THE UNIFORMISATION THEOREM OF RIEMANN SURFACES

A dissertation submitted to Birkbeck, University of London for the degree of MSc in Mathematics.

By

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Abstract

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ABSTRACT OF DISSERTATION submitted by Y. Buzoku and entitled The Uniformisation Theorem of Riemann surfaces.

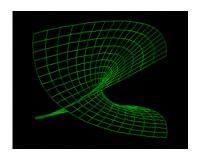
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This dissertation is an introduction to Riemann surfaces and their classification by way of the uniformisation theorem. The dissertation is split into two parts; the first being a discussion of Riemann surfaces themselves, some methods of their construction, and interesting properties and symmetries that they exhibit, concluding by putting ourselves into a position where we can present the problem of classifying Riemann surfaces naturally. The standard reference for this section is the book by Jones & Singerman on Complex Functions, [7]. The second part focuses on providing a more analytical account of the proof of the classification of Riemann surfaces. It is an analytically heavy treatment, motivated in a conceptually different way via the theory of electrostatics. We set up the problem and then diverge into considering compact Riemann surfaces, which we go into detail on and then sketch the case for non compact, simply connected Riemann surfaces as large parts of the arguments are similar in the two cases. The uniformisation theorem then follows naturally. The standard reference for the second section is the book by Donaldson on Riemann surfaces, [4], though we take a somewhat different approach, eventually converging our argument with Donaldson's.

Declaration

This dissertation is submitted under the regulations of Birkbeck, University of London as part of the examination requirements for the MSc degree in Mathematics. Any quotation or excerpt from the published or unpublished work of other persons is explicitly indicated and in each such instance a full reference of the source of such work is given. I have read and understood the Birkbeck College guidelines on plagiarism and in accordance with those requirements submit this work as my own.

Introduction



The imaginary part of the Riemann surface $f(z) = z^2$.

When given a complex function such as $f(z) = z^2$, a problem arises that that visualising this function becomes a non-trivial task as this is not a bijection between the two planes. We are taught to use branch cuts, to consider specific regions of the domain of definition, to avoid these problems. But, with the notion of a Riemann surface, we don't have leave it at that. In the case of the function w = f(z), we have the usual two branches for the w-plane. The idea of a Riemann surface conceptually comes from "stitching" together these two branches,

to form a single sheet with multiple layers, such that if one were to move around in the z-plane, the image in the w-plane should also be free to move about. Thus, as expected, a single point z gets mapped to two points w, as we can imagine that the two layers of the w-plane both cover the point in the z-plane. This example, along with functions like $w = \log(z)$, with its famous infinite tower of sheets, are what were studied by Riemann himself, and what motivated his study of what became known as Riemann surfaces. Riemann initially attempted to generalise the notion of analytic continuation, by defining the notion of a germ of a function which he then used to define Riemann surfaces. We, however, will approach this topic from a more modern perspective, using modern tools and definitions with the goal of somehow classifying all Riemann surfaces in a meaningful way. This classification that we chase is called the uniformisation theorem, and was first rigorously proven independently by Poincarè and Koebe in 1907, though its statement was conjectured much earlier. Before we begin, we wish to remark that the field of Riemann surfaces is the field of some truly remarkable mathematics, for it is one of the rare cases where notions from algebraic geometry, complex analysis and topology and group theory converge to provide equivalent descriptions of the same object, each in their own language. This is truly the domain of beautiful mathematics.

Chapter 1

Preliminaries

In this dissertation, we take \mathbb{F} to mean either the field of Real or Complex numbers. We denote open balls of a metric space, (X,d), by $B_r(x) = \{y \in X | d(x,y) < r\}$ and closed balls by $\overline{B_r(x)} = \{y \in X | d(x,y) \le r\}$. Given an open set U, its closure will be denoted as \overline{U} .

Definition 1.0.1 (Region). Given a topological space X, a region is a non-empty, connected, open subset of X. The closure of a region is called a closed region. In this dissertation, X is taken to be a Riemann surface.

Definition 1.0.2 (Holomorphic, Meromorphic, Entire and Analytic functions). Given a region U in \mathbb{C} , a function f is said to be *holomorphic* if it differentiable at all points $z \in U$. The function f is said to be *analytic* if it can be expressed as a power series, i.e. $f(z) = \sum_{j=0}^{\infty} \lambda_j z^j$ for some $\lambda_j \in \mathbb{C}$. A function whose domain of definition is \mathbb{C} is called *entire*. It is called *meromorphic*, if f(z) is undefined at a discrete subset of points $\Delta \subset U$. A meromorphic function f(z) can be written as the quotient of two coprime, entire polynomials $\frac{p(z)}{q(z)}$ with $q(z) \neq 0$. The set of discrete points Δ corresponds to the zeros of q(z).

Definition 1.0.3 (Paths and Loop spaces). Given a topological space X, a path in X is a continuous map $\gamma \colon [0,1] \to X$, where $\gamma(0)$ and $\gamma(1)$ are the endpoints of γ . If a path has starting and endpoint the same, i.e. $\gamma(0) = \gamma(1) = x$, then it is said to be a loop based at x. The set of loops based at X is called the loop space at x and is denoted by $\Omega_x X$.

Definition 1.0.4 (Fundamental group). Given a topological space X, the fundamental group $\pi_1(X, x)$ at some point $x \in X$ is defined as the quotient of the loop space $\Omega_x X$ by the relation $\gamma_1 \sim \gamma_2$ if and only if γ_1 and γ_2 are homotopic and have the same basepoint. The group operation is concatentation of loops, i.e. for two loops $\gamma_1, \gamma_2 \in \Omega_x X$, we can

concatenate them by writing

$$(\gamma_2 \cdot \gamma_1)(t) = \begin{cases} \gamma_1(2t), & t \in [0, 1/2] \\ \gamma_2(2t-1), & t \in [1/2, 1] \end{cases}$$

Recall that if this fundamental group is trivial for any $x \in X$, then the space X is said to be simply connected.

Definition 1.0.5 (Covering Spaces and the Universal Cover). A continuous map $p: Y \to X$ of topological spaces is called a covering map if there is a collection \mathcal{U} of open sets $U \subset X$, called elementary neighbourhoods with the following property: For each $x \in X$, there is a $U \in \mathcal{U}$ which contains x and such that for each $y \in p^{-1}(x)$, there is a continuous map $q: U \to p^{-1}(U)$ with $p \circ q = Id_U$ and such that q(U) is the path component of $p^{-1}(U)$ containing y. The space Y is called a covering space or cover of X, and may be written as (Y,p). If Y is simply connected, then it is said to be the universal cover of X. A proof of the existence and uniqueness of the universal cover for most topological spaces (subject to certain, technical requirements) can be found in [6, p.63].

Definition 1.0.6 (Covering Transformation). Let X be a topological space and let (Y_1, p_1) and (Y_2, p_2) be two covering spaces of X. A continuous map $F: Y_1 \to Y_2$, such that $p_1 = p_2 \circ F$ is called a covering transformation. If F is a homeomorphism, then it is called a covering isomorphism.

Definition 1.0.7 (Deck Transformation). Given a covering space $p: Y \to X$ of a topological space X, covering isomorphisms $F_i: Y \to Y$ are called covering self-isomorphisms or deck transformations. Deck transformations of a covering space (Y, p) form a group Deck(Y, p), called the deck group of the covering space (Y, p).

Definition 1.0.8 (Normal Cover). A connected covering space $p: Y \to X$ is called *normal* or regular if the deck group acts transitively on $p^{-1}(x)$ for any point $x \in X$.

Definition 1.0.9 (Support of a function). Let U be an open set. Given a function $f: U \to \mathbb{F}$, its support, supp $(f) \subset U$, is defined as the closure of the set $\{x \in U \mid f(x) \neq 0\}$. If this set is compact, then f is said to have compact support contained in U.

Definition 1.0.10 (Bump Function). Let U be an open ball in \mathbb{C} about the origin. We define, for some $\epsilon > 0$, a bump function $\beta : U \to \mathbb{R}$ to be a smooth function with support in U defined as,

$$\beta(z) = \begin{cases} 1, & |z| \le \frac{1}{4}\epsilon \\ 0, & |z| > \frac{3}{4}\epsilon \end{cases}$$

Chapter 2

Riemann Surfaces

2.1 Riemann surfaces and their properties

This chapter is intended as an overview of complex analysis and a brief introduction into Riemann surface theory. The purpose of this chapter is to motivate the problem which is solved by the uniformisation theorem, and to provide some examples of Riemann surfaces so that the reader who has never encountered them before can get an idea of some of their properties. We begin by defining and constructing some Riemann surfaces.

Definition 2.1.1 (Riemann Surface). A Hausdorff topological space X is said to be a Riemann surface if:

- There exists a collection of open sets $U_{\alpha} \subset X$, where α ranges over some index set A, such that $\bigcup_{\alpha \in A} U_{\alpha}$ covers X.
- There exists for each α , a homeomorphism, called a chart map, $\psi_{\alpha} \colon U_{\alpha} \to \widetilde{U}_{\alpha}$, where \widetilde{U}_{α} is an open set in \mathbb{C} , with the property that for all α , β , the composite map $\psi_{\alpha} \circ \psi_{\beta}^{-1}$ is holomorphic on its domain of definition, with a holomorphic inverse map, i.e. both maps are angle preserving or *conformal*. These composite maps are called transition maps or transition functions.

We call the triple $(\{U_{\alpha}\}, \{\widetilde{U}_{\alpha}\}, \{\psi_{\alpha}\})$ an atlas of charts for the Riemann surface X, though we also use the common notation (U, \widetilde{U}, ψ) to denote this, where $U = \{U_{\alpha}\}, \widetilde{U} = \{\widetilde{U}_{\alpha}\}$ and $\psi = \{\psi_{\alpha}\}.$

The condition that the transition functions need to be conformal maps is an example of an imposed structure. Such a structure is called a conformal, or Riemann surface, structure.

If we replace this condition with that of being smooth (infinitely differentiable), then the resultant objects have what is known as a smooth structure. The requirement that our surfaces have conformal transition maps is very strong and has some nice properties as we shall shortly see. Unfortunately, in practice, this definition of Riemann surfaces can be quite difficult to work with as infinite collections of open sets can be difficult to manipulate. Luckily, we can construct some very nice, basic examples using this definition to help us gain an understanding of them. Let us begin with some examples.

Example (The Complex Plane). The complex plane \mathbb{C} can be defined as a Riemann surface. To see this, we pick a cover for the complex plane, say $U = \mathbb{C}$, and just identity map the set U back into the complex plane. Hence the triple (U, U, ϕ) is an atlas for \mathbb{C} where $\phi = Id$.

Similarly, we can also define the upper half plane $\mathbb{H} = \{z \in \mathbb{C} : \Im(z) > 0\}$, as we can cover \mathbb{H} with an open cover, and identity map each open set into the upper half plane on \mathbb{C} . The same construction holds true not just for \mathbb{H} but for all open subsets in \mathbb{C} , thus giving us an infinite family of Riemann surfaces. A perhaps more interesting example of a Riemann surface is that of the Riemann sphere.

Example (The Riemann sphere). The Riemann sphere, commonly denoted as $\widehat{\mathbb{C}}$, can be defined as a topological space $\mathbb{C} \cup \{\infty\}$, where sequences of complex numbers that do not have a convergent subsequence in \mathbb{C} can be said to converge at an additional point, called the point at infinity, denoted by ∞ . This implies our space is sequentially compact and hence compact. We construct an atlas of charts for the Riemann sphere as this will show us that the Riemann sphere is indeed a Riemann surface. The following system of constructing charts on a sphere is known as stereographic projection. Let $U = \{U_1, U_2\}$, with $U_1 = \{z \in \mathbb{C} : |z| < 2\}$ and $U_2 = \{z \in \mathbb{C} : |z| > \frac{1}{2}\} \cup \{\infty\}$. Let $U_1 = \widetilde{U}_1 = \widetilde{U}_2$ and we let $\psi_1 : U_1 \to \widetilde{U}_1$, be the identity map. However, we let $\psi_2 : U_2 \to \widetilde{U}_2$ with $\psi_2(z) = \frac{1}{z}$ and $\psi_2(\infty) = 0$. Both of these maps are clearly holomorphic so all that is needed is to consider if the transition maps composed of these two maps are also holomorphic. It turns out that both $\psi_1 \circ \psi_2^{-1}$ and $\psi_2 \circ \psi_1^{-1}$ are the map $z \mapsto \frac{1}{z}$, which is also holomorphic on its domain of definition $\{z \in \mathbb{C} : \frac{1}{2} < |z| < 2\}$ and so we have our atlas for the Riemann sphere proving it is a Riemann surface.

Note the Riemann mapping theorem [1, p.221], a result that allows us to conformally map open subsets of $\mathbb C$ to one another. We particularly care about the mapping between the open disk $\mathbb D=\{z\in\mathbb C\colon |z|<1\}$ and $\mathbb H$. Both of these are simply connected regions and subsets of $\mathbb C$ so we can find a conformal map between them; such a map looks like $f:\mathbb H\to\mathbb D$ where $f(z)=\frac{z-i}{z+i}$. This map is conformal as we can write the inverse map as $f^{-1}(w)=i\frac{1+w}{1-w}$.

Since we have that, up to conformal equivalence, all simply connected, open strict subsets of \mathbb{C} are equivalent to \mathbb{D} , what is to say that there do not exist more conformal equivalences, perhaps some that equate all the spaces that we have constructed to one another?

Lemma 2.1.2. The Riemann sphere, the complex plane and the upper half plane are **not** conformally equivalent to one another.

Proof. We first distinguish $\widehat{\mathbb{C}}$ from \mathbb{C} and \mathbb{H} . In the case of the Riemann sphere, we note that it is a compact space, whereas the complex plane and upper half plane are non compact spaces. Since, conformal equivalence is a continuous bijection and the image of a compact set is compact, if we were to suppose that we had a conformal equivalence between $\widehat{\mathbb{C}}$ and X (where $X = \mathbb{C}$ or \mathbb{H}), this would imply that X is compact, which is a contradiction. Therefore, the Riemann sphere is not conformally equivalent to either the complex plane or the upper half plane. To distinguish between \mathbb{C} and \mathbb{H} , we note that \mathbb{H} is conformally equivalent to \mathbb{D} so we can pose the question for \mathbb{D} instead. In that case, by Liouville's theorem ([1, p.122]) we have that since \mathbb{D} is bounded, holomorphic functions to \mathbb{D} are therefore constant functions so for any such f we have for all f conformal equivalence exists between \mathbb{C} to be a conformal map, f must also be holomorphic, however, we have that f is not even bijective since f must also be holomorphic, however, we have that f is not even bijective since f to be a \mathbb{C} . Hence, no such conformal equivalence exists between \mathbb{C} and \mathbb{H} .

So we indeed have constructed three unique surfaces up to conformal equivalence, $\widehat{\mathbb{C}}$, \mathbb{C} and \mathbb{H} , i.e. three distinct Riemann surfaces. We now show a more abstract and special property of Riemann surfaces; that is, that their transition maps have a positive Jacobian determinant. This restricts the types of surfaces which can arise as Riemann surfaces to being only the orientable surfaces. Furthermore, since this implies that Riemann surfaces are a subset of orientable surfaces, any result which holds true for orientable surfaces, will hold true for Riemann surfaces.

Lemma 2.1.3 (Riemann surfaces are orientable). Let X be a Riemann surface. Then, the Jacobians of the transition maps all have positive determinant.

Proof. Pick an arbitrary transition map between arbitrary open sets on X, with a non-empty intersection. Let us call this map f. By the definition of transition maps, f is holomorphic on its domain of definition. Hence we can write for any point z = x + iy in the intersection, f(z) = f(x+iy) = u(x,y)+iv(x,y) for some real valued functions u, v, with u, v satisfying the Cauchy-Riemann equations $\frac{\partial u}{\partial x} = \frac{\partial v}{\partial y}$ and $\frac{\partial u}{\partial y} = -\frac{\partial v}{\partial x}$. Thus, the determinant of the Jacobian

 $J[f] \ \text{of our transition map } f \ \text{is } \det(J[f]) = \begin{vmatrix} \frac{\partial u}{\partial x} & \frac{\partial u}{\partial y} \\ \frac{\partial v}{\partial x} & \frac{\partial v}{\partial y} \end{vmatrix} = \frac{\partial u}{\partial x} \frac{\partial v}{\partial y} - \frac{\partial v}{\partial x} \frac{\partial u}{\partial y} = (\frac{\partial u}{\partial x})^2 + (\frac{\partial u}{\partial y})^2 \geq 0,$ where the last equality holds by the Cauchy-Riemann equations. This is strictly positive barring case of f = k for some $k \in \mathbb{C}$, which makes the Jacobian determinant zero. However, we discard this case as the inverse map of such an f would not be holomorphic (i.e. $f^{-1}(k)$ would not be a single point) and thus f cannot be constant. Hence, we have that the Jacobian determinant for any transition map on a Riemann surface is strictly positive. \square

So, using atlases we have constructed a few examples of Riemann surfaces, and have even grouped them into three classes up to conformal equivalence. One may ask the question, are there any more Riemann surfaces that are not conformally equivalent the three that we have constructed? How might they arise? It turns out there are a number of different and interesting methods of constructing Riemann surfaces, and indeed there are infinitely many Riemann surfaces that are not conformally equivalent to the above three. However, as we will see shortly, the above three Riemann surfaces are certainly special in their own right.

2.2 Quotients of the complex plane

One powerful method of constructing Riemann surfaces arises from quotients of topological spaces. In some cases, we can endow these quotient spaces with Riemann surface structures. We recall that in topology, we can construct quotient spaces by taking some topological space X and defining an equivalence relation \sim on it, to create a new space X/\sim . Under the projection map $p\colon X\to X/\sim$ that maps elements of X to their equivalence classes under \sim , this new space also becomes a topological space since we can define a topology on this space (called the quotient topology), by taking $U\subset X/\sim$ to be open if and only if $p^{-1}(U)$ is open.

Definition 2.2.1 (Topological group). A topological group is a topological space $\{G, *\}$ which is also a group, in which the operations

- $m: G \times G \to G$, defined as m(g,h) = g * h
- $i: G \to G$, defined as $i(g) = g^{-1}$

are continuous.

We can see easily that $\{\mathbb{C},+\}$, $\{GL(2,\mathbb{C}),+\}$ are a pair of examples of topological groups, where to see that $GL(2,\mathbb{C})$, the matrix group of 2×2 invertible matrices over \mathbb{C} , is a

topological space we identify it with \mathbb{C}^4 , where each entry of the matrix has its own basis vector. They both become topological groups when endowed with the subspace topology, as they can be considered subspaces of Euclidean space. We also define the notion of a right translation on a topological group.

Definition 2.2.2 (Right Translation). Let G be a topological group. Then for some $g \in G$, we have a map $m_g : G \to G$ such that $m_g(x) = xg$, and this map is a continuous bijection with a continuous inverse $m_{g^{-1}}$. Therefore, these maps, known as right translations, are homeomorphisms and the group of homeomorphisms of a topological group G is transitive.

Since we can map any point in G to any other using these right translations, that means that we can map any point to the identity element $e \in G$. This motivates the following definition.

Definition 2.2.3 (Discrete subgroup). Let G be a topological group. We call the subgroup $\Omega \subset G$ a discrete subgroup of G if it has the property that there exists a neighbourhood U of G in G such that $U \cap \Omega = \{e\}$.

Example (The cylinder). We can construct an infinitely long cylinder by using the notion of a quotient space from topology and endowing the quotient space with a conformal structure. We consider the (additive) subgroup $2\pi\mathbb{Z}$ in \mathbb{C} and form the quotient space $p\colon \mathbb{C}\to \mathbb{C}/2\pi\mathbb{Z}$. This space is homeomorphic to the infinite cylinder $\mathbb{R}\times S^1$. To endow a Riemann surface structure on this infinite cylinder, we consider a disc $B_{\frac{1}{2}}(z)$ around a point $z\in\mathbb{C}$. If $z_1,z_2\in B_{\frac{1}{2}}(z)$ and if $z_1=z_2+2\pi n$ for some $n\in\mathbb{Z}$, then n=0 and hence $z_1=z_2$. This occurs since the radius of the disc is less than π . Hence, the projection map p maps $B_{\frac{1}{2}}(z)$ bijectively into $\mathbb{C}/2\pi\mathbb{Z}$, where the equivalence classes are all the points whose distance is less than half away from z, i.e. collections of points which are of distance less than one away from each other. We can construct an atlas from this data by letting $U=p(B_{\frac{1}{2}}(z))$, $\widetilde{U}=B_{\frac{1}{2}}(z)$ and the local inverses of p, q_z (which map an element $p(B_{\frac{1}{2}}(z))$ to the centre point of the disc z), be the chart maps. We cover $\mathbb{C}/2\pi\mathbb{Z}$ with this data. The transition maps will have the form $(q_z \circ p)(z) = z + 2\pi n$ for some $n \in \mathbb{Z}$. This map is clearly holomorphic with holomorphic inverse and as such we have that the infinite cylinder is a Riemann surface.

Note further that we can find a conformal equivalence between $\mathbb{C} \setminus \{0\}$ and $\mathbb{C}/2\pi\mathbb{Z}$, by ways of the mapping $z \mapsto e^{iz}$. This map never attains zero and hence is holomorphic with holomorphic inverse. Thus the cylinder is conformally equivalent to the punctured plane and so, the punctured plane is also a Riemann surface, obtaining its conformal structure via the

¹To prove this, take the unit square in \mathbb{C} which is itself homeomorphic to \mathbb{R}^2 and consider the identification $[x,0] \sim [x,2\pi]$ for any point $x \in \mathbb{R}$ and show that \mathbb{C}/\sim is homeomorphic to $\mathbb{R} \times S^1$

exponential map.

Example (The complex torus). We begin by picking two different, non zero, complex numbers ω_1, ω_2 and define $\Omega = \Omega(\omega_1, \omega_2) = \{n\omega_1 + m\omega_2 | n, m \in \mathbb{Z}\}$. This object, Ω , is an example of a lattice in \mathbb{C} , an additive subgroup of \mathbb{C} . As proven in [7, p.58], such lattices are isomorphic as addititve groups to $\mathbb{Z} \times \mathbb{Z}$ and so by [6, Example 1.43], this quotient \mathbb{C}/Ω is topologically a torus. But we want to give this torus a conformal structure. We note that we can write this lattice as $\Omega(1,\omega)$ where $\omega = \frac{\omega_i}{\omega_j}$ where $i,j = \{1,2\}$ with $i \neq j$, chosen so that $\Im(\omega) > 0$. We call this number ω the modulus of the complex torus. Working with this modified lattice, we proceed by picking a point and considering a disc of radius r around z, called $B_r(z)$. The radius is yet to be determined, though we need to ensure that, in the same way as before, we can pick points to be "sufficiently close" to our chosen point z so that we can define our open sets properly and choose appropriate chart maps. Picking two different points $z_1, z_2 \in B_r(z)$ if $z_1 = z_2 + n + \omega m$ for $n, m \in \mathbb{Z} \setminus \{0\}$ then, if we want $z_1 = z_2$ (implying n=m=0) then we need that $2r < \min_{n,m} |n+m\omega|$. In fact, since for $m \neq 0$ we have $|n+m\omega| > \Im(\omega)$ and if m=0 but $n\neq 0$ then we have $|n+m\omega| \geq 1$, then we can get a better bound on r, namely that $2r < \min(1, \Im(\omega))$. Then, as before, p maps $B_r(z)$ bijectively to \mathbb{C}/Ω , so we can construct an atlas for this torus, with charts being $U = p(B_r(z))$, with $U = B_r(z)$ and the chart maps being the local inverses to the covering maps. So we see that, in much the same way as in the case of the cylinder, we can bestow a conformal structure on this torus, hence making the torus a Riemann surface. In particular this is a compact Riemann surface. If we construct a torus in such a way, we call it a complex torus, and represent it by \mathbb{T} or if the number ω is important \mathbb{T}_{ω} .

In both of these cases, our strategy was to find a discrete subgroup of the additive group \mathbb{C} . This method, by which we search for discrete subgroups of the space we are quotienting, is our main tool when it comes to studying quotient spaces of Riemann surfaces. As we will eventually show, all but three Riemann surfaces arise as a quotient of some other Riemann surface, and in all cases, the subgroups which define the Riemann surface are discrete.

Going back to our complex tori for a moment, we see that in the process of creating our conformal structure on the torus we associated a complex number ω , called the modulus of the lattice Ω , that is an element of the upper half plane. This makes us wonder, does this number hold any special meaning? We have that seemingly unrelated spaces, like $\mathbb{C}/2\pi\mathbb{Z}$ and $\mathbb{C}\setminus\{0\}$ are conformally equivalent, so perhaps complex tori are also conformally equivalent. It turns out that this is not quite the case. Showing this fact, is however a little trickier. The crucial point is that, if two tori have the same modulus, then they are conformally equivalent

by [7, Theorem 6.1.4]. An in depth discussion of this can be found in [7, Section 6.1]. What is striking is that the space of all such moduli itself forms a space called a Teichmüller or Moduli space (of genus 1 surfaces). This Moduli space obtains a Riemann surface structure, where each point of this space is an equivalence class of complex tori with a particular modulus.

So if we can have an infinite number of distinct conformal structures on the complex torus, does the same hold true on the Riemann sphere? We investigate this problem below.

2.3 Maps between Riemann surfaces

As with any mathematical theory, we want to investigate maps between our objects. Through composing our maps with the charts of the surfaces, we can define well behaved maps. However, globally, these maps may differ somewhat from their local descriptions. We begin with a powerful lemma and corollary about holomorphic maps between Riemann surfaces.

Lemma 2.3.1. Let X and Y be connected Riemann surfaces and $F: X \to Y$ a non-constant holomorphic map. For each point $x \in X$, there is a unique integer $k_x \ge 1$ such that we can find charts around x in X and F(x) in Y such that $F(x) = x^{k_x}$.

Sketch Proof. This lemma gives that, locally, holomorphic maps between Riemann surfaces can be represented by power maps. A proof of this lemma can be found in [4, Proposition 5]. To sketch the idea of the proof, we see that by letting (U, \widetilde{U}, ϕ) be the chart around $z \in X$, with $\phi(z) = 0 \in \widetilde{U} \subset \mathbb{C}$, which we can choose by just linearly translating the chart maps, and by letting (V, \widetilde{V}, ψ) be a chart about F(x) such that $\psi(F(x)) = 0 \in \widetilde{V} \subset \mathbb{C}$, similarly, gives us that the composite $(\psi \circ F \circ \phi^{-1})(x) = x^{k_x}$, as in the charts, F is represented by holomorphic function f with a power series, in which we find all terms except for one vanish.

Corollary 2.3.2. Let $F: X \to Y$ be a non-constant holomorphic map between connected Riemann surfaces. Then there is a discrete subset $R \subset X$ of points in X where $k_x > 1$.

Proof. In local coordinates about each point $p \in X$ the function F can be represented by the function f. Since f can be written as a power map, we see that for all the points p where $k_x = 1$ the derivative of f is constant and cannot be zero. So for points where $k_x > 1$, the derivative of f gives the function $k_x x^{k_x-1}$, so the zeros of the derivative of f, which is also a holomorphic function, give us the points in f which have f have that this set of points is discrete.

When considering holomorphic maps between Riemann surfaces, frequently such maps are

taken to be proper, that is, maps for which the preimage of a compact set is compact. Whilst it is not strictly necessary to do so, it is natural to study proper maps between topological spaces that may themselves not be compact, as it allows us to define and study the idea of the "ends" of a topological space². We discuss this in a bit more depth at the end of this dissertation, though for now, we accept that for a richer theory of maps between our topological spaces with Riemann surface structures on them, we wish them to be proper. Thus, we append the requirement that such maps be proper to our standard requirements from complex analysis and Riemann surface theory (such as holomorphicity).

So what does a map being proper give us? If we let $F: X \to Y$ be a non-constant, proper, holomorphic map between connected Riemann surfaces and let R represent the set of points in X where $k_x > 1$, then, rather amazingly we have that the image $\Delta = F(R)$ is discrete in Y and that for any $y \in Y$ the preimage $f^{-1}(y)$ is also a finite subset of X. The fact that Δ is a discrete subset of Y is analogous to our example in the introduction. We recall that under the mapping $w=z^2$, we had that for each value z there were two values w. These two values formed "sheets" over the z plane. We use this idea to say that F has the property that if it is a proper, non-constant holomorphic map, then it must have only finitely many sheets. Thus we note that the points in R are particularly special when studying F and are called the critical points of f, with the points in F(R) called the branch or ramification points. The value k_x is called the multiplicity of x under F. If, for example, we let $f: \widehat{\mathbb{C}} \to \widehat{\mathbb{C}}$ such that $f(z) = z + \frac{1}{z}$, then this is a map where every point has multiplicity 2, since if p is a particular branch point under f, then it is a solution to $p=z+\frac{1}{z}$, or the polynomial $z^2 - pz + 1 = 0$. Since p is a branch point, then both branches of f must converge on p so it is a repeated root. Considering the discriminant of this polynomial gives $p^2-4=0$ implying the branch points on $\widehat{\mathbb{C}}$ under the map f are $\{2, -2\}$.

We will unfortunately not go into the deep topic of covers of Riemann surfaces with ramification points, though an excellent resource on the topic of ramified covers of the Riemann sphere (and a general theory of branched covers of surfaces) can be found in [9]. However, what is important to our discussion is the following.

Definition 2.3.3. Given a proper, non-constant holomorphic map between Riemann surfaces $F: X \to Y$, with Y connected, we can define for all $y \in Y$ an integer

$$d(y) = \sum_{x \in F^{-1}(X)} k_x$$

²First introduced by Freudenthal on page 701 of his 1931 paper Über die Enden topologischer Räume und Gruppen for the journal Mathemtische Zeitschrift 33, 692-713

Since the sum runs over a finite set, by the above lemma, d(y) is finite. Furthermore, this integer d is independent of chosen point $y \in Y$, a proof of which can be found in [4, Proposition 7], and thus is an invariant of the map F. We call this d the d-egree of the map F. The theory thus developed gives us a nice extension of our definition of meromorphic functions on Riemann surfaces.

Definition 2.3.4 (Redefining Meromorphic functions). Given a Riemann surface X, a meromorphic function on X is said to be a function $F: X \to \widehat{\mathbb{C}}$, that is not always equal to ∞ . The poles of F are the points $F^{-1}(\infty)$ in X, and the order of a pole x is the integer k_x .

We can therefore say something important for our study of Riemann surfaces.

Proposition 2.3.5. Let X be a compact, connected Riemann surface. If there is a meromorphic function on X with exactly one pole, and that pole has order 1, then X is conformally equivalent to the Riemann sphere.

Proof. Let $F: X \to \widehat{\mathbb{C}}$ be the meromorphic function in the hypothesis. Since X is compact, we have that F is proper, i.e. the preimage of any compact set is compact. The degree of this map is 1 since the pole is mapped to once. Thus for all $y \in \widehat{\mathbb{C}}$, there exists exactly one $x \in X$. Thus, by the open mapping theorem, F is a bijection. The inverse map of F is continuous since for any closed set $U \subset X$ we have that F(U) is compact in $\widehat{\mathbb{C}}$, and thus closed. Therefore, F is a homeomorphism. Thus, since F has degree one, then in a chart centered on each point $p \in X$ $F'(p) \neq 0$, and so by the inverse function theorem F^{-1} is also holomorphic.

Hence, we see that any two complex structures on the Riemann sphere are compatible (meaning that we can pass from one to the other conformally) and so there exists exactly one conformal structure on $\widehat{\mathbb{C}}$ up to conformal equivalence.

2.4 Covering spaces of Riemann surfaces

We begin this section with a strong claim: All Riemann surfaces arise as a quotient of a simply connected Riemann surface by some group action. This means that if we can understand simply connected Riemann surfaces then we can perhaps construct all possible Riemann surfaces. In particular, we wish to study how these simply connected surfaces can return to us a general Riemann surface.

2.4.1 Automorphisms of Riemann surfaces

Definition 2.4.1 (Automorphisms and the Automorphism group of a Riemann surface). Let X be a Riemann surface. A function $f: X \to X$ is called an automorphism if it is a homeomorphic conformal equivalence. The set of all automorphisms of X form a group under function composition, written as Aut(X).

Lemma 2.4.2. Let $p: Y \to X$ be a covering of a Riemann surface. Then there is a unique complex structure on Y such that $p: Y \to X$ is holomorphic.

Proof. Since X is a Riemann surface, we can choose a cover of X such that the sets $\{U_{\alpha}\}$ are small enough to be mapped to homoeomorphically from sets $\{V_{\alpha}\}$ on Y (this process is described in depth in [7, Theorem 4.11.2]). We let $\{\phi_{\alpha}\}$ be the chart maps for each U_{α} from X, and hence we can extend these to define charts from Y as $\psi_{\alpha} = \phi_{\alpha} \circ p$, for each α . Hence, we define the atlas for Y to be the triple $(\{V_{\alpha}\}, \{\tilde{V}_{\alpha}\}, \{\psi_{\alpha}\})$, where $\tilde{V}_{\alpha} = \tilde{U}_{\alpha}$, the open sets in the complex plane from the atlas of X. The transition maps of this atlas are conformal, since if we take two sets V_{α} and V_{b} such that $V_{\alpha} \cap V_{b} \neq \emptyset$, then we have that the transition function is $\psi_{\alpha} \circ \psi_{b}^{-1} = (\phi_{\alpha} \circ p) \circ (\phi_{b} \circ p)^{-1} = \phi_{\alpha} \circ p \circ p^{-1} \circ \phi_{b} = \phi_{\alpha} \circ \phi_{b}^{-1}$ which is a conformal map since this is a transition function of X itself. Similarly, for any two atlases on X, we get that the atlases they induce on Y are compatible, and so we have a well defined conformal structure on Y.

To show that p is holomorphic, we apply a coordinate transformation to any set on Y, say $V \in \{V_{\alpha}\}$. Then, we have that $\psi_{V} \circ \psi_{V}^{-1} = \phi_{U} \circ p \circ p^{-1} \circ \phi_{U}^{-1} = \text{Id.}$ So these maps are conformal equivalences, implying that p must be a holomorphic map. Uniqueness of the complex structure follows from the fact that for any two atlases, the projections will map the same points in Y to the same point in X.

Corollary 2.4.3. Let $p: Y \to X$ be a covering of a Riemann surface. Then, each deck transformation $F: Y \to Y$ is an automorphism of Y as a Riemann surface.

Proof. Let $g \in \text{Deck}(Y, p)$, and points $y_1, y_2 \in Y$ such that $y_2 = g(y_1)$. Pick open sets around each point so $y_1 \in V_1$ with chart map ψ_1 and $y_2 \in V_2$ with chart map ψ_2 . Then we get the following transformation of coordinates

$$\psi_2 \circ g \circ \psi_1^{-1} = (\phi_2 \circ p) \circ g \circ (p^{-1} \circ \phi_1^{-1})$$
$$= \phi_2 \circ (p \circ g) \circ p^{-1} \circ \phi_1^{-1}$$
$$= \phi_2 \circ \phi_1^{-1} \quad \text{since } p \circ g = p$$

However, $\phi_2 \circ \phi_1^{-1}$ is a conformal transformation, and so g is holomorphic. Replacing g with $h = g^{-1}$ in the above proof gives that g^{-1} is also holomorphic and hence g is a conformal transformation of Y. Therefore, it is an automorphism of Y.

These two results give us that the deck group of a cover of a Riemann surface is precisely the group of automorphisms of the cover. This covering space is also a Riemann surface.

Lemma 2.4.4. [6, Proposition 1.39] A topological covering space $p: Y \to X$ is normal if and only if $p_*\pi_1(Y,y)$ is a normal subgroup in $\pi_1(X,x)$ for some $y \in p^{-1}(x)$.

This lemma, though not hard to prove, is a standard result and so we refer the reader to a proof. Note that the author of [6], uses the notation G(X) to mean Deck(Y, p) for the covering space $p: Y \to X$. What is important however is the following corollary.

Corollary 2.4.5. If X is a connected Riemann surface, then the universal covering surface $p: Y \to X$ is a regular covering surface of X.

Proof. Since Y is a universal cover of X that implies that it is simply connected, and hence $\pi_1(Y,y) = 1$ for all $y \in Y$. Therefore, it is trivially a normal subgroup in $\pi_1(X,x)$, for all $x \in X$, under the induced map $p_* \colon \pi_1(Y,y) \to \pi_1(X,x)$. Hence, by lemma 2.4.4, we have that the universal covering surface Y of a Riemann surface X is a regular covering of X. \square

Since we have that every Riemann surface admits a universal cover, and that it is a regular cover, then we would like to link that fact to these groups of covering transformations. In particular, we would like to show that the group of automorphisms of the universal cover of a Riemann surface can be associated with the space it covers.

Proposition 2.4.6. Let $p: Y \to X$ be a regular covering of X. Let G = Deck(Y, p) be the group of deck transformations of the cover (Y, p). Then there is a homeomorphism $q: X \to Y/G$ given by $q(x) = [\tilde{x}]_G$ where $x \in X$ and $\tilde{x} \in p^{-1}(x)$, with $[\cdot]_G$ meaning orbit under the action of G.

With this proposition, we can identify all Riemann surfaces with automorphisms of their universal covers, by setting Y to be the universal cover of X and then $G \leq \text{Deck}(Y, p)$, to be some subgroup of the group of covering transformations of the universal cover.

Proof. Let $\pi: Y \to Y/G$ be the projection, sending each point \tilde{x} to its orbit $[\tilde{x}]_G$. Hence, we define q as $q(x) = \pi(\tilde{x})$, where $p(\tilde{x}) = x$. We need to show that this only depends of x and not on our choice of $\tilde{x} \in p^{-1}(x)$.

Let $y_1, y_2 \in p^{-1}(x)$. Since we have that (Y, p) is a regular covering, the deck group acts transitively on $p^{-1}(x)$. Therefore, we have a $g \in G$ such that $g(y_1) = y_2$. Hence we have that $\pi(y_1) = \pi(g(y_1)) = \pi(y_2)$ as needed. Thus, $q: X \to Y/G$ is a well defined, surjective function.

We want q to be a bijection, so we pick two points $x_1, x_2 \in X$ such that $q(x_1) = q(x_2)$. We want to show they are the same point. Consider that for some $y_1 \in p^{-1}(x_1)$ and $y_2 \in p^{-1}(x_2)$, we have that $\pi(y_1) = \pi(y_2)$ and so we have $g \in G$ such that $y_2 = g(y_1)$. Therefore, $x_1 = p(y_1) = p(g(y_1)) = p(y_2) = x_2$. Hence we have that q is a bijection. Since p is a continuous map and $\pi: Y \to Y/G$ is a continuous map from the definition of the quotient topology on Y/G, then so is q. Since q is a surjection, the image of an open set under π is open, and hence p and q also have this property. Thus, q is a homeomorphism. \square

Furthermore, using this homeomorphism q, we can construct a conformal structure on Y/G, by defining charts on Y/G to have the form $(q(U_{\alpha}), \widetilde{U}_{\alpha}, \phi_{\alpha})$, where $(U_{\alpha}, \widetilde{U}_{\alpha}, \phi_{\alpha})$ are the charts on X. Thus, q becomes a conformal equivalence with the projection $\pi: Y \to Y/G$ being holomorphic. We conclude by showing that these subgroups act without fixed points, except under the action of the identity element. This property makes finding the appropriate subgroups that we seek a bit easier.

Definition 2.4.7 (Discontinuous group action). A group of homeomorphisms G acts discontinuously on a topological space X if for every $x \in X$, there exists a neighbourhood V of x such that $V \cap g(V) = \emptyset$ for all non identity elements $g \in G$. This implies that the action of a non identity element in G can have no fixed points in X.

Lemma 2.4.8. If (Y,p) is a connected covering space for a Riemann surface X, then Deck(Y,p) acts discontinuously on Y.

Proof. Let $g \in \text{Deck}(Y, p)$ and suppose for some $y \in Y$ we have an elementary neighbourhood V of y such that $V \cap g(V) \neq \emptyset$, i.e. there exists a $v \in V$ with $g(v) \in V$. Since $p \circ g = p$ and p is injective on V then we have that v = g(v) and thus v is a fixed point of g. Since Y is connected, and the fixed point set is both open and closed (by [7, Lemma 4.19.1]), then the whole cover Y must be the fixed point set of g. Thus g must be the identity element. \square

Thus we have shown that we can reduce the problem of classifying Riemann surfaces to the problem of studying the automorphisms of connected, simply connected Riemann surfaces. This is because, as we saw, the deck transformations of a covering space of a Riemann surface are precisely the automorphisms of the cover. These automorphisms and their subgroups play

a pivotal role in the classification of Riemann surfaces as they define exactly which Riemann surfaces can be constructed from a particular simply connected Riemann surface. Thus we need to know how many simply connected Riemann surfaces there are. This problem is solved by the uniformisation theorem. It turns out the three spaces we discovered earlier, $\widehat{\mathbb{C}}$, \mathbb{C} , \mathbb{H} are indeed the only simply connected Riemann surfaces up to conformal equivalence. Proving this point is the focus of the next portion of this dissertation. For now, let us conclude by computing the automorphism groups of these three simply connected spaces as this will give us a good idea of which spaces arise as a quotient of which spaces.

2.4.2 Automorphism groups of simply-connected Riemann surfaces

We begin by computing the automorphism group of \mathbb{C} , as it is the simplest case.

Theorem 2.4.9. The automorphism group of the complex plane is $Aut(\mathbb{C}) = \{f(z) = az + b \mid a, b \in \mathbb{C}, a \neq 0\}$

Proof. Let f(z) = az + b such that $a \neq 0$ and $a, b \in \mathbb{C}$. This is clearly a holomorphic function so we just need to find a holomorphic inverse. Let w = f(z). Then we have that if $g(w) = \frac{w-b}{a}$ then g(f(z)) = z and as such $f^{-1}(z) = g(z)$. The element g is clearly also holomorphic since $a \neq 0$ and linear functions are clearly automorphisms of \mathbb{C} . Now we wish to show that they are the only possible maps that give automorphisms of \mathbb{C} .

Since f(z) is an automorphism, it is crucially a holomorphic function and hence admits a series expansion about any point $z \in \mathbb{C}$. So, we write $f(z) = \sum_{j=0}^{\infty} \lambda_j z^j$ for $\lambda_j \in \mathbb{C}$ for all j. Writing $g(z) = f(z^{-1})$ gives that g is also holomorphic and injective but only on the annulus 0 < |z| < 1. Hence, g has a singularity (perhaps many), at 0, and hence f has one at ∞ . So we need to classify this singularity, to try and understand f. If f had an essential singularity, then by the Great Picard Theorem ([2, p.300]), f would fail to be injective, which it is because it is an automorphism. If the singularity were removable, then by Liouville's theorem ([1, p.122]), f would be a constant function, which has no holomorphic inverse. Hence, f must have a pole at ∞ . If for f is f we had that f is f then f would fail to be injective on f as every point would be mapped to by multiple other points. Hence f and cannot be constant. Therefore, f has a simple pole at infinity and f is f and f and f and f are f has a simple pole at infinity and f and f is f and f and f are f and f are f has a simple pole at infinity and f and f is f and f and f are f are f and f are f are f and f are f and f are f and f are f and f are f are f and f are f and f are f and f are f are f and f are f and f are f are f and f are f and f are f and f are f and f are f a

Theorem 2.4.10. The automorphism group of the Riemann sphere is $Aut(\widehat{\mathbb{C}}) = \{ f : \widehat{\mathbb{C}} \to \widehat{\mathbb{C}} \mid f(z) = \frac{az+b}{cz+d}, \ ad-bc \neq 0, \ a,b,c,d \in \mathbb{C} \}$

Proof. If $f(z) = \frac{az+b}{cz+d}$, for complex numbers as in the hypothesis, then f is clearly a bijective, and holomorphic map to the Riemann sphere, as it has a single pole at $z = -\frac{d}{c}$. This map also has a holomorphic inverse to the Riemann sphere, as $f^{-1}(z) = \frac{b-zd}{cz-a}$, which also has a single pole at $z = \frac{a}{c}$. So we have that $\{f(z) = \frac{az+b}{cz+d} \mid ad-bc \neq 0, \ a,b,c,d \in \mathbb{C}\} \subseteq Aut(\widehat{\mathbb{C}})$. For the converse, given an $f \in Aut(\widehat{\mathbb{C}})$, we have that f is a holomorphic bijection with a holomorphic inverse on the Riemann sphere. Such maps, are rational maps, which from complex analysis, we know are polynomials of the form f(z) = P(z)/Q(z) for coprime polynomials P and Q. But P(z) and P(z) be linear maps. Therefore, we have that $P(z) = \frac{az+b}{cz+d}$, for $P(z) = \frac{az+b}{cz+d}$, and so the sets are equal. $P(z) = \frac{az+b}{cz+d}$.

Corollary 2.4.11. We can identify automorphisms of the Riemann sphere with a multiplicative matrix group, called $PSL(2,\mathbb{C}) = \left\{ \begin{pmatrix} a & b \\ c & d \end{pmatrix} \mid ad - bc \neq 0, \ a,b,c,d \in \mathbb{C} \right\}$, the set of 2×2 , invertible, complex valued matrices.

Proof. Consider the map $\Phi: Aut(\widehat{\mathbb{C}}) \to PSL(2,\mathbb{C})$. Let $f \in Aut(\widehat{\mathbb{C}})$, so $f(z) = \frac{az+b}{cz+d}$. We have that

$$\Phi\left(\frac{az+b}{cz+d}\right) = \begin{pmatrix} a & b \\ c & d \end{pmatrix}$$

is a bijective map between the two groups with a bijective inverse. Since $Aut(\widehat{\mathbb{C}})$ is a group under function composition, we have that Φ is actually a group isomorphism. As such, $Aut(\widehat{\mathbb{C}})$ and $PSL(2,\mathbb{C})$ are isomorphic as groups and so we can identify automorphisms of the Riemann sphere with elements of $PSL(2,\mathbb{C})$.

Recall that subgroups of the automorphism groups of simply connected Riemann surfaces must act discontinuously for us to be able to create a quotient Riemann surface, that is, crucially, without fixed points. However, when we consider an automorphism of the Riemann sphere, we see that every automorphism has at least one fixed point. Thus the only subgroup of $Aut(\widehat{\mathbb{C}})$ which we can quotient $\widehat{\mathbb{C}}$ by is the trivial subgroup. Thus, $\widehat{\mathbb{C}}$ is the universal cover of only itself, and no other Riemann surfaces can arise as a quotient of the Riemann sphere by a discrete subgroup of $Aut(\widehat{\mathbb{C}})$.

When considering the automorphisms of the complex plane, recall that we showed that $Aut(\mathbb{C}) = \{f : \mathbb{C} \to \mathbb{C} \mid f(z) = az + b, \ a, b \in \mathbb{C}, \ a \neq 0\}$. We see that this has a fixed point for z = b/(1-a), if $a \neq 1$. So our subgroup of automorphisms of \mathbb{C} which act discontinuously on \mathbb{C} are of the type f(z) = z + b. Thus, our automorphisms depend on b. But we showed,

that if b is a pure real or pure imaginary number, then the action on the plane gives us the infinite cylinder, which can also be conformally mapped to the punctured plane. If b is a complex number, then the action gives us a complex torus. The last case, if b is zero, then we have the identity map, which generates the trivial subgroup of $Aut(\mathbb{C})$, which gives us the complex plane itself. Thus, in our introductory examples above, it turns out that we classified all possible Riemann surfaces which have the complex plane as a universal cover!

What remains is to study the quotients of the upper half plane. What is more striking is that, if any other Riemann surfaces exist, all such Riemann surfaces have the upper half plane as a universal cover. So what does the automorphism group of \mathbb{H} look like?

Theorem 2.4.12. The automorphism group of the upper half plane \mathbb{H} is

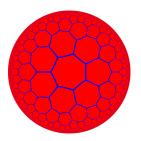
Aut(
$$\mathbb{H}$$
) = $PSL(2,\mathbb{R})$ = $\left\{ \begin{pmatrix} a & b \\ c & d \end{pmatrix} \middle| ad - bc = 1, \ a,b,c,d \in \mathbb{R} \right\}$. The composition of automorphisms of \mathbb{H} corresponds to the multiplication of elements in $PSL(2,\mathbb{R})$.

We will not prove this fact, though the proof of this can be found in [7, Theorem 4.17.3 (iii)]. In fact, by [7, Theorem 2.9.1], we can identify $PSL(2,\mathbb{R})$ with $\{f(z) = \frac{az+b}{cz+d} \mid ad-bc = 1, a,b,c,d \in \mathbb{R}\}$, using the same isomorphism as in the above corollary. So, if we quotient \mathbb{H} by discrete subgroups of $PSL(2,\mathbb{R})$, then we get not just new Riemann surfaces, but all other Riemann surfaces. The discrete subgroups of $Aut(\mathbb{H})$ are called Fuchsian groups. Before we can see some examples we must first give a topology to $PSL(2,\mathbb{R})$, by associating it with a quotient space of the subspace $X = \{(a,b,c,d) \in \mathbb{R}^4 \mid ad-bc = 1\} \subset \mathbb{R}^4$ under the homeomorphism $\delta: X \to X$ where $\delta(a,b,c,d) = (-a,-b,-c,-d)$. The group $PSL(2,\mathbb{R})$, is then given the quotient topology of X under the action of δ . We can now give some simple examples of Fuchsian groups.

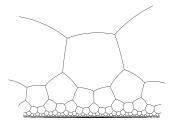
The map $z \mapsto \lambda z$ with $\lambda > 1$ gives a discrete subgroup of $PSL(2,\mathbb{R})$ as this corresponds to the subgroup of matrices where only the lead diagonal is non-zero and $\frac{a}{d} > 1$. Another example includes $z \mapsto z + 1$, which consists of the upper triangular matrix subgroup of $PSL(2,\mathbb{R})$. Both of these examples lead to discrete subspaces in the above topology of $PSL(2,\mathbb{R})$. A less trivial example is that of the modular group $PSL(2,\mathbb{Z})$. This is clearly a subgroup of $PSL(2,\mathbb{R})$. We need to show it is a discrete subgroup, though given the above topology for $PSL(2,\mathbb{R})$ we see that $PSL(2,\mathbb{Z})$ is a lattice of points in $PSL(2,\mathbb{R})$, and thus it is clearly a discrete subgroup of $PSL(2,\mathbb{R})$, and thus, a Fuchsian group.

Another important example that we shall not fully flesh out here is that of the triangle groups. A triangle group is formed by considering the reflections of a hyperbolic triangle in \mathbb{D} . These groups also form Fuchsian groups, which we denote by the group presentation $\langle X, Y | X^{m_1} = Y^{m_2} = (XY)^{m_3} = e \rangle$, where the values m_i are given by the internal angles of the hyperbolic triangle in \mathbb{D} and the generators X and Y corresponding to particular pairs of reflections across edges of the hyperbolic triangle. Thus, $\sum_{i=1}^{3} \frac{1}{m_i} < 1$. These triangles are called (m_1, m_2, m_3) triangles and the groups they generate, (m_1, m_2, m_3) groups. One example is that of the (2, 3, 7) triangle group. These triangle groups provide a nice link between hyperbolic geometry and Riemann surfaces and a simple proof that they are Fuchsian can be found in [7, p.237]. Note that these triangle groups clearly tessellate \mathbb{D} as the elements of the group reflect the triangle along its edges, so we compose these reflections to tessellate \mathbb{D} .

Whilst it was somewhat simple to describe that Fuchsian groups generate all but a few Riemann surfaces, Fuchsian groups themselves are difficult to understand and their properties vary wildly. As we saw with the triangle groups, a nice property of theirs is that they define tessellations of \mathbb{H} . We conclude this section by presenting some tessellations of both \mathbb{H} and \mathbb{D} which arise by acting on regions in these spaces by Fuchsian groups.

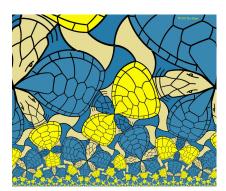






The same regular tiling of \mathbb{D} and \mathbb{H} by acting upon a heptagon on the left and right. This action can be described by the "dual" tiling by hyperbolic (2,3,7) triangles in the middle (notice that 14 triangles fit inside a heptagon). Image source: https://en.wikipedia.org/wiki/Heptagonal_tiling, Accessed: 20/08/2021.





The same regular tiling of \mathbb{D} and \mathbb{H} by the artist M.C. Escher, achieved by acting upon the turtle figure by some elements of a discrete subgroup of $PSL(2,\mathbb{R})$.

Image source: http://www.josleys.com/show_gallery.php?galid=290 , Accessed: 20/08/2021.

Chapter 3

Calculus on Riemann Surfaces

3.1 Calculus on smooth, orientable surfaces

We begin our study of calculus on Riemann surfaces by recalling some basic ideas from the theory of calculus on manifolds. This section takes the work found in Chapter 5 of [4], Chapters 2, 3, 8, 9 and 10 of [11] and Chapters 4 and 5 of [15], to provide a comprehensive set of consistent definitions and examples that should give the reader a good overview of this broad subject. The reader ideally should have a prior familiarity with at least differential forms and one homology theory, though the above quoted references make for excellent learning resources, especially [11].

In the following definitions S will denote a smooth, orientable surface. Recall that Riemann surfaces are examples of smooth, orientable surfaces.

Definition 3.1.1 (Smooth paths). Let $\epsilon > 0$. Then we say $\gamma : [-\epsilon, \epsilon] \to S$ is a smooth path on S if we can define $\frac{d\gamma}{dt}$ for all $t \in (-\epsilon, \epsilon)$.

Definition 3.1.2 (Tangent Space). Let $p \in S$ be a point on S. We define T_pS , the tangent space of S at p, to be the space of equivalence classes of smooth paths γ through p such that $\gamma(0) = p$. Two paths γ_1, γ_2 are said to be equivalent if $\frac{d\gamma_1}{dt} = \frac{d\gamma_2}{dt}$.

Definition 3.1.3 (Cotangent Space). For all $p \in S$ we define the real cotangent space at p to be $T_p^*S = \operatorname{Hom}_{\mathbb{R}}(T_pS, \mathbb{R})$ and the complex cotangent space at p to be $T_p^*S^{\mathbb{C}} = \operatorname{Hom}_{\mathbb{R}}(T_pS, \mathbb{C})$, where $\operatorname{Hom}_{\mathbb{R}}(T_pS, \mathbb{F})$ is the real vector space of linear maps from the tangent space to the field \mathbb{F} .

We state the following definitions in such a way that will be useful to us, however they all

generalise to the case of n-dimensional manifolds.

Definition 3.1.4 (Alternating Space). Let V be a \mathbb{F} vector space. A k-linear map $\omega : V^k \to \mathbb{F}$ is said to be alternating if $\omega(v_1, \ldots, v_k) = 0$ when $v_i = v_j$ where $i \neq j$. The vector space of alternating, k-linear maps is denoted by $\mathrm{Alt}^k(V)$.

Definition 3.1.5 (Exterior Product). Let V be a \mathbb{F} vector space. Let $\omega_1 \in \mathrm{Alt}^p(V)$ and $\omega_2 \in \mathrm{Alt}^q(V)$. We define the exterior product $\wedge : \mathrm{Alt}^p(V) \times \mathrm{Alt}^q(V) \to \mathrm{Alt}^{p+q}(V)$ as

$$(\omega_1 \wedge \omega_2)(v_1, \dots, v_{p+q}) = \sum_{\sigma \in S(p,q)} sign(\sigma)\omega_1(v_{\sigma(1)}, \dots, v_{\sigma(p)})\omega_2(v_{\sigma(p+1)}, \dots, v_{\sigma(q)})$$

where S(p,q) is defined as a permutation of $\{1,\ldots,p+q\}$ such that $\sigma(1)<\ldots<\sigma(p)$ and $\sigma(p+1)<\ldots<\sigma(p+q)$.

The alternating space $Alt^k(V)$ becomes an alternating \mathbb{F} algebra when considering it with the associative exterior product. Refer to [11, p.11] for a more in-depth treatment.

Definition 3.1.6 (Differential k-forms). Let U be an open set in S. A differential k-form on U is a smooth map $\omega: U \to \mathrm{Alt}^k(S)$. The vector space of all such maps, over the field $\mathbb F$ on U is denoted by $\Omega^k(U; \mathbb F)$. When $\mathbb F$ can be any field, we write $\Omega^k(U)$.

Definition 3.1.7 (Support of a k-form). Let ω be a k-form on an open set U in S. The support of ω is defined as $\operatorname{supp}(\omega) = \overline{\{p \in S | \omega(p) \neq 0\}} \subset U$.

Definition 3.1.8 (Differential k-forms of compact support). Let U be an open set in S. A differential k-form with compact support is a smooth map $\omega : U \to \operatorname{Alt}^k(S)$, where $\operatorname{supp}(\omega) \subset U$. The vector space of differential k-forms with compact support in U, is denoted as $\Omega_c^k(U)$. Note that we can extend ω onto the whole of S by setting $\omega = 0$ at all points $p \in S \setminus \operatorname{supp}(\omega)$. We call this process, extending over S by zero, or just extending by zero for short.

Note, if our space S is compact, then for all $p \in \{0, 1, 2\}$ we have that $\Omega_c^p(S) = \Omega^p(S)$. The space of differential 0-forms is also sometimes written as $C^{\infty}(U)$. We also note that if k > 2 then $\mathrm{Alt}^k(\mathbb{R}^2) = 0$ and as such, in our case, we only need to study 0,1 and 2-forms. An important note is the differential 1-forms evaluated at a point p in an open set U of \mathbb{R}^2 can be considered as elements of the cotangent space at p^{-1} . Similarly, for differential 2-forms, we can consider such objects as elements of the 1-dimensional algebra called the external

¹Such elements are sometimes called vectors of the cotangent space, or co-vectors when they need to be distinguished from vectors in the tangent space at p.

algebra of the cotangent space of U at p, represented symbolically as $\Lambda^2 T_p^*(U)$. If $T_p^*(U)$ has as basis elements dx_1, dx_2 then the basis element of $\Lambda^2 T_p^*(U)$ is $dx_1 \wedge dx_2$, with dx_i defined as the exterior derivative of the coordinate functions x_1 and x_2 .

Definition 3.1.9 (Exterior Derivative). Let U be an open set in S. We define the exterior derivative of 0-forms and 1-forms.

If $f \in \Omega^0(U)$, such that $f = f(x_1, x_2)$ then df is a 1-form on U, with $df = \frac{\partial f}{\partial x_1} dx_1 + \frac{\partial f}{\partial x_2} dx_2$. This definition is independent of parametrisation, a proof of which can be found in [11, p.69]. If $g \in \Omega^0(U)$ then d(fg) = fdg + (df)g.

If $\alpha \in \Omega^1(U)$ such that $f = \alpha_1 dx_1 + \alpha_2 dx_2$, where $\alpha_1, \alpha_2 \in \Omega^0(U)$, then we can define $d\alpha = \left(\frac{\partial \alpha_2}{\partial x_1} - \frac{\partial \alpha_1}{\partial x_2}\right) dx_1 \wedge dx_2$. The exterior derivative has some important properties. Firstly, if $f \in \Omega^0(U)$ then d(df) = 0, giving rise to the identity $d^2 = d \circ d = 0$. Secondly, if $f \in \Omega^0(U)$ and $\alpha \in \Omega^1(U)$ then we have $d(f\alpha) = f d\alpha + df \wedge \alpha$.

Definition 3.1.10 (Closed and Exact forms). A k-form α , where $k \in \{0, 1, 2\}$, is called closed if $d\alpha = 0$. A k-form α is called exact if it can be written as $d\beta = \alpha$ for some (k-1)-form β .

We see from the above definitions that every exact form is closed, since if we have that $\alpha = d\beta$ then $d\alpha = d(d\beta) = d^2\beta = 0$. The question to then ask is, is every closed form exact? That is the question that cohomology aims to answer, as we shall see below.

Definition 3.1.11 (Orientation and Volume forms). Given a smooth surface S, if we have a strictly positive 2-form, say $\rho \in \Omega^2(S)$ on S, then we can take ρ to be an orientation form on S. Two orientation forms ω, μ are equivalent if they differ by a positive smooth function f, i.e $\omega = f\mu$. If the integral (in the sense of definition 3.1.13) of an orientation form ω is 1 on the surface, then it is said to be an area form or more generally a volume form.

Orientation forms are helpful as they allow us to identify smooth functions on S with smooth 2-forms on S. By that, we mean that, ρ forms a basis element for $\Omega^2(S)$ and that we can then write every smooth 2-form, ω on S as $\omega = f\rho$, for some smooth function f on S. This concept is already familiar to us, as when we integrate over regions on \mathbb{R}^2 , we have that $dx \wedge dy$ is our cartesian volume form (often written as dxdy), and so we can integrate smooth functions f by constructing a new 2-form $fdx \wedge dy$. The main take away though is that these orientation forms are not unique and can be adapted to suit the problem at hand. However, we wish to restrict to studying area forms, rather than general orientation forms as we want to keep the values of integrals the same when changing coordinates. Since

they are supposed to represent the same unit area, by changing area form in a problem, we are essentially changing bases of $\Omega^2(S)$ and as such, must adapt the problem via the Jacobian. As we know from integration on \mathbb{R}^2 , we can go from cartesian coordinates to, for example, polar coordinates. This change might completely change the look of our integrand (the 2-form) but the end result should be the same, i.e. problems are independent of choice of area form. Though we don't prove this, we actually have a stronger result, that is proven in [11], namely that studying differential forms is independent of choice of coordinates on S. Before we move on to the integral of a differential form, we must first state the following lemma, that defines partitions of unity. Partitions of unity are a powerful tool that help us to define integration of 2-forms on orientable surfaces.

Lemma 3.1.12. [4, Lemma 7] Let K be a compact subset of S and let U_1, \ldots, U_n be open sets in S with $K \subset \bigcup_{i=1}^n U_i$. There, there exist real non-negative functions f_1, \ldots, f_n on S, each with compact support and with $supp(f_i) \subset U_i$ such that $\sum_{i=1}^n f_i = 1$.

The set of functions $\{f_i\}$ is called a partition of unity subordinate to the cover $\{U_i\}$.

Definition 3.1.13 (Integration). Integration on a surface can be defined for 1-forms and 2-forms. Given a smooth path living in a single chart on S, γ : $[0,1] \to S$, we can define the integral of a 1-form $\alpha = \alpha_1 dx_1 + \alpha_2 dx_2$ along γ , by

$$\int_{\gamma} \alpha = \int_{0}^{1} \left(\alpha_{1} \frac{d\gamma_{1}}{dt} + \alpha_{2} \frac{d\gamma_{2}}{dt} \right) dt$$

where $\gamma_1(t), \gamma_2(t)$ are the x_1, x_2 coordinate functions of the smooth path γ . If the path crosses over multiple charts, then our technique is to break up the interval [0,1] into subintervals corresponding to different charts and to sum the integrals over each subinterval to evaluate the integral. Furthermore, if we have some smooth function $\psi \colon [0,1] \to [0,1]$ such that $\psi(0) = \gamma(0)$ and $\psi(1) = \gamma(1)$, then the integral over the new path $\gamma \circ \psi$ remains the same as the integral over γ . This means that we can deform our paths to make them easier to integrate over, so long as this deformation keeps the endpoints intact.

When integrating 2-forms, in much the same way as for 1-forms, we have to be careful about whether the support of the 2-form we integrate goes over many charts or not. Suppose first that a 2-form, ρ , has its compact support contained within a single chart U. Then, using the coordinates of that chart, we can write, for some $f \in \Omega_c^0(U)$ that $\rho = f(x_1, x_2)dx_1 \wedge dx_2$ and we can integrate this 2-form as

$$\int_{U} \rho = \int_{U} f(x_1, x_2) dx_1 \wedge dx_2$$

Let us now suppose that ρ has compact support on a set K, that crosses finitely many sets coordinate charts on S (this can be assumed as the set of compact support is compact and is hence finitely covered). The above lemma gives us the existence of a partition of unity subordinate to this finite cover. Writing these functions as $\{\chi_i\}$, we have that $\sum_{i=0}^n \chi_i = 1$ and that for each chart U_i , χ_i is compactly supported within it. Hence, we define the integral ρ over K to be $\int_K \rho = \sum_{i=0}^n \int_{U_i} \chi_i \rho$ which must be finite as χ_i has compact support within each U_i and ρ also has compact support within K. Since, in both cases, we can extend ρ by 0, outside of its support on S, then we can write this as

$$\int_{S} \rho = \sum_{i=0}^{n} \int_{S} \chi_{i} \rho$$

and as such, $\rho = \sum_{i=0}^{n} \chi_i \rho$. Whilst we do not show this here, the linearity of the integral gives that this definition of integration holds true for any arbitrary partition of unity.

Whilst from the above definition it might not be clear how to integrate expressions of the type $\int_U f(x_1, x_2) dx_1 \wedge dx_2$, we can alleviate this by writing it as

$$\int_{U} f(x_1, x_2) dx_1 \wedge dx_2 = \int_{U} f(x_1, x_2) dx_1 dx_2$$

A minor, but important, point to make is that this equality arises when given an area form on either an open set, or even on the whole of S. In these cases, we can define a Lebesgue integral on our surface, which coincides with the usual notion of a surface integral in \mathbb{R}^2 , a fact of which can be found in [15, Chapter 4]. Defining the Lebesgue integral on a surface properly requires careful detail that we do not have the space to go into here and as such it shall be avoided. Henceforth, whenever mention of a measure is made, it will refer to this equality between a volume form and the Lebesgue measure, unless explicitly stated otherwise. The Lebesgue measure is frequently written in literature as $d\mu$ or $d\mu_x$ where the x represents the variable of integration. However, to keep in line with classical vector calculus, in this dissertation, whenever we make use of the Lebesgue measure, we will denote it as dS or dS_x , to represent the unit surface element. The reader is invited to read [13, Chapter 11] for a deeper introduction to measure theory.

Our final, important result is that of Stokes' theorem, a proof of which can be found in [15, Theorem 5-5].

Theorem 3.1.14 (Stokes' Theorem). If α is a compactly supported 1-form on a smooth, orientable surface S, then for any bounded set U, with boundary ∂U , in S, we have that

$$\int_{\partial U} \alpha = \int_{U} d\alpha$$

In particular, if U has no boundary, then we have that $\int_U d\alpha = 0$. Taking U to be the whole surface, we get that if our surface has no boundary, then $\int_S d\alpha = 0$. This means that any exact form on a smooth, orientable surface without boundary has integral 0 over the whole surface.

3.2 Complex structures

Below, let X denote a general Riemann surface. By the definition of the complex cotangent space, $T_p^*X^{\mathbb{C}}$ is simply the vector space of linear maps from T_pX to \mathbb{C} . This gives us that the derivatives of any complex valued function on X are elements of this space. So what do elements of this space look like locally?

Definition 3.2.1 (Complex structure). Given a vector space V over \mathbb{R} , we can define a complex structure J on V, by defining $J:V\to V$ such that $J^2=-1$ to be an \mathbb{R} -linear map on V.

Definition 3.2.2 (Complex linear and antilinear maps). Let V be a real vector space with a complex structure J. A linear map $A:V\to\mathbb{C}$ is said to be complex linear if for all $v\in V$ we have that A(Jv)=iA(v) and complex antilinear if for all $v\in V$ we have A(Jv)=-iA(v).

Using these definitions, we state the following lemma.

Lemma 3.2.3. Any \mathbb{R} -linear map from a vector space V to \mathbb{C} , can be written in a unique way as a sum of complex linear and antilinear maps.

Definition 3.2.4 (Dolbeault complexes). Given a complex coordinate on X, say z = x + iy where x, y are real coordinates, then, we have that dz = dx + idy and the conjugate is $d\overline{z} = dx - idy$. The form dz gives a basis of $\Omega^{1,0}(X)$, a real vector space of forms which may be written as $\alpha = f(z)dz$, where $f \in \Omega^0(X)$, a real, smooth function. These forms are called (1,0)-forms. Similarly, $d\overline{z}$ gives a basis of $\Omega^{0,1}(X)$, with elements $\beta = g(\overline{z})d\overline{z}$, where $g \in \Omega^0(X)$, a real, smooth function. These forms are called (0,1)-forms.

Given a function $f \in \Omega^0(X)$ we can take its exterior derivative. This yields the usual $df = \frac{\partial f}{\partial x} dx + \frac{\partial f}{\partial y} dy$. However, given the definitions of dz and $d\overline{z}$, we can rewrite df as

$$df = \frac{1}{2} \left(\frac{\partial f}{\partial x} - i \frac{\partial f}{\partial y} \right) dz + \frac{1}{2} \left(\frac{\partial f}{\partial x} + i \frac{\partial f}{\partial y} \right) d\overline{z}$$

where we have $dx = \frac{1}{2}(dz + d\overline{z})$ and $dy = \frac{1}{2i}(dz - d\overline{z})$. Thus we define two new differential operators on complex valued functions.

Definition 3.2.5 (Dolbeault differentials). We define for a smooth complex valued function f the differentials

$$\partial f = \frac{\partial f}{\partial z} dz, \qquad \overline{\partial} f = \frac{\partial f}{\partial \overline{z}} d\overline{z}$$

where the derivatives are defined as

$$\frac{\partial f}{\partial z} = \frac{1}{2} \left(\frac{\partial f}{\partial x} - i \frac{\partial f}{\partial y} \right), \qquad \frac{\partial f}{\partial \overline{z}} = \frac{1}{2} \left(\frac{\partial f}{\partial x} + i \frac{\partial f}{\partial y} \right)$$

From this, we see that the exterior derivative d has itself been decomposed and can be thus written as $d = \partial + \overline{\partial}$. If we have a holomorphic function f then $df = \partial f + \overline{\partial} f$. However, the fact that f is holomorphic gives that the second term vanishes, giving us a criterion for holomorphic functions; namely that $df = \partial f$.

The Dolbeault differentials are not necessarily restricted to just functions, but can also be applied to (1,0) and (0,1)-forms. Given a (1,0)-form $\alpha = f(z)dz$ we have that $\partial \alpha = \frac{\partial f}{\partial z}dz \wedge dz = 0$ and $\overline{\partial}\alpha = \frac{\partial f}{\partial \overline{z}}d\overline{z} \wedge dz$. Similarly, for the (0,1)-form $\beta = g(\overline{z})d\overline{z}$ we have that $\partial \beta = \frac{\partial g}{\partial \overline{z}}dz \wedge d\overline{z}$ and $\overline{\partial}\beta = \frac{\partial g}{\partial \overline{z}}d\overline{z} \wedge d\overline{z} = 0$, thus giving us the following diagram of derivatives

$$\begin{array}{ccc} \Omega^{0,1}(X) & \stackrel{\partial}{\longrightarrow} \Omega^2(X) \\ \hline \bar{\partial} & \bar{\partial} \\ \Omega^0(X) & \stackrel{\partial}{\longrightarrow} \Omega^{1,0}(X) \end{array}$$

Thus we conclude with the following definition.

Definition 3.2.6 (Holomorphic and Meromorphic 1-forms). Given $\alpha \in \Omega^{1,0}(X)$, it is said to be a holomorphic 1-form if $\overline{\partial}\alpha = 0$. It is said to be a meromorphic 1-form if it is a holomorphic 1-form on $X \setminus \Delta$, where Δ is a finite discrete subset of X, where around each $p_i \in \Delta$, α can be written as f(z)dz and f(z) is a meromorphic function with a pole at p_i .

3.3 Cohomology of surfaces

We conclude this section by briefly introducing the de Rham cohomology and the Dolbeault cohomology groups. We wish to use results from these cohomology theories to help us solve the uniformisation theorem.

Definition 3.3.1 (de Rham Cohomology). Let S be a smooth surface. We define the (additive) de Rham cohomology groups $H^i(S)$ for $i = \{0, 1, 2\}$ to be the following:

- $H^0(S) = \ker(d: \Omega^0(S) \to \Omega^1(S))$
- $H^1(S) = \ker(d:\Omega^1(S) \to \Omega^2(S)) / \operatorname{Im}(d:\Omega^0(S) \to \Omega^1(S))$
- $\bullet \ H^2(S) = \Omega^2(S)/\operatorname{Im}(d:\Omega^1(S) \to \Omega^2(S))$

Example. If our underlying surface is $S = \mathbb{R}^2$, the Euclidean plane, we see that the functions that live in the kernel of $d: \Omega^0(S) \to \Omega^1(S)$ are the constant functions, and only the constant functions. Hence, $\dim(H^0(S)) = 1$ and we say that $H^0(S) = \mathbb{R}$, as this group is generated by the constant functions. Furthermore, on the plane, we can show that all closed forms are exact² using a generalisation of the fundamental theorem of calculus, implying that $H^1(S) = H^2(S) = 0$.

We can also define the de Rham cohomology with compact support, $H_c^i(X)$, by swapping $\Omega^i(S)$ with $\Omega_c^i(S)$, in the definitions of the cohomology groups. This notion is useful on surfaces that are not compact, as $H_c^i(S) = H^i(S)$ if S is compact.

Example. Considering the example of the plane again, this time we wish to compute the cohomology groups of compact support. We must be a bit more careful as the only constant function with compact support in $\Omega_c^0(S)$ is the zero function. Thus, $H_c^0(S) = 0$. We also have that $H_c^1(S) = 0$. However, since the plane is orientable, we do have the notion of integration and integrable 2-forms in $\Omega_c^2(S)$ and as such we can extend the integration map $\int_S : \Omega_c^2(S) \to \mathbb{R}$ to be a map in cohomology, since any element $[\theta] \in H_c^2(S)$, can be written as $\theta + d\phi$, for some compactly supported 1-form ϕ and by Stokes' theorem, the integral of

²Recall this means that given a k-form, α , if $d\alpha = 0$ then there exists a k-1-form β such that $d\beta = \alpha$

 $d\phi$ vanishes so the map $\int_S : H_c^2(S) \to \mathbb{R}$ can be computed by $\int_S [\theta] = \int_S \theta$. This map defines an (group) isomorphism, we have that $H_c^2(S) = \mathbb{R}$. A relevant proof of the above results can be found in [11, Theorem 10.13].

Note that there is a symmetry in the cohomology groups of the plane when considering cohomology with and without compact support. This is an example of Poincaré duality, a powerful tool which helps relate the cohomology groups for surfaces, be they compact or not. An excellent resource on Poincaré duality can be found in [11, Chapter 13]. The theorem states that given a connected orientable surface S, then $H^p(S) \cong H_c^{2-p}(S)^*$, where the $H_c^{2-p}(S)^*$ is the dual space of $H_c^{2-p}(S)$, for $p \in \{0,1,2\}$. Crucially though for compact, connected Riemann surfaces, Poincaré duality gives that $H^1(S)$ must be even dimensional. A deeper exposé into this can be found in [11, p.130].

We state now some results now that will be useful for later. Let S be a connected, path-connected orientable surface and Σ be a connected, compact, orientable surface.

- 1. $\int_{\Sigma}: H^2(\Sigma) \to \mathbb{R}$ is an isomorphism. [11, Corollary 10.14]
- 2. $H^0(\Sigma) \cong H^2(\Sigma)$, by Poincaré duality.
- 3. If Σ has genus $g \geq 0$, then $H^1(\Sigma) \cong \mathbb{R}^{2g}$ by [4, p.69].
- 4. If S is simply connected, then $H^1(\Sigma) = 0$. In fact, by [6, Theorem 2A.1], we have that if our surface S is path connected, then $H^1(S)$ is isomorphic to the abelianisation of the fundamental group $\pi_1^{ab}(S) = \pi_1(S)/[\pi_1(S), \pi_1(S)]$ where $[\pi_1(S), \pi_1(S)]$ is the subgroup of elements of the form $aba^{-1}b^{-1}$ for $a, b \in \pi_1(S)$.

Therefore, we see that the Riemann sphere has de Rham cohomology

$$H^0(\widehat{\mathbb{C}}) = \mathbb{R}$$
 $H^1(\widehat{\mathbb{C}}) = 0$ $H^2(\widehat{\mathbb{C}}) = \mathbb{R}$

The complex torus of modulus ω , \mathbb{T}_{ω} , similarly has de Rham cohomology

$$H^0(\mathbb{T}_\omega) = \mathbb{R} \qquad H^1(\mathbb{T}_\omega) = \mathbb{R}^2 \qquad H^2(\mathbb{T}_\omega) = \mathbb{R}$$

Finally, we have that a compact, connected, orientable, genus g, surface Σ_g has de Rham cohomology

$$H^0(\Sigma_g) = \mathbb{R}$$
 $H^1(\Sigma_g) = \mathbb{R}^{2g}$ $H^2(\Sigma_g) = \mathbb{R}$

Before we move on, we also want to briefly justify that connected, compact, orientable surfaces also admit a Riemann surface structure, thus giving us that for every genus g, we have a Riemann surface. If S is a connected, smooth, oriented surface in \mathbb{R}^3 , then we can

define isothermal coordinates on S. The notion of isothermal coordinates is a tool from Riemannian geometry that gives charts on S, such that for $U \subset S$, we have $\psi : U \to \tilde{U} \subset \mathbb{R}^2$ such that ψ^{-1} takes orthogonal vectors from T_pS for some $p \in U$ to orthogonal vectors in \mathbb{R}^3 . Thus the composition of any such chart with the inverse of another chart preserves the orientation of any two vectors on S, thus giving S a conformal structure thereby making it a Riemann surface. A particularly nice proof of the fact that we can cover such surfaces with isothermal coordinates is given in [3]. Hence, we can talk about compact Riemann surfaces of genus g.

We conclude by defining a complex theory of cohomology, which works with the Dolbeault complexes and differentials more naturally, and are natural spaces to consider when studying Riemann surfaces.

Definition 3.3.2 (Dolbeault Cohomology). Let X be a Riemann surface. We define the (additive) Dolbeault cohomology groups as follows,

- $H^{0,0}(X) = \ker(\overline{\partial}: \Omega^0(X) \to \Omega^{0,1})$
- $H^{1,0}(X) = \ker(\overline{\partial}: \Omega^{1,0} \to \Omega^2(X))$
- $H^{0,1}(X) = \Omega^{0,1}/\ker(\overline{\partial}:\Omega^0(X) \to \Omega^{0,1}(X))$
- $H^{1,1}(X) = \Omega^2(X)/\ker(\overline{\partial}:\Omega^{1,0}(X) \to \Omega^2(X))$

These spaces seem rather abstract, and it is difficult to attach any meaning to them. We see that $H^{0,0}(X)$ is the space of holomorphic functions on X and similarly $H^{1,0}(X)$ is the space of holomorphic 1-forms. The other two spaces are more difficult to understand, and in the next section, we will hopefully start to understand them a bit better and how, in general, the Dolbeault cohomology groups relate to the de Rham cohomology groups of Riemann surfaces.

Chapter 4

The Uniformisation Theorem

4.1 The Laplacian and Hilbert spaces

4.1.1 The Δ operator and Harmonic functions

We begin our treatment of the uniformisation theorem, with a discussion of a particularly important and famous differential operator on Riemann surfaces; the Laplacian.

Definition 4.1.1 (The Laplacian). Let X be a Riemann surface. We define a linear map $\Delta: \Omega^0(X) \to \Omega^2(X)$, where $\Delta = 2i(\overline{\partial} \circ \partial)$. We call this linear map the Laplacian. Often, we simply write $\Delta = 2i\overline{\partial}\partial$. Where convenient, we also may use the fact that $\overline{\partial}\partial = -\partial\overline{\partial}$ which arises from the fact that $d^2 = 0$, and write $\Delta = -2i\partial\overline{\partial}$.

In local co-ordinates, Δ simply recovers the expected Laplacian of a function since for a given function $f \in \Omega^0(X)$ we get that

$$\Delta f = 2i\overline{\partial}\partial f = 2i\left(\frac{\partial}{\partial\overline{z}}\right)\left(\frac{\partial}{\partial z}\right)f(d\overline{z}\wedge dz)$$

$$= \frac{i}{2}\left(\frac{\partial}{\partial x} + i\frac{\partial}{\partial y}\right)\left(\frac{\partial}{\partial x} - i\frac{\partial}{\partial y}\right)f(2idx\wedge dy)$$

$$= -\left(\frac{\partial^2 f}{\partial x^2} + \frac{\partial^2 f}{\partial y^2}\right)dx\wedge dy$$

In this dissertation, we write ∇^2 for the classical Laplacian in local coordinates. Therefore we can say $\Delta f = -(\nabla^2 f) dx \wedge dy$.

A special class of function that is important to our theory is that of the Harmonic functions. They play an important role in studying the Laplacian operator generally and have particularly nice properties that we will rely on in the coming sections.

Definition 4.1.2 (Harmonic function). Let X be a Riemann surface. A function $f \in \Omega^0(X)$ is said to be harmonic if $\Delta f = 0$. If f is a smooth complex function, it is said to be harmonic if both its real and imaginary parts are real valued harmonic functions.

We claim that any holomorphic function is harmonic and show this using our language of differential forms. Recall that any complex function f(z) can be written as a sum of two real valued functions of (x, y), i.e. f(z) = u(x, y) + iv(x, y).

Lemma 4.1.3. Holomorphic functions are harmonic functions, i.e. they have harmonic real and imaginary parts.

Proof. Let X be a Riemann surface and f a holomorphic function defined on X. Let us consider the Laplacian applied to the sum $\frac{1}{2i}(f \pm \overline{f})$, which represents both the real and imaginary parts of f.

$$\frac{1}{2i}\Delta(f\pm\overline{f})=\overline{\partial}\partial(f\pm\overline{f})=-\partial(\overline{\partial}f)\pm\overline{\partial}(\overline{\overline{\partial}f)}=-\partial(0)\pm\overline{\partial}(0)=0\pm0=0$$

where $\overline{\partial} f = 0$ in both brackets because f is holomorphic.

Note that whilst holomorphic implies harmonic, the converse isn't always true. However, we can at least say the following.

Lemma 4.1.4. Let X be a Riemann surface, U an open set around a point $p \in X$ and let $\phi \in \Omega^0(U)$ be a harmonic function. Then there exists an open neighbourhood $V \subset U$ of the point p and a holomorphic function $f: V \to \mathbb{C}$ with $\phi = \Re(f)$.

Proof. Let $A \in \Omega^1(U)$ be the real 1-form $A = i\overline{\partial}\phi + \overline{(i\overline{\partial}\phi)} = -\frac{\partial\phi}{\partial y}dx + \frac{\partial\phi}{\partial x}dy$. Since ϕ is harmonic, $\overline{\partial}\partial\phi = 0$ and $d = \partial + \overline{\partial}$, we get that dA = 0. Hence, if V is an open, simply-connected set (such that $H^1(V) = 0$), then we can find a function $\psi \in \Omega^0(V)$ such that $d\psi = A$. By equating coefficients we get that $\partial\psi = -i\partial\psi$ and $\overline{\partial}\psi = i\overline{\partial}\phi$. So if we construct a function $f: V \to \mathbb{C}$ where $f = \phi + i\psi$, we see that $\overline{\partial}f = \overline{\partial}(\phi + i\psi) = \overline{\partial}\phi + i(i\overline{\partial}\phi) = 0$ and hence, f is a holomorphic function whose real part is ϕ .

Harmonic functions are, clearly, quite well behaved and have nice properties that are not generally exhibited by other functions. These properties, play an important role in our treatment of the uniformisation theorem. One such property is the following.

Lemma 4.1.5 (The Mean Value Theorem for Harmonic Functions). Let U be a region in \mathbb{C} around a point $a \in U$ and let $\phi \in \Omega^0(U)$ be harmonic on U. Let γ be a closed circle of

¹We note that we can relax the differentiability condition considerably to make the functions only twice differentiable. Thus harmonic functions need to be at least twice differentiable.

radius R > 0 contained within U that encircles the point a. Then, the mean value of ϕ over the set bounded by γ is equal to the value of ϕ at the point a.

We note that while the statement of this lemma talks about regions of \mathbb{C} , we could equally talk about regions on any Riemann surface as by definition of a Riemann surface, we have chart maps that map regions on a Riemann surface to regions in \mathbb{C} homeomorphically.

Proof. Since ϕ is a real-valued harmonic function, by lemma 4.1.4 we have a holomorphic function f such that the real part of f is ϕ . We now consider the statement of Cauchy's integral formula about γ :

$$f(a) = \frac{1}{2\pi i} \int_{\gamma} \frac{f(z)}{z - a} dz$$

Parametrising by γ gives us that $z = a + Re^{i\theta}$ for $\theta \in [0, 2\pi)$ on γ and that $dz = iRe^{i\theta}$. Plugging these into the above equation gives us the relation that

$$f(a) = \frac{1}{2\pi i} \int_{\gamma} \frac{f(z)}{z - a} dz = \frac{1}{2\pi} \int_{0}^{2\pi} f(a + Re^{i\theta}) d\theta$$

Since ϕ is the real part of f, we can split f into real and imaginary parts with the integrals splitting into real and imaginary parts, thus giving us that the average value of $\phi(z)$ along γ equals $\phi(a)$.

4.1.2 Hilbert Spaces and the space $\mathcal{H}(X)$

Hilbert spaces are in some ways the most familiar class of function spaces that one may study. They crop up everywhere in Mathematics and in Theoretical Physics and indeed they play a major part in the proof of the uniformisation theorem. Specifically, we will want to construct a particular Hilbert space that will contain a function that we will use to prove the uniformisation theorem. Let us provide some definitions first.

Definition 4.1.6 (Hilbert Space). An inner product space V is called a Hilbert Space if it is a complete metric space with a norm induced by the inner product.

Hilbert spaces over the real numbers are called real Hilbert spaces and those over the complex numbers are called complex Hilbert spaces. We see that Hilbert spaces are very similar to just normal inner product spaces and as such a lot of our examples are familiar vector spaces. Examples include:

1. \mathbb{R}^n with the usual Euclidean inner product is a Hilbert space as it is an inner product space in which all Cauchy sequences attain their limits. In fact, one definition of \mathbb{R} is

the completion of the space of sequences of rational numbers. In this space for example, we can define the number 1 as the limit of the sequence of partial sums $\left\{\sum_{i=1}^{n} \frac{9}{10^{i}}\right\}$, as n goes to infinity.

2. Given a space A, on which we can define integration with respect to a $measure^2 d\mu$, we define the Hilbert space $L^2(A)$ to be the space of functions $f:A\to \mathbb{F}$ such that $\int_A |f|^2 d\mu$ is finite. Such functions are called square integrable. We define an inner product of two functions $f,g\in L^2(A)$ to be $\langle f,g\rangle=\int_A \overline{g}fd\mu$, where the overline denotes conjugation. We call this inner product, the L^2 inner product on A and its associated norm $||f||^2=\int_A |f|^2 d\mu$ the L^2 norm.

We are now in a position to introduce a space of importance to our study of the Laplacian on a Riemann surface. **Note** that henceforth we consider only connected Riemann surfaces to alleviate issues related to having multiple connected components.

Lemma 4.1.7. Let X be a connected Riemann surface. We define a relation on $\Omega^0(X)$, whereby two functions are equivalent if they differ by a constant function. This relation, given the symbol \sim defines an equivalence relation.

Proof. Let $f, g, h \in \Omega^0(X)$ and let $a, b \in \mathbb{R}$. Since any function f is equal to f + 0 then we have that $f \sim f$. If $f \sim g$ implying f = g + a, then we have that g + (-a) = f and thus $g \sim f$, since $-a \in \mathbb{R}$. Finally, if $f \sim g$ and $g \sim h$ implying that f = g + a and g = h + b, then we have that $f \sim h$ since f = g + a = (h + b) + a = h + (b + a) and $b + a \in \mathbb{R}$. Hence, \sim is indeed an equivalence relation.

With this equivalence relation we can now define an important function space that is at the heart of the analysis presented in this chapter.

Definition 4.1.8 (The space $\mathcal{H}(X)$). If X is a compact, connected Riemann surface we define $\mathcal{H}(X) = \Omega^0(X)/\sim$, to be the space of cosets under the above defined equivalence relation \sim . If X is non compact, then we simply take $\mathcal{H}(X) = \Omega_c^0(X)$.

We note that $\mathcal{H}(X)$ is a vector space. This follows from the fact that $\Omega^0(X)$ and $\Omega^0_c(X)$ are themselves vector spaces, and in the compact case, the vector space structure descends through the quotient. Note more importantly that functions in $\mathcal{H}(X)$ are supported in X. This is because in the case of non compact Riemann surfaces, we are explicitly studying functions of compact support. In the case of compact Riemann surfaces, since every closed subset of a compact topological space is compact, then all continuous functions on compact

²See reference [13, Chapter 11]

Riemann surfaces have compact support implying all smooth functions do too. Therefore, $\mathcal{H}(X)$ only contains functions of compact support on X for both compact and non compact Riemann surfaces. We now can define a new inner product and associated norm on $\mathcal{H}(X)$ which will be of particular interest to our study of the Laplacian on Riemann surfaces.

Definition 4.1.9 (The Dirichlet Inner Product). Let X be a connected Riemann surface and let $f, g \in \mathcal{H}(X)$. We define the Dirichlet Inner Product $\langle \cdot, \cdot \rangle_D : \mathcal{H}(X) \times \mathcal{H}(X) \to \mathbb{R}$ as

$$\langle f, g \rangle_D = \langle df, dg \rangle = 2i \int_X \partial f \wedge \overline{\partial} g$$

where the inner product in the middle is the L^2 inner product. This inner product also defines a norm, called the Dirichlet norm, by $||f||_D^2 = \langle f, f \rangle_D$.

If the functions were to not be supported within this region, then the value of the inner product and norm might be $+\infty$, so by considering only functions of suitable support, we guarantee that we have an inner product and norm. With all of this theoretical set up, it would be nice if we could somehow get back our claimed main object of study for this section; the Laplacian. Although hidden, the Dirichlet norm (and the associated inner product) already contain the Laplacian in their definition. The following proposition shows how.

Proposition 4.1.10. Given a connected Riemann surface X, if at least one of the functions f, g is in $\mathcal{H}(X)$ (hence having compact support on X), then

$$\langle f, g \rangle_D = \int_X f \Delta g = \int_X g \Delta f$$

Proof. We first consider the following two identities:

- $\overline{\partial}(f\overline{\partial}g) = f(\overline{\partial}^2g) + \overline{\partial}f \wedge \overline{\partial}g = 0 + 0 = 0$
- $\bullet \ \partial (f\overline{\partial}g) = \partial f \wedge \overline{\partial}g + f(\partial \overline{\partial}g) = \partial f \wedge \overline{\partial}g \tfrac{1}{2i}f\Delta g$

Now let us expand the inner product. By the first identity above, we have that $f\overline{\partial}g$ is holomorphic, so that $d=\partial$. Hence, we can apply Stokes' theorem. Note that since X has no boundary, $\partial X=\emptyset$ and hence that integral vanishes.

$$\langle f, g \rangle_D = 2i \int_X \partial f \wedge \overline{\partial} g = 2i \int_X \partial (f \overline{\partial} g) + \int_X f \Delta g$$
$$= 2i \int_{\partial X} f \overline{\partial} g + \int_X f \Delta g$$
$$= \int_X f \Delta g$$

Since the Dirichlet inner product is a real inner product, it is symmetric and hence we have that $\langle f, g \rangle_D = \langle g, f \rangle_D$ and so $\int_X f \Delta g = \int_X g \Delta f$ as claimed.

4.2 The Dirichlet Energy Functional

A point that we will not be proving explicitly is that all functionals we consider hereon are "continuous", in the sense that, given a functional $F: \mathcal{H}(X) \to \mathbb{R}$, we have that for any $\epsilon > 0$ we have a $\delta > 0$ such that for all $x, y \in \mathcal{H}(X)$, $||x - y|| < \delta \Rightarrow |F(x) - F(y)| < \epsilon$. This continuity allows us to search for limits, by looking at how the functional acts on Cauchy sequences and their corresponding Cauchy sequences in \mathbb{R} .

Definition 4.2.1 (The $\hat{\rho}$ functional). Let X be a connected Riemann surface and $\rho \in \Omega_c^2(X)$ be a 2-form of compact support in X such that $\int_X \rho = 0$. Then, we define the functional $\hat{\rho} \colon \mathcal{H}(X) \to \mathbb{R}$, for some f in $\mathcal{H}(X)$ to be:

$$\hat{\rho}(f) = \frac{1}{2} \int_X f \rho$$

Since $\int_X \rho = 0$ then we have that this functional is linear on functions in $\mathcal{H}(X)$. Note that, if X is a compact Riemann surface, then all functions and 2-forms on X automatically have compact support on X, so the compact support condition really only matters in the cases of non compact Riemann surfaces. We claim that $\hat{\rho}$ is a linear functional.

Lemma 4.2.2 ($\hat{\rho}$ is a linear functional). Let X be a connected Riemann surface. Then $\hat{\rho}$ is a linear functional, i.e. $\hat{\rho}(\lambda f + \mu g) = \lambda \hat{\rho}(f) + \mu \hat{\rho}(g)$.

Recall that elements of $\mathcal{H}(X)$ are equivalence classes of smooth functions with compact support on X which differ by a constant. We denote cosets as \tilde{f} , where f is the coset representative.

Proof. Let X be a compact, connected Riemann surfaces and $\tilde{f}, \tilde{g} \in \mathcal{H}(X)$ and $\lambda, \mu \in \mathbb{R}$. Then for some real constant, $a, b \in \mathbb{R}$ we have that $\tilde{f}(z) = f(z) + a$ and $\tilde{g}(z) = g(z) + b$. Recall also that $\int_X \rho = 0$. Hence, we get that:

$$\begin{split} \hat{\rho}(\lambda \tilde{f} + \mu \tilde{g}) &= \hat{\rho}(\lambda (f + a) + \mu (g + b)) \\ &= \frac{1}{2} \int_X (\lambda (f + a) + \mu (g + b)) \rho \\ &= \frac{\lambda}{2} \int_X f \rho + \frac{\lambda a}{2} \int_X \rho + \frac{\mu}{2} \int_X g \rho + \frac{\mu b}{2} \int_X \rho \\ &= \lambda \hat{\rho}(f) + \mu \hat{\rho}(g) = \lambda \hat{\rho}(\tilde{f}) + \mu \hat{\rho}(\tilde{g}) \end{split}$$

with the last equality holding since clearly $f \in \tilde{f}$ and $g \in \tilde{g}$. If X is non compact, then the above holds by setting a = b = 0 since constant functions do not have compact support in X.

We now define a functional in terms of $\hat{\rho}$ and the Dirichlet norm on $\mathcal{H}(X)$ of crucial importance. We will be using this functional to help us set up a problem whose solution allows us to prove the uniformisation theorem.

Definition 4.2.3 (Dirichlet Energy Functional). Let X be a connected Riemann surface, and suppose $f \in \mathcal{H}(X)$. Suppose further that we have a 2-form $\rho \in \Omega^2(X)$ (or $\Omega_c^2(X)$ if X is non compact) such that $\int_X \rho = 0$. Then we can define a functional $\mathcal{L} : \mathcal{H}(X) \to \mathbb{R}$ such that $\mathcal{L}(f) = ||f||_D^2 - 2\hat{\rho}(f)$. This \mathcal{L} is called the Dirichlet energy functional, or sometimes Dirichlet energy for short.

So how can we use \mathcal{L} to set up this problem? We motivate this problem physically, using the theory of Electrostatics; a subject which mathematicians in Riemann's time were already very familiar with. The physical principle of least action gives that physical systems always act to minimise their energy. Thus we too should look for a physical description of the electrostatic energy on an arbitrary Riemann surface and minimise it. The description of the functions which minimise this energy that we thus get will allow us to prove the uniformisation theorem. Given a Riemann surface X, we imagine it is made of a electrically conductive material. We impose on X a static, smooth distribution of charge, such that the net charge on X due to ρ is zero. We call such a ρ an externally imposed charge distribution. The fact that the net charge on X is zero gives us that $\int_X \rho = 0$.

To get a concrete idea of what sort of externally imposed charge distributions have this property, let us consider $X = \widehat{\mathbb{C}}$. We can fix an area form on X, say $\sin(\theta)d\theta \wedge d\phi$, given by the standard spherical coordinates. An example distribution may therefore be $\cos(\phi)\sin(\theta)d\theta \wedge d\phi$. We have that $\int_{0}^{2\pi} \int_{0}^{\pi} \cos(\phi)\sin(\theta)d\theta d\phi = 0$ so this is a valid charge distribution.

A non-trivial example distribution on the Riemann sphere, say ρ , would be to take two small disjoint circular regions, of diameter $\epsilon > 0$, on the Riemann sphere such that our smooth charge distribution is compactly supported inside these two regions say A and B. After fixing an area form on $\widehat{\mathbb{C}}$, say $dz \wedge d\overline{z}$, let us say that the function representing the distribution on A is a bump function β , with the function for B being $-\beta$. We see clearly that the integral of our distribution splits into two. Hence, we get that since ρ is zero outside $A \cup B$ that $\int_{\widehat{\mathbb{C}}} \rho = \int_A \beta dz \wedge d\overline{z} + \int_B (-\beta) dz \wedge d\overline{z} = 0$, by the definition of ρ . So acceptable charge distributions must be particularly well behaved. In fact, this example is similar to defining a pair of opposing point charges on our Riemann surface. Such a charge distribution is similar to a type of physical system called a dipole.

The potential energy felt by a test charge³ at a point p, denoted as $\phi(p)$, may be non zero everywhere on X. This ϕ is an example of an electric potential on X. So what is the potential energy of this system? We find the energy is split into two parts. Firstly, there is the energy stored in ϕ itself, which is given by $\frac{1}{2}\|\phi\|_D^2$, half the Dirichlet norm of ϕ . The second term is the energy contribution due to ρ interacting with ϕ which is described by $-\hat{\rho}(\phi)$. An in depth discussion of the physics needed for this can be found in [5, Chapter 1]. Since we are considering the total energy of the system, we can sum up these two terms up to give that $\frac{1}{2}\|\phi\|_D^2 - \hat{\rho}(\phi)$ describes the total energy of our problem. We note though that $\frac{1}{2}\|\phi\|_D^2 - \hat{\rho}(\phi) = \frac{1}{2}\mathcal{L}(\phi)$. Thus, we seek to minimise \mathcal{L} over all functions in $\mathcal{H}(X)$, to find the minimum energy of our system. We call such an object, which encapsulates the energy of the system, the Lagrangian of the system. Hence, we can state the problem we want to solve:

Problem. Given a charge distribution ρ , of integral zero, over our Riemann surface X, which functions $\phi \in \mathcal{H}(X)$ exist that minimise $\mathcal{L}(\phi)$?

The answer, as we will find, is precisely those fields which arise as the potential due to the same charge distribution ρ . We must be careful however, as we might be looking for functions which minimise \mathcal{L} , but such functions may not exist. As we will show, this functional is bounded from below, giving us a necessary condition for a minimum to exist. However, we note that this condition is not sufficient, and we will deal with that in the latter part of this section. Let us, for a moment, suppose such a potential function exists and that \mathcal{L} is bounded from below. Which functions act to minimise \mathcal{L} ? We begin by computing $\delta \mathcal{L} = 0$, the first variation of \mathcal{L} . A formal introduction to the notion of the first variation of a functional can be found in [8, p.4], though it is sufficient to imagine it acts much the same as a derivative for functions in that its zeros give the stationary points of the functional. To compute it, we let $\epsilon > 0$ and let g be function on X. The conditions on g will be more precisely defined later, so for now, we assume such a g exists and it is, at least, integrable. Computing the first variation of \mathcal{L} gives:

$$\delta \mathcal{L} = \frac{d}{d\epsilon} \Big|_{\epsilon=0} \mathcal{L}(f + \epsilon g) = \frac{d}{d\epsilon} \Big|_{\epsilon=0} (\|f + \epsilon g\|_D^2 - 2\hat{\rho}(f + \epsilon g)) = 0$$

$$\frac{d}{d\epsilon} \Big|_{\epsilon=0} (\langle f, f \rangle_D + 2\epsilon \langle f, g \rangle_D + \epsilon^2 \langle g, g \rangle_D) - 2\hat{\rho}(g) = 0$$

$$2 \langle f, g \rangle_D - 2\hat{\rho}(g) = 0$$

$$\int_X g(\Delta f - \rho) = 0$$

³A hypothetical object which interacts with electric charge, but doesn't affect the net charge of the system.

The appearance of the Laplacian comes from proposition 4.1.10. Since we have that f and ρ have compact support on X, then for this integral to be zero, we need the following statement to be true: There exists a unique solution f, up to addition by a constant, of the equation $\Delta f = \rho$ if and only if $\int_X \rho = 0$

The equation on the left, known as Poisson's equation, is precisely the problem we are seeking. We call finding such an f, "solving the Poisson equation". It turns out, to prove the uniformisation theorem, we really only need to solve the Poisson equation on two types of Riemann surface; connected, compact Riemann surfaces and connected, simply-connected, non compact Riemann surfaces. Before we study Poisson's equation too deeply however, we must be careful and show that \mathcal{L} is really bounded from below thus motivating our study of Poisson's equation. We do this by showing that $\hat{\rho}$ is a bounded functional on $\mathcal{H}(X)$, and thus, by the positivity of $\|\cdot\|_D^2$, we get a lower bound for \mathcal{L} .

4.2.1 Showing the Dirichlet energy \mathcal{L} is bounded from below

To prove $\hat{\rho}$ is bounded, we need the following theorem and its subsequent corollary. We proceed by working in a bounded, convex, open set in the complex plane. We can do so as we can map any bounded, convex, open set from our Riemann surfaces homeomorphically to similar sets in the complex plane using chart maps. Similarly, due to conformal equivalence, we can assume such sets are circular discs in the complex plane. We omit the proofs of the following theorem and its corollary due to them being long arguments that are well described in the quoted references.

Theorem 4.2.4. [4, Theorem 11] Let Ω be a bounded, convex, open set in \mathbb{R}^2 and ψ be a smooth function on an open set containing the closure $\overline{\Omega}$ with $\overline{\psi}$ denoting the average

$$\overline{\psi} = \int_{\Omega} \psi dS$$

where A is the area of Ω . Then, for $x \in \Omega$, we have

$$|\psi(x) - \overline{\psi}| \le \frac{d^2}{2A} \int_{\Omega} \frac{1}{|x - y|} |\nabla \psi(y)| dS_y$$

Corollary 4.2.5. [4, Corollary 6] Under the hypothesis from Theorem 4.2.4, we have that

$$\int_{\Omega} |\psi(x) - \overline{\psi}|^2 dS_x \le \left(\frac{d^3 \pi}{A}\right)^2 \int_{\Omega} |\nabla \psi|^2 dS$$

where ∇ denotes the usual gradient operator in 2 dimensions.

Using these two estimates we can prove the following theorem.

Proposition 4.2.6. Let X be a connected Riemann surface. Then the functional $\hat{\rho}$: $\mathcal{H}(X) \to \mathbb{R}$ is bounded, i.e. there exists a constant C such that $|\hat{\rho}(\tilde{f})| \leq C ||\tilde{f}||_D$, for all $\tilde{f} \in \mathcal{H}(X)$.

The proof follows the proof found in [4, p.125], though we generalise the result to hold for any Riemann surface, rather than for only connected Riemann surfaces.

Proof. We split this proof into two parts. First we show that $\hat{\rho}$ is bounded in the case when the compact set $\operatorname{supp}(\rho) = K$ is contained within a single coordinate chart. Then we show that we can extend the boundedness of $\hat{\rho}$ over the whole of K if it cannot be contained in a single coordinate chart. Note that if X is compact, then K = X. We consider a bounded, convex set U in \mathbb{C} . Let $\tilde{f} \in \mathcal{H}(X)$. We write $\hat{\rho}(\tilde{f}) = \int_U (f+a)\rho$ for some $a \in \mathbb{R}$ and set the constant a to be equal to $-\overline{f}$, the average value of f on U, giving that $\hat{\rho}(\tilde{f}) = \int_U (f-\overline{f})\rho$. By fixing an area form in this coordinate chart, say dS, we can write $\rho = gdS$, where $g \in \Omega^0(X)$. This therefore allows us to write $\hat{\rho}(\tilde{f}) = \int_U (f-\overline{f})gdS$ which by the Cauchy-Schwarz inequality gives

$$|\hat{\rho}(\tilde{f})| = \left| \int_{U} (f - \overline{f}) g dS \right| \le ||g|| ||f - \tilde{f}||$$

Hence, by Corollary 4.2.5, we get that $|\hat{\rho}(\tilde{f})| \leq C \|\nabla f\| = C \|df\|$ where $C = d^3\pi A \|g\|$, and we have used the fact that the norm of the gradient and exterior derivative of a function are equal. Finally, if we now compose with a chart map to map back to X from U, and letting $\|\cdot\|_U$ indicate the L^2 norm on U, we have that $\|df\|_U \leq \|df\|_X = \|\tilde{f}\|_{D,X}$. Hence we have that if $\hat{\rho}$ is supported in a single coordinate chart on X, then $\hat{\rho}$ is a bounded operator.

So to extend this to the case where the compact set K is contained in multiple coordinate charts on X, we first fix a finite cover of K by coordinate charts of the type considered in the previous case. This is possible since K is compact. We call this family of coordinate charts $\{U_{\alpha}\}$. Since we showed that for each coordinate chart, a 2-form of compact support in that chart of integral zero is bounded, we wish to somehow stitch together these forms to form a 2-form of integral zero over the whole of K. To do this stitching, we introduce a partition of unity χ_{α} subordinate to our choice of finite cover $\{U_{\alpha}\}$. Since integration of 2-forms on a compact, connected set $V \subseteq X$ defines an isomorphism between $H^2(V)$ and \mathbb{R} , we can write the 2-form from our functional $\hat{\rho}$ as $\rho = d\theta$ for some $\theta \in \Omega^1(K)$. Now on each U_{α} we derive a 2-form of compact support from θ using our partition of unity by setting $\rho_{\alpha} = d(\chi_{\alpha}\theta)$. Each 2-form ρ_{α} is of compact support in U_{α} which it gets from χ_{α} . Thus we have for each α

$$\int_{K} \rho_{\alpha} = \int_{K} d(\chi_{\alpha} \theta) = \int_{\partial K} \chi_{\alpha} \theta = 0$$

which vanish by Stokes' theorem if K = X since X has no boundary or if X is not compact since ρ_{α} must vanish on the boundary⁴ of K. So $\int_{K} \rho_{\alpha} = 0$. Finally, since $\sum_{\alpha} \chi_{\alpha} = 1$, we have that $\rho = d\theta = d\left(\sum_{\alpha} \chi_{\alpha}\theta\right) = \sum_{\alpha} d(\chi_{\alpha}\theta) = \sum_{\alpha} \rho_{\alpha}$. As such, we have constructed a 2-form ρ such that $\int_{K} \rho = 0$. Since we showed in the first part that each $\hat{\rho}_{\alpha}$ is bounded in its coordinate chart, and since ρ is equal to a finite sum of ρ_{α} we conclude that the functional $\hat{\rho} = \sum_{\alpha} \hat{\rho}_{\alpha}$ is also bounded since it is a finite sum of bounded linear maps. Furthermore, the boundedness holds even outside K since we can extend $\hat{\rho}$ by zero over the whole of $X \setminus K$. \square

So we have that the functional $\hat{\rho}$ is bounded both when X is compact and non compact. Let us now directly show that \mathcal{L} is hence bounded from below. We have that $\mathcal{L}(f) = \|f\|_D^2 - 2\hat{\rho}(f)$ for some function $f \in \mathcal{H}(X)$. We showed above that $|\hat{\rho}(f)| \leq C\|f\|_D$ for the constant C defined above. Since $\hat{\rho}(f) \leq |\hat{\rho}(f)|$ we have that $\mathcal{L}(f) \geq \|f\|_D^2 - 2C\|f\|_D = (\|f\|_D - C)^2 - C^2 \geq -C^2$. Hence, we have that the Dirichlet energy is indeed bounded from below, irrespective of the compactness of X.

4.3 Poisson's equation on compact Riemann surfaces

Since \mathcal{L} is bounded on both compact and non compact Riemann surfaces, we need to study the two cases separately. We begin by stating our problem of solving Poisson's equation for compact Riemann surfaces.

Theorem 4.3.1 (Poisson's equation on compact Riemann surfaces). Let X be a connected, compact Riemann surface. Let $\rho \in \Omega^2(X)$ be a 2-form on X. Then, there exists a unique, smooth solution $f \in \Omega^0(X)$ to the equation $\Delta f = \rho$, up to the addition of a constant, if and only if $\int_X \rho = 0$. This equation is called the Poisson equation.

We can immediately prove two things:

- 1. If f is a smooth solution to $\Delta f = \rho$, then $\int_X \rho = 0$.
- 2. f is unique up to a constant.

Unfortunately, the *only if* direction is much tricker to prove and the majority of this section is dedicated to proving it. So we begin the proof of Theorem 4.3.1 by proving the *if* direction and uniqueness of the solution.

Proof. We wish to first prove that if f is a smooth solution to $\Delta f = \rho$ then $\int_X \rho = 0$. Recall that since $d = \partial + \overline{\partial}$ then $d(\partial f) = \partial^2 f + \overline{\partial} \partial f = \overline{\partial} \partial f$ because $\partial^2 = 0$ and that $\Delta f = 2i \, \overline{\partial} \partial f$.

⁴Technically, ρ_{α} isn't defined on ∂K but we can extend ρ out of K over the whole of X by zero.

Combining these two gives us that $\Delta f = 2i d(\partial f)$. So let us integrate ρ .

$$\int_X \rho = \int_X \Delta f = 2i \int_X \overline{\partial} \partial f = 2i \int_X d(\partial f) = 2i \int_{\partial X} \partial f = 0$$

where the last equality follows from Stokes' theorem and the fact that X has no boundary.

Now to show that such a solution is unique up to a constant. We assume f and g are both solutions to the Poisson equation, i.e. $\Delta f = \rho$ and $\Delta g = \rho$. Then we have that $\Delta(f-g) = 0$, or in other words, that f-g is a harmonic function. Thinking of the Dirichlet norm for a moment, if we consider $||f-g||_D$ we get $\int_X (f-g)\Delta(f-g) = 0$. Hence, we have $||f-g||_D = 0$. Unfortunately, the Dirichlet norm is not a norm on $\Omega^0(X)$ but on $\mathcal{H}(X)$, but elements of $\Omega^0(X)$ can easily be mapped to elements of $\mathcal{H}(X)$. So let us consider $||\tilde{f}-\tilde{g}||_D = ||df-dg||$. It follows that $||\tilde{f}-\tilde{g}||_D = 0$ since $\Delta(c) = 0$, for any constant c, so by the definition of d and the Dirichlet norm, we get that df-dg=d(f-g)=0 implying that f-g=c where c is a constant. So any solution of the Poisson equation is unique up to the addition of a constant.

To complete our proof of Theorem 4.3.1, now must now show that given a 2-form $\rho \in \Omega^2(X)$ such that $\int_X \rho = 0$, we can find a smooth function f on X such that $\Delta f = \rho$. To do so, we show that whilst minimising functions of \mathcal{L} are indeed obtained as solutions to Poisson's equation, they may live "outside" our space of all possible smooth functions on X. As we will show however, this will not be the case.

4.3.1 The completion of $\mathcal{H}(X)$ and Weyl's lemma

Though we showed earlier that \mathcal{L} is bounded from below, we must be careful as its minimum may not live in $\mathcal{H}(X)$. We begin this section by constructing the completion of the space $\mathcal{H}(X)$, an extension of the space $\mathcal{H}(X)$ that contains the limits of all *Cauchy sequences* of functions on $\mathcal{H}(X)$. What does this space look like? We construct this space abstractly as follows.

Definition 4.3.2 (The space $\overline{\mathcal{H}}(X)$). Let X be a connected Riemann surface. The space $\overline{\mathcal{H}}(X)$, the abstract completion of $\mathcal{H}(X)$, can be constructed from the inner product space $\mathcal{H}(X)$ with the Dirichlet inner product, as the space of equivalence classes of Cauchy sequences of elements in $\mathcal{H}(X)$, under the equivalence relation \sim , where two Cauchy sequences $\{\psi_i\}$ and $\{\phi_i\}$ are equivalent if $\lim_{i\to\infty} \|\psi_i - \phi_i\|_D \to 0$. Therefore, a point in $\overline{\mathcal{H}}(X)$ is an equivalence class of Cauchy sequences of functions in $\mathcal{H}(X)$ that have the same limit.

This space is clearly very complicated, but its usefulness to us comes from the fact that we can guarantee that Cauchy sequences of functions will attain their limit in $\overline{\mathcal{H}}(X)$. This space is also a vector space because we can add points of this vector space, by summing up the sequences term by term. The space $\overline{\mathcal{H}}(X)$ also has an inner product, which we define as, for any two points $\psi, \phi \in \overline{\mathcal{H}}(X)$, where $\psi = \{\psi_i\}$ and $\phi = \{\phi_i\}$, the product $(\psi, \phi)_D = \lim_{i \to \infty} \langle \psi_i, \phi_i \rangle_D$. Note that we use the notation $(\cdot, \cdot)_D$ for the Dirichlet inner product on $\overline{\mathcal{H}}(X)$ and $\langle \cdot, \cdot \rangle_D$ for the Dirichlet inner product on $\mathcal{H}(X)$. Note that the sequence of values in \mathbb{R} of the inner product is itself a Cauchy sequence in \mathbb{R} . Since we have a complete inner product space, with a norm that we can derive from our new inner product, $(\cdot, \cdot)_D$, then $\overline{\mathcal{H}}(X)$ is, in fact, a Hilbert space. We now wish to extend our previously defined functionals to $\overline{\mathcal{H}}(X)$. We identify these extended functionals with an underline, except for the extended Dirichlet inner product, which we denote as $(\cdot, \cdot)_D$. Note that this extended Dirichlet inner product also extends our norm to $\overline{\mathcal{H}}(X)$ by defining it as $\|\psi\|_D^2 = (\psi, \psi)_D$ for some $\psi \in \overline{\mathcal{H}}(X)$. So we now need to extend $\hat{\rho}$. Its extension to $\overline{\mathcal{H}}(X)$ can be accomplished similarly to the Dirichlet inner product by

$$\underline{\hat{\rho}}(\psi) = \lim_{i \to \infty} \hat{\rho}(\psi_i) = \lim_{i \to \infty} \int_X \psi_i \rho$$

for some $\psi = \{\psi_i\} \in \overline{\mathcal{H}}(X)$. We note that this functional is also Cauchy in \mathbb{R} since if we take a Cauchy sequence from $\mathcal{H}(X)$, say $\{\phi_i\}$, then the sequence $\hat{\rho}(\phi_i)$ is Cauchy in \mathbb{R} . Furthermore, $\hat{\rho}$ is bounded since we have that for some $\psi \in \overline{\mathcal{H}}(X)$:

$$|\underline{\hat{\rho}}(\psi)| = \lim_{i \to \infty} |\hat{\rho}(\psi_i)| \le \lim_{i \to \infty} C||\psi_i|| = C||\psi||$$

As such, we can define our extended version of the Dirichlet energy to be $\underline{\mathcal{L}}(f) = -2\underline{\hat{\rho}}(f) + \|\underline{f}\|_D^2$, for any $f \in \overline{\mathcal{H}}(X)$. It follows that $\underline{\mathcal{L}}$ is also bounded from below. The proof is identical to that of \mathcal{L} . Therefore, since $\underline{\mathcal{L}}$ is a functional on $\overline{\mathcal{H}}(X)$ we can now guarantee that if there is an abstract element Φ for which $\underline{\mathcal{L}}(\Phi) = -C^2$ then $\Phi \in \overline{\mathcal{H}}(X)$. To obtain such an element, we need to construct a Cauchy sequence of functions in $\mathcal{H}(X)$ for which \mathcal{L} attains this minimum, thereby guaranteeing the existence of an element of $\overline{\mathcal{H}}(X)$ for which $\underline{\mathcal{L}}$ attains its minimum.

Lemma 4.3.3. Let X be a compact, connected Riemann surface. We let \mathcal{L} be the Dirichlet energy functional, and let it be bounded from below by a number $M \in \mathbb{R}$, i.e. $\mathcal{L} \geq M$ (which holds true by the previous section with $M = -C^2$). Let $\{\phi_i\}$ be a sequence of functions such that $\lim_{i \to \infty} \mathcal{L}(\phi_i) = M$. Then $\{\phi_i\}$ is a Cauchy sequence in $\mathcal{H}(X)$ which defines our abstract element $\Phi \in \overline{\mathcal{H}}(X)$ such that $\underline{\mathcal{L}}(\Phi) = M$.

Note that the work in this lemma and proof follow the proposition in [10, p.30].

Proof. Let $\epsilon > 0$. Then there exists some $N \in \mathbb{N}$ such that for $i, j \geq N$ we have that $\mathcal{L}(\phi_i) \leq M + \frac{1}{4}\epsilon$ and $\mathcal{L}(\phi_j) \leq M + \frac{1}{4}\epsilon$. We define a function $I : [0,1] \to \mathbb{R}$, as a function of t, with $I(0) = \mathcal{L}(\phi_i)$ and $I(1) = \mathcal{L}(\phi_i)$, which interpolates between the two giving

$$I(t) = \mathcal{L}(t\phi_i + (1-t)\phi_j) = ||t\phi_i + (1-t)\phi_j||_D^2 - 2\hat{\rho}(t\phi_i + (1-t)\phi_j)$$

$$= ||\phi_j + (\phi_i - \phi_j)t||_D^2 - 2\hat{\rho}(\phi_j + (\phi_i - \phi_j)t)$$

$$= ||\phi_i - \phi_j||_D^2 t^2 + 2\langle\phi_i, \phi_j\rangle_D t + ||\phi_j||_D^2 - 2\hat{\rho}(\phi_j) - 2\hat{\rho}(\phi_i - \phi_j)t$$

So we have that I(t) is a quadratic polynomial. By expanding $2(I(0) - 2I(\frac{1}{2}) + I(1))$ we get the coefficient of t^2 . Considering that $I(0) = \mathcal{L}(\phi_j) \leq M + \frac{1}{4}\epsilon$ and $I(1) = \mathcal{L}(\phi_i) \leq M + \frac{1}{4}\epsilon$ and for all t we have $I(t) \geq M$, since $\mathcal{L} \geq M$ for all functions, then we can apply the following approximation

$$\|\phi_i - \phi_j\|_D^2 = 2(I(0) - 2I(1/2) + I(1)) \le 2(M + \frac{1}{4}\epsilon - 2M + M + \frac{1}{4}\epsilon) = \epsilon$$

Therefore $\{\phi_i\}$ is Cauchy so there exists an element $\Phi = \lim_{i \to \infty} \phi_i$ such that $\underline{\mathcal{L}}(\Phi) = M$. \square

So we have proven the existence of the minimising element Φ of $\underline{\mathcal{L}}$ in $\overline{\mathcal{H}}(X)$. But we want to understand it a bit better. As things stand, it is simply an equivalence class of Cauchy sequences with limit Φ ; not a very usable object. We want to instead show it is an element of $\mathcal{H}(X)$, that is, show it can be identified with an equivalence class of smooth functions (up to the addition of a constant). This implies Weyl's lemma; a theorem that states that every solution of the Poisson equation is a smooth function, which in our case, allows us to conclude our proof of the reverse direction of Theorem 4.3.1.

Before we embark on this final part of our proof, let us first show that indeed, we are on the right track and that this minimising element does indeed give a solution to the Poisson equation. The method mirrors what we did before for \mathcal{L} with the arbitrary function g, however now we use $\underline{\mathcal{L}}$ and Φ instead of an arbitrary function. Recall, earlier on we assumed that there existed a function g which minimised \mathcal{L} . As things stand, we still don't know that. What we do know however is that we can minimise $\underline{\mathcal{L}}$ by applying it to Φ . Let us compute the first variation of $\underline{\mathcal{L}}$. For some $f \in \mathcal{H}(X)$ we have:

$$\delta \underline{\mathcal{L}} = \frac{d}{d\epsilon} \Big|_{\epsilon=0} \underline{\mathcal{L}}(\Phi + \epsilon f) = 0$$

Using an argument, largely the same as before, we reach the following conclusion

$$(\Phi, f)_D = \underline{\hat{\rho}}(f) \iff \lim_{i \to \infty} \langle \phi_i, f \rangle_D = \underline{\hat{\rho}}(f) \iff \lim_{i \to \infty} \int_{V} f \Delta \phi_i = \underline{\hat{\rho}}(f)$$

Since $\hat{\rho}$ is a bounded linear functional, we have that the sequence of real numbers, defined by the integrals as the limit is taken to infinity, is a Cauchy sequence in \mathbb{R} and so we can precisely write that:

$$\int_{X} f \Delta \Phi = \underline{\hat{\rho}}(f) \Longleftrightarrow \int_{X} f \Delta \Phi - f \rho = 0 \Longleftrightarrow \int_{X} f (\Delta \Phi - \rho) = 0$$

In other words, Φ is indeed a solution to the Poisson equation. Such a Φ is often called a Weak Solution to the Poisson equation since it is not necessarily a smooth, or even continuous function. Now all that is left to do is show Φ can be identified with a smooth function. Before we state Weyl's lemma, we give the definition of a convolution of two functions.

Definition 4.3.4 (Convolution of functions). Let f, g be L^2 functions defined on \mathbb{R}^n . Then their convolution is defined as $(f * g)(x) = \int_{\mathbb{R}^n} f(y)g(x-y)dy$ where dy is volume form on \mathbb{R}^n .

For a detailed treatment of convolutions of functions and the theory of distributions and a more in depth discussion of the theory of Partial Differential Equations, one can refer to [12, Chapter 6]. We can now state the infamous Weyl Lemma.

Theorem 4.3.5 (Weyl's Lemma). Let X be a compact, connected Riemann surface. If $\rho \in \Omega^2(X)$ such that $\int_X \rho = 0$, then a weak solution to Poisson's equation $\phi \in \overline{\mathcal{H}}(X)$ in fact is an element of $\mathcal{H}(X)$, i.e. ϕ is a smooth function.

We also want to use the following lemma, that we call the "Local Weyl Lemma".

Lemma 4.3.6 (Local Weyl Lemma). Let U be a bounded region in \mathbb{C} and let $\rho \in \Omega^2(U)$ be a 2-form on U. Suppose $\phi \in L^2(U)$ such that, for any smooth function h of compact support in U, we have that

$$\int_{U} \Delta h \phi = \int_{U} \chi \rho$$

Then ϕ is smooth and satisfies the equation $\Delta \phi = \rho$.

We see that, if this lemma were to be true in a single coordinate chart and if we could somehow stitch up a function all over X to have this property hold, then it implies Weyl's lemma. To avoid simply copying pages from [4, Chapter 10], we instead refer the reader to the reference and sketch the steps involved in the argument.

Sketch Proof. (Theorem 4.3.5 - Weyl's Lemma)

1. We begin the proof by considering Φ on a single open coordinate chart of \mathbb{C} and showing that, in local charts, Φ can by identified with an L^2 function, by taking a sequence ϕ_i

that tends to Φ and showing that the L^2 norm $\|\phi_i - \phi_j\| \le C \|\phi_i - \phi_j\|_D$ is bounded by something which converges. Thus, the limit of ϕ_i is an L^2 function.

- 2. We then extend Φ to be L^2 over the whole of X, by considering a set A to be the set of points $x \in X$ with the property that ϕ_i converges in the L^2 sense in a coordinate chart around x. This set is open by definition, so the compliment of A is not open unless it is empty since X is connected. Then we prove by contradiction that A = X, otherwise there would be a point in the closure of A that is not in A.
- 3. We then return to a single coordinate chart and want to prove the Local Weyl lemma, as if we can show it is true, then Weyl's lemma must also be true.
- 4. We look to then simplify the problem by looking for a weak solution ψ to the Laplace equation $\Delta \psi = 0$. Such a solution is called a weak harmonic function. We construct this weak harmonic function in such a way that if we can prove that it is smooth, then it implies that our weak solution Φ is also smooth. This is done by constructing ψ as $\psi = \Phi \phi'$ where ϕ' is smooth.
- 5. We then show that the convolution $\beta(|z|) * \psi \psi$, where β is a bump function in a neighbourhood of the origin, vanishes outside an ϵ neighbourhood of supp $(\Delta \psi)$ if and only if ψ is smooth. We use this to show that ψ is smooth and hence Φ is also smooth by the Local Weyl lemma which implies that our weak solution is actually smooth on the whole of X.

Thus we have the rather remarkable fact that if we have a solution to the Poisson equation, it must be smooth. \Box

We note that it is not strictly necessary that we study Poisson's equation on all compact, connected Riemann surfaces as our goal is to classify compact, simply connected Riemann surfaces. However, as we will see in the next section, by studying compact, connected Riemann surfaces we can get some rather powerful tools to help us understand compact Riemann surfaces a bit better.

4.4 Cohomology groups of compact Riemann surfaces

The implications of having solved Poisson's equation are not immediately obvious, other than the statement of Theorem 4.3.1. So how can we use the guaranteed existence of such a function f for a given 2-form to help us classify connected, compact Riemann surfaces? How

exactly does Poisson's equation help us prove the uniformisation theorem? In the compact, connected case, we get the following theorem of relations between the Dolbeault and de Rham cohomology groups.

Theorem 4.4.1. Let X be a compact, connected Riemann surface. Then we have the following isomorphims:

- 1. The induced map $\sigma: H^{1,0}(X) \to \overline{H^{0,1}}(X)$, induced by the mapping $\alpha \mapsto \overline{\alpha}$ is an isomorphism.
- 2. The bilinear map $B: H^{1,0}(X) \times H^{0,1}(X) \to \mathbb{C}$, defined by $B(\alpha, \tilde{\theta}) = \int_X \alpha \wedge \theta$ gives an isomorphism between $H^{0,1}(X)$ and the dual space $(H^{1,0}(X))^*$.
- 3. Define the mapping $\mu: H^{1,0}(X) \times H^{0,1}(X) \to H^1(X)$, given by $\mu(\alpha,\beta) = i(\alpha) + \overline{i(\sigma^{-1}(\beta))}$, where $i: H^{1,0}(X) \to H^1(X)$ is the inclusion map, defined by mapping a holomorphic 1-form to its cohomology class in $H^1(X)$. The map μ defines an isomorphism.
- 4. The map $\nu: H^{1,1}(X) \to H^2(X)$, defined as the natural inclusion of $\operatorname{Im}(\overline{\partial}: \Omega^{1,0}(X) \to \Omega^2(X))$ in $\operatorname{Im}(d: \Omega^1(X) \to \Omega^2(X))$, is an isomorphism.

Sketch Proof. We do not prove these all here, though we do prove that the first map is surjective, as this explicitly uses the fact that we can solve the Poisson equation on X. Given a class $\tilde{\theta} \in H^{0,1}(X)$, we want to find an element $\theta' \in \Omega^{0,1}(X)$ such that $\theta' = \theta + \overline{\partial} f$ with $\partial \theta' = 0$, since this would mean that $\overline{\theta'}$ is a holomorphic 1-form and that $\tilde{\theta} = -\sigma(\overline{\theta'})$. Hence, we need to solve the equation $\partial \overline{\partial} f = -\partial \theta$, as this holds true when $\partial \theta' = 0$. Note that $\partial \overline{\partial} = \frac{i}{2} \Delta$, is just the Laplacian in disguise, so we can indeed solve this equation, given that $\int_X \partial \theta = 0$. But $\int_X \partial \theta = \int_{\partial X} \theta = 0$ by Stokes' theorem. Thus, σ is surjective. Injectivity can be proven by composing σ with the bilinear map from point 2 in the theorem giving a positive definite bilinear form, thus implying σ is injective. The other proofs follow similarly and require the solution of the Poisson equation.

These isomorphisms provide concrete relations between the de Rham and Dolbeault cohomology groups of our connected, compact Riemann surface. But what information do the Dolbeault cohomology groups capture in particular? We notice that the third isomorphism implies that the cohomology groups $H^{1,0}(X)$ and $H^{0,1}(X)$ not only have the same dimension, but when combined with the first isomorphism, we get that each group takes half the value of the dimension of $H^1(X)$. We discussed earlier that for a connected, compact, genus g smooth surface, Σ_g , we had that $H^1(\Sigma_g) \cong \mathbb{R}^{2g}$. Therefore, we have that $\dim_{\mathbb{R}} H^{1,0}(X) = \dim_{\mathbb{R}} H^{0,1}(X) = g$. So it appears that compact, connected Riemann surfaces can be classified by their cohomology groups. However, this is not the case since cohomology alone is not a strong enough invariant to distinguish between two Riemann surfaces, since two Riemann surfaces with equal cohomology groups may not be conformally equivalent. We encountered this when we introduced complex tori of different moduli ω_1 and ω_2 , as two complex tori \mathbb{T}_{ω_1} and \mathbb{T}_{ω_2} have the same cohomology groups but are not conformally equivalent.

In Chapter 2, we proved lemma 2.3.5 which gave us that we had a conformal equivalence between a compact, connected Riemann surface X and $\widehat{\mathbb{C}}$ if we could construct a meromorphic function on X such that it has a single simple pole. The following proposition generalises this result by giving us which meromorphic functions exist with simple poles on a genus g Riemann surface, but more importantly shows that for a surface of genus 0, we can always construct a meromorphic function with a single simple pole. This implies that the genus is not strictly a topological property of our surface, but has some underlying effect on the Riemann surface structures possible on such a surface.

Proposition 4.4.2. Let X be a connected, compact Riemann surface with $H^{0,1}(X)$ of finite dimension g. Then, given any distinct g+1 points on X, p_1, \ldots, p_{g+1} , there exists a meromorphic function on X with a simple pole at some (perhaps all) of the points p_1, \ldots, p_{g+1} .

Proof. We begin by considering a single point $p \in X$. This implies, by hypothesis, that we have a genus 0 Riemann surface. We now pick a coordinate chart around p such that p maps to 0, and let U be a small disc encircling p. Given an $\epsilon > 0$, we define a bump function β on U and define a function $f(z) = \beta(z)z^{-1}$, and note that whilst f is defined on U, we can extend f by 0 to the whole of X. We note further that f is holomorphic in the region $|z| < \epsilon/4$ since $f(z) = z^{-1}$ in this region. We lose this fact as soon as β stops being constant however since we can't impose that β decays to zero holomorphically. We now construct a useful (0,1)-form on X from f, namely $A = \overline{\partial} f(z) = \overline{\partial} (\beta(z))z^{-1}$. By definition, A is an element of $\Omega^{0,1}(X \setminus \{p\})$. We would like to extend A over the whole of X and understand this extension's cohomology class. Namely, we want to understand the class $[A] \in H^{0,1}(X) = \Omega^{0,1}(X)/\operatorname{Im}(\overline{\partial}:\Omega^0(X) \to \Omega^{0,1}(X))$.

We see from the definition of $H^{0,1}(X)$ that we can find a smooth function g on X such that $\overline{\partial}(g) = A$ if and only if [A] = 0. We see that A can be easily extended over the whole of X. In the annulus $\epsilon/4 < |z| \le 3\epsilon/4$, we see the A has some value corresponding to the rate of decay of β . So let us consider when $|z| \le \epsilon/4$. In this case, we see that β is constant so $\overline{\partial}\beta = 0$ for

all $|z| \leq \epsilon/4$. Thus, the limit of A as z approaches 0 vanishes and so we have extended A over the whole of X. Thus, we have that because X is a genus 0 surface, $H^{0,1}(X) = 0$ and, [A] = 0, and so there exists a smooth function $\psi \in \Omega^0(X)$ such that $\overline{\partial}\psi = A$. Note that, if we multiply the function f by some $\lambda \in \mathbb{C}$ then our (0,1)-form A becomes λA and hence, any associated cohomology class becomes $\lambda[A]$.

Now, we let ϕ be a function defined as $\phi = \psi + \lambda f$ for some $\lambda \in \mathbb{C}$. The function ψ is defined to be smooth over all of X and f is smooth on $X \setminus \{p\}$ so ϕ is also smooth on $X \setminus \{p\}$. Since we have that ϕ is smooth, if we can show that this ϕ is holomorphic near p then it is precisely the meromorphic function we need. We know that f, on the punctured disc $\overline{B_{\epsilon/4}(0)} \setminus \{0\}$, is holomorphic, therefore so is λf . Therefore, if we can show that ψ is holomorphic on this disc, then ϕ will be holomorphic on the punctured disc. But we know that on X, we have $\overline{\partial}\psi = A$ and, except for in the annulus $\epsilon/4 < |z| \le 3\epsilon/4$, we have that A = 0. In particular, on the disc of radius $\epsilon/4$ we have that A = 0 and so $\overline{\partial}\psi = 0$. Hence, ψ is holomorphic and so it follows that, in this punctured disc, ϕ is holomorphic. Thus, in this disc around p, the function ϕ is meromorphic, with a pole at p and smooth over X, as required.

We can now generalise this construction, whereby for g+1 points p_1, \ldots, p_{g+1} , we can repeat the exact same procedure, creating functions f_1, \ldots, f_{g+1} such that they extend by zero over X and have a single, simple pole at the points p_1, \ldots, p_{g+1} respectively. Then we can define forms A_i and show that since the cohomology group $H^{0,1}(X) \cong \mathbb{R}^g$ (implying it is g dimensional), then the linear combination $\lambda_1[A_1] + \ldots + \lambda_{g+1}[A_{g+1}] = 0$, is a linear dependency, implying that the equality to zero is had even if not all $\lambda_i = 0$. Then we proceed as before with this linear combination equalling zero, rather than with [A] = 0 as we did in the single point case, thereby constructing a meromorphic function on X with at most g+1 simple poles.

So the existence of the linear dependency of the cohomology classes of (0,1)-forms, over their respective poles, gives a link between the genus of X and the number of poles a meromorphic function may have on X. The number of points which have a pole over them correspond to which of the λ_i are non zero in the linear dependency. This result is incredibly powerful as it turns the topological notion of a genus into an invariant of connected, compact Riemann surfaces. It means two compact Riemann surfaces of different genus *cannot* be conformally equivalent. Therefore, we have the following:

Corollary 4.4.3. Let X be a connected, compact Riemann surface.

1. If X has genus 0, then X is conformally equivalent to the Riemann sphere.

- 2. If X has genus 1, then X is conformally equivalent to a complex torus.
- 3. If X has genus n, then X is conformally equivalent to a smooth, orientable surface of genus n.

Note especially the first point in Corollary 4.4.3: If we have a connected, simply connected, compact Riemann surface, it is conformally equivalent to the Riemann sphere $\widehat{\mathbb{C}}$, since the above lemma constructed a meromorphic function with a single simple pole on a genus 0 surface thus giving us the conformal equivalence to $\widehat{\mathbb{C}}$. Thus, crucially, we have classified all connected, simply connected, compact Riemann surfaces.

4.5 Poisson's equation on simply connected, non compact Riemann surfaces

Before we move onto the uniformisation theorem, we need to briefly also discuss the case of connected, simply connected, non compact Riemann surfaces. At the end of Chapter 2, we showed that all Riemann surfaces arise as quotients of a simply connected covering Riemann surface. We hypothesised that the three covering spaces are $\widehat{\mathbb{C}}$, \mathbb{C} , \mathbb{H} and have hence shown that the only compact, simply connected Riemann surface is $\widehat{\mathbb{C}}$, by studying compact, connected Riemann surfaces generally. We now wish to fill in the gap. The loss of compactness complicates things slightly so, to keep things simple, we discuss only simply connected non compact Riemann surfaces. This section is a discussion of the method of proving the following theorem rather than a formal analysis of the proof as the results are well known and referenced accordingly. Let us state the problem of Poisson's equation for connected, simply connected, non compact Riemann surfaces.

Theorem 4.5.1. Let X be a connected, simply connected, non compact Riemann surface. If ρ is a 2-form of compact support in X with $\int_X \rho = 0$ then there is a smooth function $\phi \in \Omega^0(X)$ such that $\Delta \phi = \rho$, with ϕ tending to 0 at infinity in X.

We see that this case introduces a subtly weaker constraint on ϕ , namely that ϕ tends to 0 at infinity. By that we mean the following.

Definition 4.5.2 (A function tending to a value at infinity on X). Let X be a connected, simply connected, non compact Riemann surface. If we have a smooth function $f \in \Omega^0(X)$, and a real number $c \in \mathbb{R}$, then we say that f tends to c at infinity in X if, for all $\epsilon > 0$ there is a compact subset $K \subset X$ such that $|f(x) - c| < \epsilon$ for $x \notin K$. Conversely, we say that

f tends to $+\infty$, at infinity in X if for all $C \in \mathbb{R}$, there is a compact set $K \subset X$ such that f(x) > C for all $x \notin K$ (and we say that f tends to $-\infty$ in X if -f tends to $+\infty$).

If one thinks about our physical interpretation of the problem, this is equivalent to saying that, at large enough scales, the effect of an electric potential on a test charge is negligible, which is in fact a real "boundary condition" used in electrostatics. So it is reassuring to see such a condition appearing in our study. So what differs in the proof of the non compact case? We let such a connected, simply connected, non compact Riemann surface be denoted as X. We see that we can still construct the Hilbert space $\overline{H}(X)$, and construct functionals $\underline{\hat{\rho}}$ and $\underline{\mathcal{L}}$ for some $\rho \in \Omega_c^2(X)$ in much the same way as for compact X. We have that $\underline{\mathcal{L}}$ is still bounded from below and so we run through the argument of Weyl's lemma to show that the minimising function is L^2 . Locally, the solution still exists. The problem now arises when we try to understand what happens over the whole surface X. Now the minimising function, whilst smooth, doesn't behave in quite the same way, leading to the following proposition.

Proposition 4.5.3. [4, Proposition 32] Let ρ be a 2-form of integral 0 supported in a coordinate disk D. Then there is a smooth function ϕ on X which satisfies the equation $\Delta \phi = \rho$, and a sequence of functions ϕ_i on X which have the following properties:

- 1. There are real numbers c_i and compact sets $B_i \subset X$ such that $\phi_i = c_i$ outside B_i .
- 2. For any 1-form $\alpha \in \Omega^1_c(X)$, the L^2 norms $\|(\phi \phi_i)\alpha\|$ tend to 0 as i tends to infinity.
- 3. The L^2 norms $||d\phi d\phi_i||$ tend to 0 as i tends to infinity.

The goal is then to show that outside of $X \setminus \text{supp}(\rho)$, we can make ϕ tend to 0. Equally, if we can make ϕ tend to a constant, then we can replace ϕ by $\phi - c$ and have that go to 0 without changing the solution to the Laplace equation. The proof then goes to show that outside of $X \setminus \text{supp}(\rho)$, we can classify that space as being the unique 'End' of the surface X. This idea, due to Freudenthal, says that an end of a (locally) connected space X is an assignment to each compact subset $K \subseteq X$, a connected component of X whose closure $X \setminus K$ is not compact, such that if $K \subseteq K'$ then $X_K \subseteq X_{K'}$. In the referenced material, it is shown henceforth that connected, simply connected, non compact Riemann surfaces have exactly one end, and that we can arrange our solution to tend to 0, thus giving us our solution to Poisson's equation. This argument is studied in depth in [4, Chapter 10], but a gentler approach is available in [10, p.32].

4.6 The Uniformisation Theorem

We are finally ready to discuss the uniformisation theorem. We state this as two points, the first of which we already proved when considering the case of compact, connected Riemann surfaces. Thus, we only need to show that the existence of a solution to Poisson's equation on connected, simply connected, non compact Riemann surfaces implies the second point of the uniformisation theorem.

Theorem 4.6.1 (The Uniformisation Theorem). Let X be a connected, simply connected Riemann surface.

- 1. If X is compact, it is conformally equivalent to $\widehat{\mathbb{C}}$, the Riemann sphere.
- 2. If X is not compact, the X is conformally equivalent to either \mathbb{C} , the complex plane or \mathbb{H} , the upper half plane.

Thus we can justifiably classify all Riemann surfaces as being quotients of $\widehat{\mathbb{C}}$, \mathbb{C} or \mathbb{H} .

Corollary 4.6.2. Any connected Riemann surface is conformally equivalent to one of the following:

- The Riemann sphere, $\widehat{\mathbb{C}}$.
- The Complex Plane \mathbb{C} , the Cylinder $\mathbb{C}/2\pi\mathbb{Z}$ or the Torus \mathbb{C}/Λ , for some lattice $\Lambda \subset \mathbb{C}$.
- The Upper Half Plane \mathbb{H} quotiented by a discrete subgroup $\Gamma \subset PSL(2,\mathbb{R})$.

Riemann surfaces may be called Elliptic, Parabolic and Hyperbolic, depending on their universal cover.

Now, we want to use the fact that we have a solution to Poisson's equation on simply connected, non compact Riemann surfaces, to prove the uniformisation theorem. We do this by showing that Poisson's equation for simply connected, non compact Riemann surfaces implies the uniformisation theorem for the non compact surfaces. Our strategy will be to try to find a meromorphic function on our simply connected, non compact Riemann surface X, with a single simple pole, at a point $p \in X$, with this function's imaginary part tending to 0 at infinity on X.

Proof. (Theorem 4.6.1 - The Uniformisation Theorem) We pick local coordinates about p, and define a (0,1)-form $A = \overline{\partial}\beta(z)/z$ for some bump function β centred on p. Thus, A is supported in an annulus around p. By choosing $\rho = \partial A$, then we have that, by Stokes' theorem, the integral of ρ on X is 0. When considering our complex valued 2-form ρ as a sum

 $\rho = \rho_r + i\rho_i$, where both ρ_r and ρ_i are real-valued 2-forms, our analysis of the Poisson equation gives us that we have a complex valued function g such that the real and imaginary parts of g tend to 0 at infinity and $\partial \overline{\partial} g = \rho$. Now let B be the real 1-form $B = (A - \overline{\partial} g) + (\overline{A} - \overline{\overline{\partial} g})$. By construction, the first bracket is holomorphic, so dB = 0. Since X is simply connected, we have that $H^1(X) = 0$, and so every closed form is exact. Thus we can find a real function ψ such that $d\psi = B$. Thus, we can write $A = \overline{\partial} g + \overline{\partial} \psi$. Therefore, in an annulus about p, we can write that $\overline{\partial} \left(\frac{\beta}{z} - (g + \psi)\right) = 0$. Thus, we have a meromorphic function $f = g + \psi$ on f0 with a simple pole at f1 whose imaginary part tends to zero at infinity on f2, since f3 is a real function. So we have a meromorphic function f3 from f4 to the Riemann sphere. Recall, a meromorphic function to the Riemann sphere is a holomorphic map.

By considering the topology⁵ of the Riemann sphere, and some subtle topological arguments relating the function f and its restrictions on open sets on $\widehat{\mathbb{C}}$, we can show that f is a proper map that maps our surface X into an open subset $U \subset H_+ \cup H_- \cup \{\infty\}$, where H_+ and $H_$ denote the open upper and lower half planes in \mathbb{C} , injectively. Thus, we get that $U = \widehat{\mathbb{C}} \setminus K$ for some compact subset $K \in \mathbb{C}$. Therefore, X is conformally equivalent to this set U. This means that K must be connected since otherwise U would not be simply connected. Thus, K is either a closed interval [a,b] for a < b or a single point $\{p\}$. In the case of K being a single point $\{p\}$, we see that this is the case of U being conformally equivalent to \mathbb{C} as it is puncturing the Riemann sphere, thereby undoing the compactification of the complex plane that we used in the first place to describe $\widehat{\mathbb{C}}$. We can treat the case of K = [a, b]by using automorphisms of $\widehat{\mathbb{C}}$, to map one endpoint of our closed interval, say b, to ∞ and mapping a to 0, with the line connecting them going along the positive real axis. With this arrangement, $\widehat{\mathbb{C}} \setminus K$ becomes a branch cut of the plane, which, when we compose with the map $z \mapsto \sqrt{z}$, gives us a conformal equivalence to \mathbb{H} , since the arguments of all points $z \in \widehat{\mathbb{C}} \setminus K$ get halved. Thus we have that, up to conformal equivalence, there are exactly 2 simply connected, non compact Riemann surfaces.

So that is it! We have shown that, up to conformal equivalence, there are exactly three connected, simply connected Riemann surfaces! Combining the result of this theorem with the conclusions reached at the end of Chapter 2, gives the rather amazing conclusion of corollary 4.6.2 that indeed every Riemann surfaces, except for the Riemann sphere, the complex plane, the punctured plane, the cylinder, and complex tori, are conformally equivalent to some quotient of the upper half plane by a discrete subgroup of $PSL(2,\mathbb{R})$.

 $^{^5}$ For an in depth proof, see [4, p.132-133]

Conclusion

The uniformisation theorem is an incredibly powerful theorem and whilst it classifies connected Riemann surfaces into three broad families, the study of Riemann surfaces by no means concludes there. Proving the uniformisation theorem required us to first show that we only needed to study simply connected Riemann surfaces. We did this by showing that every Riemann surface arises as a quotient of a simply connected Riemann surface, X, by some discrete subgroup of the group Aut(X). This led us to try and understand which simply connected Riemann surfaces exist and what their automorphism groups look like. We saw that using a problem from electrostatics, we could define a functional on our Riemann surfaces whose minimising functions were those satisfying the Poisson equation. Thus our study turned to the Poisson equation for two types of Riemann surface. In the compact case, this gave rise to a powerful relationship between the Dolbeault and de Rham cohomology groups of our compact, connected Riemann surfaces. In fact, we came just short of proving the coveted Riemann-Roch theorem, a theorem which provides a link between the types of holomorphic and meromorphic functions possible on a compact, connected Riemann surface X, the degree of said functions and the genus of the surface X. For non compact surfaces, we showed that the setup was largely the same and though we needed to take some technical detours, those technical detours gave us a solution to Poisson's equation, which we then used to prove the uniformisation theorem. However, Riemann surfaces go beyond that. As mentioned, they provide fascinating connections between different areas of mathematics with one of the most interesting, in the author's opinion, being the link between closed geodesics on compact, connected Riemann surfaces and the prime number theorem. The relationship arises due to [14], in which the author posits that there is a relationship between the values of the prime counting function $\pi(x)$ and the norms of elements in a discrete subgroup of $SL(2,\mathbb{R})$, which correspond to the lengths of simple, closed geodesics on the corresponding Riemann surface. Such connections are remarkable and though we studied in depth only one particular aspect of them, Riemann surfaces are a broad and truly wonderful field of study.

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