions of the equation should yield a conformal equivalence of the branched universal covering surface for the $2 p+2$ punctured sphere with the unit disk (at least for $p>1$ ). Thus it is really the moduli problem for hyperelliptic curves.
11.9. Schwartzian and Beltrami equations. Riemann and Schwarz used functional equations to characterize the nature of many functions. The most famous after Riemann's application to $\theta$ functions is the Schwarzian derivative. Call $\mathcal{D}$ a Schwarzian for a subgroup $G \leq \mathrm{PGL}_{2}(\mathbb{C})$ if $\mathcal{D}(g(h(z)))=\mathcal{D}(h(z))$ for $g \in G$ and any meromorphic function $h$, and conversely, the set of $g \in \mathrm{PGL}_{2}(\mathbb{C})$ for which this holds exactly for relevant meromorphic $h$ defines $G$. Use notation from Chap. 2 [9.14].
(11.25a) Show the translation group $\mathcal{T}$ has Schwarzian $\mathcal{D}_{\mathcal{T}}$ given by $h \mapsto \frac{d h}{d z}$.
(11.25b) Show the group of multiplications $\mathcal{M}$ has Schwarzian $\mathcal{D}_{\mathcal{M}}$ given by $h \mapsto$ $\frac{h^{\prime}}{h}$, the logarithmic derivative.
(11.25c) Show the affine group $\mathbb{C} \times{ }^{s} \mathbb{C}^{*}=\mathcal{A}$ has Schwarzian $\mathcal{D}_{\mathcal{A}}=\mathcal{D}_{\mathcal{M}} \circ \mathcal{D}_{\mathcal{T}}$.
(11.25d) Any element of $\mathrm{PGL}_{2}(\mathbb{C})$ has either the form $z \mapsto a z+b$ or $z \mapsto a+b /(z-c)$. Conclude (by changing $h(z)$ to $h(z)-c$ ): $\mathrm{PGL}_{2}(\mathbb{C})$ has a Schwarzian if there is a differential operator $\mathcal{D}_{\tau}(\tau: z \mapsto 1 / z)$ with

$$
\mathcal{D}_{\tau}\left(\frac{h^{\prime \prime}}{h^{\prime}}\right)=\mathcal{D}_{\tau}\left(\frac{\left(1 /(h)^{\prime \prime}\right.}{\left(1 /(h)^{\prime}\right.}\right)
$$

for any meromorphic $h$. Compute that $\mathcal{D}_{\tau}(g)=g^{\prime}-\frac{1}{2} g^{2}$ works: $\mathrm{PGL}_{2}(\mathbb{C})$ has a Schwarzian $\mathcal{D}_{\tau} \circ \mathcal{D}_{\mathcal{A}}$.
The Beltrami Equation: Irwin, P. 35 of your book. This is a loaded page, in which you take a Riemannian structure and turn it into a quasiconformal structure. I am doing exercises (Chap. 3 of a book) in which I took my own path to the Beltrami equation to suit a theory of uniformization I'm using. I wanted to do something along the lines you are doing. You have, however, the statement: "the most nontrivial $\ldots$ is the verification that $\mu f_{z}=f_{\bar{z}}$ has homeomorphic solutions."

I looked in the rest of your book, and couldn't find a proof that $\mu f_{z}=f_{\bar{z}}$ has homeomorphic solutions if $\mu$ is bounded by 1 in the neighborhood of a given point. Is it somewhere there, and if so, what is the easiest solution of this?

### 11.10. Frattini covers and half-canonical classes.

For the particular situation defining a half-canonical class from a differential satisfying Chap. 3 (5.11), the choices are defined up to chains with values in $\pm 1$. Now consider the number of boundary equivalent half-canonical classes. Every set of transition functions has a global meromorphic section (how hard would that be to prove in this case). So, each boundary equivalence class has a differential giving it by situation (5.11). So, it is enough to use the setup of Chap. 3 (5.11) to define all boundary equivalent half-canonical classes.
(11.27a) Suppose $\left\{h_{\beta, \alpha}\right\}$ and $\left\{h_{\beta, \alpha}^{\prime}\right\}$ are two sets of transition functions defining half-canonical classes by the rule (5.11). Let $h_{\alpha} / h_{\alpha}^{\prime}$ be the corresponding ratios of the functions from (5.11). Their squares form a function on $X$ and they give a homomorphism $\pi_{1}(X) \rightarrow \mathbb{Z} / 2$ by the following rule. For any closed path, form the analytic continuation of $h_{\alpha} / h_{\alpha}^{\prime}$ around the path.

This works if the formula for analytic continuation applies on a manifold which it does by $\S 6.2$. So, given any two their ratio defines a cover of $X$, and conversely.
(11.27b) Now, suppose the cover $\hat{Y} \rightarrow \hat{X}$ is unramified. We already have a homomorphism $\pi_{1}\left(U_{\boldsymbol{z}}, z_{0}\right) \rightarrow \hat{A}_{n}$ defining this. Further, the kernel of this to $A_{n}$ factors through a map $\pi_{1}(\hat{X}) \rightarrow \mathbb{Z} / 2$. Note: If you take the divisor of $d \hat{\varphi}$, it is $A_{n}$ invariant.
(11.27c)
11.11. Differential forms, orientation, area and the Laplacian. Why is an orientation forced in order to integrate a form? How would we generalize length to get area? Define $\wedge$ multiplication.
(11.28a) Pythagorian formula for area.

Show with any Riemannian manifold $X$ replacing $\mathbb{R}^{2}$, with $d s^{2}: \mathbb{T}_{X} \times \mathbb{T}_{X} \rightarrow \mathcal{C}_{X}^{\infty}$ the nondegenerate symmetric 2 -tensor.
(11.29a) Consider an open set $U$ in $\mathbb{R}^{n}$. Call a differential 1-form $\omega$ integrable if it has the form $d f$ for some $f \in \mathcal{C}_{U}^{\infty}$. The integrability condition is that $d \omega=$ 0 . $f \mapsto d f \mapsto T_{d f}$ maps to the $(1,1)$ tensor $\Delta\left(T_{d f}\right)(T, \omega)=D_{T}\left(T_{d f}\right) \otimes \omega$. Now contract back to a function.

Take a basis of differentials and recall the inner product $\left\langle\varphi_{1}, \varphi_{2}\right\rangle=$ $\int_{X} d \varphi_{1} * \bar{\varphi} \varphi_{2}$.

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