# Complex Abelian Varieties and Elliptic Curves 

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## Chapter 1

## Introduction

This text is essentially a compilation of notes I used to give lectures for a graduate class I taught at the University of Milan in late 2016 and early 2017. My original conception was of a first-semester course on elliptic curves mostly from an algebraic background as was taught to me when I was in graduate school. However, after interviewing some of my students beforehand, I decided to choose which material to include based on their interests and background. As they were very strong in complex analysis and complex algebraic geometry, I wound up spending most of my lecture time on the theory of complex abelian varieties (theta functions on complex tori, algebraization, uniformization of complex elliptic curves, and a quick survey of polarizations), with only a slight emphasis on elliptic curves. In the last several lectures I introduced the algebraic theory of elliptic curves over arbitrary fields and was able to go over certain material on the classification of their endomorphism rings. Unfortunately however, I ran out of time to get to many crucial algebraic topics such as the Weil pairing, reduction of elliptic curves over local fields, and the Mordell-Weil theorem.

I taught the course assuming prerequisite knowledge based on the areas in which my students generally seemed to have strong backgrounds. This is also reflected in the lecture notes. For instance, these notes assume basic background knowledge of complex manifolds and algebraic topology. They also assume that the reader has a reasonably strong background in abstract and linear algebra, and that they have taken a first course in algebraic geometry or at least has basic knowledge of the theory of projective varieties over fields. I went out of my way to avoid any mention of schemes, however, and in developing the theory of theta functions I chose to use the language of Cartier divisors (later changed to Weil divisors in the case of elliptic curves) rather than line bundles. Some facts from algebraic geometry which are a little less basic, such as the Riemann-Roch Theorem, are explained carefully but not proven. The same is true of basic Lie theory, which is used only for $\$ 2.1$. I have kept algebraic number theory to a minimum (although if the course had gone on longer, of course there would have been much more of it). Several ideas and results that I did not have time to teach are outlined in some of the provided exercises.

My main sources for Chapter 2 were the first chapter of David Mumford's book Abelian Varieties (5]) as well as Michael Rosen's and James Milne's articles on abelian varieties in the book Arithmetic Geometry edited by Cornell and Silverman (7] and [3] respectively), the latter of which was used particularly for the material on polarizations. For the AbelJacobi Theorem I mainly consulted Serge Lang's book Introduction to Algebraic and Abelian

Functions ([2]) while adopting some variations on his proof; for someone with a more complex analytic background, [1] might be a good alternative. For the algebraic material in Chapter 3. more or less everything came from Silverman's (first) book The Arithmetic of Elliptic Curves ([8]).

### 1.1 Three definitions for elliptic curves

We begin by giving three definitions for an elliptic curve. Eventually we will explain each definition more clearly and show that the definitions are equivalent. Our first definition presents elliptic curves as a particular case of a much more general object: abelian varieties. Our initial focus in this course will be on examining abelian varieties over the complex numbers (that is, we will assume $K=\mathbb{C}$ in the definition below) with a particular emphasis on complex elliptic curves.

Definition 1.1.1. A group variety over a field $K$ is an algebraic variety $A$ over $K$ with the property that there is a group law on the set of points $A(K)$ such that group multiplication and inversion are given by morphisms $m: A \times A \rightarrow A$ and $i: A \rightarrow A$. An abelian variety is a group variety which is complete.

An elliptic curve over a field $K$ is an abelian variety of dimension 1 over $K$.
Remark 1.1.2. a) We note that an abelian variety is smooth. Indeed, on any variety over $K$, there is a $K$-point $x_{0}$ at which the variety is smooth. If $A$ is an abelian variety, for each $y \in A$, the translation-by- $y$ map $t_{y}: A \rightarrow A$ given by $x \mapsto y x$ is an invertible morphism from $A$ to itself and therefore induces an isomorphism on the tangent spaces $\left(t_{y}\right)_{*}: T_{x_{0}} X \xrightarrow{\sim} T_{y x_{0}} X$. Since every $a \in A$ is equal to $y x_{0}$ for some $y$, the tangent spaces at all points of $A$ are isomorphic to $T_{x_{0}} A$, and so $A$ is smooth everywhere.
b) It is possible to show from Definition 1.1 .1 that the group law on $A$ is commutative, using in particular the fact that $A$ is complete (see [5, §4] or [3, §2]). We will show this for complex abelian varieties (i.e. $K=\mathbb{C}$ ) in $\$ 2.1$ but not in the general case. So every abelian variety has the structure of an abelian group. The use of the adjective "abelian" here is somewhat coincidental, as they were each independently named after Henrik Abel. However, note that a group variety may have an abelian group structure without being complete and therefore without being an abelian variety, e.g. the affine line $\mathbb{A}_{K}^{1}$ with additive group structure, or the punctured affine line $\mathbb{A}_{K} \backslash\{0\}$ with multiplicative group structure.

The next definition comes more directly from the classical setting of algebraic curves.
Definition 1.1.3. An elliptic curve over a field $K$ is a smooth projective genus-1 curve $E$ over $K$ along with a distinguished $K$-point $O \in E(K)$.

Finally, we give the most elementary definition, to provide motivation for what follows by giving a more concrete idea of the structure inherent in an elliptic curve (note that here and everywhere below, $\bar{K}$ denotes an algebraic closure of $K$ ).

Definition 1.1.4. An elliptic curve $E$ over a field $K$ is the locus $E(\bar{K})$ of points $(x, y) \in$ $\bar{K}^{2}$ satisfying an equation of the form $y^{2}=f(x):=x^{3}+a x+b$ for $a, b \in K$ such that the
discriminant $-4 a^{3}-27 b^{2}$ is nonzero (i.e. the cubic polynomial $f(x)$ doesn't have multiple roots), along with an extra "point at infinity" which we denote by $O$. It is endowed with a binary operation $(P, Q) \mapsto P+Q$ defined as follows.

For any points $P=\left(x_{P}, y_{P}\right)$ and $Q=\left(x_{Q}, y_{Q}\right)$ in $E \backslash\{O\}$, let $L$ be the line connecting them. By convention, if $P=Q$, then $L$ is the tangent line to the curve given by $y^{2}=f(x)$ at $P$. Then the line $L$ intersects a unique third point $R=\left(x_{R}, y_{R}\right)$ in $E(\bar{K})$ (if the line $L$ is vertical, then we take $R=O$ ). We define the sum of $P$ and $Q$, denoted $P+Q$ or $Q+P$, to be the point $\left(x_{R},-y_{R}\right)$ (which is again $O$ if $R=O$ ).

For any point $P \in E(\bar{K})$, we define the sum of $P$ and $O$ to be $P+O=O+P=P$.
For any algebraic extension $K^{\prime} \supseteq K$, we define $E\left(K^{\prime}\right)$ to be the set of points in $E(\bar{K})$ whose coordinates lie in $K^{\prime}$. By convention, $O \in E\left(K^{\prime}\right)$.

Remark 1.1.5. For any algebraic extension $K^{\prime}$ of $K$ the set of points $E\left(K^{\prime}\right)$ is really the set of $K^{\prime}$-points of the projective curve given by the homogenization $y^{2} z=x^{3}+a x z^{2}+b z^{3}$ of the defining equation. Viewed this way, the "point at infinity" $O$ is given by $(x: y: z)=$ ( $0: 1: 0$ ). Loosely speaking, we may visualize it as a point lying "above" the affine curve on the two-dimensional coordinate plane, with the $y$-coordinate being infinity. In this way, it makes sense that every vertical line $x=x_{0}$ intersects the curve at points $P:=\left(x_{0}, y_{0}\right)$, $\left(x_{0},-y_{0}\right)$, and $O$, thus justifying the assignment of $P+O=O+P=P$ and the convention that "if the line $L$ is vertical, then we take $R=O$ " above.

It turns out that the binary operation on the set $E(\bar{K})$ which we defined above is a commutative group law (thus justifying our use of the "+" notation).

Proposition 1.1.6. The binary operation given in Definition 1.1.4 is a group law on the set $E(\bar{K})$. More precisely,
a) it is commutative, and the point $O \in E(\bar{K})$ acts as an additive identity;
b) any point $P=\left(x_{P}, y_{P}\right) \in E(\bar{K})$ has an inverse $-P \in E(\bar{K})$ given by $-P=\left(x_{P},-y_{P}\right)$; and
c) it is associative, i.e. $(P+Q)+R=P+(Q+R)$ for any $P, Q, R \in E(\bar{K})$.

Parts (a) and (b) of the above proposition are immediate from Definition 1.1.4, while part (c) is very tedious to prove (we will prove it later in an elegant way using the concept of Picard group).

Remark 1.1.7. Note that for any algebraic extension $K^{\prime} \supseteq K, E\left(K^{\prime}\right)$ is a subgroup of $E(\bar{K})$, and if $K^{\prime \prime}$ is an algebraic extension of $K^{\prime}$, then $E\left(K^{\prime \prime}\right)$ is a subgroup of $E\left(K^{\prime}\right)$. This follows from the fact that the slope of the line $L$ from the definition must be an element of any field containing the coordinates of the points $P$ and $Q$.

In fact, it is possible to write down a formula for the addition law involving rational functions of the coordinates of the two imput points, where all coefficients lie in $K$. In the language of algebraic geometry, this is equivalent to the very important fact that the addition law is a morphism $E \times E \rightarrow E$ defined over $K$.

### 1.2 Exercises

Exercise 1.2.1. Let $E$ be the elliptic curve over $\mathbb{Q}$ given by the locus of points satisfying $y^{2}=x^{3}-x$; let $P=(0,0) \in E(\mathbb{Q}) \subset E(\overline{\mathbb{Q}})$; and let $Q:=\left(x_{0}, y_{0}\right)$ be any point in $E(\overline{\mathbb{Q}}) \backslash\{(0,0)\}$. Compute $P+Q=\left(-1 / x_{0},-y_{0} / x_{0}^{2}\right)$.

Exercise 1.2.2. Actually prove directly part (c) of Proposition 1.1 .6 - that is, show the addition law given in Definition 1.1.4 satisfies the associative property. (This is extremely tedious and I have never fully attempted it!)

Exercise 1.2.3. Given any elliptic curve $E: y^{2}=f(x)$ as in Definition 1.1.4 over a field $K$, characterize all points of order 2 in $E(K)$. Show that the 2-torsion subgroup of $E(K)$ must be finite. What are the possible structures of the 2-torsion subgroup of $E(K)$, and how do they depend on the cubic polynomial $f(x) \in K[x]$ ? What is the structure of the 2-torsion subgroup of $E(\bar{K})$ ?

### 1.3 Outline of these notes

We want to relate the three definitions given above for elliptic curves by studying the more general objects known as abelian varieties. We will show that an abelian variety of dimension 1 has to be a curve with the properties given in Definition 1.1.3 and also that any curve satisfying Definition 1.1 .3 can be expressed in the form given in Definition 1.1.4 with an additive group law on its points. In our journey through these different ways of viewing elliptic curves, we will stop to further examine many interesting results.

We will first (in Chapter 2) study the theory of complex abelian varieties from the point of view of examining complex tori and certain meromorphic functions known as theta functions which can be used to realize (some of) them as abelian varieties. We will also state and prove the Abel-Jacobi theorem (in $\$ 2.6$ ), which can be used to construct classical examples of complex abelian varieties (where the ground field $K$ is $\mathbb{C}$ ), including complex elliptic curves.

We will then (in Chapter 3) examine elliptic curves over a general ground field $K$, focusing mainly on maps between elliptic curves and endomorphism rings. In particular, we will study ordinary and supersingular elliptic curves over finite fields along with their endomorphism rings (in 3.2.4).

## Chapter 2

## Complex abelian varieties

The goal of this chapter is to describe as fully as possible what properties characterize complex abelian varieties and in particular complex elliptic curves. We start by noting that by Definition 1.1.1, a complex abelian variety is in particular a complete complex variety. Since completeness of a complex variety implies compactness as a complex manifold, and the group law as in the definition gives this complex manifold the structure of a complex Lie group, every complex abelian variety is a connected compact complex Lie group. We will therefore begin our study of complex abelian varieties by investigating connected compact complex Lie groups. We will see that in dimension 1, every connected compact complex Lie group is an abelian variety (an elliptic curve), although many connected compact complex Lie groups of dimension $\geq 2$ cannot be given an algebraic structure and therefore are not abelian varieties. Our most important result (Theorem 2.3.6) will be a criterion for when a connected compact complex Lie group is an abelian variety.

### 2.1 Connected compact complex Lie groups

We first want to determine the structure of connected compact complex Lie groups (for the moment, we are not considering whether or not they are algebraic).

Proposition 2.1.1. Let $X$ be a connected compact complex Lie group. Then $X$ is commutative.

Proof. Let $e$ denote the identity element of $X$ and let $T_{e} X$ denote the tangent space of $X$ at $e$; it is a complex vector space of dimension equal to the dimension of $X$. For any $x \in X$, write $\phi_{x}: X \rightarrow X$ for the conjugation-by- $x$ map $a \mapsto x^{-1} a x$. Then $\phi_{x}$ induces an endomorphism $\left(\phi_{x}\right)_{*}: T_{e} X \rightarrow T_{e} X$ of the tangent space. It is clear that the map $\phi: X \rightarrow \operatorname{End}_{\mathbb{C}}\left(T_{e} X\right)$ sending $x \in X$ to the endomorphism $\left(\phi_{x}\right)_{*} \in \operatorname{End}_{\mathbb{C}}\left(T_{e} X\right)$ is a homomorphism of groups (so its image lies in the group of automorphisms $\operatorname{Aut}\left(T_{e} X\right)$ ) as well as a holomorphic map. (In fact, it can be used to define the Lie bracket on $T_{e} X$.) Since $X$ is compact and $\operatorname{End}_{\mathbb{C}}\left(T_{e} X\right)$ is affine, $\phi$ must be a constant map and is therefore the trivial homomorphism; that is, $\left(\phi_{x}\right)_{*}=1 \in \operatorname{Aut}\left(T_{e} X\right)$ for all $x \in X$.

Recall that the exponential map from the Lie algebra $T_{e} X$ to the Lie group $X$ is given by $v \mapsto \gamma_{v}(1)$ where $\gamma_{v}: \mathbb{C} \rightarrow X$ is the unique holomorphic map whose differential
$\left(\gamma_{v}\right)_{*}: \mathbb{C}=T_{0} \mathbb{C} \rightarrow T_{e} X$ sends $1 \in T_{0} \mathbb{C}$ to $v \in T_{e} X$. The exponential map has the important property of being a biholomorphism when restricted to the inverse image of a small enough open neighborhood of $0 \in T_{e} X$. Moreover, exponential maps commute with holomorphic homomorphisms between Lie groups and their induced maps on the corresponding Lie algebras (which are the tangent spaces at the identity). It follows from these properties that given some small enough open neighborhood $U \ni 1 \in X, \phi_{x}$ acts as the identity on $U$ for all $x \in X$. Now we claim that any such $U$ generates $X$ as a group. Indeed, if we let $U^{\prime}$ denote the subgroup of $X$ generated by $U$, then $U^{\prime}$ is also open since for each $x \in U^{\prime}$, $x U \subseteq U^{\prime}$, and the connectedness of $X$ proves that its only open subgroup is $X$ itself. It follows from this that since each $\phi_{x}$ is a group automorphism, $\phi_{x}$ must act as the identity on all of $X$ for any $x \in X$. Thus, $X$ is commutative.

In light of the above proposition, from now on, we will use " + " to denote the group operation on $X$ and denote the identity of $X$ by 0 .
Proposition 2.1.2. Let $X$ be a connected compact complex Lie group of dimension $g$. Then $X$ is a complex torus; that is, $X \cong V / \Lambda$, where $V$ is a complex vector space of dimension $g$, $\Lambda$ is a full lattice (i.e. of rank $2 g$ ) in $V$, and $V / \Lambda$ is given the Lie group structure it inherits as a quotient of the Lie group $V$.

Proof. Let $\pi: T_{0} X \rightarrow X$ be the exponential map defined in the proof of Proposition 2.1.1. One can show directly from the definition of the exponential map that the commutativity of $X$ implies that $\pi$ is a homomorphism. Let $a \in X$ be an element lying in the image of $\pi$, and let $U$ be a small enough open neighborhood of $1 \in X$ such that $\left.\pi\right|_{W}: W \rightarrow U$ is an isomorphism for some open $W \subseteq \pi^{-1}(U)$ with $0 \in W$. Then clearly the subset $a+U$ is contained in the image of $\pi$, and so the image of $\pi$ is open. Since $X$ is connected, this implies that $\pi$ is surjective. Now if $v \in T_{0} X$ is any element of $\operatorname{ker}(\pi)$, it is clear that for $W$ as above, $(v+W) \cap \operatorname{ker}(\pi)=\{v\}$. It follows that $\operatorname{ker}(\pi)$ is a discrete subgroup of $T_{0} X$. The fact that $X$ is compact implies that $\Lambda:=\operatorname{ker}(\pi)$ is a full lattice in $V:=T_{0}(X)$, and the proposition is proved.

Remark 2.1.3. It is easy to see from the holomorphic isomorphism $X \cong V / \Lambda$ that $\pi: V \rightarrow$ $X$ is a covering map. Since $V$ is simply connected, it is in fact a universal covering space for $X$. It immediately follows that the fundamental group $\pi_{1}(X, 0)$, as well as the first singular homology group $H_{1}(X, \mathbb{Z})$, can be identified with $\Lambda \cong \mathbb{Z}^{2 g}$. Note moreover that this result shows that $X$ is homeomorphic to $\left(S^{1}\right)^{2 g}$, so it is very easy to compute the dimensions of the singular homology groups $H_{i}(X, \mathbb{Z})$ for all $i \geq 0$ in terms of the dimension of $X$.

Corollary 2.1.4. As an abstract group, $X$ is divisible; that is, for any nonzero integer $n$, the multiplication-by-n map $[n]: X \rightarrow X$ is surjective. Moreover, the kernel $X[n]$ of $[n]$ is isomorphic to the group $(\mathbb{Z} / n \mathbb{Z})^{2 g}$, where $g$ is the dimension of $X$.
Proof. We construct an isomorphism of real vector spaces $V \xrightarrow{\sim} \mathbb{R}^{2 g}$ by choosing a basis of the free $\mathbb{Z}$-module $\Lambda \subset V$ and sending each basis element to a standard basis element of $\mathbb{R}^{2 g}$. This induces an isomorphism of quotient groups $X \cong V / \Lambda \xrightarrow{\sim}(\mathbb{R} / \mathbb{Z})^{2 g}$, from which the statements immediately follow.

### 2.2 Divisors on complex tori

We now want to study the group of divisors on any complex torus $X$, with our eventual goal being a characterization of all ample divisors on $X$ which (if they exist) may be used to characterize $X$ as an algebraic variety. In this subsection, we will only use a fairly elementary definition of Cartier divisors, which will be more convenient for the moment, although the reader who has studied divisors in the context of algebraic geometry should keep in mind that this is equivalent to the notion of Weil divisors (formal sums of codimension-1 subvarieties) in the case that $X$ is a smooth complex variety.

Definition 2.2.1 (Cartier divisors). Let $Y$ be a connected complex manifold.
a) The (additive) group $\operatorname{Div}(Y)$ of "local function data" on $Y$ is defined as follows. A "local function datum" is given by $D=\left\{\left\{U_{i}\right\}_{i \in I},\left\{f_{i}\right\}_{i \in I}\right\}$, where $\left\{U_{i}\right\}_{i \in I}$ is a finite open cover of $Y$ and, for each $i \in I, f_{i}$ is a nonzero meromorphic function defined on the open subset $U_{i} \subset$ $Y$ such that for any $i, j \in I$ such that $U_{i} \cap U_{j} \neq \varphi$, the quotient $f_{i} / f_{j}$ (and $f_{j} / f_{i}$ ) is holomorphic and nonvanishing on $U_{i} \cap U_{j}$. We define the sum of two elements $D_{1}=\left(\left\{U_{i}\right\}_{i \in I},\left\{f_{i}\right\}_{i \in I}\right)$ and $\left.D_{2}=\left(\left\{V_{j}\right\}_{j \in J},\left\{g_{j}\right\}_{j \in J}\right\}\right)$ to be $D_{1}+D_{2}:=\left(\left\{U_{i} \cap V_{j}\right\}_{(i, j) \in I \times J},\left\{f_{i} g_{j}\right\}_{(i, j) \in I \times J}\right)$.
b) A "local function datum" $\left(\left\{U_{i}\right\}_{i \in I},\left\{f_{i}\right\}_{i \in I}\right)$ defined in this way is said to be effective if each $f_{i}$ is holomorphic on $U_{i}$. It is said to be trivial if each $f_{i}$ is both holomorphic and nonvanishing on $U_{i}$.
c) The group of (Cartier) divisors on $Y$, denoted $\operatorname{Div}(Y)$, is defined to be the group equivalence classes of "local function data" on $Y$ modulo trivial elements. If a divisor $D$ is identified in this way with some $\left(\left\{U_{i}\right\}_{i \in I},\left\{f_{i}\right\}_{i \in I}\right)$ for some open cover $\left\{U_{i}\right\}_{i \in I}$, then we say that $D$ can be "represented by" this data. It is an easy exercise to check that the property of effectivity is well-defined for elements of $\operatorname{Div}(Y)$.
d) A divisor $D \in \operatorname{Div}(Y)$ is said to be principal if it can be represented by $(\{Y\},\{f\})$ for some nonzero meromorphic function $f$ on $Y$. In this case, we denote $D$ by $(f)$.

We observe right away that the divisor group is commutative (hence the "+" notation) and that the trivial divisor (which we denote by 0 ) is the identity element. It is easy to verify that the notion of positivity or effectiveness of divisors induces a partial ordering on the divisor group; namely, we write $D \geq D^{\prime}$ if $D-D^{\prime}$ is positive (thus a divisor $D$ is positive if and only if $D \geq 0$ ).

We denote the field of meromorphic functions on $Y$ (resp. the set of principal divisors on $Y$ ) by $\mathcal{M}(Y)$ (resp. by $\operatorname{Prin}(Y)$ ). The following fact is obvious from the definitions.

Proposition 2.2.2. The set of principal divisors on a complex manifold $Y$ forms a subgroup of $\operatorname{Div}(Y)$ and is the image of the homomorphism $\mathcal{M}(Y)^{\times} \rightarrow \operatorname{Div}(Y)$ defined by sending any nonzero meromorphic function $f$ to the divisor represented by $(\{Y\},\{f\})$. The kernel of this homomorphism is the subgroup of nowhere-vanishing holomorphic functions on $Y$.

There is a very straightforward way to define pullbacks of divisors via surjective holomorphic maps $f: Y^{\prime} \rightarrow Y$; namely, given any divisor $D \in \operatorname{Div}(Y)$ represented by $\left(\left\{U_{i}\right\}_{i \in I},\left\{f_{i}\right\}_{i \in I}\right)$, we let $f^{*}(D) \in \operatorname{Div}\left(Y^{\prime}\right)$ be given by $\left(\left\{f^{-1}\left(U_{i}\right)\right\}_{i \in I},\left\{f_{i} \circ f\right\}_{i \in I}\right)$. In this way, $f: Y^{\prime} \rightarrow Y$ induces a homomorphism of groups $f^{*}: \operatorname{Div}(Y) \rightarrow \operatorname{Div}\left(Y^{\prime}\right)$.

In our quest to characterize the complex tori which can be embedded into projective space as algebraic varieties, it will be necessary to study positive divisors on them. It would be nice
to be able to characterize positive divisors on a complex torus $X$ as principal divisors, since then our investigation would boil down to considering the set of holomorphic functions on $X$. But unfortunately, since such an $X$ is compact, the only holomorphic functions defined everywhere on $X$ are the constant functions, so the only positive principal divisor in $\operatorname{Div}(X)$ is the trivial divisor 0 . However, we can pull back any divisor on $X$ via $\pi: V \rightarrow V / \Lambda=X$ to get a divisor on $V$, and it turns out that the resulting divisor is principal by the following classical result.

Theorem 2.2.3 (Cousins). Every divisor on the complex manifold $\mathbb{C}^{n}$ for any $n \geq 1$ is principal.

Definition 2.2.4. $A$ divisor on $V$ is said to be periodic if it lies in the image of $\pi^{*}$ : $\operatorname{Div}(X) \rightarrow \operatorname{Div}(V)$.

We now want to characterize all periodic divisors on $V$ by representing them with functions in $\mathcal{M}(V)$ that have nice properties. We observe first of all that if $f \in \mathcal{M}(V)$ has periodic divisor with respect to the lattice $\Lambda$, that means that for each $\lambda \in \Lambda, g_{\lambda}(v):=$ $f(v+\lambda) / f(v) \in \mathcal{M}(V)$ must be holomorphic and nonvanishing. (I like to say that such a function is "almost periodic" with respect to $\Lambda$; it is only actually periodic if $(f)=\pi^{*} D$ with $D \in \operatorname{Prin}(X)$.) Moreover, these functions $g_{\lambda}$ satisfy the compatibility condition that

$$
\begin{equation*}
g_{\lambda_{1}+\lambda_{2}}(v)=g_{\lambda_{1}}(v) g_{\lambda_{2}}\left(v+\lambda_{1}\right), \quad \forall \lambda_{1}, \lambda_{2} \in \Lambda \tag{2.1}
\end{equation*}
$$

Moreover, since each $g_{\lambda}$ is holomorphic and nonvanishing, we may write $g_{\lambda}=e^{2 \pi i G_{\lambda}}$ for some holomorphic function $G_{\lambda} \in \mathcal{M}(V)$. Then the condition in (2.1) becomes

$$
\begin{equation*}
G_{\lambda_{1}+\lambda_{2}}(v) \equiv G_{\lambda_{1}}(v)+G_{\lambda_{2}}\left(v+\lambda_{1}\right) \quad \bmod \mathbb{Z}, \quad \forall \lambda_{1}, \lambda_{2} \in \Lambda \tag{2.2}
\end{equation*}
$$

Recall that a function $H: V \times V \rightarrow \mathbb{C}$ is a Hermitian form on the complex vector space $V$ if it is $\mathbb{C}$-linear in the first argument and if $H(w, v)=\overline{H(v, w)}$ for all $v, w \in V$. It is an easy exercise in linear algebra to show that if a function $H: V \times V \rightarrow \mathbb{C}$ is a Hermitian form, its imaginary part $E:=\Im H: V \times V \rightarrow \mathbb{R}$ is an $\mathbb{R}$-linear alternating form (i.e. $E(w, v)=-E(v, w)$ for all $v, w \in V)$ and if $E(i v, i w)=E(v, w)$ for all $v, w \in V$. Conversely, any such $E: V \times V \rightarrow \mathbb{R}$, determines a Hermitian form $H: V \times V \rightarrow \mathbb{C}$ by $H(v, w)=E(i v, w)+i E(v, w)$. In light of this, from now on, in any context where we have a Hermitian form $H$, we will use $E$ to denote its imaginary part, and conversely, whenever we have an alternating $\mathbb{R}$-bilinear form $E$, we will denote the corresponding Hermitian form by $H$.

Definition 2.2.5. A theta function (with respect to the lattice $\Lambda$ ) is a function $\theta \in \mathcal{M}(V)$ which satisfies the property that, for all $v \in V$ and $\lambda \in \Lambda$,

$$
\theta(v+\lambda) / \theta(v)=e^{2 \pi i(L(v, \lambda)+J(\lambda))}
$$

for some map $J: \Lambda \rightarrow \mathbb{C}$ and some map $L: V \times \Lambda \rightarrow \mathbb{C}$ which is $\mathbb{C}$-linear in the first argument. Given a Hermitian form $H: V \times V \rightarrow \mathbb{C}$, we say that such a function $\theta$ is a theta function for $H$ if we have $E\left(\lambda_{1}, \lambda_{2}\right)=L\left(\lambda_{1}, \lambda_{2}\right)-L\left(\lambda_{2}, \lambda_{1}\right)$ for $\lambda_{1}, \lambda_{2} \in \Lambda$.

Note that the set of all theta functions on $V$ with respect to $\Lambda$ forms a group under multiplication, and that in fact, the map from this group to the group of divisors $\operatorname{Div}(X)$ defined by sending a theta function $\theta$ to the divisor $D \in \operatorname{Div}(X)$ such that $(\theta)=\pi^{*} D$ is a group homomorphism which we denote by

$$
\text { div : }\{\text { theta functions on } V \text { w.r.t. } \Lambda\} \longrightarrow \operatorname{Div}(X) \text {. }
$$

Its kernel is the subgroup of trivial theta divisors, i.e. theta divisors $\theta$ which are holomorphic and nonvanishing on $V$.

We will give an incomplete proof of the following result, as the only full argument I know involves techniques of sheaf cohomology which are beyond the scope of this course (see the arguments in [5, §2]).

Proposition 2.2.6. For every divisor $D \in \operatorname{Div}(X)$, there is a unique Hermitian form $H$ and a function $\theta \in \mathcal{M}(V)$ with $(\theta)=\pi^{*} D$ which is a theta function for $H$.

Proof. We omit the proof of the fact that there always exists a function $\theta \in \mathcal{M}(V)$ with $(\theta)=\pi^{*} D$ such that

$$
\theta(v+\lambda) / \theta(v)=e^{2 \pi i(L(v, \lambda)+J(\lambda))}, \quad \forall v \in V, \lambda \in \Lambda
$$

with $L: V \times \Lambda \rightarrow \mathbb{C}$ linear in the first argument, or equivalently, the map div : \{theta functions $\} \rightarrow$ $\operatorname{Div}(X)$ is surjective. Now for any $\lambda_{1}, \lambda_{2} \in \Lambda$, condition (2.2) can be written as

$$
\begin{equation*}
L\left(v, \lambda_{1}+\lambda_{2}\right)+J\left(\lambda_{1}+\lambda_{2}\right) \equiv L\left(v, \lambda_{1}\right)+L\left(v, \lambda_{2}\right)+L\left(\lambda_{1}, \lambda_{2}\right)+J\left(\lambda_{1}\right)+J\left(\lambda_{2}\right) \bmod \mathbb{Z} \tag{2.3}
\end{equation*}
$$

which forces

$$
\begin{equation*}
L\left(v, \lambda_{1}+\lambda_{2}\right)-L\left(v, \lambda_{1}\right)-L\left(v, \lambda_{2}\right)=0 \tag{2.4}
\end{equation*}
$$

(here we get equality instead of equivalence modulo $\mathbb{Z}$ because it holds when $v=0$ and therefore holds for all $v$ by continuity) and

$$
\begin{equation*}
J\left(\lambda_{1}+\lambda_{2}\right)-J\left(\lambda_{1}\right)-J\left(\lambda_{2}\right) \equiv L\left(\lambda_{1}, \lambda_{2}\right)(\bmod \mathbb{Z}) \tag{2.5}
\end{equation*}
$$

Moreover, by switching the roles of $\lambda_{1}$ and $\lambda_{2}$, we get $L\left(\lambda_{1}, \lambda_{2}\right) \equiv L\left(\lambda_{2}, \lambda_{1}\right)(\bmod \mathbb{Z})$. In particular, $L$ is $\mathbb{Z}$-linear in the second argument, and we may extend $L$ to a map $V \times V \rightarrow \mathbb{C}$ which is $\mathbb{R}$-bilinear and symmetric modulo $\mathbb{Z}$. Set $E(v, w)=L(v, w)-L(w, v)$ for all $v, w \in V$. Then $E: V \times V \rightarrow \mathbb{C}$ is clearly $\mathbb{R}$-bilinear and skew-symmetric; moreover, $E(\Lambda, \Lambda) \subseteq \mathbb{Z}$, and so $E$ takes values in $\mathbb{R}$. It is also a tricky but elementary exercise to check that $E(i v, i w)=E(v, w)$ for all $v, w \in V$. Therefore, in particular $E$ is the imaginary part of the Hermitian form $H: V \times V \rightarrow \mathbb{C}$ given by $H(v, w)=E(i v, w)+i E(v, w)$.

We now have a well-defined map from the multiplicative group of theta divisors on $V$ (with repsect to $\Lambda$ ) to the additive group of Hermitian forms on $V$ which is clearly a group homomorphism; we denote it by

$$
\mathcal{H}:\{\text { theta functions on } V \text { w.r.t. } \Lambda\} \longrightarrow\{\text { Hermitian forms on } V \text { w.r.t. } \Lambda\} .
$$

We will now show that its kernel contains the subgroup of trivial theta functions. It is easy to see that if $h \in \mathcal{M}(V)$ is any trivial theta function, then $h$ must be of the form
$e^{2 \pi i(A(v)+B(v)+C)}$, where $A: V \rightarrow \mathbb{C}$ is a quadratic form, $B: V \rightarrow \mathbb{C}$ is a linear functional, and $C \in \mathbb{C}$ is a constant. Then for any $\lambda \in \Lambda, h(v+\lambda) / h(v)$ can be written as

$$
e^{2 \pi i([A(v+\lambda)-A(v)-A(\lambda)]+[A(\lambda)+B(\lambda)])},
$$

with $A(v+\lambda)-A(v)-A(\lambda)$ symmetric and bilinear. Using this, one can check that the Hermitian form $H$ corresponding to such an $h$ is the zero form.

It follows that the map $\mathcal{H}$ factors through a map

$$
\{\text { theta functions }\} /\{\text { trivial theta functions }\} \cong \operatorname{Div}(X) \longrightarrow\{\text { Hermitian forms }\},
$$

thus showing that each divisor $D \in \operatorname{Div}(X)$ uniquely determines a Hermitian form $H$.

Note that it was shown in the above proof that if $\theta$ is a theta function for some Hermitian form $H$, the imaginary part $E$ of $H$ must be $\mathbb{Z}$-valued on $\Lambda \times \Lambda$. This motivates the following crucial definition.

Definition 2.2.7. A Riemann form associated to a complex torus $X \cong V / \Lambda$ is a Hermitian form on $V$ whose imaginary part is $\mathbb{Z}$-valued on $\Lambda \times \Lambda$.

Proposition 2.2 .6 tells us in particular that the image of the map $\mathcal{H}$ defined in its proof is contained in the additive group of Riemann forms, so that we have group homomorphisms $\{$ theta functions w.r.t. $\Lambda\} /\{$ trivial theta functions $\} \xrightarrow{\sim} \operatorname{Div}(X) \rightarrow\{$ Riemann forms w.r.t. $\Lambda\}$.

Combining Lemma 2.3.4(b) below with Exercise 2.7.1, we see that the second map is a surjection. We shall describe its kernel in $\$ 2.5$.

Proposition 2.2.8. Let $H: V \times V \rightarrow \mathbb{C}$ be the Riemann form associated to some divisor $D \in \operatorname{Div}(X)$ via Proposition 2.2.6. Then there exists a theta function for $H$, and any theta function $\theta \in \mathcal{M}(V)$ for $H$ can be written as $h \theta_{D}$, where $h$ is a trivial theta function (i.e. $(h)=0 \in \operatorname{Div}(V))$ and $\theta_{D} \in \mathcal{M}(V)$ is a theta function satisfying

$$
\theta_{D}(v+\lambda) / \theta_{D}(v)=e^{2 \pi i\left(\frac{1}{2 i} H(v, \lambda)+\frac{1}{4 i} \pi H(\lambda, \lambda)+K(\lambda)\right)}
$$

where $K: \Lambda \rightarrow \mathbb{R}$ satisfies the property that

$$
K\left(\lambda_{1}+\lambda_{2}\right)-K\left(\lambda_{1}\right)-K\left(\lambda_{2}\right) \equiv \frac{1}{2} E\left(\lambda_{1}, \lambda_{2}\right) \quad \bmod \mathbb{Z}, \quad \forall \lambda_{1}, \lambda_{2} \in \Lambda
$$

Moreover, $\theta_{D}$ is unique up to multiplication by a constant in $\mathbb{C}^{\times}$.
Proof. We first have to show that a function $\theta_{D}$ with the property given in the statement is a theta function for $H$. This follows from verifying that $L_{0}(v, \lambda):=\frac{1}{2 i} H(v, \lambda)$ satisfies $L_{0}\left(\lambda_{1}, \lambda_{2}\right)-L_{0}\left(\lambda_{2}-\lambda_{1}\right)=\frac{1}{2 i}\left(H\left(\lambda_{1}, \lambda_{2}\right)-\overline{H\left(\lambda_{2}, \lambda_{1}\right)}\right)=E\left(\lambda_{1}, \lambda_{2}\right)$.

Now we show that such a function $\theta_{D}$ exists. We know from Proposition 2.2.6 that there exists a theta function $\theta$ satisfying $\theta(v+\lambda) / \theta(v)=e^{2 \pi i(L(v, \lambda)+J(\lambda))}$ where $L$ is $\mathbb{C}$-linear in the first variable and satisfies $L\left(\lambda_{1}, \lambda_{2}\right)-L\left(\lambda_{2}, \lambda_{1}\right)=E\left(\lambda_{1}, \lambda_{2}\right)$. Then $M:=L-L_{0}$ restricted
to $\Lambda \times \Lambda$ is a symmetric $\mathbb{Z}$-bilinear form which can be extended $\mathbb{R}$-linearly to a symmetric $\mathbb{R}$-bilinear form $M: V \times V \rightarrow \mathbb{C}$ (actually $\mathbb{C}$-bilinear because $L$ and $L_{0}$ are $\mathbb{C}$-bilinear in the first argument). Then $h_{1}(v):=e^{2 \pi i\left(\frac{1}{2} L(v, v)\right)}$ is clearly a trivial theta function (see the proof of Proposition 2.2.6), and replacing $\theta$ by $\theta / h_{1}(v)$ allows us to replace $L$ by $L_{0}=\frac{1}{2 i} H$ in our expression for $\theta(v+\lambda) / \theta(v)$. Now the relation (2.5) says that $J\left(\lambda_{1}+\lambda_{2}\right)-J\left(\lambda_{1}\right)-J\left(\lambda_{2}\right) \equiv$ $\frac{1}{2 i} H\left(\lambda_{1}, \lambda_{2}\right)(\bmod \mathbb{Z})$ for $\lambda_{1}, \lambda_{2} \in \Lambda$. Setting $K(\lambda)=J(\lambda)-\frac{1}{4 i} H(\lambda, \lambda)$, we get

$$
\begin{equation*}
K\left(\lambda_{1}+\lambda_{2}\right)-K\left(\lambda_{1}\right)-K\left(\lambda_{2}\right) \equiv \frac{1}{2 i} H\left(\lambda_{1}, \lambda_{2}\right)-\frac{1}{2 i} \mathfrak{R} H\left(\lambda_{1}, \lambda_{2}\right)=\frac{1}{2 i} E\left(i \lambda_{1}, \lambda_{2}\right)(\bmod \mathbb{Z}) \tag{2.6}
\end{equation*}
$$

Since $E$ takes values in $\mathbb{R}$, the imaginary part $\Im K$ of $K$ is additive on $\Lambda$ and can therefore be extended to an $\mathbb{R}$-linear functional $\Im K: V \rightarrow \mathbb{C}$. Note that $\Im K(i v)+i \Im K(v)$ defines a $\mathbb{C}$-linear functional $V \rightarrow \mathbb{C}$ whose imaginary part is $\Im K$. Now by replacing $\theta$ with $\theta$ divided by the trivial theta function $h_{2}(v):=e^{\Im K(i v)+i \Im K(v)}$, we may assume that $K$ is takes values in $\mathbb{R}$ without affecting $K\left(\lambda_{1}+\lambda_{2}\right)-K\left(\lambda_{1}\right)-K\left(\lambda_{2}\right)$, thus fulfilling the properties given in the statement of the proposition.

The uniqueness of $\theta_{D}$ up to a constant now follows quickly from the fact that there is no nontrivial linear functional $B: V \rightarrow \mathbb{C}$ which takes values in $\mathbb{R}$.

We call a theta function satisfying the properties in the statement of the above proposition a normalized theta function for the divisor $D$.

### 2.3 Algebraization of complex tori

For any divisor $D \in \operatorname{Div}(X)$, we set

$$
\mathcal{L}(D)=\{f \in \mathcal{M}(X) \mid(f)+D \geq 0\} \cup\{0\} .
$$

(We may think of the divisor associated to the constant function 0 as greater than every other divisor.) Note that $\mathcal{L}(D)$ is a vector space.

For any divisor $D \in \operatorname{Div}(X)$, let $\theta_{0}$ be a theta function whose divisor is $\pi^{*} D$. We define $\mathcal{L}\left(\theta_{0}\right)$ denote the set of all holomorphic functions $\theta \in \mathcal{M}(V)$ which have the same "translation functions" as $\theta_{0}$ has, i.e. $\theta(v+\lambda) / \theta(v)=\theta_{0}(v+\lambda) / \theta_{0}(v)$ (by convention, $\left.0 \in \mathcal{L}\left(\theta_{0}\right)\right)$. It is clear that $\mathcal{L}\left(\theta_{0}\right)$ is also a vector space. In fact, it is easy to see from the definitions that we have an isomorphism $\mathcal{L}\left(\theta_{0}\right) \xrightarrow{\sim} \mathcal{L}(D)$ given by $\theta \mapsto \theta / \theta_{0}$.

Below we again let $\theta_{D}$ denote a normalized theta function for $D$ as given by Proposition 2.2.8, although it is clear that the statements hold when $\theta_{D}$ is replaced by any $\theta_{0}$ with the same translation functions.

Definition 2.3.1. A divisor $D \in \operatorname{Div}(X)$ is very ample if there exists a basis $\left\{\theta_{0}, \ldots, \theta_{m}\right\}$ of $\mathcal{L}\left(\theta_{D}\right)$ such that the map $\Theta: X \rightarrow \mathbb{P}_{\mathbb{C}}^{m}$ induced by the map $X \rightarrow \mathbb{P}_{\mathbb{C}}^{m}$ given by $v \mapsto$ $\left(\theta_{0}(v): \ldots: \theta_{m}(v)\right)$ is an embedding of $X$ into $\mathbb{P}_{\mathbb{C}}^{m}$, i.e. if we have the following:
i) $\Theta$ is well-defined; that is, we do not have $\theta_{0}(v)=\ldots=\theta_{m}(v)=0$ for any $v \in V$;
ii) $\Theta: X \rightarrow \mathbb{P}_{\mathbb{C}}^{m}$ is an injection; and
iii) The induced maps on tangent spaces $\Theta_{*}: T_{a} X \rightarrow T_{\Theta(a)} \mathbb{P}_{\mathbb{C}}^{m}$ are injections.
$A$ divisor $D \in \operatorname{Div}(X)$ is ample if there is an integer $n \geq 1$ such that $n D$ is very ample.

Our goal is to characterize ample divisors $D \in \operatorname{Div}(X)$ in terms of their corresponding Riemann forms. If we can find an ample divisor $D$, we can use generators of $\mathcal{L}(n D)$ for some integer $n \geq 1$ to embed $X$ into projective space over $\mathbb{C}$ and give $X$ the structure of a variety.

The following proposition shows that there is no chance of this unless the Riemann form corresponding to $D$ is positive definite.

Proposition 2.3.2. Let $D \in \operatorname{Div}(X)$ be a divisor, and let $H: V \times V \rightarrow \mathbb{C}$ be its corresponding Riemann form.
a) Let $W_{0}=\{w \in V \mid H(v, w)=0 \forall v \in V\}$ be the right kernel of the pairing $H$. Then any function $\theta \in \mathcal{L}\left(\theta_{D}\right)$ must factor through the quotient map $V \rightarrow V / W_{0}$. In particular, if $H$ is degenerate, then $\Theta$ is not one-to-one.
b) If $H$ is not positive semidefinite (i.e. if $H(v, v)<0$ for some $v \in V$ ), then we have $\mathcal{L}\left(\theta_{D}\right)=\{0\}$ and so $\Theta$ is not defined.

Proof. Choose any $\theta \in \mathcal{L}\left(\theta_{D}\right)$. We first note that since $H$ is trivial on $V \times\left(W_{0} \cap \Lambda\right)$, we have $\theta(w+\lambda) / \theta(w)=e^{2 \pi i K(\lambda)}$ for all $\lambda \in W_{0} \cap \Lambda$. Since $K$ takes values in $\mathbb{R},\left|e^{2 \pi i K(\lambda)}\right|=1$. Let $C_{0} \subset V$ be a compact subset such that $C_{0}+\left(W_{0} \cap \Lambda\right)=W_{0}$. For any fixed $v \in V$, we then have $\max _{w \in W}\{|\theta(v+w)|\}=\max _{w^{\prime} \in C_{0}}\left\{\left|\theta\left(v+w^{\prime}\right)\right|\right\}$, which is finite because $C_{0}$ is compact. Then by the maximum principal for holomorphic functions on $W_{0}$, we see that $\theta(v+w)$ must be constant as a function of $w$, proving statement (a).

Now suppose that $H$ is not positive semidefinite. Let $W \subseteq V$ be a nontrivial subspace such that $H(w, w)<0$ for all nonzero $w \in W$. Let $C \subset V$ be a compact subset such that $C+\Lambda=V$. Fix any $v \in V$ and $\theta \in \mathcal{L}\left(\theta_{D}\right)$, consider the function on $W$ given by $w \mapsto \theta(v+w)$. If we write $w=w^{\prime}+\lambda$ for $w^{\prime} \in C$ and $\lambda \in \Lambda$, then we get

$$
\begin{equation*}
|\theta(v+w)|=\left|\theta\left(v+w^{\prime}+\lambda\right)\right|=\left|\theta\left(v+w^{\prime}\right)\right| e^{\pi \Re H\left(v+w^{\prime}, \lambda\right)+\frac{1}{2} \pi H(\lambda, \lambda)} \tag{2.7}
\end{equation*}
$$

It is straightforward to compute that

$$
\begin{equation*}
\Re H\left(v+w^{\prime}, \lambda\right)+\frac{1}{2} H(\lambda, \lambda)=\frac{1}{2} H(w, w)+\Re H(v, w)-\Re H\left(v, w^{\prime}\right)-\frac{1}{2} H\left(w^{\prime}, w^{\prime}\right) . \tag{2.8}
\end{equation*}
$$

Note that as $w \rightarrow \infty$, we have $H(w, w) \rightarrow-\infty$, which dominates the other terms on the right-hand side, because $\Re H(v, w)$ is linear in $w$ and the rest of the terms depend only on $w^{\prime}$ which varies over the compact subset $C$ and are therefore bounded. Thus, $|\theta(v+w)| \rightarrow 0$ as $w \rightarrow \infty$. Then by the maximum principal for holomorphic functions on $W, \theta(v+w)=0$ for all $w \in W$ and therefore $\theta=0$, proving statement (b).

The next result, known as the Lefschetz Embedding Theorem, says that the converse is also true.

Theorem 2.3.3. Let $H$ be a positive definite Riemann form. Then the divisor $D \in \operatorname{Div}(X)$ corresponding to $H$ is ample. More precisely, for any $n \geq 3$, the divisor $n D \in \operatorname{Div}(X)$ is very ample.

Before we can prove this theorem, we need a major lemma that is due to Frobenius. Recall that for any $\mathbb{R}$-bilinear alternating form $E: V \times V \rightarrow \mathbb{R}$, the Pfaffian of $E$ is defined to be
the nonnegative real number $\sqrt{\operatorname{det}(E)}$, where $\operatorname{det}(E)$ is understood to be the determinant of the matrix $\left(E\left(\lambda_{i}, \lambda_{j}\right)\right)$ for some (any) $\mathbb{Z}$-basis $\lambda_{1}, \ldots, \lambda_{2 g}$ of $\Lambda$. We observe that $\operatorname{det}(E)$ is a positive integer if $E(\Lambda, \Lambda) \subseteq \mathbb{Z}$.

Lemma 2.3.4. Let $\Lambda \subset V$ be as above, and let $E: \Lambda \times \Lambda \rightarrow \mathbb{Z}$ be a nondegerate alternating form.
(a) There is a basis $\left\{\lambda_{1}, \ldots, \lambda_{2 g}\right\}$ of $\Lambda$ and positive integers $e_{1}, \ldots, e_{g} \in \mathbb{Z}$ with $e_{i} \mid e_{i+1}$ for $1 \leq i \leq g-1$, such that $E\left(\lambda_{i}, \lambda_{i+g}\right)=-e_{i}$ for $1 \leq i \leq g$ and $E\left(\lambda_{i}, \lambda_{j}\right)=0$ otherwise for $1 \leq i, j \leq 2 g$.
(b) Suppose that $K: \Lambda \rightarrow \mathbb{R}$ is a function satisfying $K\left(\lambda_{1}+\lambda_{2}\right)-K\left(\lambda_{1}\right)-K\left(\lambda_{2}\right) \equiv$ $\frac{1}{2} E\left(\lambda_{1}, \lambda_{2}\right)(\bmod \mathbb{Z})$. Then the $\mathbb{C}$-vector space of theta functions $\theta$ on $V$ whose translation functions are given by $\theta(v+\lambda) / \theta(v)=e^{2 \pi i\left(\frac{1}{2 i} H(v, \lambda)+\frac{1}{4 i} \pi H(\lambda, \lambda)+K(\lambda)\right)}$ has dimension equal to $\sqrt{\operatorname{det}(E)}=e_{1} \ldots e_{g}$. (In particular, there exists such a theta function which is nontrivial.)

For reasons of time, we omit the proof of the above lemma. The following corollary is immediate from the lemma and will be useful later.

Corollary 2.3.5. If $D \in \operatorname{Div}(X)$ is a divisor whose associated Riemann form $H$ is positive definite with imaginary part $E$, then for any integer $n \geq 1$, the dimension of the complex vector space $\mathcal{L}(n D)$ is $n^{g} \sqrt{\operatorname{det}(E)}$.

Proof (of Theorem 2.3.3). We will prove this for $n=3$; the argument for greater $n$ is similar.
To prove show that property (i) in Definition 2.3.1 holds for $3 D$, we observe that given any $\theta \in \mathcal{L}\left(\theta_{D}\right)$ and any choice of $a, b \in V$, the meromorphic function $\theta_{a, b} \in \mathcal{L}(V)$ defined by $\theta_{a, b}(v)=\theta(v+a+b) \theta(v-a) \theta(v-b)$ satisfies

$$
\begin{gather*}
\theta_{a, b}(v+\lambda) / \theta_{a, b}(v)=e^{2 \pi i\left(\frac{1}{2 i} H(v+a+b, \lambda)+\frac{1}{2 i} H(v-a, \lambda)+\frac{1}{2 i} H(v-b, \lambda)+\frac{3}{4 i} H(\lambda, \lambda)+3 K(\lambda)\right)} \\
=e^{2 \pi i\left(\frac{1}{2 i} 3 H(v, \lambda)+\frac{1}{4 i} 3 H(\lambda, \lambda)+3 K(\lambda)\right)}=(\theta(v+\lambda) / \theta(v))^{3} \tag{2.9}
\end{gather*}
$$

for $v \in V, \lambda \in \Lambda$. So $\theta_{a, b} \in \mathcal{L}(3 H)$, and, given any $v_{0} \in V$, by choosing appropriate $a, b \in V$ we can ensure that $\theta_{a, b}\left(v_{0}\right) \neq 0$ and therefore not all basis elements of $\mathcal{L}\left(\theta_{3 D}\right)$ vanish at $v_{0}$.

To prove that 2.3 .1 (ii) holds for $3 D$, it suffices to show that for any $\theta \in \mathcal{L}\left(\theta_{D}\right)$, for $v_{1}, v_{2} \in V$, if $\theta\left(v_{1}+a+b\right) \theta\left(v_{1}-a\right) \theta\left(v_{1}-b\right)$ is a constant multiple of $\theta\left(v_{2}+a+b\right) \theta\left(v_{2}-a\right) \theta\left(v_{2}-b\right)$ for all $a, b \in V$, then $v_{1}-v_{2} \in \Lambda$. Taking logarithmic differentials with respect to $a$, we get

$$
\begin{equation*}
\theta\left(v_{1}+a+b\right)^{-1} d \theta-\theta\left(v_{1}-a\right)^{-1} d \theta=\theta\left(v_{2}+a+b\right)^{-1} d \theta-\theta\left(v_{2}-a\right)^{-1} d \theta \tag{2.10}
\end{equation*}
$$

for any $b \in V$. It follows that the differential given by $v \mapsto \theta\left(v_{2}+v\right)^{-1} d \theta-\theta\left(v_{1}+v\right)^{-1} d \theta$ is translation-invariant (after renaming the variable $a$ as $v$ and considering that $b$ can be chosen arbitrarily) and is therefore equal to $d B(v)$ for some linear functional $B: V \rightarrow \mathbb{C}$. After integrating and exponentiating, we get

$$
\begin{equation*}
\theta\left(v+v_{2}\right) / \theta\left(v+v_{1}\right)=c e^{B(v)} \Longrightarrow \theta\left(v+\left(v_{2}-v_{1}\right)\right) / \theta(v)=c^{\prime} e^{B(v)} \tag{2.11}
\end{equation*}
$$

for some constants $c, c^{\prime} \in \mathbb{C}^{\times}$. Using the property given in Definition 2.2.5, we deduce from substituting $v+\lambda$ for $v$ in 2.11) and using our translation formula for $\theta$ that $e^{2 \pi i\left(\frac{1}{2 i} H\left(v_{2}-v_{1}, \lambda\right)\right.}=$ $e^{B(\lambda)}$ for all $\lambda \in \Lambda$. Thus, we have that $\pi H\left(v_{2}-v_{1}, \lambda\right) \equiv B(\lambda)(\bmod 2 \pi i \mathbb{Z})$. Now we have
$\pi H\left(v_{2}-v_{1}, \lambda\right)-B(\lambda)=\pi H\left(\lambda, v_{2}-v_{1}\right)+2 \pi i E\left(v_{1}-v_{2}, \lambda\right)-B(\lambda) \in 2 \pi i \mathbb{Z}$, and it follows from the fact that $E$ is real-valued that $\pi H\left(\lambda, v_{2}-v_{1}\right)-B(\lambda)$ takes purely imaginary values for $\lambda \in \Lambda$. Thus, $v \mapsto \pi H\left(v, v_{1}-v_{2}\right)-B(v)$ takes purely imaginary values, but this function is $\mathbb{C}$-linear, so we have $\pi H\left(v, v_{2}-v_{1}\right)-B(v)=0$ for $v \in V$. Then

$$
\begin{equation*}
\pi H\left(v_{2}-v_{1}, \lambda\right)-B(\lambda)=2 \pi i E\left(v_{2}-v_{1}, \lambda\right)+\pi H\left(\lambda, v_{2}-v_{1}\right)-B(\lambda)=2 \pi i E\left(v_{2}-v_{1}, \lambda\right) \in 2 \pi i \mathbb{Z} \tag{2.12}
\end{equation*}
$$

implying that $E\left(v_{2}-v_{1}, \lambda\right) \in \mathbb{Z}$ for all $\lambda \in \Lambda$. Thus, letting $\Lambda^{\prime}=\Lambda+\left(v_{2}-v_{1}\right) \mathbb{Z}$, we have $E\left(\Lambda^{\prime}, \Lambda^{\prime}\right) \subseteq \mathbb{Z}$. In turn, this implies (as another easy exercise) that the subgroup $\Lambda^{\prime} \subset V$ is another lattice containing $\Lambda$, or equivalently, that $\left[\Lambda^{\prime}: \Lambda\right]<\infty$.

Now (2.11) tells us that we have

$$
\begin{equation*}
\theta\left(v+\left(v_{2}-v_{1}\right)\right) / \theta(v)=c^{\prime} e^{\pi H\left(v, v_{2}-v_{1}\right)}=e^{2 \pi i\left(\frac{1}{2 i} H\left(v, v_{2}-v_{1}\right)+\frac{1}{4 i} H\left(v_{2}-v_{1}, v_{2}-v_{1}\right)+K^{\prime}\left(v_{2}-v_{1}\right)\right)} \tag{2.13}
\end{equation*}
$$

where $c^{\prime \prime} \in \mathbb{C}^{\times}$is a constant and $K^{\prime}: \Lambda^{\prime} \rightarrow \mathbb{C}$ is an extension of $K$ (not necessarily satisfying the perscribed properties for $K: \Lambda \rightarrow \mathbb{R}$ ). Therefore, $\theta$ is actually a theta function for $H$ with respect to the lattice $\Lambda^{\prime}$. We write $\operatorname{det}_{\Lambda^{\prime}}(E)$ for the determinant of the alternating form $E: \Lambda^{\prime} \times \Lambda^{\prime} \rightarrow \mathbb{Z}$. Note that by Lemma $2.3 .4(\mathrm{~b})$, the space $\mathcal{L}_{H, K^{\prime}}$ of theta functions $\theta^{\prime}$ with respect to $\Lambda^{\prime}$ for $H$ satisfying $\theta^{\prime}\left(v+\lambda^{\prime}\right) / \theta^{\prime}(v+\lambda)=e^{2 \pi i\left(\frac{1}{2 i} H\left(v, \lambda^{\prime}\right)+\frac{1}{4 i} H\left(\lambda^{\prime}, \lambda^{\prime}\right)+K^{\prime}\left(\lambda^{\prime}\right)\right)}$ for $\lambda^{\prime} \in \Lambda^{\prime}$ has dimension equal to $\sqrt{\operatorname{det}_{\Lambda^{\prime}}(E)}$.

Suppose that $\Lambda^{\prime} \supsetneq \Lambda$. Then one can show as an easy exercise in linear algebra that $\operatorname{det}(E)>\operatorname{det}_{\Lambda^{\prime}}(E)$. We have shown above that every theta function $\theta \in \mathcal{L}\left(\theta_{D}\right)$ belongs to the space of theta functions with respect to $\Lambda^{\prime}$ whose translation functions are as in Definition 2.2 .5 with $L(v, \lambda)=\frac{1}{2 i} H(v, \lambda)$ and $J^{\prime}(\lambda)=\frac{1}{4 i} H(\lambda, \lambda)+K(\lambda)$ for some $K^{\prime}: \Lambda^{\prime} \rightarrow \mathbb{C}$ extending $K$. There are only finitely many choices for $K^{\prime}$ since $\left[\Lambda^{\prime}: \Lambda\right]<\infty$ and the relation in (2.5) must be satisfied, and each space $\mathcal{L}_{H, K^{\prime}}$ has dimension $\sqrt{\operatorname{det}_{\Lambda^{\prime}}(E)}<\sqrt{\operatorname{det}(E)}$. We therefore a contradiction. Thus, $\Lambda^{\prime}=\Lambda$, implying that $v_{2}-v_{1} \in \Lambda$, as desired.

Finally we show that 2.3 .1 (iii) holds for $3 D$. After fixing a basis for $V$, we write $z_{1}, \ldots, z_{g}$ for the corresponding set of coordinate functions on $V$. Let $\widetilde{\Theta}: V \rightarrow \mathbb{P}_{\mathbb{C}}$ be the composition with $\Theta$ as defined in Definition 2.3 .1 with the covering map $V \rightarrow X$. Then it obviously suffices to show that at every point $v_{0} \in V$, the induced map $\widetilde{\Theta}_{*}: T_{v_{0}} V \rightarrow T_{\widetilde{\Theta}\left(v_{0}\right)} \mathbb{P}_{\mathbb{C}}^{m}$ is injective. Choose any point $v_{0} \in V$, and let $\sum_{i=1}^{g} \alpha_{i} \frac{\partial}{\partial z_{i}} \in T_{v_{0}} V$ be a tangent vector whose image under $\widetilde{\Theta}_{*}$ vanishes. If we write $\left\{\theta_{0}, \ldots, \theta_{m}\right\}$ for a basis of $\mathcal{L}\left(\theta_{3 D}\right)$ and assume without loss of generality that $\theta_{0}\left(z_{0}\right) \neq 0$, we then have

$$
\begin{equation*}
\sum_{i=1}^{g} \alpha_{i} \frac{\partial\left(\phi / \theta_{0}\right)}{\partial z_{i}}\left(v_{0}\right)=0 \tag{2.14}
\end{equation*}
$$

for every $\phi \in \mathcal{L}\left(\theta_{3 D}\right)$. Suppose that this choice of tangent vector is nonzero; we may assume without loss of generality that $\alpha_{1} \neq 0$. Then we have (using the quotient rule for differentiation)

$$
\begin{align*}
& \frac{\alpha_{1}}{\theta_{0}\left(v_{0}\right)^{2}}\left(\theta_{0}\left(v_{0}\right) \frac{\partial \phi}{\partial z_{1}}\left(v_{0}\right)-\phi\left(v_{0}\right) \frac{\partial \theta_{0}}{\partial z_{1}}\left(v_{0}\right)\right)=0 \\
\Longrightarrow & \phi\left(v_{0}\right)^{-1} \frac{\partial \phi}{\partial z_{1}}\left(v_{0}\right)=\theta_{0}\left(v_{0}\right)^{-1} \frac{d \theta_{0}(v)}{\partial z_{1}}\left(v_{0}\right)=: c \tag{2.15}
\end{align*}
$$

for every $\phi \in \mathcal{L}\left(\theta_{3 D}\right)$. Now choose $\theta \in \mathcal{L}\left(\theta_{D}\right)$ and define $\theta_{a, b}(v) \in \mathcal{L}\left(\theta_{3 D}\right)$ as before for any $a, b \in V$. Putting $\theta_{a, b}$ for $\phi$ in (2.15) gives us

$$
\begin{equation*}
\theta\left(v_{0}-a\right)^{-1} \frac{\partial \theta}{\partial z_{1}}\left(v_{0}-a\right)+\theta\left(v_{0}-b\right)^{-1} \frac{\partial \theta}{\partial z_{1}}\left(v_{0}-b\right)+\theta\left(v_{0}+a+b\right)^{-1} \frac{\partial \theta}{\partial z_{1}}\left(v_{0}+a+b\right)=c . \tag{2.16}
\end{equation*}
$$

By treating the expression in 2.16) as a (constant) function in the variable $a$ (which we rename as $v$ ) and taking partial derivatives with respect to $a$, we see that

$$
\begin{equation*}
\frac{\partial}{\partial z_{i}}\left(\left(v_{0}-v\right)^{-1} \frac{\partial \theta}{\partial z_{1}}\left(v_{0}-a\right)\right)=\frac{\partial}{\partial z_{i}}\left(\theta\left(v_{0}+a+b\right)^{-1} \frac{\partial \theta}{\partial z_{1}}\left(v_{0}+a+b\right)\right) \tag{2.17}
\end{equation*}
$$

for $1 \leq i \leq g$ and for any choice of $b$. In particular, it follows that $\theta\left(v_{0}+v\right)^{-1} \frac{\partial \theta}{\partial z_{1}}\left(v_{0}+v\right)$ is linear (but not necessarily homogeneous) in $v$; we write it as $\mu z_{1}+\nu(v)$ for some constant $\mu \in \mathbb{C}$ and linear function $\nu$ not depending on $z_{1}$. By integration and exponentiating on both sides of the equation $\theta\left(v_{0}+v\right)^{-1} \frac{\partial \theta}{\partial z_{1}}\left(v_{0}+v\right)=\mu z_{1}+\nu(v)$ with respect to $z_{1}$, we get that $\theta\left(v_{0}+z_{1}\right)=e^{\frac{1}{2} \mu z_{1}^{2}+\nu\left(v_{0}\right) z_{1}} \theta\left(v_{0}\right)$. Let $e_{1} \in V$ be the unit basis vector in the direction of $z_{1}$. Then we have

$$
\begin{equation*}
\theta\left(v_{0}+d e_{1}\right) / \theta\left(v_{0}\right)=e^{\frac{1}{2} \mu d^{2}+\nu\left(v_{0}\right) d} \tag{2.18}
\end{equation*}
$$

Similarly to our argument above, we deduce from substituting $v+\lambda$ for $v$ in (2.18) and using our translation formula for $\theta$ that $\pi H\left(z+d e_{1}, \lambda\right) \equiv d\left(\nu\left(v_{0}+\lambda\right)-\nu\left(v_{0}\right)\right)(\bmod 2 \pi i \mathbb{Z})$. Since $\nu$ is linear (possibly non-homogeneous), the map $\lambda \mapsto d\left(\nu\left(v_{0}+\lambda\right)-\nu\left(v_{0}\right)\right)$ defines a functional on $V$. Now we see using the same arguments as were used above to show that $E\left(d e_{1}, \Lambda\right) \subseteq \mathbb{Z}$ for any $d \in \mathbb{C}$. But this contradicts the nondegeneracy of $H$ and thus also our assumption that the tangent vector $\sum_{i=1}^{g} \alpha_{i} \frac{\partial}{\partial z_{i}} \in T_{v_{0}} V$ is nonzero.

We have now more or less shown the main theorem of this entire chapter, which gives a criterion for any complex torus $X$ to be algebraic. (Note that this in turn is equivalent to $X$ being an abelian variety, because a well-known result of Chow shows that the multiplication and inverse maps associated to $X$ are in fact morphisms due to the compactness of $X$.) We may think of it as a sort of "main theorem of complex tori" or "main theorem of complex abelian varieties" although, to my knowledge, it is not given such a name anywhere in the literature.

Theorem 2.3.6. A complex torus $X \cong V / \Lambda$ has the structure of an abelian variety if and only if there exists a positive definite Riemann form $H: V \times V \rightarrow \mathbb{C}$ associated to $X$.

It turns out to be the case that for $g \geq 2$, "most" complex tori of dimension $g$ do not possess any positive definite Riemann form and are therefore not abelian varieties (we will not give a proof of this here). However, as we will see immediately at the start of the next section, every 1-dimensional complex torus has a positive definite Riemann form and is therefore an elliptic curve.

### 2.4 Uniformization of elliptic curves

In this subsection we deal only with a complex torus $X$ of dimension $g=1$, which can always be realized as an elliptic curve over $\mathbb{C}$ because it always has a positive definite Riemann
form. Indeed, let $\left\{\lambda_{1}, \lambda_{2}\right\}$ any ordered basis of $\Lambda$ ordered so that $\Im\left(\lambda_{2} / \lambda_{1}\right)>0$, and define $E: \mathbb{C} \times \mathbb{C} \rightarrow \mathbb{R}$ to be the unique alternating $\mathbb{R}$-bilinear form such that $E\left(\lambda_{1}, \lambda_{2}\right)=-1$. (In fact, $E$ does not depend on the choice of ordered basis with this property - see Exercise 2.7.2. ) Then clearly $E(\Lambda, \Lambda)=\mathbb{Z}$ and it is easy to check that $E(i z, z)>0$ for any nonzero complex number $z$, and so $E$ is the imaginary part of a positive definite Riemann form $H: \mathbb{C} \times \mathbb{C} \rightarrow \mathbb{C}$. It follows from Theorem 2.3 .6 and more specifically Theorem 2.3 .3 that any divisor $D \in \operatorname{Div}(X)$ whose associated Riemann form is $H$ is ample, and in fact that $3 D$ is very ample so that a basis of $\mathcal{L}(3 D)$ can be used to embed $X$ into some projective space (this is called "uniformization"). Our goal now is to make this more explicit.

It will now be more convenient to switch from viewing elements of $\operatorname{Div}(X)$ as Cartier divisors to viewing them as Weil divisors - that is, each divisor $D \in \operatorname{Div}(X)$ is a finite formal sum of points in $X$. Since $X$ is a smooth curve, the notions of Cartier divisor and Weil divisor are equivalent. In particular, looking at Weil divisors allows us to use the notion of "degree": if an element $D \in \operatorname{Div}(X)$ can be written as $D=\sum_{i=1}^{m} n_{i}\left(P_{i}\right)$ for integers $n_{i} \in \mathbb{Z}$ and points $P_{i} \in X$, its degree is $\operatorname{deg}(D)=\sum_{i=1}^{m} n_{i}$. In this way, deg $: \operatorname{Div}(X) \rightarrow \mathbb{Z}$ is a homomorphism which preserves the partial ordering on $\operatorname{Div}(X)$.

Proposition 2.4.1. There exist meromorphic functions $x, y \in \mathcal{M}(X)$ satisfying a cubic polynomial equation $f(x, y)=0$ such that $X$ can be identified with the closure in $\mathbb{P}_{\mathbb{C}}^{2}$ of the curve defined by the relation $f(x, y)=0$.

Proof. Let $D \in \operatorname{Div}(X)$ be a divisor whose associated Riemann form is the $H$ defined with respect to an ordered basis $\left\{\lambda_{1}, \lambda_{2}\right\} \subset \Lambda$ as above. We note that the $2 \times 2$ matrix given by $\left(E\left(\lambda_{i}, \lambda_{j}\right)\right)$ has determinant 1 , so $\sqrt{\operatorname{det}(E)}=1$, and so Corollary 2.3.5 says that $\operatorname{dim}_{\mathbb{C}} \mathcal{L}(n D)=n$ for all $n \geq 1$.

We first observe that $\operatorname{deg}(D) \geq 0$. Indeed, the fact that $\mathcal{L}(D)$ is nontrivial implies that there exists a nonzero function $f_{0} \in \mathcal{M}(X)$ such that $\left(f_{0}\right)+D \geq 0$, which implies that $\operatorname{deg}\left(\left(f_{0}\right)\right) \geq-\operatorname{deg}(D)$. But every nonzero meromorphic function on a compact Riemann manifold has as an equal number of poles and zeros (counting multiplicity), so $\operatorname{deg}\left(\left(f_{0}\right)\right)=0$ and $\operatorname{deg}(D) \geq 0$. Now note that if we subtract any principal divisor from $D$, the associated Riemann form is still $H$, so in particular we may and will replace $D$ with $D-\left(f_{0}\right)$. Then we get $1 \in \mathcal{L}(D)$. Since $(1)=0 \in \operatorname{Div}(X)$, in fact we have $D \geq 0$. If $D=0$ then $\operatorname{dim}_{\mathbb{C}} \mathcal{L}(n D)=\operatorname{dim}_{\mathbb{C}} \mathcal{L}(0)=1$ for all $n \geq 1$, a contradiction, so $D>0$.

Since $\mathcal{L}(D)$ has dimension 1 and contains the constant functions, there are no nonconstant meromorphic functions on $X$ whose associated divisors are $\geq-D$.

Since $\mathcal{L}(2 D)$ has dimension 2 , it must be generated by $\{1, x\}$, where $x \in \mathcal{M}(X)$ is some nonconstant function such that $(x) \geq-2 D$.

Since $\mathcal{L}(3 D)$ has dimension 3 and $\langle 1, x\rangle=\mathcal{L}(2 D) \subsetneq \mathcal{L}(3 D)$, there is a function $y \in \mathcal{M}(X)$ such that $\{1, x, y\}$ is a basis of $\mathcal{L}(3 D)$ and $(y) \geq-3 D$. We claim that $y$ cannot be a polynomial function of $x$. Indeed, Theorem 2.3 .3 says that $3 D$ is very ample and so the function $X \rightarrow \mathbb{P}_{\mathbb{C}}^{2}$ given by $[1: x: y]$ is a projective embedding. Then if $y$ is a rational function of $x$, the injectivity of $[1: x: y]: X \hookrightarrow \mathbb{P}_{\mathbb{C}}^{2}$ implies that $[1: x]: X \rightarrow \mathbb{P}_{\mathbb{C}}^{1}$ is also injective. Note that the only possible images of $[1: x]$ are constant or the entire Riemann sphere since $X$ is compact and connected. But $X$ has nontrivial fundamental group, so this is a contradiction. Therefore, $y \notin \mathbb{C}(x)$. Moreover, $\mathcal{M}(X)$ is the fraction field of $\mathbb{C}[x, y]$ modulo some polynomial relation $f(x, y)=0$, or equivalently, $[1: x: y]$ realizes $X$ as the
closure in $\mathbb{P}_{\mathbb{C}}^{2}$ of the curve defined by $f(x, y)=0$. It remains only to show that $f \in \mathbb{C}[x, y]$ has degree 3 .

Since $\mathcal{L}(4 D)$ has dimension 4 and $\left(x^{2}\right)=2(x) \geq-4 D \in \operatorname{Div}(X)$, we get $\left\{1, x, y, x^{2}\right\}$ as a basis.

Since $\mathcal{L}(5 D)$ has dimension 5 and $(x y)=(x)+(y) \geq-2 D-3 D=-5 D \in \operatorname{Div}(X)$, we get $\left\{1, x, y, x^{2}, x y\right\}$ as a basis.

Since $\mathcal{L}(6 D)$ has dimension 6 , any set of 7 elements in $\mathcal{L}(6 D)$ is linearly dependent. But $1, x, y, x^{2}, x y, x^{3}, y^{2} \in \mathcal{L}(6 D)$ because $\left(x^{3}\right)=3(x) \geq-6 D \in \operatorname{Div}(X)$ and $\left(y^{2}\right)=2(y) \geq$ $-6 D \in \operatorname{Div}(X)$. Since $\left\{1, x, y, x^{2}, x y\right\}$ is linearly independent, there must be some relation among the functions $1, x, y, x^{2}, x y, x^{3}, y^{2}$ where the coefficients of $x^{3}$ and $y^{2}$ are both nonzero, so $x$ and $y$ satisfy a cubic relation $f(x, y)=0$.

In fact, with a little ingenuity it is possible to write down such meromorphic functions explicitly as functions $\mathbb{C} \rightarrow \mathbb{C} \cup\{\infty\}$ which are periodic with respect to a given lattice $\Lambda$. Consider the below function discovered by Weierstrass.

$$
\wp(z)=\frac{1}{z^{2}}+\sum_{\lambda \in \Lambda \backslash\{0\}}\left(\frac{1}{(z-\lambda)^{2}}-\frac{1}{\lambda^{2}}\right)
$$

This is called the "Weierstrass P-function". It is clear upon inspection that $\wp(z+\lambda)=\wp(z)$ for any $z \in \mathbb{C}$ and $\lambda \in \Lambda$, and so $\wp$ is periodic with respect to $\Lambda$. It is a little less trivial to see that the expression for $\wp(z)$ converges for every $z \in \mathbb{C} \backslash \Lambda$, and that in fact the function $\wp$ is meromorphic, with poles of order 2 at every element of $\Lambda$. Therefore, $\wp$ can also be considered as a function in $\mathcal{M}(X)$. We finally note that $\wp$ is an even function; i.e. $\wp(-z)=\wp(z)$ for all $z \in \mathbb{C}$. Therefore its derivative $\wp^{\prime}$ is an odd meromorphic function on $\mathbb{C}$ (i.e. $\wp^{\prime}(-z)=-\wp^{\prime}(z)$ for all $z \in \mathbb{C}$ ) which is also periodic with respect to $\Lambda$ (thus, it may also be considered as a function in $\mathcal{M}(X)$ ) and whose only poles are poles of order 3 at every element of $\Lambda$.

Proposition 2.4.2. The field $\mathcal{M}(X)$ of meromorphic functions on $X$ is given by $\mathbb{C}\left(\wp, \wp^{\prime}\right)$.
Proof. We want to show that every meromorphic function on $\mathbb{C}$ which is periodic with respect to the lattice $\Lambda$ is a rational function of $\wp$ and $\wp^{\prime}$. Now every meromorphic function is the sum of an odd function and an even function each periodic with respect to $\Lambda$ (easy exercise), and every such odd function divided by the odd function $\wp^{\prime}$ becomes an even function, so it suffices to show that every even nonzero meromorphic function $f$ on $\mathbb{C}$ which is periodic with respect to $\Lambda$ is a rational function of $\wp$.

Let $(f)=\sum_{P \in X} n_{P}(P) \in \operatorname{Div}(X)$ be the Weil divisor associated to $f$ considered as a function on $X$. We first note that for each $P \in X$, we have $n_{-P}=n_{P}$. Moreover, we claim that if $P \in X$ with $P=-P$, then $n_{P}$ is even. To prove this, let $z \in \mathbb{C}$ lie in the inverse image of such a $P$ under the map $\pi: \mathbb{C} \rightarrow \mathbb{C} / \Lambda=X($ so $2 z \in \Lambda)$, and let $f^{(n)}$ denote the $n$th derivative of $f$ for $n \geq 0$. Then we check that $f^{(i)}(z)=f^{(i)}(-z)=(-1)^{i-1} f^{(i)}(z)$, which implies that $f^{(i)}(z)=0$ for $i$ even. It follows that the order of $f$ at $z$ is even. Thus, the divisor of $f$ on $X$ is given by $\sum_{P \in X} n_{P}(P)+n_{P}(-P) \in \operatorname{Div}(X)$ with $n_{P} \in \mathbb{Z}$ almost all 0 .

Let $D \subset \mathbb{C}$ be a fundamental domain of the lattice $\Lambda$; that is, $D+\Lambda=\mathbb{C}$ and $(D+\lambda) \cap D=$ $\varnothing$ for any $\lambda \in \Lambda$ (note that there is an obvious bijection between $D$ and $X$ ). Let $D^{\prime} \subset D$ be a subset such that $\left(D^{\prime}+\Lambda\right) \cup\left(-D^{\prime}+\Lambda\right)=\mathbb{C}$ and $\left(D^{\prime}+\Lambda\right) \cap\left(-D^{\prime}+\Lambda\right)=\frac{1}{2} \Lambda$. Assume that $0 \in D^{\prime} \subset D$. Now consider the meromorphic function $g$ on $\mathbb{C}$ given by

$$
g(z)=\prod_{w \in D^{\prime} \backslash\{0\}}(\wp(z)-\wp(w))^{n_{\pi(w)}} .
$$

We check that for each $z \in D^{\prime} \backslash\{0\}$, the divisor associated to $\wp(z)-\wp(w)$ considered as a function on $X$ is $(P)+(-P)-2(0) \in \operatorname{Div}(X)$. It follows that for each $P \in X \backslash 0$, the functions $f$ and $g$ have the same order; thus, when considered as functions in $\mathcal{M}(X)$, they have the same associated divisor except possibly for the coefficient at $0 \in X$. But then they must have the same order at 0 as well, since $(f / g) \in \operatorname{Div}(X)$ is a principal divisor and must have degree 0 . So $f / g$ has no zeros or poles on $X$ and therefore must be a constant function since $X$ is compact. Thus, since $g \in \mathbb{C}(\wp)$, we have $f \in \mathbb{C}(\wp)$ as desired.

Corollary 2.4.3. a) The divisor $(0) \in \operatorname{Div}(X)$ is ample, and $3(0) \in \operatorname{Div}(X)$ is very ample (i.e. there is a basis of $\mathcal{L}(3(0))$ which can be used to embed $X$ into some projective space).
b) The imaginary part $E$ of the positive definite Riemann form associated to the divisor $(0) \in \operatorname{Div}(X)$ satisfies $\sqrt{\operatorname{det}(E)}=1$, so that it is given by $E\left(\lambda_{1}, \lambda_{2}\right)=-1$ for some (any) basis $\left\{\lambda_{1}, \lambda_{2}\right\}$ of $\Lambda$ with $\Im\left(\lambda_{2} / \lambda_{1}\right)>0$.
c) The set $\left\{1, \wp, \wp^{\prime}\right\}$ of meromorphic functions on $\mathbb{C}$ viewed as functions in $\mathcal{M}(X)$ is a basis of $\mathcal{L}(3(0))$, and the map $\left[1: \wp: \wp^{\prime}\right]: \mathbb{C} \rightarrow \mathbb{P}_{\mathbb{C}}^{2}$ induces an embedding of $X$ as the projective curve in $\mathbb{P}_{\mathbb{C}}^{2}$ defined by a cubic relation $f\left(\wp, \wp^{\prime}\right)=0$.
Proof. We already know that $X$ has the structure of an projective curve over $\mathbb{C}$, and Proposition 2.4.2 tells us that the function field of the curve $X$ is $\mathcal{M}(X)=\mathbb{C}\left(\wp, \wp^{\prime}\right)$. Since the function field of a complex curve must have transcendence degree 1 over $\mathbb{C}$ and $\wp$ and $\wp^{\prime}$ are obviously transcendental over $\mathbb{C}$, it is then clear that there must be some algebraic relation $f\left(\wp, \wp^{\prime}\right)=0$. This means that $X$ is the closure in projective space of the variety given by the relation $f(x, y)=0$, where $x$ and $y$ are the functions induced on $X=\mathbb{C} / \Lambda$ by $\wp$ and $\wp^{\prime}$ respectively. Thus, $\left[1: \wp: \wp^{\prime}\right]$ induces an embedding $X \hookrightarrow \mathbb{P}_{\mathbb{C}}^{2}$.

Now $x$ has a double pole at $0 \in X$ and no poles anywhere else while $y$ has a triple pole at $0 \in X$ and no poles anywhere else, so $x, y \in \mathcal{L}(3(0))$. Then $\{1, x, y\}$ can be extended to a basis of $\mathcal{L}(3(0))$ whose elements clearly induce an embedding of $X$ into some projective space, so the divisor $3(0) \in \operatorname{Div}(X)$ is very ample and $(0) \in \operatorname{Div}(X)$ is ample, implying (a).

To prove part (b), we only need to show that the imaginary part $E$ of the positive definite Riemann form associated to the ample divisor $(0) \in \operatorname{Div}(X)$ satisfies $\sqrt{\operatorname{det}(E)}=1$, since then it is easy to see that the only such $E$ is defined as given in the statement. Lemma 2.3.4 tells us that $\operatorname{dim}_{\mathbb{C}} \mathcal{L}((0))=\sqrt{\operatorname{det}(E)}$, and clearly every constant function lies in $\mathcal{L}((0))$, so it suffices to show that $\mathcal{L}((0))=\mathbb{C}$. If a function $g \in \mathcal{L}((0))$ has no pole at $0 \in X$, then it is holomorphic and therefore a constant function because $X$ is compact. Thus, we assume that there exists $g \in \mathcal{L}((0))$ with a simple pole at $0 \in X$ and no poles anywhere else. Then $g: X \rightarrow \mathbb{C} \cup\{\infty\}$ is actually a degree-1 morphism $X \rightarrow \mathbb{P}_{\mathbb{C}}^{1}$ and thus identifies $X$ with $\mathbb{P}_{\mathbb{C}}^{1}$ itself. This is a contradiction since $X$ has nontrivial fundamental group, so we get the desired statement.

Now part (c) results from the same argument as in the proof of Proposition 2.4.1(in fact, $x$ and $y$ play the same roles as in that proof).

Again using a little ingenuity, it is possible to write down a cubic polynomial equation relating $\wp$ and $\wp^{\prime}$, whose existence is guaranteed by Corollary 2.4.3. For any $k \in \mathbb{Z}$, define

$$
G_{2 k}=\sum_{\lambda \in \Lambda \backslash\{0\}} \lambda^{-2 k} .
$$

One can show that this sum converges for any $k \geq 2$. Then it is possible to show using Laurent approximations of $\wp$ and $\wp^{\prime}$ that the following relation holds:

$$
\begin{equation*}
f\left(\wp, \wp^{\prime}\right):=\wp^{\prime}(z)^{2}-4 \wp(z)^{3}+60 G_{4} \wp(z)+140 G_{6} \equiv 0 . \tag{2.19}
\end{equation*}
$$

This is remarkable in showing that in fact, $X$ can be viewed as the projective closure of a variety given by a cubic equation of the rather simple form $y^{2}=4 x^{3}+B x+C$ for constants $B, C \in \mathbb{C}$. One checks easily that the identity element $0 \in X$ corresponds to the unique point at infinity on this projective variety. We have now partially proven that Definition 1.1.1 implies Definition 1.1.4. That is, any complex abelian variety (in fact, any complex torus!) of dimension one can be defined algebraically by a cubic polynomial of the form given in Definition 1.1.4 (note that we can easily get rid of the " 4 " in the equation by scaling $x$ by $4^{1 / 3}$ to get the equation in exactly the right form). We have not yet seen that the group operation can be defined by the formula given in that definition, however.

### 2.5 Isogenies, polarizations, and duals

This subsection will provide only a brief introduction to the construction of duals in the world of complex abelian varieties which should lend some intuition for some of the algebraic results on elliptic curves that will come later.

### 2.5.1 Isogenies of complex abelian varieties

Definition 2.5.1. An isogeny of abelian varieties $X_{1}$ and $X_{2}$ of the same dimension is a surjective morphism $X_{1} \rightarrow X_{2}$ which is a homomorphism of groups.

In order to understand isogenies of complex abelian varieties, the following proposition is essential.

Proposition 2.5.2. For $i=1,2$, let $V_{i} / \Lambda_{i}$ be a complex torus, where $V_{i}$ is a complex vector space and $\Lambda_{i} \subset V_{i}$ is a lattice. Then any holomorphic homomorphism $\varphi: V_{1} / \Lambda_{1} \rightarrow V_{2} / \Lambda_{2}$ can be lifted to a unique $\mathbb{C}$-linear map $\tilde{\varphi}: V_{1} \rightarrow V_{2}$ satisfying $\tilde{\varphi}\left(\Lambda_{1}\right) \subseteq \Lambda_{2}$. Moreover, if $\varphi$ is an isogeny, then $\tilde{\varphi}$ is an isomorphism.
Proof. Recall that each quotient map $V_{i} \rightarrow V_{i} / \Lambda_{i}$ is homeomorphic onto its image when restricted to a small enough open neighborhood of the identity $0 \in V_{i}$. Then the proof is an easy exercise in complex analysis.

Corollary 2.5.3. Let $X_{1}$ and $X_{2}$ be abelian varieties of the same dimension, and let $\varphi$ : $X_{1} \rightarrow X_{2}$ be a morphism which is a homomorphism of groups. Then $\varphi$ is an isogeny if and only if it has finite kernel.

Proof. Suppose that $\varphi$ is an isogeny. Then its lifting $\tilde{\varphi}$ is an isomorphism from the covering space $V_{1}$ of $X_{1} \cong V_{1} / \Lambda_{1}$ to the covering space $V_{2}$ of $X_{2} \cong V_{2} / \Lambda_{2}$, so we may identify $V_{1}$ and $V_{2}$ and consider $\tilde{\varphi}$ to be an automorphism of a vector space $V$ such that $\tilde{\varphi}\left(\Lambda_{1}\right) \subseteq \Lambda_{2}$, where $\Lambda_{1}$ and $\Lambda_{2}$ are rank- $2 g$ lattices in $V$. Then to prove the claim, it suffices to show that the quotient $\Lambda_{2} / \tilde{\varphi}\left(\Lambda_{1}\right)$ is finite. But since $\tilde{\varphi}$ is an $\mathbb{R}$-linear isomorphism, the lattice $\tilde{\varphi}\left(\Lambda_{1}\right)$ has rank $2 g$, as does $\Lambda_{2}$, and so the induced quotient is finite.

Now suppose that $\varphi$ is not an isogeny - that is, $\varphi$ is not surjective. Then clearly $\tilde{\varphi}$ is not a surjection onto $V_{2}$. Since $V_{1}$ and $V_{2}$ have the same dimension, this implies that $\tilde{\varphi}$ is not injective either, so that $\operatorname{ker}(\tilde{\varphi}) \subseteq V_{1}$ is a subspace of positive dimension. Since $\Lambda_{1} \subset V_{1}$ is discrete, the image of $\operatorname{ker}(\tilde{\varphi})$ in $X_{1}$ is positive-dimensional and in particular not finite. Since this is contained in the kernel of $\varphi$, we see that $\varphi$ does not have finite kernel.

The above corollary gives an equivalent definition of "isogeny" which is used in many texts. In the case of elliptic curves, the "surjection" and "finite kernel" conditions are also equivalent to the condition that the homomorphism is nontrivial. (Actually, Silverman's book [8] includes the trivial map as an isogeny between elliptic curves, but this is inconsistent with our definition and most definitions of the term.)

Corollary 2.5.4. Let $X$ be a complex abelian variety, and let $N$ be a finite subgroup of $X$. Then there exists a complex abelian variety $X^{\prime}$ and an isogeny $\varphi: X \rightarrow X^{\prime}$ whose kernel is $N$, so that $X^{\prime} \cong X / N$. The isogeny $\varphi$ is unique up to automorphism of $X^{\prime}$.

Proof. The inverse image of $N$ under $V \rightarrow V / \Lambda \cong X$ is a lattice $\Lambda^{\prime} \subset V$ which contains $\Lambda$. If we let $X^{\prime}=V / \Lambda^{\prime}$, it is clear that the identity automorphism on $V$ induces a surjective map $\varphi: X \rightarrow X^{\prime}$ whose kernel coincides with $N$. The uniqueness up to automorphism is clear from considering the liftings of two such isogenies $X \rightarrow X^{\prime}$ to automorphisms of $V$. In order to show that $X^{\prime}$ also has the structure of abelian variety, it is enough to find a positive definite Riemann form for $X^{\prime}$. It is easy to check that if $H: V \times V \rightarrow \mathbb{C}$ is a positive definite Riemann form with respect to $\Lambda$ and $n=\left[\Lambda^{\prime}: \Lambda\right]$ is the degree of $\varphi$, then $n^{2} H$ is a positive definite Riemann form with respect to $\frac{1}{n} \Lambda$. Since clearly $\Lambda^{\prime} \subseteq \frac{1}{n} \Lambda$, we are done.

Definition 2.5.5. The degree of an isogeny $\varphi: X_{1} \rightarrow X_{2}$ is the order of its kernel.
We now briefly discuss the most important example of an isogeny.
Example 2.5.6. If $X$ is a complex abelian variety of dimension $g$, then Corollary 2.1.4 says that the multiplication-by- $n$ map $[n]: X \rightarrow X$ is surjective, and that it has finite kernel isomorphic to $(\mathbb{Z} / n \mathbb{Z})^{2 g}$, for any integer $n \geq 1$. Therefore, $[n]$ is an isogeny of degree $n^{2 g}$.

Proposition 2.5.7. Let $\varphi: X_{1} \rightarrow X_{2}$ be an isogeny of degree $n$ of complex abelian varieties of dimension $g$. Then there is an isogeny $\varphi^{\prime}: X_{2} \rightarrow X_{1}$ of degree $n^{2 g-1}$ such that $\varphi^{\prime} \circ \varphi$ and $\varphi \circ \varphi^{\prime}$ are the multiplication-by-n maps $[n]_{X_{1}}$ and $[n]_{X_{2}}$ on $X_{1}$ and $X_{2}$ respectively.

Proof. We again identify the covering spaces of $X_{1}$ and $X_{2}$, so that $X_{1}=V / \Lambda_{1}$ and $X_{2}=$ $V / \Lambda_{2}$ for a $g$-dimensional complex vector space $V$ and rank- $2 g$ lattices $\Lambda_{1} \subseteq \Lambda_{2} \subset V$. Let $X_{1}[n]$ denote the kernel of the multiplication-by- $n$ map $X_{1} \rightarrow X_{1}$. Then the inverse image of $X_{1}[n]$ under $V \rightarrow V / \Lambda_{1} \cong X_{1}$ clearly coincides with $\frac{1}{n} \Lambda_{1}$. It is easy to see that we have the inclusions of lattices $\Lambda_{1} \subseteq \Lambda_{2} \subseteq \frac{1}{n} \Lambda_{1}$. Then $N^{\prime}:=\frac{1}{n} \Lambda_{1} / \Lambda_{2}$ has order $n^{2 g-1}$, and by Corollary 2.5.4, there is an abelian variety $X_{3}$ and an isogeny $\varphi^{\prime}: X_{2} \rightarrow X_{3}$ with kernel $N^{\prime}$. But clearly $X_{3}=V / \frac{1}{n} \Lambda_{1}$, which is isomorphic to $X / \Lambda_{1}$ via the multiplication-by- $n$ homothety on $V$, so $\varphi^{\prime}: X_{2} \rightarrow X_{1}$ is an isogeny of degree $n^{2 g-1}$. Moreover, the kernel of $\varphi^{\prime} \circ \varphi: X_{1} \rightarrow X_{1}$ is $\frac{1}{n} \Lambda_{1} / \Lambda_{1} \cong X_{1}[n]$. Since the multiplication-by- $n$ map $[n]_{X_{1}}: X_{1} \rightarrow X_{1}$ has this kernel, by the uniqueness given in the statement of Corollary 2.5.4, after composing $\varphi^{\prime}$ with an automorphism of $X_{1}$ we can assume that $\varphi^{\prime} \circ \varphi=[n]_{X_{1}}$. The claim that $\varphi \circ \varphi^{\prime}=[n]_{X_{2}}$ now follows from a similar argument.

Remark 2.5.8. The above proposition shows that if $X_{1}$ and $X_{2}$ are abelian varieties with an isogeny $X_{1} \rightarrow X_{2}$, then there is also an isogeny from $X_{2}$ to $X_{1}$. Therefore, it makes sense to simply say " $X_{1}$ and $X_{2}$ are isogenous", with "isogenous" being an equivalence relation.

### 2.5.2 Polarizations of complex abelian varieties

We will define a particular type of holomorphic homomorphism $\varphi$ from a complex abelian variety $X$ to another complex torus $X^{\vee}$, and we will prove that $X^{\vee}$ is also an abelian variety and $\varphi$ is an isogeny. In order to do this, we first need to further develop the theory of Riemann forms associated to divisors on $X$; for this purpose, we will only assume for the moment that $X$ is a complex torus which is not necessarily an abelian variety.

In $\S 2.2$ we constructed a homomorphism $\mathcal{H}$ from $\operatorname{Div}(X)$ to the group of Hermitian forms on the covering space $V$ (and we showed that in fact, the image coincides with the group of Riemann form for $X$ ). We shall take this notion further. We first observe that the set of normalized theta functions is closed under multiplication, and in fact, $\theta_{D} \theta_{D^{\prime}}=\theta_{D+D^{\prime}}$ for any divisors $D, D^{\prime} \in \operatorname{Div}(X)$. Thus, $D \mapsto \theta_{D}$ defines a homomorphism from $\operatorname{Div}(X)$ to the group of theta functions. Moreover, if we associate to each normalized theta function $\theta_{D}$ the corresponding Riemann form $H: V \times V \rightarrow \mathbb{C}$ and map $K: \Lambda \rightarrow \mathbb{R}$ as in the statement of Proposition 2.2.8, we see that this association respects addition of divisors in $\operatorname{Div}(X)$. In other words, the map $\Phi: D \mapsto\left(H_{D}, K_{D}\right)$, where $H_{D}: V \times V \rightarrow \mathbb{C}$ and $K_{D}: \Lambda \rightarrow \mathbb{R}$ are the Riemann form and map associated to $\theta_{D}$, defines a homomorphism
$\operatorname{Div}(X) \rightarrow\{$ additive group of Herm. forms $V \times V \rightarrow \mathbb{C}\} \times\{$ additive group of maps $\Lambda \rightarrow \mathbb{R}\}$.
The kernel of $\Phi$ is the subgroup of principle divisors $\operatorname{Prin}(X) \subset \operatorname{Div}(X)$, because a normalized theta function $\theta_{D}$ has trivial $H$ and $K$ if and only if $\theta_{D}$ is periodic with respect to $\Lambda \subset V$, which means that it induces a meromorphic function on $X$ whose associated divisor is $D$.

Since $\mathcal{H}$ is simply the composition of $\Phi$ with the projection to the group of Riemann forms on $V$, its kernel is some subgroup of $\operatorname{Div}(X)$ containing $\operatorname{Prin}(X)$. We let $\operatorname{Div}^{0}(X)$ denote this kernel; in other words, $\operatorname{Div}^{0}(X) \subset \operatorname{Div}(X)$ is the subgroup of all divisors whose associated Riemann form is 0 . We let $\operatorname{Pic}(X):=\operatorname{Div}(X) / \operatorname{Prin}(X)$ denote the Picard group
of $X$ and write $\operatorname{Pic}^{0}(X)$ for the subgroup $\operatorname{Div}^{0}(X) / \operatorname{Prin}(X)$. From now on, we consider $\Phi$ as a homomorphism on $\operatorname{Pic}(X)$ (whose kernel is $\operatorname{Pic}^{0}(X)$ ).

Note that $\Phi$ associates to each divisor class $[D] \in \operatorname{Pic}^{0} \subseteq \operatorname{Pic}(X)$ a map $K: \Lambda \rightarrow \mathbb{R}$ which satisfies the property that $K\left(\lambda_{1}+\lambda_{2}\right)-K\left(\lambda_{1}\right)-K\left(\lambda_{2}\right) \equiv \frac{1}{2} E\left(\lambda_{1}, \lambda_{2}\right)=0$ modulo $\mathbb{Z}$, so that $\theta_{D}(v+\lambda) / \theta_{D}(v)=e^{2 \pi i K(\lambda)}$ is actually a group homomorphism $\Lambda \rightarrow \mathbb{C}^{\times}$which does not depend our choice of a fixed $v \in V$. In other words, for $[D] \in \operatorname{Pic}^{0}(X)$ and $K: \Lambda \rightarrow \mathbb{R}$ the associated map, the function $\lambda \mapsto e^{2 \pi i K(\lambda)}$ is a complex character $\chi_{[D]}$ on $\Lambda$. It is easy to check that $\chi_{\left[D_{1}+D_{2}\right]}=\chi_{\left[D_{1}\right]} \chi_{\left[D_{2}\right]}$ for $D_{1}, D_{2} \in \operatorname{Div}(X)$, so $\Phi$ induces an injective homomorphism from $\operatorname{Pic}^{0}(X)$ to the group of complex characters on $\Lambda$ (we will soon see that this is an isomorphism. A "complex character" is understood by definition to take values in the unit circle.)

Now in order to define a polarization, for any $a \in X$, we denote the translation-by- $a$ morphism taking $b \in X$ to $b+a$ by $t_{a}: X \rightarrow X$. It induces an pullback automorphism $t_{a}^{*}: \operatorname{Div}(X) \rightarrow \operatorname{Div}(X)$. Note that $t_{a}^{*}$ stabilizes $\operatorname{Prin}(X)$ (for any $(f) \in \operatorname{Prin}(X)$, check that $t_{a}^{*}(f)=\left(f \circ t_{a}\right)$, so it may be considered as an automorphism of $\operatorname{Pic}(X)$ as well.

Definition 2.5.9. A polarization is a map $X \rightarrow \operatorname{Pic}(X)$ of the form

$$
\varphi_{D}: a \mapsto\left[t_{-a}^{*} D-D\right]
$$

for some ample divisor $D \in \operatorname{Div}(X)$.
Remark 2.5.10. This definition as well as the results below would all still be valid with " $a$ " in place of " $-a$ " above, but we include the minus sign in order to get a nicer description of principal polarizations on elliptic curves later.

### 2.5.3 The dual of a complex abelian variety

Our main goal now is to show that $\operatorname{Pic}^{0}(X)$ has the structure of a complex abelian variety and that any polarization $\varphi_{D}$ is an isogeny. We will first show that a polarization is surjective with finite kernel. In order to show this, we need a lemma.

Lemma 2.5.11. For any $D \in \operatorname{Div}(X)$ and $a \in X$, we have the following.
a) The divisor $t_{a}^{*} D-D$ lies in $\operatorname{Div}^{0}(X)$; therefore, the image of any polarization is contained in $\operatorname{Pic}^{0}(X)$.
b) Under the homomorphism from $\operatorname{Pic}^{0}(X)$ to the group of complex characters on $\Lambda$ defined above, $\left[t_{a}^{*} D-D\right]$ goes to the character $\lambda \mapsto e^{2 \pi i E(w, \lambda)}$, where $E$ is the imaginary part of the Riemann form associated to $D$ and $w$ is some (any) element of $V$ whose image modulo $\Lambda$ is $a \in X \cong V / \Lambda$.

Proof. Let $H: V \times V \rightarrow \mathbb{C}$ be the Riemann form and $K: \Lambda \rightarrow \mathbb{R}$ be the map such that $\Phi(D)=(H, K)$. Fix a point $w \in V$ whose image modulo $\Lambda$ is $a \in X$. Let $\theta_{D}$ be a normalized theta function for $D$. Then it is clear that $\theta^{\prime}:=\left(\theta_{D} \circ t_{w}\right)=\pi^{*} t_{a}^{*} D$, where $t_{w}: V \rightarrow V$ is the translation-by- $w$ function on $V$. For $v \in V$ and $\lambda \in \Lambda$, we compute

$$
\begin{equation*}
\theta^{\prime}(v+\lambda) / \theta^{\prime}(v)=e^{\pi H(w+v, \lambda)+\frac{1}{2} H(\lambda, \lambda)+2 \pi i K(\lambda)} \tag{2.20}
\end{equation*}
$$

Let $\theta^{\prime \prime}(v)=e^{-\pi H(v, w)} \theta^{\prime}(v)$, which is a theta function with the same divisor as $\theta^{\prime}$ since $v \mapsto e^{-\pi H(v, w)}$ is a trivial theta function. This contributes a factor of $e^{-\pi H(\lambda, w)}$ to the functional equation above, so that now we have, for $v \in V$ and $\lambda \in \Lambda$,

$$
\begin{gather*}
\theta^{\prime \prime}(v+\lambda) / \theta^{\prime \prime}(v)=e^{2 \pi i\left(\frac{1}{2 i} H(v+w, \lambda)+\frac{1}{4 i} H(\lambda, \lambda)+K(\lambda)-\frac{1}{2 i} H(\lambda, w)\right.} \\
=e^{2 \pi i\left(\frac{1}{2 i} H(w, \lambda)-\frac{1}{2 i} H(\lambda, w)\right.} \cdot e^{2 \pi i\left(\frac{1}{2 i} H(v, \lambda)+\frac{1}{4 i} H(\lambda, \lambda)+K(\lambda)\right.} \\
=e^{2 \pi i E(w, \lambda)} \cdot e^{2 \pi i\left(\frac{1}{2 i} H(v, \lambda)+\frac{1}{4 i} H(\lambda, \lambda)+K(\lambda)\right.}=e^{2 \pi i E(w, \lambda)} \cdot \theta_{D}(v+\lambda) / \theta_{D}(v) . \tag{2.21}
\end{gather*}
$$

Therefore we see that $\theta^{\prime \prime} / \theta_{D}$ is a normalized theta function with $\left(\theta^{\prime \prime} / \theta_{D}\right)=\pi^{*}\left(t_{a}^{*} D-D\right) \in$ $\operatorname{Div}(V)$, so $\theta^{\prime \prime} / \theta_{D}=\theta_{t_{a}^{*} D-D}$. Moreover, we have $\theta_{t_{a}^{*} D-D}(v+\lambda) / \theta_{t_{a}^{*} D-D}(v)=e^{2 \pi i E(w, \lambda)}$ for $v \in V$ and $\lambda \in \Lambda$, and so the associated Riemann form is 0 while the associated map $\Lambda \rightarrow \mathbb{R}$ is $\lambda \mapsto E(w, \lambda)$ (note that this is independent modulo $\mathbb{Z}$ of our choice of $w$ since $E(\Lambda, \Lambda) \subseteq \mathbb{Z}$ ). Then (a) and (b) both follow from the definition of $\operatorname{Div}^{0}(X)$ and constructions in the discussion above.

We will now, for the rest of this section, assume that $X$ is an abelian variety (so $X$ has an ample divisor).

Proposition 2.5.12. Let $\varphi_{D}: X \rightarrow \operatorname{Pic}^{0}(X)$ be a polarization as defined in Definition 2.5.9. where $D \in \operatorname{Div}(X)$ is an ample divisor. Then $\varphi_{D}$ is a surjective homomorphism of abstract groups with kernel of order $\operatorname{det}(E)$, where $E$ is the imaginary part of the Riemann form associated to $D$. Moreover, the injection of $\operatorname{Pic}^{0}(X)$ into the group of complex characters of $\Lambda$ given by $\mathfrak{D} \mapsto \chi_{\mathfrak{D}}$ as constructed above is an isomorphism.

Proof. We first observe that any complex character $\chi: \Lambda \rightarrow \mathbb{C}^{*}$ must be of the form $\lambda \mapsto$ $e^{2 \pi i B(\lambda)}$ for some $\mathbb{R}$-linear function $B: \Lambda \rightarrow \mathbb{R}$. Meanwhile, by Proposition 2.3.2, since $D$ is ample, the $\mathbb{R}$-bilinear form $E: V \times V \rightarrow \mathbb{R}$ is nondegenerate. Therefore, any $\mathbb{R}$-linear function $B: \Lambda \rightarrow \mathbb{R}$ is given by $E(w, \cdot)$ for some $w \in V$. We therefore have $\chi_{\mathfrak{D}}=e^{2 \pi i E(w, \cdot)}$ for some $w \in V$, and it follows from Lemma 2.5.11(b) that $\varphi_{D}(-a)=\mathfrak{D}$ for any divisor class $\mathfrak{D} \in \operatorname{Pic}^{0}(X)$, where $a=\pi(w)$. Thus, $\mathfrak{D} \mapsto \chi_{\mathfrak{D}}$ is actually an isomorphism from $\operatorname{Pic}^{0}(X)$ to the group of complex characters on $\Lambda$, and $\varphi_{D}$ is a surjective map. The fact that $\varphi_{D}$ is also a group homomorphism follows easily from our formula $\chi_{\varphi_{D}(a)}=e^{2 \pi i E(w, \cdot)}$.

It is now clear that the kernel of $\varphi_{D}$ consists of the images modulo $\Lambda$ of elements $w \in V$ with $e^{2 \pi i E(w, \lambda)}=1$ and thus $E(w, \lambda) \in \mathbb{Z}$ for all $\lambda \in \Lambda$. It follows quickly from Lemma 2.3.4(a) that the subset $\Lambda^{\prime} \subset V$ consisting of all such $w$ is a lattice containing $\Lambda$ such that the quotient $\Lambda^{\prime} / \Lambda$ has order $\operatorname{det}(E)$.

The next proposition shows that $\operatorname{Pic}^{0}(X)$ is not only a complex torus but a complex abelian variety.

Proposition 2.5.13. Let $V^{*}$ be the complex vector space of all $\mathbb{C}$-antilinear functions $\xi$ : $V \rightarrow \mathbb{C}$, and let $\Lambda^{*}$ be the subset of all $\xi \in V^{*}$ satisfying $\Im \xi(\Lambda) \subseteq \mathbb{Z}$. Then $\Lambda^{*}$ is a lattice of maximal rank in $V^{*}$, and there is a canonical isomorphism of groups $\operatorname{Pic}^{0}(X) \xrightarrow{\sim} V^{*} / \Lambda^{*}$, thus giving $\operatorname{Pic}^{0}(X)$ the structure of a complex torus. Moreover, $\operatorname{Pic}^{0}(X) \cong V^{*} / \Lambda^{*}$ has a positive definite Riemann form and is therefore an abelian variety by Theorem 2.3.6.

Proof. It is an elementary exercise to check that $V^{*}$ and $\Lambda^{*}$ have the same dimension and rank as $V$ and $\Lambda$ respectively, so that $V^{*} / \Lambda^{*}$ is a complex torus. Fix an ample divisor $D \in \operatorname{Div}(X)$, and let $E$ be the imaginary part of its associated Riemann form. Now we have already shown that $\operatorname{Pic}^{0}(X)$ is isomorphic to the group of complex characters on $\Lambda$, each of which is of the form $\lambda \mapsto e^{2 \pi i B(\lambda)}$ for some $\mathbb{R}$-linear function $B: V \rightarrow \mathbb{R}$. Since for such a $B$ there always exists a $\mathbb{C}$-antilinear function on $V$ whose imaginary part is $-B$ (take $B(i v)-i B(v))$, we get a surjection $V^{*} \rightarrow \operatorname{Pic}^{0}(X)$ given by $\xi \mapsto e^{-2 \pi i \Im \xi(\cdot)}$ whose kernel is clearly $\Lambda^{*}$. Note that if $w \in V$ is an element such that $B(\cdot) \equiv E(w, \cdot)$ modulo $\mathbb{Z}$, then we have $B(i v)-i B(v)=E(-i w, \cdot)-i E(w, \cdot)=H(-w, \cdot)$. It now follows from what we have shown before that the map $\varphi_{D}: X \rightarrow \operatorname{Pic}^{0}(X) \xrightarrow{\sim} V^{*} / \Lambda^{*}$ is given by

$$
a \mapsto t_{-a}^{*} D-D \mapsto H(w, \cdot)+\Lambda^{*} \in V^{*} / \Lambda^{*},
$$

where $w$ is some (any) element of $V$ mapping to $a \in X$. Therefore, it lifts to the isomorphism $\widetilde{\varphi_{D}}: V \xrightarrow{\sim} V^{*}$ given by $w \mapsto H(w, \cdot) \in V^{*}$. Note that $\widetilde{\varphi_{D}}(\Lambda) \subseteq \Lambda^{*}$ is a sublattice of (finite) index equal to the order of $\operatorname{ker}\left(\varphi_{D}\right)$, which is $\operatorname{det}(E)$, implying that $\Lambda \subseteq \Lambda^{*} \subseteq \frac{1}{\operatorname{det}(E)} \widetilde{\varphi_{D}}(\Lambda)$. Now define $H^{*}: V^{*} \times V^{*} \rightarrow \mathbb{C}$ by $H^{*}\left(\xi_{1}, \xi_{2}\right)=H\left(\widetilde{\varphi_{D}}-1\left(\xi_{1}\right), \widetilde{\varphi_{D}}{ }^{-1}\left(\xi_{2}\right)\right)$. It is immediate to check that $H^{*}$ is a positive definite Riemann form on $V^{*}$ with $H^{*}\left(\widetilde{\varphi_{D}}(\Lambda), \widetilde{\varphi_{D}}(\Lambda)\right) \subseteq$ $\mathbb{Z}$. Therefore, we have $H^{*}\left(\Lambda^{*}, \Lambda^{*}\right) \subseteq \operatorname{det}(E)^{-2} \mathbb{Z}$, so that $\operatorname{det}(E)^{2} H^{*}$ is a positive definite Riemann form for $V^{*} / \Lambda^{*}$.

We denote the group $\operatorname{Pic}^{0}(X)$ with its structure as an abelian variety by $X^{\vee}$ and call it the dual abelian variety of $X$. Any polarization $\varphi_{D}: X \rightarrow X^{\vee}$ is an isogeny of degree $\operatorname{det}(E)$ (note that this is always a perfect square). If there exists an ample divisor $D \in \operatorname{Div}(X)$ whose associated Riemann form satisfies $\operatorname{det}(E)=1$, then this isogeny is an isomorphism $\varphi_{D}: X \xrightarrow{\sim} X^{\vee}$. In this case, we say that $\varphi_{D}$ is a principal polarization and that $X$ is "selfdual". It is not too deep to show that there is always a natural isomorphism $\left(X^{\vee}\right)^{\vee} \xrightarrow{\sim} X$, as one expects, but we will not do it here.

Remark 2.5.14. For any complex torus $X$ (even one which is not an abelian variety), what we have shown implies that $\operatorname{Pic}^{0}(X)$ still has the structure of a complex torus, called the dual complex torus of $X$. This intuitively makes sense if one identifies $\operatorname{Pic}^{0}(X)$ with the group of complex characters on $X$ as above. (Note in particular that as a real Lie group, $X \cong(\mathbb{R} / \mathbb{Z})^{2 g}$ and a complex character on $X$ is a homomorphism of real Lie groups $\chi: \mathbb{Z}^{2 g} \rightarrow \mathbb{R} / \mathbb{Z}$ which lifts to an $\mathbb{R}$-linear functional $\tilde{\chi}: \mathbb{R}^{2 g} \rightarrow \mathbb{R}$ which is unique up to functionals which take values in $\mathbb{Z}$ on the standard basis of $\mathbb{R}^{2 g}$.)

However, we need an ample divisor $D$ with a positive definite Riemann form $H$ in order to construct a positive definite Riemann form on the dual complex torus, as well as to get a $\operatorname{map} \varphi_{D}$ which is surjective.

### 2.5.4 Polarizations and self-duality of elliptic curves

We now apply the theory of polarizations to the elliptic curve case. Assume now that $X$ has dimension 1. Then we know from Corollary 2.4.3 that the divisor $(0) \in \operatorname{Div}(X)$ is ample and that the associated Riemann form is positive definite and satisfies $\operatorname{det}(E)=1$. Thus,
$\varphi_{(0)}: X \rightarrow X^{\vee}$, defined using the above notation, is a principal polarization and so $X \cong X^{\vee}$ is self-dual. Explicitly, $\varphi_{(0)}$ is given by $P \mapsto[(P)-(0)] \in \operatorname{Pic}^{0}(X)$. We will see this map again in other contexts. (This is the reason for the minus sign in Definition 2.5.9.)

We also get a simple description of $\operatorname{Div}^{0}(X)$ as well as of ample divisors on $X$ in the elliptic curve case.

Proposition 2.5.15. If $X$ is an elliptic curve, then $\operatorname{Div}^{0}(X) \subset \operatorname{Div}(X)$ coincides with the subgroup of divisors of degree 0. Moreover, a divisor on $X$ is ample if and only if it has positive degree.

Proof. The subgroup of divisors of degree 0 is generated by divisors of the form $(P)-$ $(0) \in \operatorname{Div}(X)$ for a point $P \in X$. But $(P)-(0)=t_{-P}^{*}(0)-(0) \in \operatorname{Div}^{0}(X)$ by Lemma 2.5.11. Therefore, any two divisors in $\operatorname{Div}(X)$ with the same degree have the same associated Riemann form. Let $H$ be the positive definite Riemann form associated to the ample divisor $(0) \in \operatorname{Div}(X)$. For any divisor $D \in \operatorname{Div}(X)$ with $\operatorname{deg}(D)=n$, the divisors $D$ and $n(0)$ have the same associated Riemann form, which is $n H$. Since $n H=0$ (resp. $n H$ is positive definite) if and only if $n=0$ (resp. $n \geq 1$ ), we get both claims of the proposition.

### 2.6 Jacobians of compact Riemann surfaces

In this section we discuss the classical theory which first led to the discovery of elliptic curves and abelian varieties over $\mathbb{C}$ in the early 19th century. Given any given compact Riemann surface $C$ of genus $g$, we will construct a complex abelian variety $J$ of dimension $g$ called the Jacobian of $C$, which as an abstract group is isomorphic to $\operatorname{Pic}^{0}(C)$. This will help to provide some more concrete examples of complex abelian varieties of dimension $\geq 2$ as well as a deeper understanding of complex elliptic curves.

Throughout this section, we will resume the same notation of $\operatorname{Div}(C)$ (group of Weil divisors), $\operatorname{Div}^{0}(C)$ (subgroup of degree- 0 divisors), $\operatorname{Prin}(C), \operatorname{Pic}(C), \operatorname{Pic}^{0}(C)$, etc. exactly as they were defined for elliptic curves in \$2.4.

### 2.6.1 Differential forms and the Riemann-Roch theorem

Since we will only be dealing with differential forms on smooth curves, we will follow the definition of Silverman: the complex vector space of differential forms on a complex curve $C$, denoted $\Omega(C)$, is the set of all symbols $d f$ for any $f$ in the field of meromorphic functions $\mathcal{M}(C)$ satisfying the following relations: $d(f+g)=d f+d g ; d(c f)=c(d f)$; and $d(f g)=$ $f(d g)+g(d f)$ for any $f, g \in \mathcal{M}(C)$ and $c \in \mathbb{C}$.

It is easy to show that if $\omega_{1}, \omega_{2} \in \Omega(C)$, then $\omega_{2}=f \omega_{1}$ for some unique $f \in \mathcal{M}(C)$, so $\Omega(C)$ may also be viewed as a 1-dimensional $\mathcal{M}(C)$-vector space. This fact enables us to define the order of a differential form at any point on the curve. Given any differential form $\omega \in \Omega(C)$ and any point $P \in C$, the order of $\omega$ at $P$ is the order of the function $\omega / d t \in \mathcal{M}(C)$ at $P$, where $t \in \mathcal{M}(C)$ is some (any) uniformizer at $P$. So there is a divisor associated to each differential form $\omega \in \Omega(C)$, denoted $(\omega) \in \operatorname{Div}(C)$, given by $\sum_{P \in C} n_{P}(P)$ where each $n_{P}$ is the order of $\omega$ at $P$. In fact, for $\omega_{1}, \omega_{2} \in \Omega(C)$, if we let $f=\omega_{2} / \omega_{1} \in \mathcal{M}(C)$,
then we see that $\left(\omega_{2}\right)=\left(\omega_{1}\right)+(f) \in \operatorname{Div}(C)$, so the divisors of all differential forms on $C$ lie in the same divisor class in $\operatorname{Pic}(C)$. We denote this divisor class by $[K] \in \operatorname{Pic}(C)$ and call any representative $K \in \operatorname{Div}(C)$ a canonical divisor of $C$.

The algebraic definition of the genus of a curve is the dimension over $\mathbb{C}$ of the vector space of all holomorphic differential forms on $C$ (differential forms whose associated divisor is positive; we will see from the Riemann-Roch theorem below that this dimension is always finite). One can show, by taking triangulations and computing Euler characteristics, that this definition of "genus" agrees with the topological definition, where the genus of a compact Riemann surface is the "number of holes".

Example 2.6.1. (hyperelliptic curves)
We define a hyperelliptic curve over $\mathbb{C}$ to be a smooth projective curve defined by an equation of the form $y^{2}=f(x)$ for some polynomial $f \in \mathbb{C}[x]$ of degree $d \geq 1$ which has $d$ distinct roots. If $d$ is odd, then there will be one added "point at infinity" $\infty$, and if $d$ is even, there will be two added "points at infinity" denoted $\infty_{1}$ and $\infty_{2}$. Note that a hyperelliptic curve is a conic if $d=1,2$, and an elliptic curve is a hyperelliptic curve for $d=3$ by Definition 1.1.4. (In fact, we will see below that the genus of a hyperelliptic curve $C$ of degree $d=4$ is 1 , and so by Definition 1.1.3, such a curve $C$ with a distinguished point is also an elliptic curve.)

In order to visualize $C$ as a complex manifold, we consider the function $x \mapsto y=\sqrt{f(x)}$ on the Riemann sphere $\mathbb{C} \cup\{\infty\}$. In fact, $\sqrt{f(x)}$ can't be defined as a meromorphic function on the Riemann sphere, but it can be defined on the compliment of some branch cuts which connect the zeroes of the polynomial $f$. Define the subset $B \subset \mathbb{C} \cup\{\infty\}$, called the subset of branch points, to be the set of roots of $f$ if $d$ is even, and let $B \subset \mathbb{C} \cup\{\infty\}$ be the set of roots of $f$ along with the point $\infty$ if $d$ is odd (note that the cardinality of $B$ is always even and equal to either $d$ or $d+1$ ). We partition $B$ into cardinality- 2 subsets and for each such subset draw a line connecting the corresponding two points in $B$, so that none of the $\# B / 2$ lines intersect. Then there is a well-defined meromorphic function $\sqrt{f(x)}$ whose square is $f(x)$ defined on the compliment of these branch cuts. Of course, $-\sqrt{f(x)}$ is another such function. With a little visual intuition, one can see that the Riemann surface defined by the equation $y^{2}=f(x)$ can be constructed by taking two copies of the Riemann sphere, one on which $y$ takes the value $\sqrt{f(x)}$ and the other on which $y$ takes the value $-\sqrt{f(x)}$, "opening" their branch cuts, and gluing them together along their opened branch cuts. Since each copy of the Riemann sphere had $\# B / 2$ branch cuts, the resulting Riemann surface has $\# B / 2-1$ holes. Since $\# B \in\{d, d+1\}$, it follows that the genus of $C$ is equal to $\lfloor(d-1) / 2\rfloor$.

We can show that this agrees with the algebraic definition of genus as follows. One can compute that the differential form $d x / y \in \Omega(C)$ has associated divisor $(d-3)(\infty) \in \operatorname{Div}(C)$ (resp. $\left.(d / 2-2)\left(\infty_{1}\right)+(d / 2-2)\left(\infty_{2}\right) \in \operatorname{Div}(C)\right)$ if $d$ is odd (resp. if $d$ is even). Moreover, one can compute that a nonzero polynomial function $h \in \mathbb{C}[x] \subset \mathcal{M}(C)$ of degree $d^{\prime} \geq 0$ has only a pole at $\infty$ of order $2 d^{\prime}$ (resp. only poles at $\infty_{1}$ and $\infty_{2}$, each of order $d^{\prime}$ ) if $d$ is odd (resp. if $d$ is even). Thus, the space of all holomorphic differentials on $C$ consists of all differentials of the form $h \cdot d x / y$ where $h \in \mathbb{C}[x] \subset \mathcal{M}(C)$ is a polynomial function of degree $d^{\prime} \leq\lfloor(d-1) / 2\rfloor-1$, which has dimension $\lfloor(d-1) / 2\rfloor$. Thus, again we see that the genus of $C$ is equal to $g:=\lfloor(d-1) / 2\rfloor$. So if a hyperelliptic curve has genus $g$, then the polynomial $f$ used to define it has degree $2 g+1$ or $2 g+2$, and there are always $2 g+2$ branch points.

In fact, we may easily describe a basis for the singular homology of $C$ as follows. Order the points in $B$ as $\left\{z_{1}, z_{2}, \ldots, z_{2 g \pm 2}\right\}$ so that the $j$ th branch cut connects the point $z_{2 i-1}$ to $z_{2 i}$ for $1 \leq i \leq g+1$. Let $\bar{a}_{i}$ (resp. $\bar{b}_{i}$ ) be a simple closed loop surrounding only the points $z_{2 i-1}$ and $z_{2 i}$ (resp. the points $z_{2 i}, z_{2 i+1}, \ldots, z_{2 g+1}$ ) for $1 \leq i \leq g$. Then the $\bar{a}_{i}$ 's and $\bar{b}_{i}$ 's lift to simple closed loops $a_{i}$ and $b_{i}$ on the compact Riemann surface $C$. This set of loops is a basis for the homology group $H_{1}(C, \mathbb{Z})$, which is a free abelian group generated by $\left\{a_{1}, \ldots, a_{g}, b_{1}, \ldots, b_{g}\right\}$.

Note that the case of $d=3$ shows that an elliptic curve given by $y^{2}=f(x)$ as in Definition 1.1.4 always has genus 1 , and its only holomorphic differentials are constant multiples of $d x / y$. In fact, these holomorphic differentials have trivial associated divisor, so that the class of canonical divisors for an elliptic curve $C$ is $[K]=0 \in \operatorname{Div}(C)$. Recall that an elliptic curve of this form can be identified with $\mathbb{C} / \Lambda$ for some rank- 2 lattice $\Lambda \subset \mathbb{C}$, and that there is a periodic function $\wp$ defined on $\mathbb{C}$ inducing $x \in \mathcal{M}(C)$ such that its derivative $\wp^{\prime}$ induces $y \in \mathcal{M}(C)$. Then note that the differential form $d x / y \in \Omega(C)$ lifts to $\wp^{\prime} d z / \wp^{\prime}=d z \in \Omega(\mathbb{C})$, and clearly the only holomorphic differentials on $\mathbb{C}$ are constant multiples of $d z$ (which has trivial divisor). Any constant multiple of $d x / y$ is called an invariant differential of the elliptic curve $C$ and will be important later.

The curious reader may find more details in [6, §1].
The Riemann-Roch theorem gives a formula which can be used to compute (among other things) the dimension of the vector space $\mathcal{L}(D)$ for any divisor $D$ on a compact smooth curve $C$. From now on, for any divisor $D \in \operatorname{Div}(C)$, we define $l(D)$ to be the dimension of the vector space $\mathcal{L}(D)$. Note that $l(D)$ does not depend on the choice of representative $D$ of the divisor class $[D] \in \operatorname{Pic}(C)$.

Theorem 2.6.2. (Riemann-Roch)
For any divisor $D \in \operatorname{Div}(C)$ and canonical divisor $K$ of $C$, there is an integer $g \geq 0$ which depends only on $C$, such that

$$
\operatorname{deg}(D)=l(D)-l(K-D)+g-1
$$

This is a standard theorem in algebraic geometry, and we do not prove it in these notes. A proof can be found in [2, §2].
Corollary 2.6 .3 .
a) The constant $g$ given in the statement of Theorem 2.6.2 is the genus of $C$.
b) The degree of any canonical divisor of $C$ is $2 g-2$.

Proof. First we observe that there is an isomorphism from the vector space of holomorphic differentials on $C$ to $\mathcal{L}(K)$, given by $\omega \mapsto \omega / \omega_{0}$ for some $\omega_{0} \in \Omega(C)$ such that $\left(\omega_{0}\right)=K \in$ $\operatorname{Div}(C)$. Thus, the genus of $C$ is equal to $l(K)$.

Part (a) is proven by putting $D=0$ into the formula given by Theorem 2.6.2. Then part (b) is proven by instead putting $D=K$ into that formula.

The following corollary is useful in proving the Abel-Jacobi theorem below.
Corollary 2.6.4. Let $O \in C$ be an arbitrary point. Then each divisor class in $\operatorname{Pic}^{0}(C)$ can be represented by a divisor of the form $\sum_{i=1}^{g} n_{i}\left(P_{i}\right)-g(O)$ for some points $P_{i} \in C$ (not necessarily distinct) and integers $n_{i}$.

Proof. Choose any divisor $D \in \operatorname{Div}^{0}(C)$. We deduce from Theorem 2.6.2 that $\ell(D+g(O))=$ $g+\ell(K-(D+g(O)))-g+1 \geq 1$, so there is a nonzero meromorphic function $h \in \mathcal{L}(D+g(O))$. Then we have $D^{\prime}:=D+(h)>-g(O)$, so $D^{\prime}$ is a divisor of the desired form which lies in the same class as $D$ in $\operatorname{Pic}^{0}(C)$.

### 2.6.2 The Abel-Jacobi map

The singular homology group $H_{1}(C, \mathbb{Z})$ of our compact genus- $g$ Riemann surface $C$ is freely generated by a cardinality- $2 g$ set $\left\{a_{1}, \ldots, a_{g}, b_{1}, \ldots, b_{g}\right\}$ of simple closed loops which look exactly like the loops explicitly constructed in the case that $C$ is a hyperelliptic curve (Example 2.6.1). In fact, as a topological space, $C$ may be constructed by taking a $4 g$-sided polygon, called a fundamental polygon of $C$ and denoted by $\mathcal{P}$, and identifying the sides by considering each one as a loop in $H_{1}(C, \mathbb{Z})$ in the order $\left\{a_{1}, b_{1},-a_{1},-b_{1}, \ldots, a_{g}, b_{g},-a_{g},-b_{g}\right\}$. Now choose a basis $\left\{\omega_{1}, \ldots, \omega_{g}\right\}$ of the $g$-dimensional vector space of holomorphic differential forms on $C$.

The Abel-Jacobi theorem arose from attempts to calculate integrals of certain differential forms along paths on hyperelliptic curves. As we know from complex analysis, such a path integral is only well-defined up to integrals along closed loops in the homology group: in general, an integral of a differential form along a nontrivial closed loop will not equal 0 . Let $\Lambda \in \mathbb{C}^{g}$ be the additive subgroup generated by

$$
\left.\left\{\left(\int_{a_{i}} \omega_{1}, \ldots, \int_{a_{i}} \omega_{g}\right),\left(\int_{b_{i}} \omega_{1}, \ldots, \int_{b_{i}} \omega_{g}\right)\right\}_{1 \leq i \leq g}\right\} \subset \mathbb{C}^{g}
$$

(We will later show that $\Lambda$ is a rank- $2 g$ lattice.)
Now we define the Abel-Jacobi map as follows. Fix a point $O \in C$. Now for any point $P \in C$, path integrals of the form $\int_{O}^{P} \omega$ for some holomorphic differential form $\omega \in \Omega(C)$ are not well-defined, because composing any path from $O$ to $P$ with a nontrivial closed loop of base $O$ may change the value of the integral. However, if we let AJ (named after Abel and Jacobi) be the function on $C$ given by

$$
P \mapsto\left(\int_{O}^{P} \omega_{1}, \ldots, \int_{O}^{P} \omega_{g}\right),
$$

we see that this takes values in $\mathbb{C}^{g}$ determined up to elements of $\Lambda$. Thus, AJ : $C \rightarrow \mathbb{C}^{g} / \Lambda$ is a well-defined function, which can be extended $\mathbb{Z}$-linearly to a function AJ : $\operatorname{Div}(C) \rightarrow \mathbb{C}^{g} / \Lambda$.

Remark 2.6.5. In the case that $C$ is an elliptic curve given by an equation of the form $y^{2}=f(x)$ where $f \in \mathbb{C}[x]$ is a cubic polynomial in $x$, we have already seen from Example 2.6.1 that all holomorphic differentials in $\Omega(C)$ are of scalar multiples of $d x / y$. In this case, the map AJ : $C \rightarrow \mathbb{C} / \Lambda$ is a biholomorphism (see Exercise 2.7.8) which is essentially an inverse to the map $\mathbb{C} \rightarrow C$ given by $z \mapsto\left(\wp(z), \wp^{\prime}(z)\right)$.

In fact, this entire area of mathematics first came to light when mathematicians in the early 1800's were attempting to solve integrals of the form $\int d x / \sqrt{f(x)}$ with $f(x)$ a cubic. Such integrals naturally came up in calculations of the arclegnth of an ellipse and were therefore called "elliptic integrals". Researchers soon realized that functions like $1 / \sqrt{f(x)}$
could not be defined everywhere on the complex plane and path integrals of them should instead by studied on the curve given by $y^{2}=f(x)$, and that is why this curve was called an "elliptic curve".

Definition 2.6.6. The Jacobian of a compact Riemann surface $C$ of genus $g$ is the $g$ dimensional complex manifold $J:=\mathbb{C}^{g} / \Lambda$.

The function AJ : $\operatorname{Div}(C) \rightarrow J$ is called the Abel-Jacobi map. One easily verifies that when restricted to $\operatorname{Div}^{0}(C)$, this map does not depend on our choice of basepoint $O \in C$. From now on we will mainly consider the Abel-Jacobi map restricted to $\operatorname{Div}^{0}(C)$. We are finally ready to present the Abel-Jacobi Theorem, which will be proven in the next two subsections.

Theorem 2.6.7. (Abel-Jacobi)
(Abel) The kernel of the map AJ : $\operatorname{Div}^{0}(C) \rightarrow J$ coincides with the subgroup of principal divisors $\operatorname{Prin}(C) \subseteq \operatorname{Div}^{0}(C)$.
(Jacobi) The map AJ is surjective.
Thus, the Abel-Jacobi map induces an isomorphism of abstract groups $\operatorname{Pic}^{0}(C) \xrightarrow{\sim} J$.

### 2.6.3 Proof of Abel's Theorem

Our goal in this section is to prove Abel's part of Theorem 2.6.7, which is that the kernel of the Abel-Jacobi map coincides with $\operatorname{Prin}(C) \subseteq \operatorname{Div}^{0}(C)$. For the proof, we fix a fundamental polygon $\mathcal{P}$ of $C$ which contains $O$ in its interior, whose sides form a homology basis $\left\{a_{1}, \ldots, a_{g}, b_{1}, \ldots, b_{g}\right\}$ as in 2.6.2. We first need some lemmas.

Lemma 2.6.8. Let $\omega$ and $\eta$ be two differential forms on $C$. Then we have the identities

$$
\int_{C} \omega \wedge \eta=\int_{\partial \mathcal{P}} F \eta=\sum_{i=1}^{g}\left(\int_{a_{i}} \omega \int_{b_{i}} \eta-\int_{a_{i}} \eta \int_{b_{i}} \omega\right)
$$

where $F(P)$ is defined to be the integral $\int_{O}^{P} \omega$ along any path from $O$ to $P$ contained in the interior of the polygon $\mathcal{P}$ described at the start of \$2.6.2.

Proof. The first equality is Stokes' Theorem. To prove the second equality, we first note that

$$
\begin{equation*}
\int_{\partial \mathcal{P}} F \eta=\sum_{i=1}^{g}\left(\int_{a_{i}} F \eta+\int_{-a_{i}} F \eta+\int_{b_{i}} F \eta+\int_{-b_{i}} F \eta,\right) \tag{2.22}
\end{equation*}
$$

where the $a_{i}$ 's, $-a_{i}$ 's, $b_{i}$ 's, and $-b_{i}$ 's are the edges of the polygon $\mathcal{P}$. We note from the configuration of the edges of $\mathcal{P}$ as described above that for each $i$, we have

$$
\begin{equation*}
\int_{a_{i}} F \eta+\int_{-a_{i}} F \eta=\int_{a_{i}}\left(\int_{O}^{P} \omega\right) \eta-\int_{a_{i}}\left(\int_{O}^{\bar{P}} \omega\right) \eta=\int_{a_{i}}\left(\int_{\bar{P}}^{P} \omega\right) \eta, \tag{2.23}
\end{equation*}
$$

where $\bar{P}$ is the point which is identified with $P$ along the identification of the side $a_{i}$ with the side $-a_{i}$. Recall that the edge of $\mathcal{P}$ bridging $a_{i}$ and $-a_{i}$ is $b_{i}$; let $R$ (resp. $\bar{R}$ ) denote the endpoint connecting $b_{i}$ to $a_{i}$ (resp. $-a_{i}$ ). Then we have

$$
\begin{equation*}
\int_{\bar{P}}^{P} \omega=\int_{\bar{P}}^{\bar{R}} \omega+\int_{\bar{R}}^{R} \omega+\int_{R}^{P} \omega=-\int_{R}^{P} \omega-\int_{b_{i}} \omega+\int_{R}^{P} \omega=-\int_{b_{i}} \omega . \tag{2.24}
\end{equation*}
$$

It follows that we have

$$
\begin{equation*}
\int_{a_{i}} F \eta+\int_{-a_{i}} F \eta=\int_{a_{i}}\left(-\int_{b_{i}} \omega\right) \eta=-\int_{a_{i}} \eta \int_{b_{i}} \omega . \tag{2.25}
\end{equation*}
$$

A similar reasoning (keeping in mind that the edge of $\mathcal{P}$ bridging $b_{i}$ and $-b_{i}$ is $-a_{i}$ ) yields $\int_{b_{i}} F \eta+\int_{-b_{i}} F \eta=\int_{b_{i}} \eta \int_{a_{i}} \omega$. Putting these relations into 2.22 , we get the desired formula.

Corollary 2.6.9. The set of vectors $\left\{\left(\int_{a_{i}} \omega_{1}, \ldots, \int_{a_{i}} \omega_{g}\right)\right\}_{1 \leq i \leq g} \subset \mathbb{C}^{g}$ is linearly independent over $\mathbb{C}$ and thus spans $\mathbb{C}^{g}$.

Proof. The claim is equivalent to the statement that if $\omega \in \Omega(C)$ is a holomorphic differential form with $\int_{a_{i}} \omega=0$ for $1 \leq i \leq g$ then $\omega=0$. Suppose that $\omega \in \Omega(C)$ is a holomorphic differential satisfying this hypothesis. Then, writing $\bar{\omega}$ for the complex conjugate of $\omega$, Lemma 2.6.8 yields

$$
\begin{equation*}
\int_{C} \omega \wedge \bar{\omega}=\sum_{i=1}^{g}\left(\int_{a_{i}} \omega \int_{b_{i}} \bar{\omega}-\int_{b_{i}} \omega \int_{a_{i}} \bar{\omega}\right)=0 . \tag{2.26}
\end{equation*}
$$

This implies that $\omega=0$ as desired, since it is fairly intuitive to verify using local charts that $\int_{C} \omega \wedge \bar{\omega} \neq 0$ if and only if $\omega \neq 0$.

Note that we may use Lemma 2.6 .8 to "normalize" our basis $\left\{\omega_{1}, \ldots, \omega_{g}\right\}$ so that $\int_{a_{j}} \omega_{i}=$ $\delta_{i, j}$ for $1 \leq i, j \leq g$. From now on, we will assume that our basis of holomorphic differentials satisfies this property. We write $\Pi$ for the $g \times g$ matrix given by $\Pi_{i, j}=\int_{b_{j}} \omega_{i}$; we observe that now our lattice $\Lambda$ is generated by $\mathbb{Z}^{g}$ and the columns of $\Pi$, so $\Lambda=\mathbb{Z}^{g}+\Pi \mathbb{Z}^{g}$. We now note that $\Pi$ is symmetric (this is the first of what are called "Riemann's relations").

Corollary 2.6.10. The matrix $\Pi$ is symmetric; that is, $\int_{b_{j}} \omega_{i}=\int_{b_{i}} \omega_{j}$ for $1 \leq i, j \leq g$.
Proof. It is clear that since $\omega_{i}$ and $\omega_{j}$ are both holomorphic and therefore one is a meromorphic function times the other, we have $\omega_{i} \wedge \omega_{j}=0$. Then Lemma 2.6.8 gives us

$$
\begin{equation*}
0=\int_{C} \omega_{i} \wedge \omega_{j}=\sum_{k=1}^{g}\left(\int_{a_{k}} \omega_{i} \int_{b_{k}} \omega_{j}-\int_{a_{k}} \omega_{j} \int_{b_{k}} \omega_{i}\right)=\int_{b_{i}} \omega_{j}-\int_{b_{j}} \omega_{i} . \tag{2.27}
\end{equation*}
$$

Lemma 2.6.11. For any degree-0 divisor of the form $D:=\sum_{P \in C} n_{P}(P) \in \operatorname{Div}^{0}(C)$, there is a differential form $\omega_{D} \in \Omega(C)$ whose only poles are simple poles with residue $n_{P}$ at each point $P \in C$ such that $n_{P} \neq 0$. We may further guarantee that $\omega_{D}$ satisfies $\int_{a_{i}} \omega_{D}=0$ for $1 \leq i \leq g$ (this uniquely determines $\omega_{D}$ ).

Proof. It is clear that since $\operatorname{Div}^{0}(C)$ is generated by divisors of the form $(P)-(Q)$ for points $P, Q \in \operatorname{Div}^{0}(C)$, it suffices to prove this claim for $(P)-(Q)$. It suffices to show the existence of a differential form $\omega \neq 0$ whose only poles are a simple pole at $P$ (with residue 1) and a simple pole at $Q$ (with residue -1 ). One checks that the space of all differential forms whose only poles are possibly simple poles at $P$ and $Q$ is isomorphic to $\mathcal{L}(K+P+Q)$. Putting $D=-P-Q$ into the formula given in Theorem 2.6.2 gives

$$
\begin{equation*}
\operatorname{deg}(-P-Q)=l(-P-Q)-l(K+P+Q)+g-1 \Rightarrow-2=0-l(K+P+Q)+g-1 \tag{2.28}
\end{equation*}
$$

so $l(K+P+Q)=g+1$ and there exists a function in $\mathcal{L}(K+P+Q) \backslash \mathcal{L}(K)$ and thus a differential form $\omega_{D}$ which is not holomorphic but whose only poles are simple poles at $P$ and/or $Q$. Because the sum of the residues of a differential form must be 0 , we know that this $\omega_{D}$ must have simple poles at both $P$ and $Q$ with opposite residues. After scaling by a suitable constant, we may assume that the residues at $P$ and $Q$ are 1 and -1 respectively.

Since Corollary 2.6 .9 tells us that $\left\{\left(\int_{a_{i}} \omega_{1}, \ldots, \int_{a_{i}} \omega_{g}\right)\right\}_{1 \leq i \leq g}$ spans $\mathbb{C}^{g}$, there is some holomorphic differential $\omega \in \Omega(C)$ satisfying $\left(\int_{a_{1}} \omega, \ldots, \int_{a_{g}} \omega\right)=\left(\int_{a_{1}} \omega_{D}, \ldots, \int_{a_{g}} \omega_{D}\right)$. By subtracting this $\omega$ from $\omega_{D}$, we guarantee that $\omega_{D}$ satisfies $\int_{a_{i}} \omega_{D}=0$ for $1 \leq i \leq g$, and uniqueness also follows quickly from Corollary 2.6.9.

Lemma 2.6.12. Choose any $D \in \operatorname{Div}^{0}(C)$ and let $\omega_{D} \in \Omega(C)$ be the corresponding differential form given by Lemma 2.6.11. Then we have the identity

$$
\begin{equation*}
\operatorname{AJ}(D)=\frac{1}{2 \pi \sqrt{-1}}\left(\int_{b_{1}} \omega_{D}, \ldots, \int_{b_{g}} \omega_{D}\right)+\Lambda \in J \tag{2.29}
\end{equation*}
$$

Proof. First note that Lemma 2.6.8 says that for a fixed $j$,

$$
\begin{equation*}
\int_{\partial \mathcal{P}}\left(\int_{O}^{P} \omega_{j}\right) \omega_{D}=\sum_{i=1}^{g}\left(\int_{a_{i}} \omega_{j} \int_{b_{i}} \omega_{D}-\int_{a_{i}} \omega_{D} \int_{b_{i}} \omega_{j}\right)=\int_{b_{j}} \omega_{D} \tag{2.30}
\end{equation*}
$$

where the second inequality follows from the properties $\int_{a_{j}} \omega_{i}=\delta_{i, j}$ and $\int_{a_{i}} \omega_{D}=0$ for $1 \leq i \leq g$. Meanwhile, we may evaluate the integral on the left-hand side by calculating residues, as follows. Note that the function $P \mapsto \int_{O}^{P} \omega_{j}$ has no poles anywhere because $\omega_{j}$ is holomorphic, but if we write $D=\sum_{k=1}^{r} n_{k}\left(P_{k}\right)$ with $n_{k} \neq 0$ then $\omega_{D}$ has only a simple pole at each $P_{k}$ with residue $n_{k}$. The residue formula then yeilds

$$
\begin{equation*}
\int_{\partial \mathcal{P}}\left(\int_{O}^{P} \omega_{j}\right) \omega_{D}=2 \pi \sqrt{-1} \sum n_{k} \int_{O}^{P_{k}} \omega_{j} \tag{2.31}
\end{equation*}
$$

where the paths from $O$ and $P$ lie in the interior of $\mathcal{P}$. The statement of the lemma now follows from $(2.30)$ and $(2.31)$ and the definition of AJ.

Lemma 2.6.13. If a function $f \in \mathcal{M}(C)$ is holomorphic on a loop $a \in H_{1}(C, \mathbb{Z})$, then the value of the integral $\int_{a} d f / f$ is an integer multiple of $2 \pi \sqrt{-1}$.

Proof. This is a standard fact in complex analysis and arises intuitively from the observation that $d f / f=d(\log f)$ with $\log f$ defined up to an integer multiple of $2 \pi \sqrt{-1}$.

We are finally ready to begin the proof of Abel's Theorem. We have to show that a divisor $D \in \operatorname{Div}^{0}(C)$ is in the kernel of AJ if and only if $D$ is the divisor associated to some meromorphic function on $C$. First choose any $f \in \mathcal{M}(C)$ and write $(f)=: D=\sum n_{k}\left(P_{k}\right) \in$ $\operatorname{Div}^{0}(C)$ with integers $n_{k} \neq 0$. Now we easily verify that the only poles of $d f / f \in \Omega(C)$ are simple poles at each $P_{k}$ with residue $n_{k}$. Then $d f / f$ has the same poles with the same residues as the differential $\omega_{D} \in \Omega(C)$ as given by Lemma 2.6.11, so we must have $d f / f=\omega_{D}+\omega$ for some holomorphic differential $\omega$. We write $\frac{1}{2 \pi \sqrt{-1}} \omega=\sum_{i=1}^{g} m_{i} \omega_{i} \in \Omega(C)$ with $m_{i} \in \mathbb{C}$. For each $i$, we get that $2 \pi \sqrt{-1} m_{i}=2 \pi \sqrt{-1} \sum_{j=1}^{g} m_{j} \int_{a_{i}} \omega_{j}=\int_{a_{i}} \omega=\int_{a_{i}} \frac{d f}{f}-\int_{a_{i}} \omega_{D}=$ $\int_{a_{i}} \frac{d f}{f}-0 \in 2 \pi \sqrt{-1} \mathbb{Z}$ by Lemmas 2.6 .11 and 2.6.13. so in fact the $m_{i}$ 's are integers.

Now, using Lemma 2.6.12, we compute that $\operatorname{AJ}(D)$ is equal to

$$
\begin{gather*}
\frac{1}{2 \pi \sqrt{-1}}\left(\int_{b_{1}} \omega_{D}, \ldots, \int_{b_{g}} \omega_{D}\right)+\Lambda=\frac{1}{2 \pi \sqrt{-1}}\left(\int_{b_{1}} \frac{d f}{f}, \ldots, \int_{b_{g}} \frac{d f}{f}\right)-\frac{1}{2 \pi \sqrt{-1}}\left(\int_{b_{1}} \omega, \ldots, \int_{b_{g}} \omega\right)+\Lambda \\
=\frac{1}{2 \pi \sqrt{-1}}\left(\int_{b_{1}} \frac{d f}{f}, \ldots, \int_{b_{g}} \frac{d f}{f}\right)-\sum_{i=1}^{g} m_{i}\left(\int_{b_{1}} \omega_{i}, \ldots, \int_{b_{g}} \omega_{i}\right)+\Lambda \tag{2.32}
\end{gather*}
$$

Note that Lemma 2.6.13 implies that the first term in (2.32) lies in $\mathbb{Z}^{g} \subset \Lambda$. Moreover, the second term in 2.32$)$ is $\sum_{i=1}^{g} m_{i}\left(\Pi_{i, 1}, \ldots, \Pi_{i, g}\right)=\sum_{j=1}^{g} m_{i}\left(\Pi_{1, j}, \ldots, \Pi_{g, j}\right)$ by Corollary 2.6.10. Moreover, since we have $m_{i} \in \mathbb{Z}$ for $1 \leq i \leq g$, the second term lies in $\Lambda$. Therefore, we have $\operatorname{AJ}(D)=0 \in \mathbb{C}^{g} / \Lambda=J$, as desired.

We now prove the converse - if $D=\sum n_{k}\left(P_{k}\right) \in \operatorname{Div}^{0}(C)$ is a divisor such that $\operatorname{AJ}(D)=$ $0 \in J$, then $D$ is the divisor associated to some function $f \in \mathcal{M}(C)$ - by essentially running the above argument backwards. We again define $\omega_{D} \in \Omega(C)$ to be the differential form given by Lemma 2.6.11. We would like to be able to set $f(P)=e^{\int_{O}^{P} \omega_{D}}$, noting that $d f / f=\omega_{D}$ and so $(f)=D$, but this function is not well defined unless the integral $\int_{O}^{P} \omega_{D}$ is well-defined up to an integer multiple of $2 \pi \sqrt{-1}$. We claim that there is a holomorphic differential $\omega \in \Omega(C)$ such that $\int_{a_{i}}\left(\omega_{D}+\omega\right), \int_{b_{i}}\left(\omega_{D}+\omega\right) \in 2 \pi \sqrt{-1} \mathbb{Z}$ for $1 \leq i \leq g$. Indeed, since $\operatorname{AJ}(D)=0$, we have $\left(\sum_{k} n_{k} \int_{O}^{P_{k}} \omega_{1}, \ldots, \sum_{k} n_{k} \int_{O}^{P_{k}} \omega_{g}\right) \in \Lambda$, where each integral is defined along a path lying in the interior of $\mathcal{P}$. Keeping in mind that $\Lambda=\mathbb{Z}^{g}+\Pi \mathbb{Z}^{g}=\mathbb{Z}^{g}+\mathbb{Z}^{g} \Pi$ by Lemma 2.6.10, for each $j$, we may write $\sum_{k} n_{k} \int_{O}^{P_{k}} \omega_{j}$ as $s_{j}-\sum_{i=1}^{g} m_{i} \Pi_{i, j}$ for some integers $s_{j}, m_{1}, \ldots, m_{g} \in \mathbb{Z}$. Now we let $\omega=2 \pi \sqrt{-1} \sum_{i=1}^{g} m_{i} \omega_{i}$. Then for each $j$, we have $\frac{1}{2 \pi \sqrt{-1}} \int_{a_{j}}\left(\omega_{D}+\omega\right)=\frac{1}{2 \pi \sqrt{-1}} \int_{a_{j}} \omega=m_{j} \in \mathbb{Z}$. Moreover, using Lemma 2.6.12, we have
$\frac{1}{2 \pi \sqrt{-1}} \int_{b_{j}}\left(\omega_{D}+\omega\right)=\sum_{k} n_{k} \int_{O}^{P_{k}} \omega_{j}+\sum_{i=1}^{g} m_{i} \int_{b_{j}} \omega_{i}=s_{j}-\sum_{i=1}^{g} m_{i} \Pi_{i, j}+\sum_{i=1}^{g} m_{i} \Pi_{i, j}=s_{j} \in \mathbb{Z}$.

Therefore, we have $\int_{a_{i}}\left(\omega_{D}+\omega\right), \int_{b_{i}}\left(\omega_{D}+\omega\right) \in 2 \pi \sqrt{-1} \mathbb{Z}$, and so the expression $\int_{O}^{P}\left(\omega_{D}+\omega\right)$ is defined up to integer multiples of $2 \pi \sqrt{-1}$. Thus, the meromorphic function $f(P):=$ $e^{\int_{O}^{P}\left(\omega_{D}+\omega\right)}$ is well-defined. Moreover, since $\omega$ is holomorphic, $\omega_{D}+\omega$ has the same poles with the same residues as $\omega_{D}$ does, and it follows that $(f)=D$. Abel's Theorem is proved.

### 2.6.4 Proof of Jacobi's Theorem

We now prove Jacobi's part of Theorem 2.6.7. This proof is much simpler and (in my opinion) a little more intuitive. The idea is that by Corollary 2.6.4, given a basepoint $O \in C$, every divisor in $\operatorname{Div}^{0}(C)$ can be represented by a divisor of the form $\sum_{i=1}^{g}\left(P_{i}\right)-g(O)$. Therefore, Abel's Theorem suggests that it makes sense to study the map AJ ${ }^{(g)}: C^{g} \rightarrow J$ given by

$$
\Psi:\left(P_{1}, \ldots, P_{g}\right) \mapsto \sum_{i=1}^{g}\left(\int_{O}^{P_{i}} \omega_{1}, \ldots, \int_{O}^{P_{i}} \omega_{g}\right)+\Lambda
$$

and to show that it is surjective. (In fact, it is possible to show in the algebraic case that $J$ is a $g$-dimensional variety and $\mathrm{AJ}^{(g)} i$ induces a birational map from the $g$-fold symmetric product of $C$ to $J$, but we will not show this here.)

Lemma 2.6.14. There exist distinct points $Q_{1}, \ldots, Q_{g} \in C$ such that there is no nonzero holomorphic differential form $\omega$ with $(\omega) \geq \sum_{i=1}^{g}\left(Q_{i}\right) \in \operatorname{Div}(C)$.

Proof. Choose a holomorphic differential form $\omega_{0} \in \Omega(C)$, and choose a point $Q_{1} \in C$ at which $\omega_{0}$ does not have a zero. Since by Corollary 2.6.3, the dimension of the space of holomorphic differentials has dimension $g$, it follows that the space of holomorphic differentials $\omega$ with $(\omega) \geq\left(Q_{1}\right) \in \operatorname{Div}(C)$ has dimension strictly less than $g$. But it follows from putting $D=\left(Q_{1}\right)$ into the formula in Theorem 2.6.2 that $l\left(K-\left(Q_{1}\right)\right) \geq g-1$, and thus the dimension of the latter space, which is isomorphic to $\mathcal{L}\left(K-\left(Q_{1}\right)\right)$, is equal to $g-1$. Now choose a holomorphic differential form $\omega_{1} \in \Omega(C)$ with $\left(\omega_{1}\right) \geq\left(Q_{1}\right)$, and choose a point $Q_{2} \in C$ at which $\omega_{1}$ does not have a zero. By a similar argument, we get $l\left(K-\left(Q_{1}\right)-\left(Q_{2}\right)\right)=g-2$. We may repeat this process until we have distinct points $Q_{1}, \ldots, Q_{g}$ such that $l\left(K-\left(Q_{1}\right)-\ldots-\left(Q_{g}\right)\right)=g-g=0$, which implies the desired property.

Lemma 2.6.15. Choose distinct points $Q_{1}, \ldots, Q_{g} \in C$ satisfying the property given in the statement of Lemma 2.6.14. Then there is an open neighborhood $U$ of $\left(Q_{1}, \ldots, Q_{g}\right) \in C^{g}$ such that the map $\Psi: C^{g} \rightarrow J$ given by

$$
\left(P_{1}, \ldots, P_{g}\right) \mapsto \sum_{i=1}^{g}\left(\int_{Q_{i}}^{P_{i}} \omega_{1}, \ldots, \int_{Q_{i}}^{P_{i}} \omega_{g}\right)+\Lambda
$$

when restricted to $U$, is a biholomorphism of $U$ onto an open neighborhood of $0 \in J$.
Proof. Choose an open neighborhood $U^{\prime}$ of $\left(Q_{1}, \ldots, Q_{g}\right) \in C^{g}$ such that each projection of $U^{\prime}$ is simply connected. Now consider the map $\tilde{\Psi}: U^{\prime} \rightarrow \mathbb{C}^{g}$ given by

$$
\left(P_{1}, \ldots, P_{g}\right) \mapsto \sum_{i=1}^{g}\left(\int_{Q_{i}}^{P_{i}} \omega_{1}, \ldots, \int_{Q_{i}}^{P_{i}} \omega_{g}\right)
$$

which is well-defined because of simple connecteness. Let $t_{i}$ be a local parameter of $Q_{i}$ for $1 \leq i \leq g$. Then locally $\tilde{\Psi}=\left(\tilde{\Psi}_{1}, \ldots, \tilde{\Psi}_{g}\right)$ can be expressed as a map $\mathbb{C}^{g} \rightarrow \mathbb{C}^{g}$ given by $\left(t_{i}\right)_{i=1}^{g} \mapsto\left(\tilde{\Psi}_{j}\left(t_{1}, \ldots, t_{g}\right)\right)_{j=1}^{g}$. Thus the Jacobian matrix of the holomorphic function $\Psi$ at the point $\left(Q_{1}, \ldots, Q_{g}\right) \in U^{\prime}$ is given by $\left(d \tilde{\Psi}_{i} / d t_{j}\right)_{1 \leq i, j \leq g}=\left(\omega_{i} / d t_{j}\right)_{1 \leq i, j \leq g}$. This Jacobian matrix can be viewed as a linear map from the space of holomorphic differential forms on $C$ to $\mathbb{C}^{g}$, where each holomorphic $\omega=\sum_{i=1}^{g} m_{i} \omega_{i} \in \Omega(C)$ is considered as the row vector $\left(m_{1}, \ldots, m_{g}\right) \in \mathbb{C}^{g}$ being multiplied to the left of the matrix. So if the kernel of this map were nontrivial, then there would exist a nonzero holomorphic $\omega \in \Omega(C)$ with $\omega / d t_{i}=0$ for $1 \leq i \leq g$. In other words, $\omega$ would have a zero at all points $Q_{i}$, which contradicts the property of the $Q_{i}$ 's given in the hypothesis. Thus, the Jacobian matrix is invertible, and by the Inverse Function Theorem there is a neighborhood $U \subseteq U^{\prime}$ of $\left(Q_{1}, \ldots, Q_{g}\right) \in \mathbb{C}^{g}$ such that the map given in the statement is a biholomorphism of $U$ onto an open neighborhood of $0 \in J$.

Remark 2.6.16. Obviously it is no coincidence that in determining the group structure of what we now call the "Jacobian variety" of a curve, the mathematician Jacobi used what we now call a "Jacobian matrix". Referring simply to "Jacobians" can occasionally cause confusion in my experience, because many mathematicians and physicists often assume that one is talking about the somewhat better-known Jacobian matrices.

Now we can easily prove Jacobi's theorem. We choose distinct points $Q_{1}, \ldots, Q_{g}$ satisfying the property given in Lemma 2.6.14 which we know exist by that lemma. Then it follows from Lemma 2.6 .15 and the above discussion that the map AJ : $\operatorname{Div}^{0}(C) \rightarrow J$ is surjective onto an open neighborhood of $v:=\sum_{i=1}^{g}\left(\int_{O}^{Q_{i}} \omega_{1}, \ldots, \int_{O}^{Q_{i}} \omega_{g}\right)+\Lambda \in J$. Since $v=\operatorname{AJ}\left(\sum_{i=1}^{g}\left(Q_{i}\right)-g(O)\right)$, we see that AJ is surjective onto a translation of that open neighborhood by $-v$, which is an open neighborhood of $0 \in J$ which we denote by $U^{\prime \prime}$. It is easy to see that given any $a \in J$, we have $a / N \in U^{\prime \prime}$ for some large enough integer $N$. Then there is some divisor $D \in \operatorname{Div}^{0}(C)$ with $\operatorname{AJ}(D)=a / N$ and we get $\operatorname{AJ}(N \cdot D)=a$. Therefore, the Abel-Jacobi map AJ is surjective onto all of $J$, and Jacobi's Theorem is proved.

The following easy corollary directly implies that $J$ is actually a complex torus.
Corollary 2.6.17. The lattice $\Lambda \in \mathbb{C}^{g}$ has rank $2 g$.
Proof. Theorem 2.6.7 tells us that any $a \in J$ is the image under AJ of some $[D] \in \operatorname{Pic}^{0}(C)$, and in fact, Corollary 2.6.4 tells us that we can choose a representative of $[D]$ of the form $D=\sum_{i=1}^{g}\left(P_{i}\right)-g(O)$ (here $O$ is the basepoint used to define the map AJ). Thus, the image $A \subset \mathbb{C}^{g}$ of the holomorphic map on $C^{g}$ given by $\left(P_{1}, \ldots, P_{g}\right) \mapsto \sum_{i=1}^{g}\left(\int_{O}^{P_{i}} \omega_{1}, \ldots, \int_{O}^{P_{i}} \omega_{g}\right) \in \mathbb{C}^{g}$, where again we integrate along any path from $O$ to each $P_{i}$ lying inside $\mathcal{P}$, satisfies $A+\Lambda=\mathbb{C}^{g}$. Note that $A$ is the image of the compact space $C$ and is thus a compact subset of $\mathbb{C}^{g}$. The fact that $\mathbb{C}^{g}=A+\Lambda$ for a compact subset $A \subset \mathbb{C}^{g}$ is equivalent to the statement of the corollary.

Remark 2.6.18. a) Note that Theorem 2.6.7 and Corollary 2.6.17 give us a canonical isomorphism $H_{1}(C, \mathbb{Z}) \xrightarrow{\sim} \Lambda=H_{1}(J, \mathbb{Z})$ of $\mathbb{Z}$-lattices of rank $2 g$.
b) It can be shown using Riemann's relations (Corollary 2.6.10 and Exercise 2.7.10) that there is a positive definite Riemann form on $\mathbb{C}^{g}$ with respect to $\Lambda$ induced by the intersection pairing on $H_{1}(C, \mathbb{Z}) \cong \Lambda$ (which is given by $\left\langle a_{i}, a_{j}\right\rangle=\left\langle b_{i}, b_{j}\right\rangle=0$ and $\left\langle a_{i}, b_{i}\right\rangle=-1$ for $1 \leq i<j \leq g)$. It follows that $J$ is in fact a complex abelian variety of dimension $g$. Moreover, since the determinant of the intersection pairing is 1 , we see that $J$ has a canonical principal polarization so that $J$ is self-dual.
c) It is possible also for any smooth, proper curve $C$ over any field $K$ to show that the group $\operatorname{Pic}^{0}(C)$ also has the structure of an abelian variety over $K$ and thus to define the Jacobian of a curve in a completely algebraic setting. This is done by Milne in [4].

### 2.7 Exercises

Exercise 2.7.1. Show that given any alternating form $E: \Lambda \times \Lambda \rightarrow \mathbb{Z}$, there is a function $K: \Lambda \rightarrow \mathbb{Q}$ satisfying $K\left(\lambda_{1}+\lambda_{2}\right)-K\left(\lambda_{1}\right)-K\left(\lambda_{2}\right) \equiv \frac{1}{2} E\left(\lambda_{1}, \lambda_{2}\right)(\bmod \mathbb{Z})$ for $\lambda_{1}, \lambda_{2} \in \Lambda$. (Hint: start by turning $\frac{1}{2} E$ into a symmetric $\mathbb{Z}$-bilinear function on $\Lambda$.)

Exercise 2.7.2. Show that the Riemann form defined at the beginning of $\$ 2.4$ is well-defined regardless of choice of basis of $\Lambda$.

Exercise 2.7.3. By directly using the formula for $\wp$, find the zeroes of the derivative $\wp^{\prime}$. Relate this to the characterization of the 2-torsion subgroup of an elliptic curve coming from Exercise 1.2.3.

Exercise 2.7.4. Show that the subgroup $\operatorname{Div}^{0}(X) \subseteq \operatorname{Div}(X)$ may instead be defined as consisting of all divisors $D \in \operatorname{Div}(X)$ such that $t_{a}^{*} D-D \in \operatorname{Prin}(X)$ for all $a \in X$ (this is how $\operatorname{Div}^{0}(X)$ is defined in a purely algebraic setting).

Exercise 2.7.5. Let $\varphi: X_{1} \rightarrow X_{2}$ be an isogeny of complex abelian varieties. In a natural way, construct an isogeny $\varphi^{\vee}: X_{2}^{\vee} \rightarrow X_{1}^{\vee}$, called the dual isogeny. This shows that taking duals is in some sense a contravariant functor from the category of abelian varieties to itself. (Hint: in fact, the map $\varphi^{\vee}: \operatorname{Pic}^{0}\left(X_{2}\right) \rightarrow \operatorname{Pic}^{0}\left(X_{1}\right)$ is given by pulling back divisors as discussed earlier; the hard part is showing that $\varphi^{\vee}$ is an isogeny. This exercise is rather long and difficult but can be done by unwinding several definitions and constructions.)

Exercise 2.7.6. Show that if $\varphi: X_{1} \rightarrow X_{2}$ is an isogeny of complex elliptic curves, and we identify each elliptic curve with its dual, then the dual isogeny $\varphi^{\vee}: X_{2} \rightarrow X_{1}$ defined in the above exercise is actually the isogeny $\varphi^{\prime}$ guaranteed by (and constructed in the proof of) Proposition 2.5.7.

Exercise 2.7.7. Use the Riemann-Roch formula (Theorem 2.6.2) to verify some of what was shown in Example 2.6.1.

Exercise 2.7.8. Show that the Abel-Jacobi map defined on points of $C$, which was denoted AJ : $C \rightarrow \mathbb{C}^{g} / \Lambda$ in 2.6 .2 is always injective as long as $g \geq 1$. Conclude that AJ is an isomorphism if $g=1$ (this is another way to see that a complex curve which has genus 1 according to the algebraic definition can be realized as $\mathbb{C}$ modulo a lattice).

Exercise 2.7.9. Formulate a shorter argument for one direction of Abel's Theorem - the claim that $\operatorname{Prin}(C) \subseteq \operatorname{ker}(\mathrm{AJ})$ - based on the following ideas. For any $f \in \mathcal{M}(C)$, one may define a holomorphic map $\mathbb{C} \cup\{\infty\} \rightarrow J$ by composing the map $\mathbb{C} \cup\{\infty\} \rightarrow \operatorname{Div}(C)$ given by $z \mapsto f^{*}((z)) \in \operatorname{Div}(C)$ with AJ. First show that this map is constant. Then use this to prove that $\operatorname{AJ}((f))=0 \in J$.

Exercise 2.7.10. Show using Riemann-Roch that a set of points $Q_{1}, \ldots, Q_{g} \in C$ satisfies the property described in Lemma 2.6 .14 if and only if $\sum_{i=1}^{g}\left(Q_{i}\right)-g(O)$ is the "shortest" representative of its class in $\operatorname{Pic}^{0}(C)$; that is, there is no divisor in that class which can be written as $\sum_{i=1}^{g-1}\left(Q_{i}^{\prime}\right)-(g-1)(O)$

Exercise 2.7.11. Use Lemma 2.6 .8 to prove the second of "Riemann's relations", which says that the imaginary part of $\Pi$ is positive definite. (Hint: use the fact that $\frac{-1}{2 i} \int_{C} \omega \wedge \bar{\omega}>0$ for any holomorphic $\omega \neq 0$.)

## Chapter 3

## Elliptic curves and arithmetic

We now turn to the second part of the course, where we focus solely on elliptic curves (rather than abelian varieties) and study them from a purely algebraic point of view. This means that we will no longer assume that a given elliptic curve is defined over the complex numbers, and we will instead specify a field $K$ over which the curve is defined. Of course, if in many cases, such as when $K$ is finitely generated over $\mathbb{Q}$, there is obviously an embedding $K \hookrightarrow \mathbb{C}$ and our elliptic curve can still be defined over $\mathbb{C}$, which at times will allow us to draw from the theory of complex elliptic curves.

### 3.1 Defining equations and the group law

We start this chapter by finally uniting the three definitions of "elliptic curve" given in \$1.1. For the most part, we have already shown (at least in the complex case) that they are equivalent. We sum up most of the relevant facts in the following proposition. Note that we will only prove part (b) below in the case that $K=\mathbb{C}$.

Proposition 3.1.1. a) Given a smooth projective curve $C$ over a field $K$, if the characteristic of $K$ is different from 2 (resp. if $K$ has characteristic 2), the following properties are equivalent (resp. the following properties (i) and (ii) are equivalent):
i) $C$ has genus 1 and possesses a $K$-point $O \in C(K)$;
ii) $C$ can be defined by an equation of the form $y^{2}+a_{1} x y+a_{3} y=x^{3}+a_{2} x^{2}+a_{4} x+a_{6}$ with $a_{i} \in K$ such that there are no singular points; and
iii) $C$ can be defined by an equation of the form $y^{2}=f(x)$ where $f \in K[x]$ is a monic cubic polynomial without multiple roots in $\bar{K}$ (this is called Weierstrass form).

Moreover, if these equivalent conditions hold, then distinguishing a $K$-point $O \in C$ gives $C$ the structure of a 1-dimensional abelian variety over $K$ whose identity element is $O$.
b) Let $A$ be an abelian variety of dimension 1. Then conversely, $A$ is a smooth projective curve of genus 1 and can therefore be defined using an equation of a form given in (a)(ii) or in (a) (iii).

Proof. For part (a), we first note that (iii) immediately implies (ii), while given an equation of the form in (ii), as long as $K$ doesn't have characteristic 2 one can get rid of $a_{1}$ and $a_{3}$ by a simple completing-the-square operation which replaces $y$ by $y-\left(a_{1} x+a_{3}\right) / 2$. Therefore,
(ii) and (iii) are equivalent in the non-characteristic-2 case, and it remains to show that both are equivalent to (i).

Given a smooth projective curve $C$ of genus $g=1$, recall that we showed using the Riemann-Roch theorem (which holds over general fields as well as $\mathbb{C}$ ) that any canonical divisor $K \in \operatorname{Div}(C)$ has degree $2 g-2=0$. It follows from an application of the formula in Theorem 2.6 .2 that $l(n(O))=n$ for $n \geq 1$. Then we may apply the same argument in the proof of Proposition 2.4.1 to show that $C$ may be defined by a polynomial relation in two variables $x$ and $y$ which is a linear combination of $1, x, y, x^{2}, x y, x^{3}, y^{2} \in K[x, y]$, where the coefficients of both $x^{3}$ and $y^{2}$ must be nonzero - we denote them by $r, s \in K^{\times}$respectively. By scaling $x$ and $y$ by $-r s$ and $r^{2} s$ respectively, we get a defining equation of the form given in (ii), and so (i) implies (ii).

Conversely, suppose that $K$ has characteristic different from 2 and $C$ is defined by an equation of the form given in (iii). Then it is straightforward to verify that all holomorphic differentials on $C$ are scalar multiples of $d x / y$ (see the explanation in Example 2.6.1) and therefore the genus is 1 . Thus, (iii) implies (i); one can do a similar computation of the holomorphic differentials to prove that (ii) implies (i) in the characteristic-2 case.

To prove the second statement of (a), we note that in the complex case, a genus 1 curve $C$ is isomorphic via the Abel-Jacobi map AJ (which we define with respect to the chosen point $O$ ) to its Jacobian $J$, which is an 1-dimensional abelian variety (see Exercise 2.7.8). For general $K$, first let $C$ be defined by an equation of the form in (a)(iii). We may assume that the distinguished $K$-point $O$ is the infinity point $\infty$, because otherwise we may "subtract by $O$ " using the operation described in Definition 1.1.4, which is clearly algebraic and defined over $K$. Now, in order to show that $C$ is an abelian variety it suffices to show that this law is in fact a group law as claimed, which we will do below (Proposition 3.1.2).

To prove (b) in the complex case, we use the results in $\$ 2.1$ as well as Proposition 2.4.1, which show that every 1 -dimensional abelian variety is given by $\mathbb{C}$ modulo a lattice, which can be uniformized as a complex curve defined by a cubic polynomial of the form in (a)(ii).

In order to finish proving that our three Definitions 1.1.1, 1.1.3, and 1.1.4 are all equivalent, all we need to do is to show that the law described in the last definition is in fact a group law. Actually, all the group axioms were proven for this law except for the associative property. In order to show this, we will show that if $E$ is a curve of genus 1 over any field $K$ with a distinguished $K$-point $O \in E$, then there is a bijection of sets $E \xrightarrow[\rightarrow]{\sim} \operatorname{Pic}^{0}(E)$ (which in the complex case is just the Abel-Jacobi map) such that the group law on $\operatorname{Pic}^{0}(E)$ can be defined for the points of the curve $E$ as in Definition 1.1.4. Before stating this next proposition, we note that the law given in the definition is essentially equivalent to saying that the sum of three points on the cubic add up to 0 if and only if there is a line in $\mathbb{P}_{K}^{2}$ which intersects $E$ at exactly those three points (counting multiplicities).

Proposition 3.1.2. Let $E$ be a smooth projective curve of genus 1 over a field $K$, and let $O$ be a $K$-point of $E$.
a) The map on $\bar{K}$-points of $E$ given by $P \mapsto[(P)-(O)] \in \operatorname{Pic}^{0}(E)$ defines a bijection of sets $\phi: E(\bar{K}) \xrightarrow{\sim} \operatorname{Pic}^{0}(E)$. In this way, $E(\bar{K})$ inherits a group law from $\operatorname{Pic}^{0}(E)$ in which $O$ is the identity element (which we also denote by $0 \in E(\bar{K})$ ).
b) Now let $E$ be given by an equation of the form in Definition 1.1 .4 and let $O=\infty$. The group law on $E(\bar{K})$ which we obtain as in part (a) can be defined by the following rule: for any points $P, Q, R \in E(\bar{K})$, we have $P+Q+R=0$ if and only if there is a line in $\mathbb{P}_{K}^{2}$ which intersects $E$ at exactly $P, Q$, and $R$ (counting multiplicities).

Proof. We note that in the complex case, (a) is given by Exercise 2.7.8. For general $K$, note first that there cannot be a function $h \in \bar{K}(E)$ with a pole of order 1 at one point $P \in E(\bar{K})$ and no poles anywhere else because then $h: E \rightarrow \mathbb{P}_{K}^{1}$ would be a degree1 morphism which would force $E=\mathbb{P}_{K}^{1}$ (this can also be checked using Riemann-Roch). Therefore $(P)-(\infty) \in \operatorname{Div}^{0}(E)$ is not principal for $P \neq \infty$ and the map is injective. Surjectivity follows from Corollary 2.6.4, which says that every divisor class in $\operatorname{Pic}^{0}(E)$ can be represented by a divisor of the form $(P)-(\infty) \in \operatorname{Div}^{0}(E)$.

Now let $P, Q, R \in E(\bar{K})$ be any points, and assume that $P, Q, R \neq \infty$ (the full argument in the case that some of these points are $\infty$ is very similar). Suppose first that there is a line in $\mathbb{P}_{K}^{2}$ as in the statement of part (b). That is to say, if this line is given by a linear equation $h(x, y)=0$ for some $h \in \bar{K}(E)$, then the linear function $h$ has associated divisor $(P)+(Q)+(R)-3(\infty)$ (note that a line intersects a cubic at 3 points by Bézout's Theorem, so $h$ cannot have any other zeros). Then $0=[(P)+(Q)+(R)-3(\infty)]=\phi(P)+\phi(Q)+\phi(R) \in$ $\operatorname{Pic}^{0}(E)$ and therefore $P+Q+R=0$.

Conversely, suppose that $P+Q+R=0$. Then $[(P)+(Q)+(R)-3(\infty)]=\phi(P)+\phi(Q)+$ $\phi(R)=0 \in \operatorname{Pic}^{0}(E)$, so there is some function $h \in \bar{K}(E)$ with $(h)=(P)+(Q)+(R)-3(\infty)$. Then $h \in \mathcal{L}(3(\infty))$, which has dimension 3 by Riemann-Roch. Since $1, x, y \in \mathcal{L}(3(\infty))$ with no linear relations between them, $\mathcal{L}(3(\infty))=\langle 1, x, y\rangle$ and $h$ is a linear function in $x$ and $y$. Therefore, the line in $\mathbb{P}_{K}^{2}$ given by $h(x, y)=0$ intersects $E$ at exactly the points $P, Q$, and $R$ (counting multiplicities).

### 3.2 Isogenies and endomorphisms of elliptic curves

For this entire section, we make the simplifying assumption that the characteristic of the ground field $K$ is not 2. For a treatment over characteristic 2, see [8, Appendix A].

### 3.2.1 Isomorphisms and the $j$-invariant

We assume for this subsection that every elliptic curve $E$ over $K$ is defined using a model in Weierstrass form $\left(y^{2}=f(x)\right.$ for some cubic polynomial $f \in K[x]$ with no multiple roots in $\bar{K}$ ). We now investigate how to tell from any cubic polynomial $f \in K[x]$ whether the curve given by $y^{2}=f(x)$ is smooth or has a particular type of singularity, and whether or not it is isomorphic to another curve of this form. First we define some important constants associated to a Weierstrass cubic $y^{2}=f(x)$.

Definition 3.2.1. Let $f(x)=x^{3}+b x^{2}+c x+d$ with $b, c, d \in K$. Then the discriminant of the (projective) cubic curve $E$ given by $y^{2}=f(x)$ is defined to be 16 times the discriminant of the cubic $f$ (it's given by $\Delta(E)=16\left(b^{2} c^{2}-4 c^{3}-27 d^{2}+18 b c d\right) \in K$ ). If $\Delta \neq 0$, then the $j$-invariant of $E$ is given by $j(E)=\left(16 b^{2}-48 c\right)^{3} / \Delta(E) \in K$.

Proposition 3.2.2. a) Let $E$ be the projective curve over $K$ given by an equation of the form $y^{2}=f(x)$ as in Definition 3.2.1. Then $E$ is an elliptic curve (i.e. $E$ is smooth) if and only if $\Delta(E) \neq 0$. If $\Delta(E)=0$, then $E$ has a node (resp. a cusp) if the cubic $f$ has a double root and a single root (resp. has a triple root) in $\bar{K}$.
b) If $E_{1}$ and $E_{2}$ are elliptic curves over $K$ given by equations of the form $y^{2}=f_{1}(x)$ and $y^{2}=f_{2}(x)$ respectively, then the following are equivalent:
i) $E_{1}$ and $E_{2}$ are isomorphic over some algebraic extension $K^{\prime} / K$;
ii) $f_{2}(x)=f_{1}\left(u^{2} x+v\right)$ for some elements $u, v$ in some algebraic extension $K^{\prime} / K$; and iii) $j\left(E_{1}\right)=j\left(E_{2}\right)$.
c) For any $j_{0} \in \bar{K}$, there is an elliptic curve $E$ defined over $K\left(j_{0}\right)$ with $j(E)=j_{0}$. (In particular, for any $j_{0} \in K$, there is an elliptic curve $E$ defined over $K$ with $j(E)=j_{0}$.)

Proof. Part (a) can be proven through an elementary computation using algebraic geometry, noting that $\Delta(E)=0$ if and only if the cubic $f(x)$ has a multiple root $\alpha \in \bar{K}$. If this is the case, then the point $(\alpha, 0) \in E(\bar{K})$ is a singular point, and it is a node (resp. a cusp) if $\alpha$ is a double root (resp. a triple root).

For (b), we first check that any isomorphism of elliptic curves $E_{1} \xrightarrow{\sim} E_{2}$ must fix the point at infinity and therefore projects to an affine transformation on the $x$-line, so $x \mapsto u^{\prime} x+v$ for some $u, v \in K^{\prime}$. Since scaling $x$ by $u^{\prime}$ forces $y$ to be scaled by $\pm u^{\prime 3 / 2}$, it is clear that $u^{\prime}=u^{2}$ for some $u \in K^{\prime}$. Now it is straightforward to check that $x \mapsto u^{2} x+v$ multiplies the discriminant of the cubic $f$, and therefore also $\Delta(E)$, by $u^{12}$. However, it also multiplies the expression $\left(16 b^{2}-48 c\right)^{3}$ by $u^{12}$ and therefore doesn't change $j(E)$. Meanwhile, translating $x$ by $v$ changes neither the discriminant $\Delta(E)$ nor the expression $\left(16 b^{2}-48 c\right)^{3}$, so that also doesn't change $j(E)$. Finally, assume that $E_{1}$ and $E_{2}$ two elliptic curves such that $j\left(E_{1}\right)=j\left(E_{2}\right)$. Then if we assume that the characteristic of $K$ is different from 3 (for fields of characteristic 3, this is a separate exercise), we can reduce the cubics in the equations defining $E_{1}$ and $E_{2}$ so that the $x^{2}$-terms vanish, and directly find some $u \in \bar{K}$ such that the $\operatorname{map}(x, y) \mapsto\left(u^{2} x, u^{3} y\right)$ is an isomorphism $E_{1} \xrightarrow{\sim} E_{2}$.

For (c), choose any $j_{0} \in \bar{K}$. Then if $j_{0} \neq 0,1728$, consider the elliptic curve $E_{j_{0}}$ over the field $K\left(j_{0}\right)$ defined by the equation

$$
\begin{equation*}
y^{2}+x y=x^{3}-\frac{36}{j_{0}-1728} x-\frac{1}{j_{0}-1728} . \tag{3.1}
\end{equation*}
$$

Then we check that this equation describes a smooth curve in these cases and compute $j\left(E_{j_{0}}\right)=j_{0}$. If $j_{0}=0$, then we let $E_{0}$ be given by $y^{2}=x^{3}-1$ and check that $j\left(E_{0}\right)=0$. If $j_{0}=1728$, then we let $E_{1728}$ be given by $y^{2}=x^{3}-x$ and check that $j\left(E_{1728}\right)=1728$.

Remark 3.2.3. The reader may wonder why quantities like $\Delta(E)$ and $j(E)$ are defined with "extra" factors like 1728 and $16^{3}$ which actually play no role in the statement of the above proposition nor in its proof. The main explanation is that when considering reduction of elliptic curves over local fields of residue characteristic 2 and 3 , it becomes useful to define these quantities to have certain divisibilities by 2 and by 3 so that results can be stated more elegantly.

### 3.2.2 Isogenies and their duals

We recall the definition of an "isogeny" of abelian varieties which was given as Definition 2.5.1] an isogeny of elliptic curves $E_{1}$ and $E_{2}$ is a surjective morphism $E_{1} \rightarrow E_{2}$ which is a homomorphism of groups. We now present an equivalent definition in the case of elliptic curves.

Definition 3.2.4. An isogeny of elliptic curves $E_{1}$ and $E_{2}$ is a nonconstant morphism $\varphi: E_{1} \rightarrow E_{2}$ with $\varphi(0)=0$.

Proposition 3.2.5. The definition of "isogeny" of elliptic curves given as Definition 2.5.1 is equivalent to Definition 3.2.4.

Proof. Clearly Definition 2.5.1 implies Definition 3.2.4. Now let $\varphi: E_{1} \rightarrow E_{2}$ be a nonconstant morphism satisfying $\varphi(0)=0$. Let $\varphi_{*}: \operatorname{Pic}^{0}\left(E_{1}\right) \rightarrow \operatorname{Pic}^{0}\left(E_{2}\right)$ be the induced homomorphism on divisor classes given by $\left[\sum_{i} n_{i}\left(P_{i}\right)\right] \mapsto\left[\sum_{i} n_{i}\left(\varphi\left(P_{i}\right)\right)\right]$ (one checks that this is well-defined by noting that the pushforward by $\varphi$ of the divisor of a function $f \in K\left(E_{1}\right)$ is the divisor of the norm in $K\left(E_{2}\right)$ of $\left.f\right)$. Then if we let $\phi_{j}: E_{j} \xrightarrow{\sim} \operatorname{Pic}^{0}\left(E_{j}\right)$ be the isomorphism given in Proposition 3.1 .2 for $j=1,2$, it is easy to check that the maps $\varphi: E_{1} \rightarrow E_{2}$ and $\varphi_{*}: \operatorname{Pic}^{0}\left(E_{1}\right) \rightarrow \operatorname{Pic}^{0}\left(E_{2}\right)$ commute with the $\phi_{j}$ 's (using the fact that $\varphi(0)=0$ ). Therefore, $\varphi$ is also a group homomorphism. Moreover, since any nonconstant map between proper curves is surjective, $\varphi$ is a surjection.

Note that one can easily verify using the Hurwitz Formula that every separable isogeny (in particular, every isogeny of degree prime to the characteristic of the ground field) is an unramified morphism (see Exercise 3.3.1).

Proposition 3.2.6. Let $E$ be an elliptic curve over a field $K$, and let $N \subset E(\bar{K})$ be a finite subgroup. Then there is an elliptic curve $E^{\prime}$ defined over some algebraic extension $K^{\prime}$ of $K$ and a separable isogeny $\varphi: E \rightarrow E^{\prime}$ defined over $K^{\prime}$ whose kernel coincides with $N$. In other words, the quotient $E / N$ is an elliptic curve.

Remark 3.2.7. Note that this was already shown for complex elliptic curves in Corollary 2.5.4. In fact, the extension $K^{\prime}$ is the subfield of $\bar{K}$ fixed by the elements of the absolute Galois group of $K$ which stabilize $N \subset E(\bar{K})$, but we will not prove this here.

Proof. For each point $P \in E(\bar{K})$, let $t_{P}: E \rightarrow E$ be the translation-by- $P$ morphism defined over $\bar{K}$ which is given by $Q \mapsto P+Q$, and let $t_{P}^{*}: \bar{K}(E) \rightarrow \bar{K}(E)$ be the pullback of this morphism to an automorphism of the function field. Let $\Phi$ be the (finite) group of automorphisms of $\bar{K}(E)$ consisting of elements $t_{P}^{*}$ for each $P \in N$, and let $\bar{K}(E)^{\Phi}$ be the subfield of $\bar{K}(E)$ fixed by this automorphism group. Then clearly $\bar{K}(E)$ is a finite Galois extension of $\bar{K}(E)^{\Phi}$ with Galois group $\Phi$. Since $\bar{K}(E)^{\Phi}$ also has transcendence degree 1 over $\bar{K}$, there is a smooth curve $E^{\prime}$ over $\bar{K}$ with $\bar{K}\left(E^{\prime}\right)=\bar{K}(E)^{\Phi}$, and the inclusion $\bar{K}\left(E^{\prime}\right) \hookrightarrow \bar{K}(E)$ corresponds to a morphism $\varphi: E \rightarrow E^{\prime}$ (we can always find a subfield $K^{\prime} \subseteq \bar{K}$ finite over $K$ over which $E^{\prime}$ and $\varphi$ are defined). It is easy to check that the inverse image of any point in $E^{\prime}(\bar{K})$ is a coset of the subgroup $N \subset E(\bar{K})$ by showing that $f \circ \varphi \circ t_{P}=f \circ \varphi$ for every $f \in \bar{K}\left(E^{\prime}\right)$. Thus, the morphism $\varphi$ is unramified. It follows
from an application of the Hurwitz Formula that $E^{\prime}$ has genus 1, and therefore, by choosing the distinguished point $\varphi(0) \in E^{\prime}(\bar{K})$, we see that $E^{\prime}$ is an elliptic curve and that $\varphi$ is an isogeny.

As we saw before in $\$ 2.5 .1$, an isogeny $\varphi$ always has finite degree (which we denote by $\operatorname{deg}(\varphi))$ and finite kernel. It is clear from basic algebraic principles that the degree is equal to the order of the kernel if and only if the isogeny is separable (which is always the case when $K$ has characteristic 0 ). More generally, since an isogeny $\varphi: E_{1} \rightarrow E_{2}$ induces an inclusion of function fields $\bar{K}\left(E_{2}\right) \subseteq \bar{K}\left(E_{1}\right)$ which has a maximal separable subextension over which $\bar{K}\left(E_{2}\right)$ is purely inseparable of degree $\operatorname{deg}_{\text {insep }}(\varphi)$, we have $\operatorname{deg}(\varphi)=\operatorname{deg}_{\text {insep }}(\varphi) \cdot \# \operatorname{ker}(\varphi)$.

Proposition 3.2.8. For any elliptic curve $E$ over $K$, the multiplication-by-n map $[n]: E \rightarrow$ $E$ is an isogeny for any integer $n \neq 0$.

Proof. We have to show that $[n]$ does not take everything in $E(\bar{K})$ to $0 \in E(\bar{K})$. Assume now that $K$ does not have characteristic 2 (if it does, then there is a similar argument using a trickier computation). One can check that this is true for $n=2$ using an explicit characterization of the (finite number of) 2-torsion points of $E(\bar{K})$ given by Exercise 1.2.3. Then given a point $P \in E(\bar{K})$ of order 2 , we have $n P=P \neq 0$ for any odd $n$, and the claim follows.

We denote by $\operatorname{Hom}\left(E_{1}, E_{2}\right)$ the set of all isogenies from the elliptic curve $E_{1}$ to the elliptic curve $E_{2}$, along with the trivial map $0: E_{1} \rightarrow E_{2}$. Then it is clear from the abelian group structure on $E_{2}$ that $\operatorname{Hom}\left(E_{1}, E_{2}\right)$ has the structure of a $\mathbb{Z}$-module. We adopt the convention that the degree of the trivial map is 0 .

Proposition 3.2.9. The $\mathbb{Z}$-module $\operatorname{Hom}\left(E_{1}, E_{2}\right)$ defined above is free of rank at most 4 .
Proof (sketch). We first show that the torsion subgroup of $\operatorname{Hom}\left(E_{1}, E_{2}\right)$ is trivial. Indeed, if there exists an isogeny $\varphi: E_{1} \rightarrow E_{2}$ and an integer $n \geq 1$ such that $[n] \circ \varphi=0$, then it follows from the surjectiveness (on $\bar{K}$-points) of $\varphi$ that $[n]$ is the trivial endomorphism on $E_{2}$. But Proposition 3.2 .8 says that this is impossible.

Now the proof that $\operatorname{Hom}\left(E_{1}, E_{2}\right)$ is in fact finitely generated and has rank at most 4 follows from showing that for some prime $\ell$ there is an injection of $\operatorname{Hom}\left(E_{1}, E_{2}\right) \otimes \mathbb{Z}_{\ell}$ into the group of $\mathbb{Z}_{\ell}$-homomorphisms between certain free rank- $2 \mathbb{Z}_{\ell}$-modules associated to $E_{1}$ and $E_{2}$ (this $\mathbb{Z}_{\ell}$-homomorphism group is obviously itself a free rank- $4 \mathbb{Z}_{\ell}$-module) which are called the $\ell$-adic Tate modules of $E_{1}$ and $E_{2}$. For reasons of time, we do not develop this construction in these notes.

In $\S 2.5$, given an isogeny of abelian varieties $\varphi: X_{1} \rightarrow X_{2}$, we came up with two notions of a dual isogeny: the isogeny $\varphi^{\prime}: X_{2} \rightarrow X_{1}$ given by Proposition 2.5.7 and the isogeny $\varphi^{\vee}: X_{2}^{\vee} \rightarrow X_{1}^{\vee}$ constructed in Exercise 2.7.5. One can show (Exercise 2.7.6 that these two notions coincide in the case of complex elliptic curves, and we will now show that dual isogenies can be constructed in the purely algebraic situation as well. Note that this will allow us to say that two elliptic curves are "isogenous" (satisfying an equivalence relation)
if there is an isogeny from one to the other, as we did for complex abelian varieties (see Remark 2.5.8.

Definition 3.2.10. Let $\varphi: E_{1} \rightarrow E_{2}$ be an isogeny of elliptic curves. Then we define the dual isogeny $\varphi^{\vee}: E_{2} \rightarrow E_{1}$ as the map given by $\phi_{1}^{-1} \circ \varphi^{*} \circ \phi_{2}$, where $\varphi^{*}: \operatorname{Pic}^{0}\left(E_{2}\right) \rightarrow$ $\operatorname{Pic}^{0}\left(E_{1}\right)$ is the pullback map on divisor classes induced by $\varphi$, and where $\phi_{j}: E_{j} \xrightarrow{\sim} \operatorname{Pic}^{0}\left(E_{j}\right)$ is the isomorphism given in Proposition 3.1.2 for $j=1,2$.

It might not be immediately obvious from the definition of "dual isogeny" that this map $\varphi^{\vee}$ is actually an isogeny. The next proposition claims this and much more.

Proposition 3.2.11. Let $\varphi: E_{1} \rightarrow E_{2}$ be an isogeny of elliptic curves.
a) We have $\varphi^{\vee} \circ \varphi=[m]: E_{1} \rightarrow E_{1}$ and $\varphi \circ \varphi^{\vee}=[m]: E_{2} \rightarrow E_{2}$, where $m=\operatorname{deg}(\varphi)$.
b) If $\psi: E_{1} \rightarrow E_{2}$ is another isogeny, then $(\varphi+\psi)^{\vee}=\varphi^{\vee}+\psi^{\vee}$ and $(\varphi \circ \psi)^{\vee}=\psi^{\vee} \circ \varphi^{\vee}$.
c) We have $[n]^{\vee}=[n]: E \rightarrow E$ and $\operatorname{deg}([n])=n^{2}$ for any integer $n$.
d) The map $\varphi^{\vee}$ is an isogeny with $\operatorname{deg}\left(\varphi^{\vee}\right)=\operatorname{deg}(\varphi)$.
e) We have $\left(\varphi^{\vee}\right)^{\vee}=\varphi: E_{1} \rightarrow E_{2}$.

Proof. Choose some $Q \in E_{2}(\bar{K})$ and $P_{0} \in \varphi^{-1}(Q)$, and compute

$$
\begin{gathered}
\varphi^{\vee}(Q)=\phi_{1}^{-1}\left(\varphi^{*}([(Q)-(O)])\right)=\phi_{1}^{-1}\left(\left[\operatorname{deg}_{\text {insep }}(\varphi)\left(\sum_{P \in \varphi^{-1}(Q)}(P)-\sum_{T \in \varphi^{-1}(O)}(T)\right)\right]\right) \\
=\phi_{1}^{-1}\left(\operatorname{deg}_{\text {insep }}(\varphi)\left[\sum_{T \in \varphi^{-1}(O)}\left(P_{0}+T\right)-\sum_{T \in \varphi^{-1}(O)}(T)\right]\right)=\phi_{1}^{-1}\left(\operatorname{deg}_{\text {insep }}(\varphi) \sum_{T \in \varphi^{-1}(O)}\left(\phi_{1}\left(P_{0}+T\right)-\phi_{1}\left(P_{0}\right)\right)\right) \\
=\phi_{1}^{-1}\left(\operatorname{deg}_{\text {insep }}(\varphi) \cdot \# \varphi^{-1}(Q) \cdot \phi_{1}\left(P_{0}\right)\right)=\operatorname{deg}(\varphi) \phi_{1}^{-1}\left(\phi_{1}\left(P_{0}\right)\right)=m P_{0} .
\end{gathered}
$$

Therefore $\varphi^{\vee} \circ \varphi=[m]$. Now we verify

$$
[m] \circ \varphi=\varphi \circ[m]=\varphi \circ\left(\varphi^{\vee} \circ \varphi\right)=\left(\varphi \circ \varphi^{\vee}\right) \circ \varphi
$$

from which it follows that $\varphi \circ \varphi^{\vee}=[m]$ since $\varphi$ is surjective, proving (a).
We omit the proof of the first claim in (b) (which is surprisingly difficult), and the second claim follows immediately from the definition of the dual isogeny.

To prove (c), we can see directly that $[0]^{\vee}=[0]$ and $[1]^{\vee}=[1]$, and the fact that $[n]^{\vee}=[n]$ for general $n \geq \mathbb{Z}$ follows from an obvious inductive argument using the additivity given in (b). Now let $m=\operatorname{deg}([n])$ and note that by (a), $[m]=[n]^{\vee} \circ[n]=[n] \circ[n]=\left[n^{2}\right]$. Since then $[0]=[m]-\left[n^{2}\right]=\left[m-n^{2}\right]$ and multiplication by any nonzero integer is a nonconstant map, we have $m-n^{2}=0$ and so $\operatorname{deg}([n])=n^{2}$.

To prove (d), it suffices to show that there always exists an isogeny $\varphi^{\prime}: E_{2} \rightarrow E_{1}$ with the property that $\varphi^{\prime} \circ \varphi=[m]$, from which it will follow from (a) that as maps from $E_{2}$ to $E_{1}, \varphi^{\vee}=\varphi^{\prime}$ and therefore $\varphi^{\vee}$ is an isogeny. If $\varphi$ is separable, this claim follows from an algebraic argument which we omit here concerning inclusions among the function fields $\bar{K}\left(E_{1}\right)$ and $\bar{K}\left(E_{2}\right)$. We will defer the proof of the claim in the case that $\varphi$ is not separable to the next subsection.

Now using (a) and the fact that $\operatorname{deg}\left(\phi^{\vee}\right)=m$ by (d), we have

$$
\varphi \circ[m]=[m] \circ \varphi=\left(\left(\varphi^{\vee}\right)^{\vee} \circ \varphi^{\vee}\right) \circ \varphi=\left(\varphi^{\vee}\right)^{\vee} \circ\left(\varphi^{\vee} \circ \varphi\right)=\left(\varphi^{\vee}\right)^{\vee} \circ[m] .
$$

Since $\left[m\right.$ ] is surjective, it follows that $\varphi=\left(\varphi^{\vee}\right)^{\vee}$, proving (e).

For any elliptic curve $E$ over a field $K$ and integer $n \geq 1$, we denote by $E[n] \subset E[\bar{K}]$ the kernel of the multiplication-by- $n$ map $[n]: E(\bar{K}) \rightarrow E(\bar{K})$.

Corollary 3.2.12. If the characteristic of $K$ does not divide $n$, then $E[n]$ has order $n^{2}$, and we have $E[n] \cong \mathbb{Z} / n \mathbb{Z} \oplus \mathbb{Z} / n \mathbb{Z}$.

Proof. By Proposition 3.2 .11 (c), we have $\operatorname{deg}([n])=n^{2}$, and $\# E[n]=\operatorname{deg}([n])$ since $[n]$ is not divisible by the characteristic of $K$ and therefore separable. Now by decomposing the finite abelian group $E[n]$ into a direct sum of cyclic subgroups and considering that $\# E[d]=d^{2}$ for each $d$ dividing $n$, we see that the only possible group structure for $E[n]$ is given by $\mathbb{Z} / n \mathbb{Z} \oplus \mathbb{Z} / n \mathbb{Z}$.

Corollary 3.2.13. For elliptic curves $E_{1}$ and $E_{2}$ over $K$, the degree map deg : $\operatorname{Hom}\left(E_{1}, E_{2}\right) \rightarrow$ $\mathbb{Z}$ is a positive definite quadratic form.

Proof. We already know that $\operatorname{deg}(\varphi) \geq 0$ for all $\varphi \in \operatorname{Hom}\left(E_{1}, E_{2}\right)$ with equality if and only if $\varphi=0$, and it is clear that $\operatorname{deg}(n \varphi)=n^{2} \operatorname{deg}(\varphi)$ from Proposition 3.2.11(c). It remains to prove that $(\varphi, \psi) \mapsto \operatorname{deg}(\varphi+\psi)-\operatorname{deg}(\varphi)-\operatorname{deg}(\psi)$ is a bilinear form. To do this, we compute $[\operatorname{deg}(\varphi+\psi)-\operatorname{deg}(\varphi)-\operatorname{deg}(\psi)]=(\varphi+\psi)^{\vee} \circ(\varphi+\psi)-\varphi^{\vee} \circ \varphi-\psi^{\vee} \circ \psi=\varphi^{\vee} \circ \psi-\psi^{\vee} \circ \varphi$,
which is bilinear (here we have used (a) as well as the additivity given by (b) in Proposition 3.2.11).

### 3.2.3 The endomorphism ring

Following the notation of $\$ 3.2 .2$, for any elliptic curve $E$, the $\mathbb{Z}$-module $\operatorname{Hom}(E, E)$ clearly has the structure of a ring via composition of elements. We denote this ring by $\operatorname{End}(E)$ and call it the endomorphism ring of $E$. Our goal in this subsection is to investigate the possible structure of this ring. Since each multiplication-by- $n$ map is an endomorphism, there is always a natural inclusion of rings $\mathbb{Z} \subseteq \operatorname{End}(E)$ and we may consider $\operatorname{End}(E)$ as a $\mathbb{Z}$-algebra. In order to find all possible ring structures of $\operatorname{End}(E)$, we recall the following facts that were shown in $\$ 3.2 .2$ (i) as an additive group, $\operatorname{End}(E)$ is a free $\mathbb{Z}$-module of rank at most 4; (ii) there is an anti-involution $(\cdot)^{\vee}: \operatorname{End}(E) \rightarrow \operatorname{End}(E)$ given by $\varphi \mapsto \varphi^{\vee}$ which preserves addition and reverses the order of multiplication and (iii) $\varphi^{\vee} \varphi=\varphi \varphi^{\vee}=: m \in \mathbb{Z}_{\geq 0}$ for any $\varphi \in \operatorname{End}(E)$ with $m=0$ if and only if $\varphi=0$.

The following lemma is purely a statement in abstract algebra, and we leave it as Exercise 3.3.4.

Lemma 3.2.14. Suppose that $R$ is a $\mathbb{Z}$-algebra satisfying the following properties:
(i) as an additive group, $R$ is a free $\mathbb{Z}$-module of rank at most 4;
(ii) there is an anti-involution $(\cdot)^{\vee}: R \rightarrow R$ fixing $\mathbb{Z} \subseteq R$ and satisfying $(\varphi+\psi)^{\vee}=$ $\varphi^{\vee}+\psi^{\vee},(\varphi \psi)^{\vee}=\psi^{\vee} \varphi^{\vee}$, and $\left(\varphi^{\vee}\right)^{\vee}=\varphi$ for all $\varphi, \psi \in R$; and
(iii) $\varphi^{\vee} \varphi=\varphi \varphi^{\vee}=: m \in \mathbb{Z}_{\geq 0}$ with $m=0$ if and only if $\varphi=0$.

Then we have $R=\mathbb{Z}$; or $R \otimes \mathbb{Q}$ is a quadratic imaginary number field; or $R \otimes \mathbb{Q}$ is a definite quaternion algebra over $\mathbb{Q}$.

Proposition 3.2.15. Let $E$ be an elliptic curve over a field $K$. Then the $\mathbb{Z}$-algebra $\operatorname{End}(E)$ either coincides with $\mathbb{Z}$, is an order in an imaginary quadratic number field, or is an order in a definite quaternion algebra. Moreover, if $K$ has characteristic 0, $\operatorname{End}(E)$ must be commutative (and therefore must be $\mathbb{Z}$ or an order in an imaginary quadratic field).

Proof. The first claim comes from the observations summarized in the discussion above along with Lemma 3.2.14. Now if $K$ has characteristic 0 , we may assume $K \subseteq \mathbb{C}$ by the Lefschetz Principle (see [8, Chapter VI, §6]), and it suffices to show that $\operatorname{End}(E)$ is commutative for any elliptic curve $E$ over $\mathbb{C}$. Indeed, it follows from Proposition 2.5.2 that any endomorphism $\varphi: E \rightarrow E$ can be lifted uniquely to a $\mathbb{C}$-linear map $\tilde{\varphi}: \mathbb{C} \rightarrow \mathbb{C}$ which is an isomorphism if and only if $\varphi \neq 0$. Obviously $\tilde{\varphi}$ is just a homothety $z \mapsto a z$ for some $a \in \mathbb{C}$ with $a=0$ if and only if $\varphi=0$, so we have an injection $\operatorname{End}(E) \hookrightarrow \mathbb{C}$ and $\operatorname{End}(E)$ is commutative.

In general, elliptic curves over fields of characteristic 0 have endomorphism rings equal to $\mathbb{Z}$, and it is considered to be exceptional if $\operatorname{End}(E)$ is instead an order in an imaginary quadratic field. In this case, we say that $E$ "has complex multiplication" or is "a CM elliptic curve".

Example 3.2.16. Let $E$ be the elliptic curve over any field $K$ of characteristic different from 2 given by the equation $y^{2}=x^{3}-x$. Then we check that the map $\varphi: E \rightarrow E$ given by $(x, y) \mapsto(-x, i y)$ is an endomorphism of $E$ by noting that it is a morphism with $\varphi(0)=0$. Now we check that $\varphi^{2}: E \rightarrow E$, which is given by $(x, y) \mapsto(x,-y)$, is in fact the map $[-1]: E \rightarrow E$. Therefore, we may identify $\varphi$ with a square root of -1 in the endomorphism ring, and it follows from Proposition 3.2.15 that $\operatorname{End}(E) \cong \mathbb{Z}[i]$ and $E$ has complex multiplication. In fact, it is possible to show that over $\mathbb{C}$, this elliptic curve $E$ is isomorphic to $\mathbb{C} /\langle 1, i\rangle$, from which it is apparent that the endomorphism ring is identified with the ring of homotheties $\mathbb{Z}[i] \subset \mathbb{C}$.

### 3.2.4 The endomorphism ring in positive characteristic

For this subsection, we assume that $E$ is an elliptic curve defined over a field $K$ of characteristic $p>0$. For simplicity, we will assume that $K$ is perfect, although most of the below results would hold anyway with slight modifications. Our main goal is to investigate the possible structures of $\operatorname{End}(E)$ in this case. First it is essential to define the Frobenius isogeny.

Definition 3.2.17. For any integer $r \geq 0$, let $E^{(r)}$ be the elliptic curve defined over $K$ by the equation given by raising all the coefficients of the equation defining $E$ to the $p^{r}$ th
power. Then the (rth power) Frobenius isogeny $\mathrm{Fr}^{r}: E \rightarrow E^{(r)}$ is the map given by $(x, y) \mapsto\left(x^{p^{r}}, y^{p^{r}}\right)$. In the case that $K=\mathbb{F}_{p^{r}}$, we have $E=E^{(r)}$ and $\operatorname{Fr}^{r} \in \operatorname{End}(E)$ is called the Frobenius endomorphism.

Note that each $\operatorname{Fr}^{r}$ is obviously a nonconstant morphism with $\operatorname{Fr}^{r}(0)=0$ and therefore an isogeny; moreover, as the notation suggests, each $\mathrm{Fr}^{r}$ is the $r$ th power of the Frobenius map $\mathrm{Fr}:=\mathrm{Fr}^{1}: E \rightarrow E^{(1)}$. It is clear that each Frobenius isogeny $\mathrm{Fr}^{r}$ is purely inseparable of degree $p^{r}$ from checking that $\left(\mathrm{Fr}^{r}\right)^{*}: K\left(E^{(r)}\right) \rightarrow K(E)$ is the inclusion $K(E)^{p^{r}} \hookrightarrow K(E)$. It is also clear from basic field theory that if $E^{\prime}$ is another elliptic curve over $K$, then every isogeny $\varphi: E \rightarrow E^{\prime}$ factors as the composition of the purely inseparable map $\mathrm{Fr}^{r}: E \rightarrow E^{(r)}$ for some $r \geq 0$ and a separable map $\varphi^{\text {sep }}: E^{(r)} \rightarrow E^{\prime}$. In other words, $\mathrm{Fr}^{r}$ is the inseparable part $\varphi^{\text {insep }}$ of $\varphi$.

The next result functions as a lemma both to Proposition 3.2.11 above (it almost immediately implies the statement of part (d) for the inseparable case) and to Proposition 3.2.19 below.

Lemma 3.2.18. The multiplication-by-p map $[p]: E \rightarrow E$ is inseparable, and there exists an isogeny $\varphi^{\prime}: E^{(1)} \rightarrow E$ such that $\varphi^{\prime} \circ \mathrm{Fr}=\mathrm{Fr} \circ \varphi^{\prime}=[p]$.

Proof. It can be shown either through direct computations or through a more abstract argument using functoriality that for any elliptic curve over a field of any characteristic, if $\omega \in \Omega(E)$ is a holomorphic differential, we have $[n]^{*}(\omega)=n \omega$ for $n \geq \mathbb{Z}$. It follows that since $K$ has characteristic $p$, we have $[p]^{*}(\omega)=0$, and therefore, $[p]: E \rightarrow E$ is inseparable. By the above discussion, we much have $[p]=[p]^{\text {sep }} \circ \mathrm{Fr}^{r}$ for some $r \geq 1$ and some separable $\operatorname{map}[p]^{\text {sep }}: E^{(r)} \rightarrow E$. The claim then follows by taking $\varphi^{\prime}=[p]^{\text {sep }} \circ \operatorname{Fr}_{r-1}$.

One can see from the proof of Proposition 3.2 .11 (d) that the isogeny $\varphi^{\prime}$ above is the dual isogeny $\mathrm{Fr}^{\vee}$. Our main tool in proving the following proposition will be studying this dual isogeny.

Proposition 3.2.19. a) The following are equivalent:
i) $E\left[p^{r}\right]=0$ for any (every) $r \geq 1$ (i.e. the map $[p]: E \rightarrow E$ is purely inseparable);
ii) $\mathrm{Fr}^{\vee}$ is (purely) inseparable;
iii) $\operatorname{End}(E)$ is an order in a quaternion algebra.

Moreover, these statements also imply
iv) $j(E) \in \mathbb{F}_{p^{2}}$.
b) If the above equivalent conditions do not hold, then $E\left[p^{r}\right] \cong \mathbb{Z} / p^{r} \mathbb{Z}$ for all $r \geq 0$, and if $j(E) \in \overline{\mathbb{F}}_{p}$, then $\operatorname{End}(E)$ is an order in an imaginary quadratic field.

Remark 3.2.20. If the equivalent conditions in part (a) above do not hold and $j(E)$ is transcendental over $\mathbb{F}_{p}$, then it is possible to show that in $\operatorname{fact} \operatorname{End}(E)=\mathbb{Z}$.

Proof. We first show that (i) and (ii) are equivalent. It is immediately clear that $E[p]=0$ implies that $E\left[p^{r}\right]=0$ for every $r \geq 1$. Now Lemma 3.2.18 and the discussion below it show that $[p]: E \rightarrow E$ factors as the composition $\mathrm{Fr}^{\vee} \circ \mathrm{Fr}$. Note that Fr obviously has degree $p$ and so $\mathrm{Fr}^{\vee}$ also has degree $p$ by Proposition 3.2.11(d). Moreover, the kernel of Fr is trivial
since Fr is purely inseparable, so the kernel of $[p]$ is trivial if and only if $\mathrm{Fr}^{\vee}$ is also purely inseparable.

Now we show that (ii) implies (iv). Assume that $\mathrm{Fr}^{\vee}$ is inseparable. Then the inseparable part of $\mathrm{Fr}^{\vee}$ is some power of the Frobebius map Fr : $E^{(1)} \rightarrow E^{(2)}$. By comparing degrees, we see that the inseparable part of $\mathrm{Fr}^{\vee}$ is this Frobenius map, while the separable part is an isomorphism $E^{(2)} \xrightarrow{\sim} E$. We note from the formula for the $j$-invariant that $j\left(E^{(2)}\right)=j(E)^{p^{2}}$. But $j(E)=j\left(E^{(2)}\right)$ by Proposition 3.2 .2 (b), so $j(E)=j(E)^{p^{2}}$ and it follows that $j(E) \in \mathbb{F}_{p^{2}}$.

Now assume (ii), and we will show (iii) by assuming that $\operatorname{End}(E)$ is commutative and deriving a contradiction. If $\operatorname{End}(E)$ is commutative, then by Proposition 3.2.15 we have that either $\operatorname{End}(E)=\mathbb{Z}$ or $\operatorname{End}(E)$ is an order in an imaginary quadratic field. One can show that for any isogeny $\varphi: E \rightarrow E^{\prime}$, we have $\operatorname{End}\left(E^{\prime}\right) \otimes \mathbb{Q}=\operatorname{End}(E) \otimes \mathbb{Q}$ (see Exercise 3.3.5) and moreover that inseparability of $[p]: E \rightarrow E$ implies inseparability of $[p]: E^{\prime} \rightarrow E^{\prime}$ (by checking separable and inseparable degrees of $[p] \circ \varphi=\varphi \circ[p])$. We therefore get $j\left(E^{\prime}\right) \in \mathbb{F}_{p^{2}}$ since (ii) implies (iv). It follows (using Proposition 3.2.2(b) and the fact that $\mathbb{F}_{p^{2}}$ is finite) that there are only finitely many isomorphism classes of elliptic curves over $\bar{K}$ which are isogenous to $E$.

We claim that there exists a prime integer $\ell$ such that the ideal $(\ell)$ remains prime in $\operatorname{End}\left(E^{\prime}\right)$ for every elliptic curve $E^{\prime}$ isogenous to $E$. If $\operatorname{End}(E)=\mathbb{Z}$, then each $\operatorname{End}\left(E^{\prime}\right)$ is also $\mathbb{Z}$ and then the claim is trivial (we can take any $\ell$ ). If $\operatorname{End}(E)$ is an order in an imaginary quadratic field, then each $\operatorname{End}\left(E^{\prime}\right)$ is an order in that same imaginary quadratic field. Then it is an elementary exercise in algebraic number theory to show that there always exists a prime integer $\ell$ such that $(\ell)$ is a prime ideal in all of a given finite set of orders in a fixed imaginary quadratic field. Note that this $\ell$ is different from $p$, since $(p)$ decomposes in $\operatorname{End}(E)$ as $\left(\mathrm{Fr}^{\vee}\right)(\mathrm{Fr})$.

Choose such a prime $\ell$, and note from Corollary 3.2 .12 that $E[\ell] \cong \mathbb{Z} / \ell \mathbb{Z} \oplus \mathbb{Z} / \ell \mathbb{Z}$. Thus, we may choose a sequence of subgroups $0=N_{0} \subset N_{1} \subset N_{2} \subset \ldots \subset E(\bar{K})$ with $N_{i} \cong \mathbb{Z} / \ell^{i} \mathbb{Z}$ for $i \geq 0$, with none of these subgroups containing $E\left[\ell^{r}\right]$ for any $r \geq 1$. By Proposition 3.2.6, we see that this sequence of subgroups corresponds to a sequence of isogenies

$$
E=E_{0} \xrightarrow{\varphi_{7}} E_{1} \xrightarrow{\varphi_{2}} E_{2} \xrightarrow{\varphi_{3}} \ldots
$$

with $E_{i} \cong E / N_{i}$ and $\operatorname{ker}\left(\varphi_{i}\right) \cong N_{i} / N_{i-1}$ for each $i \geq 0$. But as was observed above, there are only finitely many elliptic curves isogenous to $E$ up to isomorphism, so for some $i, j \geq 1$ we have $E_{i+j} \cong E_{i}$ and so the composition $\varphi_{i+j} \circ \ldots \circ \varphi_{i+1}: E_{i} \rightarrow E_{i+j} \cong E_{i}$ is an endomorphism $\psi \in \operatorname{End}\left(E_{i}\right)$ with kernel isomorphic to $N_{i+j} / N_{i}$. Since $\operatorname{deg}(\psi)=\ell^{j}$ and $(\ell)$ is prime in $\operatorname{End}\left(E^{\prime}\right)$, we have that $j$ is even and $\psi$ is the composition of $\left[\ell^{j / 2}\right]$ with some isomorphism. But this implies that $E\left[\ell^{j / 2}\right] \subseteq N_{i+j}$, a contradiction.

We finish the proof of (a) by showing that (iii) implies (ii). To do this, we assume that (ii) is false (i.e. $\mathrm{Fr}^{\vee}$ is separable) and show that then (iii) is false (i.e. $\operatorname{End}(E)$ is commutative). Indeed, since $\mathrm{Fr}^{\vee}$ is separable and has order $p$ while Fr is purely inseparable, we have $\# E\left[p^{r}\right]=\# \operatorname{ker}\left(\left(\mathrm{Fr}^{\vee}\right)^{r}\right) \cdot \# \operatorname{ker}\left(\mathrm{Fr}^{r}\right)=p^{r}$ for any $r \geq 0$. By a similar argument to the one used to prove Corollary 3.2.12, we see that the only possible group structure for $E\left[p^{r}\right]$ is given by $\mathbb{Z} / p^{r} \mathbb{Z}$.

We now claim that for any nonzero endomorphism $\varphi \in \operatorname{End}(E)$, there is some $r \geq 1$ such that $\varphi\left(E\left[p^{r}\right]\right) \neq 0$. Indeed, since $\varphi$ has finite kernel, we get $\varphi\left(E\left[p^{r}\right]\right) \neq 0$ for any $r$
such that $p^{r}>\# \operatorname{ker}(\varphi)$. It follows that the obvious map from $\operatorname{End}(E)$ to the ring of $\mathbb{Z}_{p^{-}}$ endomorphisms of $\lim _{\leftarrow r} E\left[p^{r}\right]$ is injective. (Here $\mathbb{Z}_{p}$ denotes the ring of $p$-adic integers, while the inverse limit above is taken with respect to multiplication-by- $p$ maps $E\left[p^{r}\right] \rightarrow E\left[p^{r-1}\right]$ and is obviously a $\mathbb{Z}_{p}$-module.) But since $E\left[p^{r}\right] \cong \mathbb{Z} / p^{r} \mathbb{Z}$, we have $\lim _{\leftarrow r} E\left[p^{r}\right] \cong \mathbb{Z}_{p}$ and so $\operatorname{End}_{\mathbb{Z}_{p}}\left(\lim _{\leftarrow r} E\left[p^{r}\right]\right) \cong \mathbb{Z}_{p}$ is commutative, so $\operatorname{End}(E)$ is commutative, as desired.

To prove (b), we assume that the equivalent conditions (i), (ii), (iii) in (a) do not hold. It was already shown above that then $E\left[p^{r}\right] \cong \mathbb{Z} / p^{r} \mathbb{Z}$ for $r \geq 0$. Now we assume in addition that $j(E) \in \overline{\mathbb{F}}_{p}$ and proceed to show that $\operatorname{End}(E) \supsetneq \mathbb{Z}$. Clearly there is some $r \geq 1$ such that $j(E) \in \mathbb{F}_{p^{r}}$, and so it follows from Proposition 3.2.2(c) that we may assume that $E$ is defined over $\mathbb{F}_{p^{r}}$. Then we have the Frobenius endomorphism $\operatorname{Fr}^{r} \in \operatorname{End}(E)$. Suppose that $\operatorname{Fr}^{r}=[n] \in \mathbb{Z}$. Then since $p^{r}=\operatorname{deg}\left(\operatorname{Fr}^{r}\right)=\operatorname{deg}([n])=n^{2}$ by Proposition 3.2.11(c), we have that $r$ is even and $n= \pm p^{r / 2}$. Since Fr is purely inseparable, this implies that $\left[p^{r / 2}\right]$ and hence also $[p]$ are purely inseparable, which contradicts condition (i) of (a). Therefore, $\operatorname{Fr}^{r} \notin \mathbb{Z}$ and $\operatorname{End}(E) \supsetneq \mathbb{Z}$.

The fact that any elliptic curve $E$ satisfying the equivalent conditions in part (a) of the above proposition must have $j(E) \in \mathbb{F}_{p^{r}}$ suggests that such elliptic curves are rather exceptional (since $\mathbb{F}_{p^{r}}$ is a small subfield of any algebraically closed field of characteristic $p>0)$. This is reflected in the terminology defined as follows.

Definition 3.2.21. If $E$ is an elliptic curve over a field of charaacteristic $p>0$, then $E$ is said to be supersingular if $E$ satisfies the equivalent conditions in Proposition 3.2.19(a) and is said to be ordinary otherwise.

We end by noting that the notion of supersingularity has nothing to do with singularities in the context of algebraic geometry (after all, every elliptic curve is nonsingular). This potentially confusing terminology came about because the word singular was commonly used to mean "unusual" or "special", and certainly supersingular elliptic curves are unusual among elliptic curves over fields of characteristic $p$ since their $j$-invariants lie in the small field $\mathbb{F}_{p^{2}}$.

### 3.3 Exercises

For all of the following exercises, we retain the assumption that $K$ is not a field of characteristic 2 .

Exercise 3.3.1. Let $C_{1}$ and $C_{2}$ be smooth curves of genus $g_{1}$ and $g_{2}$ respectively and $\varphi$ : $C_{1} \rightarrow C_{2}$ is a nonconstant separable morphism of degree $d$, and assume that the characteristic of $K$ doesn't divide any of the ramification indices. Recall the Hurwitz Formula

$$
2 g_{1}-2=d\left(2 g_{2}-2\right)+\sum_{P \in C_{1}(\bar{K})}\left(e_{P}-1\right),
$$

where $e_{P}$ is the ramification index of $\varphi$ at the point $P$ (note that this is equal to 1 for all but finitely many $P \in C_{1}(\bar{K})$ ).
a) Suppose that $C_{1}$ is an elliptic curve. Show that $C_{2}$ can also be given the structure of an elliptic curve if and only if $\varphi$ is unramified eveerywhere.
b) Show that every elliptic curve $E$ over $K$ has a degree- 2 morphism $\varphi: E \rightarrow \mathbb{P}_{K}^{1}$ which is ramified at exactly 4 points. Show furthermore that these four points are the elements of the 2-torsion subgroup $E[2]$.
c) Prove the converse statement that if $\varphi: C \rightarrow \mathbb{P}_{K}^{1}$ is a morphism of degree 2 which is ramified at exactly 4 points, then $C$ is isomorphic over $\bar{K}$ to an elliptic curve $E$ over $K$ such that these four ramification points correspond to the elements of $E[2]$. (Hint: choose one of the ramification points to be the distinguished point $O \in E(K)$ and study the only nontrivial automorphism $\iota: C \rightarrow C$ satisfying $\varphi \circ \iota=\varphi$. It might be easier to assume that $K$ does not have characteristic 3 , which allows us to reduce the cubic polynomial defining $E$ by getting rid of the $x^{2}$-term.)

Exercise 3.3.2. Let $E$ is an elliptic curve over $K$ given by an equation of the form $y^{2}=$ $f(x)=\left(x-\alpha_{1}\right)\left(x-\alpha_{2}\right)\left(x-\alpha_{3}\right)$ for distinct $\alpha_{1}, \alpha_{2}, \alpha_{3} \in \bar{K}$. For $n \geq 1$, write $K(E[n])$ for the (algebraic) extension generated over $K$ by coordinates of all of the points in $E[n] \subset E(\bar{K})$.
a) Show that elements of the absolute Galois group $\operatorname{Gal}(\bar{K} / K)$ permute the points in each $E[n]$. Conclude that each field extension $K(E[n]) / K$ is Galois.
b) Show that $K(E[2])=K\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right)$, so that the Galois group of the extension $K(E[2]) / K$ is canonically isomorphic to the Galois group of the cubic polynomial $f \in K[x]$.
c) Compute $K(E[4])=K\left(\alpha_{1}, \alpha_{2}, \alpha_{3}, \sqrt{\alpha_{1}-\alpha_{2}}, \sqrt{\alpha_{1}-\alpha_{3}}, \sqrt{\alpha_{2}-\alpha_{3}}\right)$. (This is difficult but can be proven either by direct computations using doubling formulas, by directly considering divisor classes in $\operatorname{Pic}^{0}(E)$, or by formulas for computing 2-descent given in [8, §X.1]. For formulas for generators of $K(E[8]) / K$, see my article [9]!)

Exercise 3.3.3. Let $\beta, \gamma \in \bar{K}$ be distinct and nonzero and consider the elliptic curves $E$ and $E^{\prime}$ over $\bar{K}$ given by the following Weierstrass equations:

$$
E: y^{2}=x(x-\beta)(x-\gamma), \quad E^{\prime}: y^{2}=x^{3}+2(\beta+\gamma) x^{2}+(\beta-\gamma)^{2} x
$$

Let $\varphi: E \rightarrow E^{\prime}$ be the morphism over $K$ given by $(x, y) \mapsto\left(\frac{y^{2}}{x^{2}}, y\left(\frac{\beta \gamma}{x^{2}}-1\right)\right)$.
a) Show that $\varphi$ is an isogeny of degree 2 whose kernel is cyclically generated by $(0,0) \in$ $E[2]$.
b) Define the elliptic curve $E^{\prime \prime}$ and the isogeny $\varphi^{\prime}: E^{\prime} \rightarrow E^{\prime \prime}$ in the same way that $E^{\prime}$ and $\varphi$ were defined respectively, except with respect to $\beta^{\prime}:=-(\beta+\gamma)+2 \sqrt{\beta \gamma}$ and $\gamma^{\prime}:=-(\beta+\gamma)-2 \sqrt{\beta \gamma}$. Check that the kernel of $\varphi^{\prime} \circ \varphi$ is $E[2]$ and that there is an isomorphism $E^{\prime \prime} \xrightarrow{\sim} E$ defined over $K$. Conclude that this isomorphism composed with $\varphi^{\prime}$ is the dual isogeny $\varphi^{\vee}$ (at least up to an automorphism of $E$ ).
c) Use these formulas to construct an elliptic curve $E$ over $\mathbb{Q}$ with $\operatorname{End}(E) \cong \mathbb{Z}[\sqrt{-2}]$.

Exercise 3.3.4. Prove Lemma 3.2 .14 (this is a lengthy but elementary algebraic argument). As an immediate corollary, show that if $\operatorname{End}(E)$ is commutative, then the dual of any endomorphism $\varphi \in \operatorname{End}(E)$ is its complex conjugate.

Exercise 3.3.5. If $\varphi: E_{1} \rightarrow E_{2}$ is an isogeny of elliptic curves, show that there is a canonical $\mathbb{Q}$-algebra isomorphism $\operatorname{End}\left(E_{1}\right) \otimes \mathbb{Q} \xrightarrow{\sim} \operatorname{End}\left(E_{2}\right) \otimes \mathbb{Q}$. (Hint: use the dual isogeny.)

Exercise 3.3.6. Assume for this exercise that $K$ doesn't have characteristic 3.
a) Prove that the automorphism group $\operatorname{Aut}(E)$ of an elliptic curve $E$ over $K$ is always finite and cyclic and of order 2,4 , or 6 . Show that there is only one elliptic curve $E$ up to isomorphism with $\operatorname{Aut}(E) \cong \mathbb{Z} / 4 \mathbb{Z}$ and show the same for $\operatorname{Aut}(E) \cong \mathbb{Z} / 6 \mathbb{Z}$. (Hint: refer to the proof of Proposition 3.2 .2 (c).)
b) Note that the absolute Galois $\operatorname{group} \operatorname{Gal}(\bar{K} / K)$ acts on $\operatorname{Aut}(E)$. Show that as a $\operatorname{Gal}(\bar{K} / K)$-module, $\operatorname{Aut}(E)$ is isomorphic to the multiplicative group $\left\langle\zeta_{d}\right\rangle \subset K^{\times}$where $d$ is its order and $\zeta_{d}$ is a primitive $d$ th root of unity. Check that we have $\varphi^{\vee}=\varphi^{-1}$ for any $\varphi \in \operatorname{Aut}(E)$.

Exercise 3.3.7. Let $E$ be an elliptic curve over a finite field $\mathbb{F}_{p^{r}}$. Clearly the group $E\left(\mathbb{F}_{p^{r}}\right)$ is finite, and it is natural to ask how large it is. The purpose of this exercise is to prove the Hasse Bound, which says

$$
\left|\# E\left(\mathbb{F}_{p^{r}}\right)-p^{r}-1\right| \leq 2 p^{r / 2}
$$

You may assume without proof that the endomorphism $\operatorname{Fr}^{r}-1 \in \operatorname{End}(E)$ is separable (but this can be proven by direct computation using the invariant differential).
a) Show that $\# E\left(\mathbb{F}_{p^{r}}\right)=\operatorname{deg}\left(\mathrm{Fr}^{r}-1\right)$.
b) Verify using Corollary 3.2 .13 that for any endomorphisms $\varphi, \psi \in \operatorname{End}(E)$, we have $|\operatorname{deg}(\varphi-\psi)-\operatorname{deg}(\varphi)-\operatorname{deg}(\psi)| \leq 2 \sqrt{\operatorname{deg}(\varphi) \operatorname{deg}(\psi)}$.
c) Conclude that the Hasse Bound holds.

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