# A FORMULA FOR THE JACOBIAN OF A GENUS ONE CURVE OF ARBITRARY DEGREE 

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#### Abstract

We extend the formulae of classical invariant theory for the Jacobian of a genus one curve of degree $n \leq 4$ to curves of arbitrary degree. To do this, we associate to each genus one normal curve of degree $n$, an $n \times n$ alternating matrix of quadratic forms in $n$ variables, that represents the invariant differential. We then exhibit the invariants we need as homogeneous polynomials of degrees 4 and 6 in the coefficients of the entries of this matrix.


## Introduction

Let $C$ be a smooth curve of genus one defined over a field $K$. Its Jacobian is an elliptic curve $E$ defined over the same field $K$. However it is only if $C$ has a $K$ rational point that $C$ and $E$ are isomorphic over $K$. Starting with equations for $C$ we would like to compute a Weierstrass equation for $E$.

Let $D$ be a $K$-rational divisor on $C$ of degree $n \geq 1$. It is natural to split into cases according to the value of $n$. If $n=1$ then $C$ has a $K$-rational point, and our task is that of writing an elliptic curve in Weierstrass form. If $n \geq 2$ then the complete linear system $|D|$ defines a morphism $C \rightarrow \mathbb{P}^{n-1}$. Explicitly, the map is given by $\left(f_{1}: \ldots: f_{n}\right)$ where $f_{1}, \ldots, f_{n}$ are a basis for the Riemann-Roch space $\mathcal{L}(D)$. If $n=2$ then $C$ is a double cover of $\mathbb{P}^{1}$ and is given by an equation of the form $y^{2}=F\left(x_{1}, x_{2}\right)$ where $F$ is a binary quartic. In this case Weil [W1], [W2] showed that the classical invariants of the binary quartic $F$ give a formula for the Jacobian.

If $n \geq 3$ then the morphism $C \rightarrow \mathbb{P}^{n-1}$ is an embedding. The image is a genus one normal curve of degree $n$. The word normal refers to the fact $C$ is projectively normal (see for example [H, Proposition IV.1.2]), i.e. if $H$ is the divisor of a hyperplane section then the natural map

$$
\begin{equation*}
S^{d} \mathcal{L}(H) \rightarrow \mathcal{L}(d H) \tag{1}
\end{equation*}
$$

is surjective for all $d \geq 1$. If $n=3$ then $C \subset \mathbb{P}^{2}$ is a smooth plane cubic, say with equation $F\left(x_{1}, x_{2}, x_{3}\right)=0$. The invariants of a ternary cubic $F$ were computed by Aronhold [A], and again Weil (in the notes to [W1] in his collected papers) showed that these give a formula for the Jacobian. If $n=4$ then $C \subset \mathbb{P}^{3}$ is the complete intersection of two quadrics. If we represent these quadrics by $4 \times 4$ symmetric matrices $A$ and $B$, then $F\left(x_{1}, x_{2}\right)=\operatorname{det}\left(A x_{1}+B x_{2}\right)$ is a binary quartic. The invariants of this binary quartic again give a formula for the Jacobian. For further details of these formulae in the cases $n=2,3,4$ see [ $\mathrm{AKM}^{3} \mathrm{P}$ ], [ARVT] or [F1].

If $n=5$ then $C \subset \mathbb{P}^{4}$ is no longer a complete intersection, and indeed the homogeneous ideal is generated by 5 quadrics. The Buchsbaum-Eisenbud structure theorem [BE1], [BE2] shows that these quadrics may be written as the $4 \times 4$ Pfaffians of a $5 \times 5$ alternating matrix of linear forms. The space of all such matrices is a 50 -dimensional

[^0]affine space, with a natural action of $\mathrm{GL}_{5} \times \mathrm{GL}_{5}$. In [F1] we computed generators for the ring of invariants and showed that they again give a formula for the Jacobian. In fact the invariants are too large to write down as explicit polynomials, so instead we gave a practical algorithm for evaluating them (based in part on the case $n=5$ of Proposition 9.3). More recently, B. Gross [G] gave a uniform description of the invariants in the cases $n=2,3,4,5$, using results of Vinberg, although this does not appear to give any way of evaluating the invariants in the case $n=5$.

In this paper we extend these formulae for the Jacobian to genus one normal curves of arbitrary degree.

Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve of degree $n \geq 3$. Since $C$ has genus one, the space of regular differentials on $C$ has dimension 1 , say spanned by $\omega$. We call $\omega$ an invariant differential, since geometrically it is invariant under all translation maps. There is a linear map

$$
\begin{equation*}
\wedge^{2} \mathcal{L}(H) \rightarrow \mathcal{L}(2 H) ; \quad f \wedge g \mapsto \frac{f d g-g d f}{\omega} \tag{2}
\end{equation*}
$$

Since (1) is surjective for $d=2$, we may represent this map by an $n \times n$ alternating matrix of quadratic forms in $x_{1}, \ldots, x_{n}$. This matrix $\Omega$ represents $\omega$ in the sense that

$$
\omega=\frac{x_{j}^{2} d\left(x_{i} / x_{j}\right)}{\Omega_{i j}\left(x_{1}, \ldots, x_{n}\right)} \quad \text { for all } i \neq j
$$

However if $n \geq 4$ then there are quadrics vanishing on $C \subset \mathbb{P}^{n-1}$ and so this description does not determine $\Omega$ uniquely. Nonetheless we show, by proving [F3, Conjecture 7.4], that there is a canonical choice of $\Omega$. We then define polynomials $c_{4}$ and $c_{6}$ of degrees 4 and 6 in the coefficients of the entries of $\Omega$, and show that the Jacobian has Weierstrass equation

$$
y^{2}=x^{3}-27 c_{4}(\Omega) x-54 c_{6}(\Omega) .
$$

These main results are stated in Section 1. In the next two sections we show that $c_{4}$ and $c_{6}$ are invariants for the appropriate action of $\mathrm{GL}_{n}$, and that they reduce to the previously known formulae for $n \leq 5$. At this point the proof of our results for any given value of $n$ is a finite calculation. However finding a proof that works for all $n$ is more challenging.

In Section 4 we show that if we can find a matrix $\Omega$ satisfying some apparently weaker hypotheses, then it will satisfy the properties claimed in Theorem 1.1. For the actual construction of $\Omega$ in Section 5 we reduce to the case $C$ is an elliptic curve $E$ embedded in $\mathbb{P}^{n-1}$ via the complete linear system $\left|n .0_{E}\right|$. At first we specify $\Omega$ as a linear map $\wedge^{2} \mathcal{L}\left(n .0_{E}\right) \rightarrow S^{2} \mathcal{L}\left(n .0_{E}\right)$, and use this in Section 6 to complete the proof of Theorem 1.1. Then in Section 7 we make a specific choice of basis for $\mathcal{L}\left(n .0_{E}\right)$, so that $\Omega$ becomes an alternating matrix of quadratic forms. We compute this matrix explicitly and, in Section 8 , prove the formula for the Jacobian by computing $c_{4}(\Omega)$ and $c_{6}(\Omega)$. Much of the work here is in checking that the invariants $c_{4}$ and $c_{6}$ are scaled correctly for all $n$.

The description of $\Omega$ in Theorem 1.1 involves higher secant varieties. We quote any general results we need about these as required. Proofs, or references to the literature, are then given in Section 9.

In future work we plan to study the space of all matrices $\Omega$. This appears to be defined by $d_{1}+d_{2}$ quadrics in $\mathbb{P}^{N-1}$ where $N=\left(n^{2}-1\right)\left(n^{2}-4\right) / 4$ and

$$
\begin{aligned}
& d_{1}=\left(n^{2}-1\right)\left(n^{2}-4\right)\left(n^{2}-9\right) / 36 \\
& d_{2}=\left(n^{2}-1\right)^{2}\left(n^{2}-9\right) / 9
\end{aligned}
$$

The numbers $N, d_{1}$ and $d_{2}$ are dimensions of irreducible representations for $\mathrm{GL}_{n}$. Moreover, as suggested by Manjul Bhargava, we expect that $d_{2}$ of the quadrics can be explained by an associative law, similar to that used in [B, Section 4].

We work throughout over a field $K$ of characteristic 0 , although it would in fact be sufficient that the characteristic is not too small compared to $n$. Except at the end of Section 1, where we give the application to computing Jacobians, we will assume that $K$ is algebraically closed. For a projective variety $X$ we write $I(X)$ for its homogeneous ideal, and $T_{P} X$ for the tangent space at $P \in X$. A Magma script containing some of the formulae in this paper is available from the author's website.

## 1. Statement of Results

Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve of degree $n \geq 3$. For any integer $r \geq 1$ the $r$ th higher secant variety $\operatorname{Sec}^{r} C$ is the Zariski closure of the locus of all $(r-1)$-planes through $r$ points on $C$. For example, if $r=1$ then $\operatorname{Sec}^{1} C=C$. The codimension of $\operatorname{Sec}^{r} C$ in $\mathbb{P}^{n-1}$ is $\max (n-2 r, 0)$. So according as $n$ is odd or even there is a higher secant variety of codimension 1 or 2 . If $n=2 r+1$ then $\operatorname{Sec}^{r} C$ is a hypersurface of degree $n$, whereas if $n=2 r+2$ then $\operatorname{Sec}^{r} C$ is the complete intersection of two forms of degree $r+1$. In Section 9 we give references for these facts about higher secant varieties, and also explain how to compute equations for $\mathrm{Sec}^{r} C$ from equations for $C$.

We give the polynomial ring $R=K\left[x_{1}, \ldots, x_{n}\right]$ its usual grading by degree, say $R=\oplus_{d} R_{d}$, and write $R(d)$ for the graded $R$-module with $e$ th graded piece $R_{d+e}$. Maps between graded free $R$-modules are required to have relative degree 0 , and are labelled by the matrices of forms that represent them. Our first main result is

Theorem 1.1. Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve of degree $n \geq 3$.
(i) If $n$ is odd, say $n=2 r+1$, and $\operatorname{Sec}^{r} C=\{F=0\}$ then there is a minimal free resolution

$$
0 \longrightarrow R(-2 n) \xrightarrow{\nabla^{T}} R(-n-1)^{n} \xrightarrow{\Omega} R(-n+1)^{n} \xrightarrow{\nabla} R
$$

where $\Omega$ is an $n \times n$ alternating matrix of quadratic forms and

$$
\nabla=\nabla(F)=\left(\begin{array}{lll}
\frac{\partial F}{\partial x_{1}} & \cdots & \frac{\partial F}{\partial x_{n}}
\end{array}\right) .
$$

(ii) If $n$ is even, say $n=2 r+2$, and $\operatorname{Sec}^{r} C=\left\{F_{1}=F_{2}=0\right\}$ then there is a minimal free resolution

$$
0 \longrightarrow R(-n)^{2} \xrightarrow{\nabla^{T}} R\left(\frac{-n-2}{2}\right)^{n} \xrightarrow{\Omega} R\left(\frac{-n+2}{2}\right)^{n} \xrightarrow{\nabla} R^{2}
$$

where $\Omega$ is an $n \times n$ alternating matrix of quadratic forms and

$$
\nabla=\nabla\left(F_{1}, F_{2}\right)=\left(\begin{array}{lll}
\frac{\partial F_{1}}{\partial x_{1}} & \cdots & \frac{\partial F_{1}}{\partial x_{n}} \\
\frac{\partial F_{2}}{\partial x_{1}} & \cdots & \frac{\partial F_{2}}{\partial x_{n}}
\end{array}\right)
$$

We remarked in [F3, Section 7] that Theorem 1.1(i) follows from the BuchsbaumEisenbud structure theorem for Gorenstein ideals of codimension 3. In this paper we give a different proof, not only so that it runs in parallel with our proof of Theorem 1.1(ii), but also because this is needed for the proof of Theorem 1.2.

If the matrix $\Omega$ exists then, by the uniqueness of minimal free resolutions (see for example [ E , Section 20.1] or [ P, Section 7]), it is uniquely determined up to scalars. Moreover starting from equations for $\operatorname{Sec}^{r} C$ we can solve for $\Omega$ by linear algebra. The details are very similar to those in [F4, Section 4].

Let $\Omega=\left(\Omega_{i j}\right)$ be as specified in Theorem 1.1. We put

$$
\begin{equation*}
M_{i j}=\sum_{r, s=1}^{n} \frac{\partial \Omega_{i r}}{\partial x_{s}} \frac{\partial \Omega_{j s}}{\partial x_{r}} \quad \text { and } \quad N_{i j k}=\sum_{r=1}^{n} \frac{\partial M_{i j}}{\partial x_{r}} \Omega_{r k} \tag{3}
\end{equation*}
$$

We then define

$$
\begin{equation*}
c_{4}(\Omega)=\frac{3(n-2)^{2}}{2^{4} n\binom{n+3}{5}} \sum_{i, j, r, s=1}^{n} \frac{\partial^{2} M_{i j}}{\partial x_{r} \partial x_{s}} \frac{\partial^{2} M_{r s}}{\partial x_{i} \partial x_{j}} \tag{4}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{6}(\Omega)=\frac{-(n-2)^{3}}{2^{6} n\binom{n+5}{7}} \sum_{i, j, k, r, s, t=1}^{n} \frac{\partial^{3} N_{i j k}}{\partial x_{r} \partial x_{s} \partial x_{t}} \frac{\partial^{3} N_{r s t}}{\partial x_{i} \partial x_{j} \partial x_{k}} . \tag{5}
\end{equation*}
$$

Let $C_{1}$ and $C_{2}$ be genus one curves with invariant differentials $\omega_{1}$ and $\omega_{2}$. An isomorphism $\gamma:\left(C_{1}, \omega_{1}\right) \rightarrow\left(C_{2}, \omega_{2}\right)$ is an isomorphism of curves $\gamma: C_{1} \rightarrow C_{2}$ with $\gamma^{*} \omega_{2}=\omega_{1}$.

Theorem 1.2. Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve of degree $n \geq 3$, and let $\Omega$ be an alternating matrix of quadratic forms as specified in Theorem 1.1. Then
(i) There is an invariant differential $\omega$ on $C$ such that

$$
\omega=\frac{x_{j}^{2} d\left(x_{i} / x_{j}\right)}{\Omega_{i j}\left(x_{1}, \ldots, x_{n}\right)} \quad \text { for all } i \neq j .
$$

(ii) The pair $(C, \omega)$ is isomorphic (over $K=\bar{K}$ ) to

$$
\left(y^{2}=x^{3}-27 c_{4}(\Omega) x-54 c_{6}(\Omega), 3 d x / y\right)
$$

The following corollary gives the application of Theorem 1.2 to computing Jacobians. For this result only we drop our assumption that $K$ is algebraically closed.

Corollary 1.3. Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve defined over a field $K$. Suppose we scale the matrix $\Omega$ in Theorem 1.1 so that the coefficients of its entries are in K. Then $C$ has Jacobian elliptic curve $y^{2}=x^{3}-27 c_{4}(\Omega) x-54 c_{6}(\Omega)$.

Proof: Let $E$ be the elliptic curve $y^{2}=x^{3}-27 c_{4}(\Omega) x-54 c_{6}(\Omega)$. By Theorem 1.2 there is an isomorphism $\gamma: C \rightarrow E$ with $\gamma^{*}(3 d x / y)=\omega$. Let $\xi_{\sigma}=\sigma(\gamma) \gamma^{-1}$ for $\sigma \in \operatorname{Gal}(\bar{K} / K)$. Since $3 d x / y$ and $\omega$ are both $K$-rational it follows that $\xi_{\sigma}^{*}(3 d x / y)=$ $3 d x / y$. This implies, as explained for example in [F1, Lemma 2.4], that $\xi_{\sigma}: E \rightarrow E$ is a translation map. Then $C$ is the twist of $E$ by the class of $\left\{\xi_{\sigma}\right\}$ in $H^{1}(K, E)$. It follows by Theorems 3.6 and 3.8 in [S, Chapter X$]$ that $C$ is a principal homogeneous space under $E$, and $E$ is the Jacobian of $C$.

Remark 1.4. Although we will not need it for the proofs of Theorems 1.1 and 1.2, it is natural to ask whether $C \subset \mathbb{P}^{n-1}$ is uniquely determined by $\Omega$. The answer is that it is. Indeed by the minimal free resolutions in Theorem 1.1 we can recover $\nabla$ from $\Omega$. Then by Euler's identity we obtain equations for $\operatorname{Sec}^{r} C$ where $n-2 r=1$ or 2 . This then determines $\operatorname{Sec}^{1} C=C$ by Theorem 9.1(v).

## 2. Changes of co-ordinates

We show that the constructions in Section 1 behave well under all changes of coordinates. First we define an action of $\mathrm{GL}_{n}$ on the space of all $n \times n$ alternating matrices of quadratic forms in $x_{1}, \ldots x_{n}$. For $g \in \mathrm{GL}_{n}$ we put

$$
g \star \Omega=g^{-T}\left(\Omega\left(\sum_{i=1}^{n} g_{i 1} x_{i}, \ldots, \sum_{i=1}^{n} g_{i n} x_{i}\right)\right) g^{-1}
$$

where $g^{-T}$ is the inverse transpose of $g$. Since the scalar matrices act trivially, this could equally be viewed as an action of $\mathrm{PGL}_{n}$.
Lemma 2.1. Let $C \subset \mathbb{P}^{n-1}$ and $C^{\prime} \subset \mathbb{P}^{n-1}$ be genus one normal curves. Let $\Omega$ and $\Omega^{\prime}$ be alternating matrices of quadratic forms that satisfy the conclusions of Theorem 1.1, and define invariant differentials $\omega$ and $\omega^{\prime}$ on $C$ and $C^{\prime}$. If $\gamma: C^{\prime} \rightarrow C$ is an isomorphism given by

$$
\left(x_{1}: \ldots: x_{n}\right) \mapsto\left(\sum_{i=1}^{n} g_{i 1} x_{i}: \ldots: \sum_{i=1}^{n} g_{i n} x_{i}\right)
$$

for some $g \in \mathrm{GL}_{n}$ then there exists $\lambda \in K^{\times}$such that $g \star \Omega=\lambda \Omega^{\prime}$ and $\gamma^{*} \omega=\lambda^{-1} \omega^{\prime}$.
Proof: Suppose $n$ is odd, say $n=2 r+1$ and $\operatorname{Sec}^{r} C=\{F=0\}$. Then $\operatorname{Sec}^{r} C^{\prime}$ is defined by

$$
F^{\prime}\left(x_{1}, \ldots, x_{n}\right)=F\left(y_{1}, \ldots, y_{n}\right)
$$

where $y_{j}=\sum_{i=1}^{n} g_{i j} x_{i}$. By the chain rule

$$
\nabla\left(F^{\prime}\right)\left(x_{1}, \ldots, x_{n}\right)=\nabla(F)\left(y_{1}, \ldots, y_{n}\right) g^{T}
$$

Then

$$
\nabla(F) \Omega=0 \Longrightarrow \nabla\left(F^{\prime}\right)(g \star \Omega)=0 .
$$

It follows by the uniqueness of minimal free resolutions that $g \star \Omega=\lambda \Omega^{\prime}$ for some $\lambda \in K^{\times}$. The case $n$ is even is similar.

We also have $\gamma^{*} \omega=\mu \omega^{\prime}$ for some $\mu \in K^{\times}$. If $y_{j}=\sum_{i=1}^{n} g_{i j} x_{i}$ then

$$
y_{s}^{2} d\left(y_{r} / y_{s}\right)=\sum_{i, j=1}^{n} g_{i r} g_{j s} x_{j}^{2} d\left(x_{i} / x_{j}\right)
$$

Dividing by $\gamma^{*} \omega=\mu \omega^{\prime}$ gives

$$
\Omega\left(y_{1}, \ldots, y_{n}\right)=\mu^{-1} g^{T} \Omega^{\prime}\left(x_{1}, \ldots, x_{n}\right) g
$$

Hence $g \star \Omega=\mu^{-1} \Omega^{\prime}$ and so $\mu=\lambda^{-1}$.
Lemma 2.2. The polynomials $c_{4}$ and $c_{6}$ are invariants for the action of $\mathrm{GL}_{n}$, i.e. $c_{4}(g \star \Omega)=c_{4}(\Omega)$ and $c_{6}(g \star \Omega)=c_{6}(\Omega)$ for all $g \in \mathrm{GL}_{n}$.

Proof: Let $\Omega^{\prime}=g \star \Omega$, i.e.

$$
\Omega_{i j}^{\prime}\left(x_{1}, \ldots, x_{n}\right)=\sum_{a, b=1}^{n}\left(g^{-1}\right)_{a i}\left(g^{-1}\right)_{b j} \Omega_{a b}\left(y_{1}, \ldots, y_{n}\right)
$$

where $y_{j}=\sum_{i=1}^{n} g_{i j} x_{i}$. Direct calculation using (3) shows that

$$
\begin{aligned}
M_{i j}^{\prime}\left(x_{1}, \ldots, x_{n}\right) & =\sum_{a, b=1}^{n}\left(g^{-1}\right)_{a i}\left(g^{-1}\right)_{b j} M_{a b}\left(y_{1}, \ldots, y_{n}\right), \\
N_{i j k}^{\prime}\left(x_{1}, \ldots, x_{n}\right) & =\sum_{a, b, c=1}^{n}\left(g^{-1}\right)_{a i}\left(g^{-1}\right)_{b j}\left(g^{-1}\right)_{c k} N_{a b c}\left(y_{1}, \ldots, y_{n}\right) .
\end{aligned}
$$

Then

$$
\begin{aligned}
\frac{\partial^{2} M_{i j}^{\prime}}{\partial x_{r} \partial x_{s}} & =\sum_{a, b, c, d=1}^{n}\left(g^{-1}\right)_{a i}\left(g^{-1}\right)_{b j} g_{r c} g_{s d} \frac{\partial^{2} M_{a b}}{\partial x_{c} \partial x_{d}}, \\
\frac{\partial^{2} M_{r s}^{\prime}}{\partial x_{i} \partial x_{j}} & =\sum_{A, B, C, D=1}^{n}\left(g^{-1}\right)_{C r}\left(g^{-1}\right)_{D s} g_{i A} g_{j B} \frac{\partial^{2} M_{C D}}{\partial x_{A} \partial x_{B}} .
\end{aligned}
$$

Multiplying these together and summing gives

$$
\sum_{i, j, r, s=1}^{n} \frac{\partial^{2} M_{i j}^{\prime}}{\partial x_{r} \partial x_{s}} \frac{\partial^{2} M_{r s}^{\prime}}{\partial x_{i} \partial x_{j}}=\sum_{a, b, c, d=1}^{n} \frac{\partial^{2} M_{a b}}{\partial x_{c} \partial x_{d}} \frac{\partial^{2} M_{c d}}{\partial x_{a} \partial x_{b}} .
$$

Thus $c_{4}\left(\Omega^{\prime}\right)=c_{4}(\Omega)$. A similar argument shows that $c_{6}\left(\Omega^{\prime}\right)=c_{6}(\Omega)$.
The following corollary shows that to prove Theorems 1.1 and 1.2 for a fixed value of $n$, it suffices to prove them for a family of curves covering the $j$-line.

Corollary 2.3. Let $\Omega_{1}$ and $\Omega_{2}$ correspond to pairs $\left(C_{1}, \omega_{1}\right)$ and $\left(C_{2}, \omega_{2}\right)$. If there is an isomorphism $\gamma: C_{1} \rightarrow C_{2}$ with $\gamma^{*} \omega_{2}=\lambda \omega_{1}$ then $c_{4}\left(\Omega_{1}\right)=\lambda^{4} c_{4}\left(\Omega_{2}\right)$ and $c_{6}\left(\Omega_{1}\right)=\lambda^{6} c_{6}\left(\Omega_{2}\right)$.

Proof: Let $C_{1}$ and $C_{2}$ have hyperplane sections $H_{1}$ and $H_{2}$. Then $H_{1}$ and $\gamma^{*} H_{2}$ are degree $n$ divisors on $C_{1}$. After composing the isomorphism $\gamma$ with a translation map, we may suppose (see [S, III.3.5]) that $H_{1} \sim \gamma^{*} H_{2}$. Then $\gamma$ is given by a change of co-ordinates on $\mathbb{P}^{n-1}$. The case $\lambda=1$ is immediate from Lemmas 2.1 and 2.2. In general we use that $c_{4}$ and $c_{6}$ are homogeneous polynomials of degrees 4 and 6 .

## 3. Curves of small degree

We compare our general formula for the Jacobian with the formulae previously known for genus one normal curves of degrees 3,4 and 5 .

For curves of degrees 3 and 4 it is easy to write down a matrix $\Omega$ satisfying the conclusions of Theorems 1.1 and 1.2(i). Indeed for $C=\left\{F\left(x_{1}, x_{2}, x_{3}\right)=0\right\} \subset \mathbb{P}^{2}$ a plane cubic we put

$$
\Omega=\left(\begin{array}{ccc}
0 & \frac{\partial F}{\partial x_{3}} & -\frac{\partial F}{\partial x_{2}} \\
-\frac{\partial F}{\partial x_{3}} & 0 & \frac{\partial F}{\partial x_{1}} \\
\frac{\partial F}{\partial x_{2}} & -\frac{\partial F}{\partial x_{1}} & 0
\end{array}\right),
$$

and for $C=\left\{F_{1}=F_{2}=0\right\} \subset \mathbb{P}^{3}$ a quadric intersection we let $\Omega$ be the $4 \times 4$ alternating matrix with entries

$$
\Omega_{i j}=\frac{\partial F_{1}}{\partial x_{k}} \frac{\partial F_{2}}{\partial x_{l}}-\frac{\partial F_{1}}{\partial x_{l}} \frac{\partial F_{2}}{\partial x_{k}},
$$

where $(i, j, k, l)$ is an even permutation of $(1,2,3,4)$. To prove Theorem 1.2(ii) in these cases we may check by direct computation that $c_{4}(\Omega)$ and $c_{6}(\Omega)$ are the classical invariants of a ternary cubic or quadric intersection, as scaled in [F1, Section 7]. We note that these are polynomials of degrees 4 and 6 in the coefficients of $F$, respectively of degrees 8 and 12 in the coefficients of $F_{1}$ and $F_{2}$.

As described for example in [F4, Section 4], a genus one normal curve of degree $n=5$ is defined by the $4 \times 4$ Pfaffians $p_{1}, \ldots, p_{5}$ of a $5 \times 5$ alternating matrix of linear forms on $\mathbb{P}^{4}$. We call the matrix of linear forms $\Phi$ a genus one model of degree 5 , and note that there is a natural action of $\mathrm{GL}_{5} \times \mathrm{GL}_{5}$ on the space of all such models. It is shown in [H, Proposition VIII.2.5] that the secant variety $\mathrm{Sec}^{2} C$ is a hypersurface of degree 5 with equation $F=0$ where $F$ is the determinant of the Jacobian matrix of $p_{1}, \ldots, p_{5}$. In [F3, Section 7] we proved that there is a degree 5 covariant $\Omega$ satisfying the conclusions of Theorems 1.1 and 1.2(i). We gave an explicit formula for this covariant in [FS, Section 2].

We claim that $c_{4}(\Omega)$ and $c_{6}(\Omega)$ are invariants for the action of $\mathrm{SL}_{5} \times \mathrm{SL}_{5}$. For the action of $\mathrm{SL}_{5}$ via changes of co-ordinates on $\mathbb{P}^{4}$ this follows from Lemma 2.2. For the action of $\mathrm{SL}_{5}$ via $\Phi \mapsto A \Phi A^{T}$ it turns out that the coefficients of the entries of $\Omega$ are already invariants. Since $\Omega$ is a covariant of degree 5 , the invariants $c_{4}(\Omega)$ and $c_{6}(\Omega)$ have degrees 20 and 30 in the coefficients of the entries of $\Phi$. Computing a single numerical example (to check the scaling) shows that $c_{4}(\Omega)$ and $c_{6}(\Omega)$ are the same as the invariants $c_{4}(\Phi)$ and $c_{6}(\Phi)$ constructed in [F1].

## 4. Minimal free resolutions

Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve of degree $n \geq 3$. Let $\Omega$ be an $n \times n$ alternating matrix of quadratic forms in $x_{1}, \ldots, x_{n}$. In Sections 5 and 6 we exhibit $\Omega$ satisfying the following three hypotheses.
(H1) If $n-2 r \geq 1$ and $f \in I\left(\operatorname{Sec}^{r} C\right)$ then $\sum_{i=1}^{n} \frac{\partial f}{\partial x_{i}} \Omega_{i j} \in I\left(\operatorname{Sec}^{r} C\right)$ for all $j$.
(H2) If $n-2 r=2$ and $\operatorname{Sec}^{r} C=\left\{F_{1}=F_{2}=0\right\}$ then $\sum_{i, j=1}^{n} \frac{\partial F_{1}}{\partial x_{i}} \Omega_{i j} \frac{\partial F_{2}}{\partial x_{j}}=0$.
(H3) If $n-2 r \geq 1$ then there exists $P \in \operatorname{Sec}^{r} C$ with $\operatorname{rank} \Omega(P)=2 r$.
In this section we prove:
Theorem 4.1. Let $\Omega$ be an $n \times n$ alternating matrix of quadratic forms, satisfying the hypotheses $(\mathbf{H} \mathbf{1}),(\mathbf{H} \mathbf{2})$ and $(\mathbf{H} 3)$. Then there is a minimal free resolution as described in Theorem 1.1.

The next two propositions are proved in Section 9. By abuse of notation we write $P$ both for a point in $\mathbb{P}^{n-1}$ and for a vector of length $n$ representing this point.

Proposition 4.2. If $n-2 r \geq 1$ and $P=\sum_{i=1}^{r} \xi_{i} P_{i}$ for some $P_{1}, \ldots, P_{r} \in C$ distinct and $\xi_{1}, \ldots, \xi_{r} \neq 0$ then the tangent space $T_{P} \operatorname{Sec}^{r} C$ is the linear span of the tangent lines $T_{P_{1}} C, \ldots, T_{P_{r}} C$.

Proposition 4.3. Let $\nabla(F)$ and $\nabla\left(F_{1}, F_{2}\right)$ be as defined in Theorem 1.1.
(i) If $n-2 r=1$ and $\operatorname{Sec}^{r} C=\{F=0\}$ then the entries of $\nabla(F)$ define a variety in $\mathbb{P}^{n-1}$ of codimension 3.
(ii) If $n-2 r=2$ and $\operatorname{Sec}^{r} C=\left\{F_{1}=F_{2}=0\right\}$ then the $2 \times 2$ minors of $\nabla\left(F_{1}, F_{2}\right)$ define a variety in $\mathbb{P}^{n-1}$ of codimension 3.

Proof: (i) Theorem 9.1 tells us that $\operatorname{Sec}^{r} C$ has singular locus $\operatorname{Sec}^{r-1} C$, and that this has codimension 3.
(ii) This is proved in Section 9.3.

We start the proof of Theorem 4.1 with the following lemma.
Lemma 4.4. Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve. Suppose that $n-2 r \geq 1$ and $\ell_{1}, \ldots, \ell_{n}$ are linear forms in $x_{1}, \ldots, x_{n}$ such that

$$
\begin{equation*}
\sum_{i=1}^{n} \ell_{i} \frac{\partial f}{\partial x_{i}} \in I\left(\operatorname{Sec}^{r} C\right) \quad \text { for all } f \in I\left(\operatorname{Sec}^{r} C\right) \tag{6}
\end{equation*}
$$

Then there exists $\lambda \in K$ such that $\ell_{i}=\lambda x_{i}$ for all $1 \leq i \leq n$.
Proof: The coefficients of $\ell_{1}, \ldots, \ell_{n}$ form an $n \times n$ matrix. Let $V \subset \operatorname{Mat}_{n}(K)$ be the subspace of all solutions to (6). We must show that $V$ consists only of scalar matrices. Let $E$ be the Jacobian of $C$. Translation by $T \in E[n]$ is an automorphism of $C$ that extends to an automorphism of $\mathbb{P}^{n-1}$, say given by a matrix $M_{T}$. Now $V$ is stable under conjugation by each $M_{T}$. By considering the standard representation of the Heisenberg group (see for example [F2, Section 3]) it follows that $V$ has a basis $\left\{M_{T}: T \in X\right\}$ for some subset $X \subset E[n]$.

We suppose for a contradiction that $M_{T} \in V$ for some $0_{E} \neq T \in E[n]$. Then translation by $T$ on $C$ extends to an automorphism of $\mathbb{P}^{n-1}$ that sends each point $P \in \operatorname{Sec}^{r} C$ to a point in the tangent space $T_{P} \operatorname{Sec}^{r} C$. Let $H$ be the divisor of a hyperplane section on $C$. For $D$ an effective divisor on $C$ we write $\bar{D} \subset \mathbb{P}^{n-1}$ for the linear subspace cut out by $\mathcal{L}(H-D) \subset \mathcal{L}(H)$. For example, if $D$ is a sum of distinct points on $C$ then $\bar{D}$ is the linear span of these points. We also write $D_{T}$ for $D$ translated by $T$. We choose $D=P_{1}+\ldots+P_{r}$ an effective divisor of degree $r$ such that
(i) $P_{1}, \ldots, P_{r} \in C$ are distinct,
(ii) $D$ and $D_{T}$ have disjoint support,
(iii) $2 D+D_{T} \nsim H$.

Proposition 4.2 shows that for generic $P \in \bar{D}$ we have $T_{P} \operatorname{Sec}^{r} C=\overline{2 D}$. It follows from our assumption $M_{T} \in V$ that $\overline{D_{T}} \subset \overline{2 D}$, equivalently $\mathcal{L}(H-2 D) \subset \mathcal{L}\left(H-D_{T}\right)$. Then by (ii) we have

$$
\mathcal{L}(H-2 D)=\mathcal{L}(H-2 D) \cap \mathcal{L}\left(H-D_{T}\right)=\mathcal{L}\left(H-2 D-D_{T}\right) .
$$

However by (iii) and the Riemann-Roch theorem these spaces do not have the same dimension. Indeed, since $r \geq 1$ and $n-2 r \geq 1$ we have

$$
\operatorname{dim} \mathcal{L}(H-2 D)=n-2 r \neq \max (n-3 r, 0)=\operatorname{dim} \mathcal{L}\left(H-2 D-D_{T}\right) .
$$

This is the required contradiction.
We show that the resolution in Theorem 1.1 is a complex.

Lemma 4.5. Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve, and let $\Omega$ be an alternating matrix of quadratic forms satisfying the hypotheses (H1) and (H2).
(i) If $n=2 r+1$ and $\operatorname{Sec}^{r} C=\{F=0\}$ then

$$
\sum_{i=1}^{n} \frac{\partial F}{\partial x_{i}} \Omega_{i j}=0 \quad \text { for all } 1 \leq j \leq n
$$

(ii) If $n=2 r+2$ and $\operatorname{Sec}^{r} C=\left\{F_{1}=F_{2}=0\right\}$ then

$$
\sum_{i=1}^{n} \frac{\partial F_{1}}{\partial x_{i}} \Omega_{i j}=\sum_{i=1}^{n} \frac{\partial F_{2}}{\partial x_{i}} \Omega_{i j}=0 \quad \text { for all } 1 \leq j \leq n
$$

Proof: (i) By the hypothesis (H1) we have

$$
\sum_{i=1}^{n} \frac{\partial F}{\partial x_{i}} \Omega_{i j}=\ell_{j} F \quad \text { for all } 1 \leq j \leq n
$$

for some linear forms $\ell_{1}, \ldots, \ell_{n}$. We multiply by $\frac{\partial F}{\partial x_{j}}$ and sum over $j$. Since $\Omega$ is alternating the left hand side is zero. Therefore

$$
\sum_{j=1}^{n} \ell_{j} \frac{\partial F}{\partial x_{j}}=0
$$

By Lemma 4.4 and Euler's identity it follows that $\ell_{1}=\ldots=\ell_{n}=0$ as required.
(ii) By the hypothesis (H1) we have

$$
\begin{equation*}
\sum_{i=1}^{n} \frac{\partial F_{1}}{\partial x_{i}} \Omega_{i j}=\ell_{j} F_{1}+m_{j} F_{2} \quad \text { for all } 1 \leq j \leq n \tag{7}
\end{equation*}
$$

for some linear forms $\ell_{1}, \ldots, \ell_{n}$ and $m_{1}, \ldots, m_{n}$. We multiply by $\frac{\partial F_{1}}{\partial x_{j}}$ and sum over $j$. Since $\Omega$ is alternating the left hand side is zero. Since $F_{1}$ and $F_{2}$ are forms defining a variety of codimension 2 they must be coprime. Therefore

$$
\sum_{j=1}^{n} \ell_{j} \frac{\partial F_{1}}{\partial x_{j}}=\xi F_{2} \quad \text { and } \quad \sum_{j=1}^{n} m_{j} \frac{\partial F_{1}}{\partial x_{j}}=-\xi F_{1}
$$

for some $\xi \in K$. If instead we multiply (7) by $\frac{\partial F_{2}}{\partial x_{j}}$ and sum over $j$ then using the hypothesis (H2) we find that

$$
\sum_{j=1}^{n} \ell_{j} \frac{\partial F_{2}}{\partial x_{j}}=\eta F_{2} \quad \text { and } \quad \sum_{j=1}^{n} m_{j} \frac{\partial F_{2}}{\partial x_{j}}=-\eta F_{1}
$$

for some $\eta \in K$.
By Lemma 4.4 there exist $\lambda, \mu \in K$ such that $\ell_{i}=\lambda x_{i}$ and $m_{i}=\mu x_{i}$ for all $1 \leq i \leq n$. By Euler's identity and the linear independence of $F_{1}$ and $F_{2}$ it follows that $\lambda=\mu=0$. Therefore

$$
\sum_{i=1}^{n} \frac{\partial F_{1}}{\partial x_{i}} \Omega_{i j}=0 \quad \text { for all } 1 \leq j \leq n
$$

The corresponding result for $F_{2}$ follows by symmetry.
To complete the proof of Theorem 4.1 we must show that the complex is exact. First we need some linear algebra. If $B$ is an $n \times n$ matrix and $S \subset\{1, \ldots, n\}$ then we write
$B^{S}$ for the $(n-|S|) \times(n-|S|)$ matrix obtained by deleting the rows and columns indexed by $S$. The Pfaffian $\operatorname{pf}(M)$ of an alternating matrix $M$ is a polynomial in the matrix entries with the property that $\operatorname{det}(M)=\operatorname{pf}(M)^{2}$.
Lemma 4.6. (i) Let $A=\left(a_{i}\right)$ be a $1 \times n$ matrix and $B$ an $n \times n$ alternating matrix over a field $K$. Suppose that $\operatorname{rank} A=1, \operatorname{rank} B=n-1$ and $A B=0$. Then there exists $\lambda \in K^{\times}$such that

$$
(-1)^{i} \operatorname{pf}\left(B^{\{i\}}\right)=\lambda a_{i}
$$

for all $1 \leq i \leq n$.
(ii) Let $A=\left(a_{i j}\right)$ be a $2 \times n$ matrix and $B$ an $n \times n$ alternating matrix over a field $K$. Suppose that $\operatorname{rank} A=2, \operatorname{rank} B=n-2$ and $A B=0$. Then there exists $\lambda \in K^{\times}$such that

$$
(-1)^{i+j} \operatorname{pf}\left(B^{\{i, j\}}\right)=\lambda\left(a_{1 i} a_{2 j}-a_{1 j} a_{2 i}\right)
$$

for all $1 \leq i<j \leq n$.
Proof: (i) It is well known that the vector with $i$ th entry $(-1)^{i} \operatorname{pf}\left(B^{\{i\}}\right)$ belongs to the kernel of $B$. See for example [BH, Section 3.4]. Since $\operatorname{rank} B=n-1$ this vector is non-zero and the kernel is 1-dimensional. The result follows.
(ii) We first claim there exist $\lambda_{1}, \ldots, \lambda_{n} \in K$ such that

$$
(-1)^{i+j} \operatorname{pf}\left(B^{\{i, j\}}\right)= \begin{cases}\lambda_{i}\left(a_{1 i} a_{2 j}-a_{1 j} a_{2 i}\right) & \text { if } i<j, \\ -\lambda_{i}\left(a_{1 i} a_{2 j}-a_{1 j} a_{2 i}\right) & \text { if } i>j .\end{cases}
$$

Indeed taking $a_{2 i}$ times the first row of $A$ minus $a_{1 i}$ times the second row of $A$, gives a non-zero vector in the kernel of $B^{\{i\}}$. If rank $B^{\{i\}}=n-2$ then we argue as in (i). Otherwise we can simply take $\lambda_{i}=0$. This proves the claim.

Now let $C=\left(a_{1 i} a_{2 j}-a_{1 j} a_{2 i}\right)_{i, j=1, \ldots, n}$ and let $D$ be the diagonal matrix with entries $\lambda_{1}, \ldots, \lambda_{n}$. We must show that if $C D=D C$ then $C D$ is a scalar multiple of $C$. More generally this is true for any rank 2 alternating matrix $C$ and diagonal matrix $D$. Indeed we may re-order the rows and columns so that the diagonal entries of $D$ which are equal are grouped together. Then $C$ is in block diagonal form. Since $C$ is alternating of rank 2 , exactly one of these blocks is non-zero. The result is then clear.

Lemma 4.7. Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve, and let $\Omega$ be an alternating matrix of quadratic forms satisfying the hypotheses (H1), (H2) and (H3).
(i) If $n=2 r+1$ and $\operatorname{Sec}^{r} C=\{F=0\}$ then the $(n-1) \times(n-1)$ Pfaffians of $\Omega$ are (scalar multiples of) the partial derivatives of $F$.
(ii) If $n=2 r+2$ and $\operatorname{Sec}^{r} C=\left\{F_{1}=F_{2}=0\right\}$ then the $(n-2) \times(n-2)$ Pfaffians of $\Omega$ are (scalar multiples of) the $2 \times 2$ minors of $\nabla\left(F_{1}, F_{2}\right)$.

Proof: We apply Lemma 4.6 over the function field $K\left(x_{1}, \ldots, x_{n}\right)$.
(i) By Lemma 4.5 we have $\sum_{i=1}^{n} \frac{\partial F}{\partial x_{i}} \Omega_{i j}=0$. By the hypothesis (H3) the generic rank of $\Omega$ is $n-1$. So by Lemma 4.6(i) there exists $\lambda \in K\left(x_{1}, \ldots, x_{n}\right)$ such that

$$
(-1)^{i} \operatorname{pf}\left(\Omega^{\{i\}}\right)=\lambda \frac{\partial F}{\partial x_{i}} \quad \text { for all } 1 \leq i \leq n .
$$

Since $\operatorname{pf}\left(\Omega^{\{i\}}\right)$ and $\frac{\partial F}{\partial x_{i}}$ are forms of degree $n-1$, we can write $\lambda=u / v$ where $u$ and $v$ are coprime forms of the same degree. Then $v$ divides $\frac{\partial F}{\partial x_{i}}$ for all $i$, and so must be a
constant by Proposition 4.3(i). Therefore $\lambda$ is a constant.
(ii) By Lemma 4.5 we have $\sum_{i=1}^{n} \frac{\partial F_{1}}{\partial x_{i}} \Omega_{i j}=\sum_{i=1}^{n} \frac{\partial F_{2}}{\partial x_{i}} \Omega_{i j}=0$. By the hypothesis (H3) the generic rank of $\Omega$ is $n-2$. So by Lemma $4.6\left(\right.$ (ii) there exists $\lambda \in K\left(x_{1}, \ldots, x_{n}\right)$ such that

$$
(-1)^{i+j} \operatorname{pf}\left(\Omega^{\{i, j\}}\right)=\lambda \frac{\partial\left(F_{1}, F_{2}\right)}{\partial\left(x_{i}, x_{j}\right)} \quad \text { for all } 1 \leq i<j \leq n
$$

Since $\operatorname{pf}\left(\Omega^{\{i, j\}}\right)$ and $\frac{\partial\left(F_{1}, F_{2}\right)}{\partial\left(x_{i}, x_{j}\right)}$ are forms of degree $n-2$, we can write $\lambda=u / v$ where $u$ and $v$ are coprime forms of the same degree. Then $v$ divides $\frac{\partial\left(F_{1}, F_{2}\right)}{\partial\left(x_{i}, x_{j}\right)}$ for all $i, j$, and so must be a constant by Proposition 4.3(ii). Therefore $\lambda$ is a constant.

Let $R=K\left[x_{1}, \ldots, x_{n}\right]$. Consider a complex of graded free $R$-modules

$$
\begin{equation*}
0 \rightarrow F_{m} \xrightarrow{\varphi_{m}} F_{m-1} \longrightarrow \ldots \longrightarrow F_{1} \xrightarrow{\varphi_{1}} F_{0} . \tag{8}
\end{equation*}
$$

We write $V_{k} \subset \mathbb{P}^{n-1}$ for the subvariety defined by the $r_{k} \times r_{k}$ minors of $\varphi_{k}$ where $r_{k}=$ $\operatorname{rank}\left(\varphi_{k}\right)$. The Buchsbaum-Eisenbud acyclicity criterion (see [BH, Theorem 1.4.13] or [E, Theorem 20.9]) states that (8) is exact if and only if $\operatorname{rank} F_{k}=\operatorname{rank} \varphi_{k}+$ $\operatorname{rank} \varphi_{k+1}$ and codim $V_{k} \geq k$ for all $1 \leq k \leq m$.
Proof of Theorem 4.1: We already saw in Lemma 4.5 that the resolution in Theorem 1.1 is a complex. We must prove it is exact. If $n$ is odd then the free $R$-modules have ranks $1, n, n, 1$ and the maps have ranks $1, n-1,1$. If $n$ is even then the free $R$-modules have ranks $2, n, n, 2$ and the maps have ranks $2, n-2,2$. By Lemma 4.7 we have $V_{1}=V_{2}=V_{3}$ and Proposition 4.3 shows that this variety has codimension 3 . We now apply the Buchsbaum-Eisenbud acyclicity criterion.

## 5. A BASIS-FREE CONSTRUCTION

The results of Section 2 show that for the proof of Theorems 1.1 and 1.2 we are free to make changes of co-ordinates on $\mathbb{P}^{n-1}$. Since we are working over an algebraically closed field we can therefore reduce to the following situation. Let $E$ be the elliptic curve

$$
y^{2}+a_{1} x y+a_{3} y=x^{3}+a_{2} x^{2}+a_{4} x+a_{6}
$$

with point at infinity $0_{E}$ and invariant differential

$$
\omega=d x /\left(2 y+a_{1} x+a_{3}\right)=d y /\left(3 x^{2}+2 a_{2} x+a_{4}-a_{1} y\right) .
$$

Let $C \subset \mathbb{P}^{n-1}$ be the image of $E$ embedded via the complete linear system $\left|n .0_{E}\right|$. The embedding depends on a choice of basis for the Riemann-Roch space $\mathcal{L}\left(n .0_{E}\right)$, but the only effect of changing this is to make a change of co-ordinates on $\mathbb{P}^{n-1}$. In this section we define a linear map $\Omega: \wedge^{2} \mathcal{L}\left(n .0_{E}\right) \rightarrow S^{2} \mathcal{L}\left(n .0_{E}\right)$. In the next section we show that the corresponding alternating matrix of quadratic forms satisfies the hypotheses (H1), (H2) and (H3).

For $f \in \mathcal{L}\left(n .0_{E}\right)$ we put $\dot{f}=d f / \omega \in \mathcal{L}\left((n+1) .0_{E}\right)$. Motivated by (2) we define a linear map

$$
\left.\begin{array}{rl}
A: \wedge^{2} \mathcal{L}\left(n .0_{E}\right) & \rightarrow S^{2} \mathcal{L}\left((n+1) .0_{E}\right) \\
f & \wedge g
\end{array}\right) f \otimes \dot{g}-g \otimes \dot{f} .
$$

Lemma 5.1. Let $f, g \in \mathcal{L}\left(n .0_{E}\right)$. Then the rational function on $E \times E$ given by

$$
(P, Q) \mapsto \frac{y_{P}+y_{Q}+a_{1} x_{Q}+a_{3}}{x_{P}-x_{Q}}(f(Q) g(P)-f(P) g(Q))
$$

belongs to $\mathcal{L}\left((n+1) .0_{E}\right) \otimes \mathcal{L}\left((n+1) .0_{E}\right)$.
Proof: (i) If we fix $Q=\left(x_{Q}, y_{Q}\right)$ then as rational functions of $P=(x, y)$,

$$
\frac{y+y_{Q}+a_{1} x_{Q}+a_{3}}{x-x_{Q}} \in \mathcal{L}\left(0_{E}+Q\right) \quad \text { and } \quad f(Q) g-g(Q) f \in \mathcal{L}\left(n .0_{E}-Q\right) .
$$

Therefore the product belongs to $\mathcal{L}\left((n+1) \cdot 0_{E}\right)$.
(ii) If we fix $P=\left(x_{P}, y_{P}\right)$ then as rational functions of $Q=(x, y)$,

$$
\frac{y_{P}+y+a_{1} x+a_{3}}{x_{P}-x} \in \mathcal{L}\left(0_{E}+P\right) \quad \text { and } \quad g(P) f-f(P) g \in \mathcal{L}\left(n .0_{E}-P\right) .
$$

Therefore the product belongs to $\mathcal{L}\left((n+1) \cdot 0_{E}\right)$.
We define a second linear map

$$
\left.\begin{array}{rl}
B: \wedge^{2} \mathcal{L}\left(n .0_{E}\right) & \rightarrow S^{2} \mathcal{L}\left((n+1) \cdot 0_{E}\right) \\
f & \mapsto g
\end{array}\right)\left.\frac{y_{P}+y_{Q}+a_{1} x_{Q}+a_{3}}{x_{P}-x_{Q}}(f(Q) g(P)-f(P) g(Q))\right|_{P=Q}
$$

where $\left.\right|_{P=Q}$ is our notation for the natural map

$$
\mathcal{L}\left((n+1) .0_{E}\right) \otimes \mathcal{L}\left((n+1) .0_{E}\right) \rightarrow S^{2} \mathcal{L}\left((n+1) .0_{E}\right) .
$$

We show that $A$ and $B$ both represent the invariant differential $\omega$, in the sense of Theorem 1.2(i).

Lemma 5.2. As rational functions on $E$ we have

$$
A(f \wedge g)=B(f \wedge g)=f \dot{g}-g \dot{f}=\frac{f d g-g d f}{\omega}
$$

Proof: This is clear for $A$. For $B$ we apply l'Hôpital's rule to get

$$
\left.\frac{f(Q) g-g(Q) f}{x-x_{Q}}\right|_{P=Q}=\left.\frac{f(Q) \dot{g}-g(Q) \dot{f}}{\dot{x}}\right|_{P=Q}
$$

and then use that $\dot{x}=2 y+a_{1} x+a_{3}$.
If we pick bases for $\mathcal{L}\left(n .0_{E}\right)$ and $\mathcal{L}\left((n+1) .0_{E}\right)$ then $A$ and $B$ are (represented by) $n \times n$ alternating matrices of quadratic forms in $n+1$ variables. However the matrix $\Omega$ we seek is an $n \times n$ alternating matrix of quadratic forms in $n$ variables. It turns out that the correct choice of $\Omega$ is a linear combination of $A$ and $B$.

We may expand rational functions on $E$ as Laurent power series in the local parameter $t=x / y$ at $0_{E}$. Let $\phi$ be the linear map that reads off the coefficient of $t^{-n-1}$. There are exact sequences

$$
0 \rightarrow \mathcal{L}\left(n .0_{E}\right) \rightarrow \mathcal{L}\left((n+1) .0_{E}\right) \xrightarrow{\phi} K \rightarrow 0
$$

and

$$
\begin{equation*}
0 \rightarrow S^{2} \mathcal{L}\left(n .0_{E}\right) \rightarrow S^{2} \mathcal{L}\left((n+1) .0_{E}\right) \xrightarrow{\phi_{\mathcal{Z}}} \mathcal{L}\left((n+1) .0_{E}\right) \rightarrow 0 \tag{9}
\end{equation*}
$$

where $\phi_{2}(f \otimes g)=\phi(f) g+\phi(g) f$.

Lemma 5.3. Let $f, g \in \mathcal{L}\left(n .0_{E}\right)$ be rational functions whose coefficients of $t^{-n}$ (when expanded as Laurent power series in $t$ ) are 0,1 respectively. Then

$$
\phi_{2}(A(f \wedge g))=n f \quad \text { and } \quad \phi_{2}(B(f \wedge g))=2 f
$$

Proof: (i) We have $x=t^{-2}+\ldots$ and $y=t^{-3}+\ldots$. Then $\dot{x}=2 y+a_{1} x+a_{3}=2 t^{-3}+$ $\ldots$ and $\dot{y}=3 x^{2}+2 a_{2} x+a_{4}-a_{1} y=3 t^{-4}+\ldots$. Writing $g$ as a polynomial in $x$ and $y$ it follows that $g=t^{-n}+\ldots$ and $\dot{g}=n t^{-n-1}+\ldots$. Therefore $\phi(f)=\phi(g)=\phi(\dot{f})=0$ and $\phi(\dot{g})=n$. We compute

$$
\phi_{2}(A(f \wedge g))=\phi_{2}(f \otimes \dot{g}-g \otimes \dot{f})=n f
$$

(ii) If we fix $Q=\left(x_{Q}, y_{Q}\right)$ then as rational functions of $P=(x, y)$,

$$
\frac{y+y_{Q}+a_{1} x_{Q}+a_{3}}{x-x_{Q}}=t^{-1}+\ldots \quad \text { and } \quad f(Q) g-g(Q) f=f(Q) t^{-n}+\ldots
$$

with product $f(Q) t^{-n-1}+\ldots$.
If we fix $P=\left(x_{P}, y_{P}\right)$ then as rational functions of $Q=(x, y)$,

$$
\frac{y_{P}+y+a_{1} x+a_{3}}{x_{P}-x}=-t^{-1}+\ldots \quad \text { and } \quad g(P) f-f(P) g=-f(P) t^{-n}+\ldots
$$

with product $f(P) t^{-n-1}+\ldots$. In both cases the leading coefficient is $f$. Adding these together gives $\phi_{2}(B(f \wedge g))=2 f$.

Corollary 5.4. Let $\Omega=n B-2 A$. Then $\Omega$ is a linear map $\wedge^{2} \mathcal{L}\left(n .0_{E}\right) \rightarrow S^{2} \mathcal{L}\left(n .0_{E}\right)$.
Proof: This follows from Lemma 5.3 and the exact sequence (9).

## 6. Proof of Theorem 1.1

If we pick a basis for $\mathcal{L}\left(n .0_{E}\right)$ then the linear map defined in Corollary 5.4 is represented by an $n \times n$ alternating matrix of quadratic forms in $n$ variables. In this section we complete the proof of Theorem 1.1 by showing that this matrix $\Omega$ satisfies the hypotheses (H1), (H2) and (H3), as stated at the start of Section 4.

For $0_{E} \neq P \in E$ we write $\mathbf{P}$ and $d \mathbf{P}$ for the linear maps $f \mapsto f(P)$ and $f \mapsto \dot{f}(P)$ in the dual space $\mathcal{L}\left(n .0_{E}\right)^{*}$. For example, if $\mathcal{L}\left(n .0_{E}\right)$ has basis $1, x, y, x^{2}, x y, \ldots$ then

$$
\begin{aligned}
\mathbf{P} & =\left(1, x_{P}, y_{P}, x_{P}^{2}, x_{P} y_{P}, \ldots\right), \\
d \mathbf{P} & =\left(0,2 y_{P}+a_{1} x_{P}+a_{3}, 3 x_{P}^{2}+2 a_{2} x_{P}+a_{4}-a_{1} y_{P}, \ldots\right) .
\end{aligned}
$$

We note that $[\mathbf{P}]$ is a point on $C \subset \mathbb{P}^{n-1}=\mathbb{P}\left(\mathcal{L}\left(n .0_{E}\right)^{*}\right)$, with tangent line passing through $[d \mathbf{P}]$. The square brackets indicate that we are taking the 1 -dimensional subspaces spanned by these vectors, i.e. the corresponding points in projective space. For $0_{E} \neq Q \in E$ we likewise define $\mathbf{Q}$ and $d \mathbf{Q}$.

For $P, Q \in E$ let $\lambda_{P, Q}$ be the slope of the chord (or tangent line if $P=Q$ ) joining $P$ and $Q$. In the following lemma the vectors $\mathbf{P}, \mathbf{Q}, d \mathbf{P}, d \mathbf{Q}$ in $\mathcal{L}\left(n .0_{E}\right)^{*}$ are extended to $\mathcal{L}\left((n+1) .0_{E}\right)^{*}$ using exactly the same definition. Evaluating $A$ or $B$ at a linear combination $\xi \mathbf{P}+\eta \mathbf{Q}$ gives an element of $\left(\wedge^{2} \mathcal{L}\left(n .0_{E}\right)\right)^{*}=\wedge^{2}\left(\mathcal{L}\left(n .0_{E}\right)^{*}\right)$.

Lemma 6.1. Let $0_{E} \neq P, Q \in E$ and $\xi, \eta \in K$. Then

$$
\begin{aligned}
& A(\xi \mathbf{P}+\eta \mathbf{Q})=\xi^{2}(\mathbf{P} \wedge d \mathbf{P})+\xi \eta(\mathbf{P} \wedge d \mathbf{Q}+\mathbf{Q} \wedge d \mathbf{P})+\eta^{2}(\mathbf{Q} \wedge d \mathbf{Q}) \\
& B(\xi \mathbf{P}+\eta \mathbf{Q})=\xi^{2}(\mathbf{P} \wedge d \mathbf{P})+\xi \eta\left(\lambda_{Q,-P}-\lambda_{P,-Q}\right)(\mathbf{P} \wedge \mathbf{Q})+\eta^{2}(\mathbf{Q} \wedge d \mathbf{Q})
\end{aligned}
$$

Proof: (i) For $f, g \in \mathcal{L}\left(n .0_{E}\right)$ we compute

$$
A(\mathbf{P})(f \wedge g)=(f \dot{g}-g \dot{f})(P)=(\mathbf{P} \wedge d \mathbf{P})(f \wedge g)
$$

The formula for $A(\xi \mathbf{P}+\eta \mathbf{Q})$ follows by bilinearity.
(ii) For $f, g \in \mathcal{L}\left(n .0_{E}\right)$ we write

$$
B(\xi \mathbf{P}+\eta \mathbf{Q})(f \wedge g)=\xi^{2} B_{0}+\xi \eta B_{1}+\eta^{2} B_{2}
$$

By Lemma 5.2 we have

$$
\begin{aligned}
& B_{0}=(f \dot{g}-g \dot{f})(P)=(\mathbf{P} \wedge d \mathbf{P})(f \wedge g), \\
& B_{2}=(f \dot{g}-g \dot{f})(Q)=(\mathbf{Q} \wedge d \mathbf{Q})(f \wedge g)
\end{aligned}
$$

Since for $s, t \in \mathcal{L}\left((n+1) .0_{E}\right)$ we have

$$
\begin{aligned}
(s \otimes t)(\xi \mathbf{P}+\eta \mathbf{Q}) & =s(\xi \mathbf{P}+\eta \mathbf{Q}) t(\xi \mathbf{P}+\eta \mathbf{Q}) \\
& =\xi^{2} s(P) t(P)+\xi \eta(s(P) t(Q)+s(Q) t(P))+\eta^{2} s(Q) t(Q)
\end{aligned}
$$

it follows from the definition of $B$ that

$$
\begin{aligned}
B_{1} & =\lambda_{P,-Q}(f(Q) g(P)-f(P) g(Q))+\lambda_{Q,-P}(f(P) g(Q)-f(Q) g(P)) \\
& =\left(\lambda_{Q,-P}-\lambda_{P,-Q}\right)(\mathbf{P} \wedge \mathbf{Q})(f \wedge g) .
\end{aligned}
$$

We pick a basis for $\mathcal{L}\left(n .0_{E}\right)$, so that now $\Omega(\mathbf{P})$ is an $n \times n$ alternating matrix, and $\mathbf{P}, \mathbf{Q}, d \mathbf{P}, d \mathbf{Q}$ are column vectors.
Lemma 6.2. Let $0_{E} \neq P_{1}, \ldots, P_{r} \in E$ distinct and $\xi_{1}, \ldots, \xi_{r} \in K$. Then

$$
\Omega\left(\sum_{i=1}^{r} \xi_{i} \mathbf{P}_{i}\right)=\Pi\left(\begin{array}{cc}
* & \Xi \\
-\Xi & 0
\end{array}\right) \Pi^{T}
$$

where

$$
\Xi=\left(\begin{array}{cccc}
(n-2) \xi_{1}^{2} & -2 \xi_{1} \xi_{2} & \ldots & -2 \xi_{1} \xi_{r}  \tag{10}\\
-2 \xi_{1} \xi_{2} & (n-2) \xi_{2}^{2} & \ldots & -2 \xi_{2} \xi_{r} \\
\vdots & \vdots & \ddots & \vdots \\
-2 \xi_{1} \xi_{r} & -2 \xi_{2} \xi_{r} & \ldots & (n-2) \xi_{r}^{2}
\end{array}\right)
$$

and $\Pi$ is the $n \times 2 r$ matrix with columns $\mathbf{P}_{1}, \ldots, \mathbf{P}_{r}, d \mathbf{P}_{1}, \ldots, d \mathbf{P}_{r}$.
Proof: Recall that $\Omega=n B-2 A$. The case $r=2$ is immediate from Lemma 6.1. Since the entries of $\Omega$ are quadratic forms the general case follows.

We now check the hypotheses (H1), (H2) and (H3).
Proof of (H1) and (H3): Suppose $n-2 r \geq 1$. A generic point $P \in \operatorname{Sec}^{r} C$ may be written $P=\left[\sum_{i=1}^{r} \xi_{i} \mathbf{P}_{i}\right]$ for some $0_{E} \neq P_{1}, \ldots, P_{r} \in E$ distinct and $\xi_{1}, \ldots, \xi_{r} \neq 0$. By Proposition 4.2 the tangent space $T_{P} \operatorname{Sec}^{r} C \subset \mathbb{P}^{n-1}$ is spanned by $\mathbf{P}_{1}, \ldots, \mathbf{P}_{r}, d \mathbf{P}_{1}, \ldots d \mathbf{P}_{r}$. In particular these $2 r$ vectors are linearly independent.

For $f \in I\left(\operatorname{Sec}^{r} C\right)$ we have $\sum_{i=1}^{n} \frac{\partial f}{\partial x_{i}}(P) \mathbf{v}_{i}=0$ for any $\mathbf{v}$ in the linear span of $\mathbf{P}_{1}, \ldots, \mathbf{P}_{r}, d \mathbf{P}_{1}, \ldots d \mathbf{P}_{r}$. By Lemma 6.2 the columns of $\Omega$ are linear combinations of these vectors. So for each $1 \leq j \leq n$ the form $\sum_{i=1}^{n} \frac{\partial f}{\partial x_{i}} \Omega_{i j}$ vanishes at $P$. Since $P \in \operatorname{Sec}^{r} C$ is generic, this proves (H1). Since $n \notin\{0,2 r\}$ and $\xi_{1}, \ldots, \xi_{r} \neq 0$, the matrix (10) is non-singular. Therefore $\operatorname{rank} \Omega(P)=2 r$ and this proves (H3).

Proof of (H2): We write $n=2 r$ and $\operatorname{Sec}^{r-1} C=\left\{F_{1}=F_{2}=0\right\}$ where $F_{1}$ and $F_{2}$ are forms of degree $r$. We must show that the form

$$
\begin{equation*}
\sum_{i, j=1}^{n} \frac{\partial F_{1}}{\partial x_{i}} \Omega_{i j} \frac{\partial F_{2}}{\partial x_{j}} \tag{11}
\end{equation*}
$$

is identically zero. A generic point $P \in \operatorname{Sec}^{r} C=\mathbb{P}^{n-1}$ may be written $P=$ [ $\sum_{i=1}^{r} \xi_{i} \mathbf{P}_{i}$ ] for some $0_{E} \neq P_{1}, \ldots, P_{r} \in E$ distinct and $\xi_{1}, \ldots, \xi_{r} \neq 0$. In addition we may assume that $2\left(P_{1}+\ldots+P_{r}\right) \nsucc H$ where $H$ is the hyperplane section. This ensures that the vectors $\mathbf{P}_{1}, \ldots, \mathbf{P}_{r}, d \mathbf{P}_{1}, \ldots d \mathbf{P}_{r}$ are linearly independent. We choose co-ordinates on $\mathbb{P}^{n-1}$ so that $\left[\mathbf{P}_{1}\right]=(1: 0: \ldots: 0),\left[\mathbf{P}_{2}\right]=(0: 1: 0: \ldots: 0), \ldots$, $d \mathbf{P}_{r}=(0: \ldots: 0: 1)$. Since $F_{1}$ and $F_{2}$ vanish on Sec ${ }^{r-1} C$ they vanish on the linear span of any $r-1$ of the $\left[\mathbf{P}_{i}\right]$. Replacing $F_{1}$ and $F_{2}$ by suitable linear combinations we may assume

$$
\begin{aligned}
& F_{1}\left(x_{1}, \ldots, x_{r}, 0, \ldots, 0\right)=0 \\
& F_{2}\left(x_{1}, \ldots, x_{r}, 0, \ldots, 0\right)=x_{1} x_{2} \ldots x_{r} .
\end{aligned}
$$

Therefore at $P=\left(\xi_{1}: \ldots: \xi_{r}: 0: \ldots: 0\right)$ we have

$$
\begin{aligned}
& \left(\frac{\partial F_{1}}{\partial x_{1}}(P), \ldots, \frac{\partial F_{1}}{\partial x_{n}}(P)\right)=(0, \ldots, 0, *, \ldots, *), \\
& \left(\frac{\partial F_{2}}{\partial x_{1}}(P), \ldots, \frac{\partial F_{2}}{\partial x_{n}}(P)\right)=\left(\prod_{i \neq 1} \xi_{i}, \ldots, \prod_{i \neq r} \xi_{i}, *, \ldots, *\right) .
\end{aligned}
$$

By Lemma 6.2 we have

$$
\Omega(P)=\left(\begin{array}{cc}
* & \Xi \\
-\Xi & 0
\end{array}\right)
$$

where $\Xi$ is the matrix (10). Since $n=2 r$ the coefficients in each row and column of $\Xi$ sum to zero. Therefore the form (11) vanishes at $P$. Since $P \in \mathbb{P}^{n-1}$ is generic, this shows that the form is identically zero.

This completes the proof of Theorem 1.1.

## 7. Explicit formulae

In this section we give an explicit formula for the matrix $\Omega$ defined in Section 5. As before $E$ is the elliptic curve

$$
y^{2}+a_{1} x y+a_{3} y=x^{3}+a_{2} x^{2}+a_{4} x+a_{6}
$$

with invariant differential $\omega=d x /\left(2 y+a_{1} x+a_{3}\right)$. We embed $E$ in $\mathbb{P}^{n-1}$ via

$$
\begin{equation*}
\left(x_{0}: x_{2}: x_{3}: \ldots: x_{n}\right)=\left(1, x, y, x^{2}, x y, x^{3}, x^{2} y, x^{4}, \ldots\right) \tag{12}
\end{equation*}
$$

Notice there is no $x_{1}$. The indicator function of a set $X$ is denoted $\mathbf{1}_{X}$. We define linear forms in indeterminates $\left\{x_{m}: m \in \mathbb{Z}\right\}$ as follows. The relation to the notation $\dot{f}=d f / \omega$ used in Section 5 will be explained below.

$$
\begin{aligned}
& \dot{x}_{m}=\frac{m}{2}\left(2 x_{m+1}+a_{1} x_{m}+a_{3} x_{m-2}\right)+\mathbf{1}_{\text {odd }}(m) \sum_{i=1}^{6}(-1)^{i}\left(m-\frac{i}{2}\right) a_{i} x_{m+1-i} \\
& \bar{x}_{m}=\frac{1}{2}\left(2 x_{m+1}+a_{1} x_{m}+a_{3} x_{m-2}\right)+\mathbf{1}_{\text {odd }}(m) \sum_{i=1}^{6}(-1)^{i} a_{i} x_{m+1-i}
\end{aligned}
$$

where by convention $a_{5}=0$. For $x \in \mathbb{R}$ we let $\operatorname{sign}(x)=-1,0,1$ according as $x$ is negative, zero or positive. For $r, s \in \mathbb{Z}$ we define

$$
\begin{aligned}
& A_{r s}=x_{r} \dot{x}_{s}-x_{s} \dot{x}_{r}, \\
& B_{r s}=\sum_{k=-\infty}^{\infty} \operatorname{sign}\left(k+\frac{1}{2}\right)\left(x_{r+2 k} \bar{x}_{s-2 k}-x_{s+2 k} \bar{x}_{r-2 k}\right) .
\end{aligned}
$$

Theorem 7.1. Let $C \subset \mathbb{P}^{n-1}$ be the image of $E$ under the embedding (12).
(i) $A=\left(A_{r s}\right)_{r, s=0,2, \ldots, n}$ and $B=\left(B_{r s}\right)_{r, s=0,2, \ldots, n}$ are $n \times n$ alternating matrices of quadratic forms in $x_{0}, x_{2}, \ldots, x_{n+1}$.
(ii) $\Omega=n B-2 A$ is an $n \times n$ alternating matrix of quadratic forms in $x_{0}, x_{2}, \ldots$, $x_{n}$. It satisfies the conclusions of Theorem 1.1 and

$$
\begin{equation*}
(n-2) \omega=\frac{x_{j}^{2} d\left(x_{i} / x_{j}\right)}{\Omega_{i j}\left(x_{1}, \ldots, x_{n}\right)} \quad \text { for all } i \neq j \tag{13}
\end{equation*}
$$

Proof: It is part of the theorem that the indeterminates $x_{m}$ for $m \notin\{0,2,3, \ldots, n\}$ cancel from the formula for $\Omega$. So when applying the theorem we simply set them to be zero. However we will not do this in the proof. Since $\bar{x}_{m}$ is a linear combination of $x_{m+1}, x_{m}, \ldots, x_{m-5}$ each $B_{r s}$ is of the form $\sum_{i j} c_{i j} x_{i} x_{j}$ where each $c_{i j}$ is a finite sum. But it is not immediately clear that the $B_{r s}$ are polynomials, i.e. that $c_{i j}=0$ for all but finitely many pairs $(i, j)$. We check this first.

If $r \equiv s(\bmod 2)$ and $r<s$ then

$$
\begin{equation*}
B_{r s}=2\left(x_{r} \bar{x}_{s}+x_{r+2} \bar{x}_{s-2}+\ldots+x_{s-2} \bar{x}_{r+2}\right) \tag{14}
\end{equation*}
$$

whereas if $r$ is even and $s$ is odd then

$$
\begin{equation*}
B_{r s}=-a_{1} x_{r} x_{s}+Q_{r, s+1}+a_{2} Q_{r, s-1}+a_{4} Q_{r, s-3}+a_{6} Q_{r, s-5}-Q_{s, r+1} \tag{15}
\end{equation*}
$$

where

$$
Q_{i j}= \begin{cases}x_{i} x_{j}+x_{i+2} x_{j-2}+\ldots+x_{j} x_{i} & \text { if } i<j+2, \\ 0 & \text { if } i=j+2, \\ -\left(x_{i-2} x_{j+2}+x_{i-4} x_{j+4}+\ldots+x_{j+2} x_{i-2}\right) & \text { if } i>j+2\end{cases}
$$

Since $B_{s r}=-B_{r s}$ this proves that the $B_{r s}$ are polynomials.
We show that the matrices $A$ and $B$ defined in the statement of the theorem represent the linear maps $A$ and $B$ defined in Section 5. The theorem then follows from the results of Sections 4, 5 and 6. In particular (13) follows from Lemma 5.2.

In the statement of the theorem the $\left\{x_{m}: m \in \mathbb{Z}\right\}$ are indeterminates. However for the proof they will be the following rational functions on $E$,

$$
x_{m}= \begin{cases}x^{m / 2} & \text { if } m \text { is even }, \\ x^{(m-3) / 2} y & \text { if } m \text { is odd }\end{cases}
$$

As rational functions on $E$, we claim that $\dot{x}_{m}=d x_{m} / \omega$ (in agreement with the notation in Section 5) and $\bar{x}_{m}=\frac{1}{2} x_{m-2}\left(2 y+a_{1} x+a_{3}\right)$. In checking these claims, we start with the right hand sides, since this also serves to motivate the definitions of $\dot{x}_{m}$ and
$\bar{x}_{m}$. For $m$ even we have

$$
\begin{aligned}
d x_{m} / \omega & =\frac{m}{2} x^{(m-2) / 2}(d x / \omega) \\
& =\frac{m}{2} x^{(m-2) / 2}\left(2 y+a_{1} x+a_{3}\right) \\
& =\frac{m}{2}\left(2 x_{m+1}+a_{1} x_{m}+a_{3} x_{m-2}\right)
\end{aligned}
$$

and

$$
\frac{1}{2} x_{m-2}\left(2 y+a_{1} x+a_{3}\right)=\frac{1}{2}\left(2 x_{m+1}+a_{1} x_{m}+a_{3} x_{m-2}\right)
$$

For $m$ odd we have

$$
\begin{aligned}
d x_{m} / \omega= & \frac{m-3}{2} x^{(m-5) / 2} y(d x / \omega)+x^{(m-3) / 2}(d y / \omega) \\
= & \frac{m-3}{2} x^{(m-5) / 2}\left(2 y^{2}+a_{1} x y+a_{3} y\right)+x^{(m-3) / 2}\left(3 x^{2}+2 a_{2} x+a_{4}-a_{1} y\right) \\
= & \frac{m-3}{2} x^{(m-5) / 2}\left(-a_{1} x y-a_{3} y+2 x^{3}+2 a_{2} x^{2}+2 a_{4} x+2 a_{6}\right) \\
& \quad+x^{(m-5) / 2}\left(3 x^{3}+2 a_{2} x^{2}+a_{4} x-a_{1} x y\right) \\
= & x^{(m-5) / 2}\left(m x^{3}-\frac{m-1}{2} a_{1} x y-\frac{m-3}{2} a_{3} y+\sum_{i=1}^{3}(m-i) a_{2 i} x^{3-i}\right) \\
= & \frac{m}{2}\left(2 x_{m+1}+a_{1} x_{m}+a_{3} x_{m-2}\right)+\sum_{i=1}^{6}(-1)^{i}\left(m-\frac{i}{2}\right) a_{i} x_{m+1-i}
\end{aligned}
$$

and

$$
\begin{aligned}
\frac{1}{2} x_{m-2}\left(2 y+a_{1} x+a_{3}\right) & =\frac{1}{2} x^{(m-5) / 2}\left(2 y^{2}+a_{1} x y+a_{3} y\right) \\
& =\frac{1}{2} x^{(m-5) / 2}\left(-a_{1} x y-a_{3} y+2 x^{3}+2 a_{2} x^{2}+2 a_{4} x+2 a_{6}\right) \\
& =\frac{1}{2}\left(2 x_{m+1}+a_{1} x_{m}+a_{3} x_{m-2}\right)+\sum_{i=1}^{6}(-1)^{i} a_{i} x_{m+1-i} .
\end{aligned}
$$

It is now clear that $A\left(x_{r} \wedge x_{s}\right)=A_{r s}$ for all $r, s \in \mathbb{Z}$. It remains to prove the same for $B$. By definition of $B$ we have

$$
B\left(x_{r} \wedge x_{s}\right)=\left.\frac{y_{P}+y_{Q}+a_{1} x_{Q}+a_{3}}{x_{P}-x_{Q}}\left(x_{r}(Q) x_{s}(P)-x_{r}(P) x_{s}(Q)\right)\right|_{P=Q}
$$

where $P, Q$ are points on $E$. Since $\bar{x}_{m}=\frac{1}{2} x_{m-2}\left(2 y+a_{1} x+a_{3}\right)$ we have

$$
\begin{aligned}
2 x_{r}(P) \bar{x}_{s}(Q) & =\left(2 y_{Q}+a_{1} x_{Q}+a_{3}\right) x_{r}(P) x_{s-2}(Q) \\
& =\frac{2 y_{Q}+a_{1} x_{Q}+a_{3}}{x_{P}-x_{Q}}\left(x_{r+2}(P) x_{s-2}(Q)-x_{r}(P) x_{s}(Q)\right)
\end{aligned}
$$

Adding this to the same expression with $(r, s)$ replaced by $(s-2, r+2)$ and then setting $P=Q$ gives

$$
\begin{equation*}
B_{r s}-B_{r+2, s-2}=2\left(x_{r} \bar{x}_{s}+x_{s-2} \bar{x}_{r+2}\right)=B\left(x_{r} \wedge x_{s}\right)-B\left(x_{r+2} \wedge x_{s-2}\right) \tag{16}
\end{equation*}
$$

Rather more obviously, replacing $(r, s)$ by $(r+2, s+2)$ changes $B_{r s}$ and $B\left(x_{r} \wedge x_{s}\right)$ in the same way, that is, by shifting the subscripts up by 2 . So to prove $B\left(x_{r} \wedge x_{s}\right)=B_{r s}$ for all $r, s \in \mathbb{Z}$ it suffices to prove it for all $r \in\{0,1\}$ and $s \in\{0,1,2,3\}$. This is a
finite calculation. We give two examples:

$$
\begin{aligned}
B\left(x_{0} \wedge x_{3}\right) & =\left.\frac{y_{P}+y_{Q}+a_{1} x_{Q}+a_{3}}{x_{P}-x_{Q}}\left(y_{P}-y_{Q}\right)\right|_{P=Q} \\
& =\frac{\left(y_{P}^{2}+a_{1} x_{P} y_{P}+a_{3} y_{P}\right)-\left(y_{Q}^{2}+a_{1} x_{Q} y_{Q}+a_{3} y_{Q}\right)}{x_{P}-x_{Q}}-\left.a_{1} y_{P}\right|_{P=Q} \\
& =\left.\left(x_{P}^{2}+x_{P} x_{Q}+x_{Q}^{2}-a_{1} y_{P}+a_{2}\left(x_{P}+x_{Q}\right)+a_{4}\right)\right|_{P=Q} \\
& =2 x_{0} x_{4}+x_{2}^{2}-a_{1} x_{0} x_{3}+2 a_{2} x_{0} x_{2}+a_{4} x_{0}^{2}
\end{aligned}
$$

and

$$
\begin{aligned}
B\left(x_{2} \wedge x_{3}\right) & =\left.\frac{y_{P}+y_{Q}+a_{1} x_{Q}+a_{3}}{x_{P}-x_{Q}}\left(y_{P}\left(x_{Q}-x_{P}\right)+x_{P}\left(y_{P}-y_{Q}\right)\right)\right|_{P=Q} \\
& =\left.\left(-y_{P}\left(y_{P}+y_{Q}+a_{1} x_{Q}+a_{3}\right)+x_{P}\left(x_{P}^{2}+x_{P} x_{Q}+\ldots+a_{4}\right)\right)\right|_{P=Q} \\
& =\left.\left(x_{P}^{2} x_{Q}+x_{P} x_{Q}^{2}-y_{P} y_{Q}-a_{1} x_{Q} y_{P}+a_{2} x_{P} x_{Q}-a_{6}\right)\right|_{P=Q} \\
& =2 x_{2} x_{4}-x_{3}^{2}-a_{1} x_{2} x_{3}+a_{2} x_{2}^{2}-a_{6} x_{0}^{2} .
\end{aligned}
$$

It is easy to check using (15) that these are equal to $B_{03}$ and $B_{23}$. The other cases we need can then be checked using (16) and the fact that $B$ is alternating.

## 8. Proof of Theorem 1.2

Let $\Omega=n B-2 A$ be as in Theorem 7.1. Then $c_{4}(\Omega)=f_{n}\left(a_{1}, \ldots, a_{6}\right)$ and $c_{6}(\Omega)=$ $g_{n}\left(a_{1}, \ldots, a_{6}\right)$ for some polynomials $f_{n}$ and $g_{n}$. We consider the effect of a change of Weierstrass equation, with notation as in [S, Chapter III].

Lemma 8.1. Let $a_{1}, \ldots, a_{6}$ and $a_{1}^{\prime}, \ldots, a_{6}^{\prime}$ be the coefficients of two Weierstrass equations related by $x=u^{2} x^{\prime}+r$ and $y=u^{3} y^{\prime}+u^{2} s x^{\prime}+t$. Then

$$
\begin{aligned}
& f_{n}\left(a_{1}, \ldots, a_{6}\right)=u^{4} f_{n}\left(a_{1}^{\prime}, \ldots, a_{6}^{\prime}\right) \\
& g_{n}\left(a_{1}, \ldots, a_{6}\right)=u^{6} g_{n}\left(a_{1}^{\prime}, \ldots, a_{6}^{\prime}\right)
\end{aligned}
$$

Proof: This follows from Corollary 2.3 and $u^{-1} \omega^{\prime}=\omega$.
It follows by Lemma 8.1, and the standard procedure for converting a Weierstrass equation to the shorter form $y^{2}=x^{3}+a x+b$, that $f_{n}$ and $g_{n}$ are scalar multiples of the usual polynomials $c_{4}$ and $c_{6}$ in $a_{1}, \ldots, a_{6}$. Explicitly,

$$
\begin{align*}
& f_{n}\left(a_{1}, \ldots, a_{6}\right)=\xi_{n}\left(b_{2}^{2}-24 b_{4}\right)=\xi_{n}\left(a_{1}^{4}+\ldots\right), \\
& g_{n}\left(a_{1}, \ldots, a_{6}\right)=\eta_{n}\left(-b_{2}^{3}+36 b_{2} b_{4}-216 b_{6}\right)=\eta_{n}\left(-a_{1}^{6}+\ldots\right) \tag{17}
\end{align*}
$$

where $b_{2}=a_{1}^{2}+4 a_{2}, b_{4}=2 a_{4}+a_{1} a_{3}$ and $b_{6}=a_{3}^{2}+4 a_{6}$.
To complete the proof of Theorem 1.2 we must compute the constants $\xi_{n}$ and $\eta_{n}$. For any given value of $n$ these can be read off from a single numerical example. However we need to compute these constants for all $n$. We write

$$
\Omega=\Omega^{(0)}+a_{1} \Omega^{(1)}+a_{2} \Omega^{(2)}+a_{3} \Omega^{(3)}+a_{4} \Omega^{(4)}+a_{6} \Omega^{(6)} .
$$

Since $c_{4}(\Omega)$ and $c_{6}(\Omega)$ have degrees 4 and 6 in the coefficients of the entries of $\Omega$, we see by (17) that it suffices to compute the invariants of $\Omega^{(1)}$.

We put

$$
\gamma_{r s}=(-1)^{\max (r, s)} \operatorname{sign}(s-r) n-2\left((-1)^{s}\left\lfloor\frac{s}{2}\right\rfloor-(-1)^{r}\left\lfloor\frac{r}{2}\right\rfloor\right)
$$

Lemma 8.2. The alternating matrix $\Omega^{(1)}$ has entries above the diagonal

$$
\begin{equation*}
\gamma_{r s} x_{r} x_{s}+(-1)^{s} n \mathbf{1}_{\text {even }}(r+s) \sum_{k=1}^{(s-r) / 2-1} x_{r+2 k} x_{s-2 k} \tag{18}
\end{equation*}
$$

Proof: Since $\Omega=n B-2 A$ we have $\Omega^{(1)}=n B^{(1)}-2 A^{(1)}$ where the superscripts indicate that we are taking the coefficient of $a_{1}$. Then $A^{(1)}$ has $(r, s)$ entry

$$
\left((-1)^{s}\left\lfloor\frac{s}{2}\right\rfloor-(-1)^{r}\left\lfloor\frac{r}{2}\right\rfloor\right) x_{r} x_{s}
$$

whereas (14) and (15) show that if $r<s$ then $B^{(1)}$ has $(r, s)$ entry

$$
\begin{cases}(-1)^{s}\left(x_{r} x_{s}+x_{r+2} x_{s-2}+\ldots x_{s-2} x_{r+2}\right) & \text { if } r \equiv s \quad(\bmod 2) \\ (-1)^{s} x_{r} x_{s} & \text { if } r \not \equiv s \quad(\bmod 2)\end{cases}
$$

Lemma 8.3. The matrices $\Omega^{(1)}, \Omega^{\prime}=\left(\gamma_{r s} x_{r} x_{s}\right)_{r, s=0,2,3, \ldots, n}$ and

$$
\Lambda=\left((\operatorname{sign}(j-i) n-2(j-i)) x_{i} x_{j}\right)_{i, j=0,1, \ldots, n-1}
$$

all have the same invariants $c_{4}$ and $c_{6}$.
Proof: We first explain why $\Omega^{(1)}$ and $\Omega^{\prime}$ have the same invariants, despite the "extra terms" in (18). We start with $\Omega^{(1)}$. The only entries involving $x_{0}$ are in the first row and column. We replace $x_{0}$ by $\lambda^{-1} x_{0}$ and multiply the first row and column by $\lambda$. By Lemma 2.2 this does not change the invariants, but setting $\lambda=0$ removes the extra terms from the first row and column. Now the only entries involving $x_{2}$ are in the second row and column. We replace $x_{2}$ by $\lambda^{-1} x_{2}$ and multiply the second row and column by $\lambda$. This does not change the invariants, but setting $\lambda=0$ removes the extra terms from the second row and column. We repeat this procedure for all subsequent rows and columns. In the end we remove all the extra terms, and are left with the matrix $\Omega^{\prime}$.

We define a bijection $\pi:\{0,1, \ldots, n-1\} \rightarrow\{0,2,3, \ldots, n\}$ by

$$
\pi(i)= \begin{cases}2 i & \text { if } i \leq n / 2 \\ 2(n-i)+1 & \text { if } i>n / 2\end{cases}
$$

We then compute

$$
\gamma_{\pi(i), \pi(j)}= \begin{cases}\operatorname{sign}(j-i) n-2(j-i) & \text { if } i \leq n / 2 \text { and } j \leq n / 2 \\ -n-2(-(n-j)-i) & \text { if } i \leq n / 2 \text { and } j>n / 2 \\ n-2(j+(n-i)) & \text { if } i>n / 2 \text { and } j \leq n / 2 \\ \operatorname{sign}(j-i) n-2(-(n-j)+(n-i)) & \text { if } i>n / 2 \text { and } j>n / 2\end{cases}
$$

In all cases we have $\gamma_{\pi(i), \pi(j)}=\operatorname{sign}(j-i) n-2(j-i)$. Therefore $\Omega^{\prime}$ and $\Lambda$ are related by a permutation matrix. It follows by Lemma 2.2 that they have the same invariants.

## Lemma 8.4. The alternating matrix of quadratic forms

$$
\Lambda=\left(\begin{array}{cccccc}
0 & (n-2) x_{1} x_{2} & (n-4) x_{1} x_{3} & (n-6) x_{1} x_{4} & \cdots & (2-n) x_{1} x_{n} \\
& 0 & (n-2) x_{2} x_{3} & (n-4) x_{2} x_{4} & \cdots & (4-n) x_{2} x_{n} \\
& & 0 & (n-2) x_{3} x_{4} & \cdots & (6-n) x_{3} x_{n} \\
& - & & & & \vdots \\
& & & \ddots & & (n-2) x_{n-1} x_{n} \\
& & & & & 0
\end{array}\right)
$$

has invariants $c_{4}(\Lambda)=(n-2)^{4}$ and $c_{6}(\Lambda)=-(n-2)^{6}$.
Proof: We have $\Lambda=\left(\lambda_{r s} x_{r} x_{s}\right)_{r, s=1, \ldots, n}$ where $\lambda_{r s}=\operatorname{sign}(s-r) n-2(s-r)$. Following the definitions of $c_{4}$ and $c_{6}$ in Section 1 we put

$$
\begin{aligned}
M_{i j} & =\sum_{r, s=1}^{n} \frac{\partial \Lambda_{i r}}{\partial x_{s}} \frac{\partial \Lambda_{j s}}{\partial x_{r}}=\mu_{i j} x_{i} x_{j} \\
N_{i j k} & =\sum_{r=1}^{n} \frac{\partial M_{i j}}{\partial x_{r}} \Lambda_{r k}=\nu_{i j k} x_{i} x_{j} x_{k}
\end{aligned}
$$

where $\mu_{i j}=\left(\sum_{r=1}^{n} \lambda_{i r} \lambda_{j r}\right)-\lambda_{i j}^{2}$ and $\nu_{i j k}=\mu_{i j}\left(\lambda_{i k}+\lambda_{j k}\right)$. It is not hard to show that

$$
\begin{aligned}
\sum_{r=1}^{n} \operatorname{sign}(i-r) \operatorname{sign}(j-r) & =n-2|i-j|-\delta_{i j} \\
\sum_{r=1}^{n}(i-r) \operatorname{sign}(j-r) & =2 i j-j^{2}-(n+1) i+n(n+1) / 2 \\
\sum_{r=1}^{n}(i-r)(j-r) & =n i j-(i+j) n(n+1) / 2+n(n+1)(2 n+1) / 6
\end{aligned}
$$

We use these to compute

$$
\sum_{r=1}^{n} \lambda_{i r} \lambda_{j r}=2 n|i-j|^{2}-2 n^{2}|i-j|-\delta_{i j} n^{2}+\left(n^{3}+2 n\right) / 3
$$

and then subtract off

$$
\lambda_{i j}^{2}=4|i-j|^{2}-4 n|i-j|+\left(1-\delta_{i j}\right) n^{2}
$$

to get

$$
\mu_{i j}=2(n-2)\left(|i-j|^{2}-n|i-j|\right)+n(n-1)(n-2) / 3 .
$$

Noting the symmetries $\mu_{i j}=\mu_{j i}$ and $\nu_{i j k}=\nu_{j i k}$, and using computer algebra to check our calculations, we find

$$
\sum_{i, j, r, s=1}^{n} \frac{\partial^{2} M_{i j}}{\partial x_{r} \partial x_{s}} \frac{\partial^{2} M_{r s}}{\partial x_{i} \partial x_{j}}=4 \sum_{i \leq j} \mu_{i j}^{2}=(16 / 3) n(n-2)^{2}\binom{n+3}{5}
$$

and

$$
\begin{aligned}
\sum_{i, j, k, r, s, t=1}^{n} & \frac{\partial^{3} N_{i j k}}{\partial x_{r} \partial x_{s} \partial x_{t}} \frac{\partial^{3} N_{r s t}}{\partial x_{i} \partial x_{j} \partial x_{k}} \\
& =4 \sum_{i \leq j \leq k}\left(\nu_{i j k}+\nu_{j k i}+\nu_{k i j}\right)^{2} \\
& =4 \sum_{i \leq j \leq k}\left(\lambda_{i j}\left(\mu_{i k}-\mu_{j k}\right)+\lambda_{j k}\left(\mu_{i j}-\mu_{i k}\right)+\lambda_{i k}\left(\mu_{i j}-\mu_{j k}\right)\right)^{2} \\
& =64(n-2)^{2} \sum_{i \leq j \leq k}(i-2 j+k)^{2}(n+i+j-2 k)^{2}(n+2 i-j-k)^{2} \\
& =64 n(n-2)^{3}\binom{n+5}{7} .
\end{aligned}
$$

The final sums are evaluated using the standard formulae for $\sum_{i=1}^{n} i, \sum_{i=1}^{n} i^{2}$, etc. In practice it is simpler to observe that the answer is a polynomial in $n$, say of degree at most $d$, and then check the result for $d+1$ distinct values of $n$.

Finally scaling by the constants included in the definitions (4) and (5) it follows that $c_{4}(\Lambda)=(n-2)^{4}$ and $c_{6}(\Lambda)=-(n-2)^{6}$.

The last two lemmas show that $\xi_{n}=(n-2)^{4}$ and $\eta_{n}=(n-2)^{6}$. Therefore $c_{4}(\Omega)=(n-2)^{4} c_{4}(E)$ and $c_{6}(\Omega)=(n-2)^{6} c_{6}(E)$. Let $\omega=d x /\left(2 y+a_{1} x+a_{3}\right)$. By the formulae in [S, Chapter III] we have

$$
(E, \omega) \cong\left(y^{2}=x^{3}-27 c_{4}(E) x-54 c_{6}(E), 3 d x / y\right)
$$

Therefore

$$
(E,(n-2) \omega) \cong\left(y^{2}=x^{3}-27 c_{4}(\Omega) x-54 c_{6}(\Omega), 3 d x / y\right)
$$

Recalling from Theorem 7.1 that $\Omega=n B-2 A$ represents the invariant differential $(n-2) \omega$, this completes the proof of Theorem 1.2.

## 9. Higher secant varieties

In this final section we give references and proofs for the facts about higher secant varieties we used earlier in the paper.

Theorem 9.1. Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve of degree $n \geq 3$.
(i) $\mathrm{Sec}^{r} C \subset \mathbb{P}^{n-1}$ is an irreducible variety of codimension $\max (n-2 r, 0)$.
(ii) The vector space of forms of degree $r+1$ vanishing on $\mathrm{Sec}^{r} C$ has dimension $\beta(r+1, n)$, where

$$
\beta(r, n)=\binom{n-r}{r}+\binom{n-r-1}{r-1}
$$

is the number of ways of choosing $r$ elements from $\mathbb{Z} / n \mathbb{Z}$ such that no two elements are adjacent.
(iii) If $n-2 r \geq 2$ then the homogeneous ideal $I\left(\operatorname{Sec}^{r} C\right)$ is generated by forms of degree $r+1$.
(iv) if $n-2 r=1$ then $\operatorname{Sec}^{r} C$ is a hypersurface of degree $n$.
(v) If $n-2 r \geq 1$ then $\operatorname{Sec}^{r} C$ has singular locus $\operatorname{Sec}^{r-1} C$.

Proof: (i) This is a general fact about curves. See for example [L, Section 1]. (ii), (iii), (iv). More generally the minimal free resolution for $I\left(\mathrm{Sec}^{r} C\right)$ was computed in [vBH, Section 8]. See [GP, Section 5] for the cases $r=1,2$, and [F2, Section 4] for further discussion.
(v) This is [vBH, Proposition 8.15].
9.1. Computing equations for higher secant varieties. The following two propositions may be used to compute equations for $\operatorname{Sec}^{r} C$ from equations for $C$. We say that a form $f$ vanishes on $C$ with multiplicity $r$ if (passing to affine co-ordinates) the Taylor expansion of $f$ at each point $P \in C$ begins with terms of order greater than or equal to $r$.

Proposition 9.2. Let $C \subset \mathbb{P}^{n-1}$ be a variety contained in no hyperplane. Let $f$ be a form of degree $r+1$.
(i) If $r \geq 1$ then

$$
f \in I\left(\operatorname{Sec}^{r} C\right) \Longleftrightarrow f \text { vanishes on } C \text { with multiplicity } r .
$$

(ii) If $r \geq 2$ then

$$
f \in I\left(\operatorname{Sec}^{r} C\right) \Longleftrightarrow \frac{\partial f}{\partial x_{i}} \in I\left(\operatorname{Sec}^{r-1} C\right) \text { for all } i=1, \ldots, n .
$$

Proof: (i) We choose $P_{1}, \ldots, P_{n} \in C$ spanning $\mathbb{P}^{n-1}$. By a change of co-ordinates we may assume $P_{1}=(1: 0: \ldots: 0), P_{2}=(0: 1: 0: \ldots: 0), \ldots, P_{n}=(0: 0: \ldots: 1)$. If $f \in I\left(\mathrm{Sec}^{r} C\right)$ then it vanishes on the linear span of any $r$ of the $P_{i}$. Therefore the monomials appearing in $f$ involve at least $r+1$ of the $x_{i}$, and since $f$ has degree $r+1$ must be squarefree. But then $f$ vanishes at $P_{1}$ with multiplicity $r$. Since $P_{1} \in C$ was arbitrary it follows that $f$ vanishes on $C$ with multiplicity $r$.

Conversely, suppose $f$ vanishes on $C$ with multiplicity $r$. Let $\Pi$ be an $(r-1)$ plane spanned by points $P_{1}, \ldots, P_{r} \in C$. By a change of co-ordinates we may assume $P_{1}=(1: 0: \ldots: 0), P_{2}=(0: 1: 0: \ldots: 0), \ldots$ Then $f\left(x_{1}, \ldots, x_{r}, 0, \ldots, 0\right)$ has total degree $r+1$, but has degree at most 1 in each of the variables. It follows that $f$ vanishes on $\Pi$. By definition $\operatorname{Sec}^{r} C$ is the Zariski closure of the union of all such $(r-1)$-planes. Therefore $f \in I\left(\operatorname{Sec}^{r} C\right)$ as required.
(ii) Since $\operatorname{char}(K)=0$ this follows from (i).

Now let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve. Taking $r=1$ in Theorem 9.1 shows that the homogeneous ideal $I(C)$ is generated by a vector space of quadrics of dimension $n(n-3) / 2$. Suppose we know a basis for this space. Then by repeatedly applying Proposition 9.2 (ii) we can find a basis for the space of forms of degree $r+1$ vanishing on $\mathrm{Sec}^{r} C$. Theorem 9.1(iii) tells us that if $n-2 r \geq 2$ then these forms define $\mathrm{Sec}^{r} C$. The following proposition covers the remaining case.

Proposition 9.3. Suppose $n-2 r=1$. Let $f$ be a form of degree $n$. If $r \geq 2$ then

$$
f \in I\left(\operatorname{Sec}^{r} C\right) \Longleftrightarrow \frac{\partial f}{\partial x_{i}} \in I\left(\operatorname{Sec}^{r-1} C\right)^{2} \text { for all } i=1, \ldots, n
$$

Proof: " $\Rightarrow$ " Let $H$ be the divisor of a hyperplane section, and let $P \in C$ be any point. Let $C_{+} \subset \mathbb{P}^{n}$ and $C_{-} \subset \mathbb{P}^{n-2}$ be the images of $C$ embedded via the linear systems
$|H \pm P|$. We choose co-ordinates so that the isomorphisms $C_{+} \rightarrow C \rightarrow C_{-}$are given by

$$
\left(x_{1}: \ldots: x_{n+1}\right) \mapsto\left(x_{1}: \ldots: x_{n}\right) \mapsto\left(x_{1}: \ldots: x_{n-1}\right)
$$

In particular $P$ is the point $\left(x_{1}: \ldots: x_{n}\right)=(0: \ldots: 0: 1)$. By Theorem 9.1 we know that $I\left(\mathrm{Sec}^{r-1} C_{-}\right)$is generated by forms $g_{1}, g_{2} \in K\left[x_{1}, \ldots, x_{n-1}\right]$ of degree $r$. By [F2, Corollary 2.3] there exist forms $h_{1}, h_{2} \in K\left[x_{1}, \ldots, x_{n}\right]$ of degree $r+1$ such that $f_{i}=x_{n+1} g_{i}+h_{i} \in I\left(\operatorname{Sec}^{r} C_{+}\right)$for $i=1,2$. Then $F=g_{1} h_{2}-g_{2} h_{1}$ belongs to

$$
I\left(\operatorname{Sec}^{r} C_{+}\right) \cap K\left[x_{1}, \ldots, x_{n}\right]=I\left(\operatorname{Sec}^{r} C\right)
$$

Since $g_{1}, g_{2}$ are coprime and $f_{1}, f_{2}$ are irreducible it is clear that $F$ is non-zero. By Theorem 9.1(iv) we have $I\left(\operatorname{Sec}^{r} C\right)=(F)$. We compute

$$
\frac{\partial F}{\partial x_{n}}=\frac{\partial f_{1}}{\partial x_{n+1}} \frac{\partial f_{2}}{\partial x_{n}}-\frac{\partial f_{1}}{\partial x_{n}} \frac{\partial f_{2}}{\partial x_{n+1}}
$$

On the other hand, for $i=1,2$ and $j=n, n+1$ we have

$$
\frac{\partial f_{i}}{\partial x_{j}} \in I\left(\mathrm{Sec}^{r-1} C_{+}\right) \cap K\left[x_{1}, \ldots, x_{n}\right]=I\left(\mathrm{Sec}^{r-1} C\right)
$$

Therefore $\frac{\partial F}{\partial x_{n}} \in I\left(\operatorname{Sec}^{r-1} C\right)^{2}$. Since $P \in C$ was arbitrary, and $C$ spans $\mathbb{P}^{n-1}$, the result follows.
" $\Leftarrow$ " Let $P_{1}, \ldots, P_{r}$ be $r$ distinct points on $C$. By a change of co-ordinates we may assume $P_{1}=(1: 0: \ldots: 0), P_{2}=(0: 1: 0: \ldots: 0), \ldots$ By Proposition 9.2 we know that $f$ vanishes on $C$ with multiplicity $2(r-1)+1=n-2$. Therefore $f\left(x_{1}, \ldots, x_{r}, 0, \ldots, 0\right)$ has total degree $n$, but has degree at most 2 in each of the variables. Since $2 r<n$ it follows that $f$ vanishes on the linear span of $P_{1}, \ldots, P_{r}$. By definition $\operatorname{Sec}^{r} C$ is the Zariski closure of the union of all such $(r-1)$-planes. Therefore $f \in I\left(\operatorname{Sec}^{r} C\right)$ as required.
9.2. Proof of Proposition 4.2. Let $C \subset \mathbb{P}^{n-1}$ be a genus one normal curve of degree $n$. Let $H$ be the divisor of a hyperplane section. We identify $\mathcal{L}(H)$ with the space of linear forms on $\mathbb{P}^{n-1}$. For $D$ an effective divisor on $C$ we write $\bar{D} \subset \mathbb{P}^{n-1}$ for the linear subspace cut out by $\mathcal{L}(H-D) \subset \mathcal{L}(H)$. We have

$$
\operatorname{Sec}^{r} C=\bigcup_{\operatorname{deg} D=r} \bar{D}
$$

We also put $D^{\circ}=\bar{D} \backslash \cup_{D^{\prime}<D} \overline{D^{\prime}}$. The gcd and lcm of divisors $\sum m_{P} P$ and $\sum m_{P}^{\prime} P$ are $\sum \min \left(m_{P}, m_{P}^{\prime}\right) P$ and $\sum \max \left(m_{P}, m_{P}^{\prime}\right) P$.
Lemma 9.4. Let $D, D_{1}, D_{2}$ be effective divisors on $C$.
(i) If $\operatorname{deg} D<n$ then $\operatorname{dim} \bar{D}=\operatorname{deg} D-1$.
(ii) The linear span of $\overline{D_{1}}$ and $\overline{D_{2}}$ is $\overline{\operatorname{lcm}\left(D_{1}, D_{2}\right)}$.
(iii) If $\operatorname{deg}\left(\operatorname{lcm}\left(D_{1}, D_{2}\right)\right)<n$ then $\overline{D_{1}} \cap \overline{D_{2}}=\overline{\operatorname{gcd}\left(D_{1}, D_{2}\right)}$.

Proof: (i) By Riemann-Roch we have $\operatorname{dim} \mathcal{L}(H-D)=n-\operatorname{deg} D$.
(ii) We have $\mathcal{L}\left(H-D_{1}\right) \cap \mathcal{L}\left(H-D_{2}\right)=\mathcal{L}\left(H-\operatorname{lcm}\left(D_{1}, D_{2}\right)\right)$.
(iii) The inclusion " $\supset$ " is clear. Equality follows by counting dimensions using (i) and (ii).

With the above notation, Proposition 4.2 becomes

Proposition 9.5. Suppose $n-2 r \geq 1$. Let $D=P_{1}+\ldots+P_{r}$ be an effective divisor of degree $r$ with $P_{1}, \ldots, P_{r} \in C$ distinct. Then for any $P \in D^{\circ}$ we have $T_{P} \operatorname{Sec}^{r} C=\overline{2 D}$.
Proof: If $P \in \overline{D^{\prime}}$ for $D^{\prime}$ an effective divisor of degree at most $r$, then by Lemma 9.4(iii) we have $D=D^{\prime}$. In particular $P \notin \operatorname{Sec}^{r-1} C$. It follows by Theorem 9.1(v) that $P$ is a smooth point on $\mathrm{Sec}^{r} C$. The next lemma shows that $\overline{2 D} \subset T_{P} \operatorname{Sec}^{r} C$, and equality follows by comparing dimensions, using Lemma 9.4(i) and Theorem 9.1(i).

Lemma 9.6. Let $X$ be an affine variety and $P_{1}, \ldots, P_{r} \in X$. Let $P=\sum \xi_{i} P_{i}$ where $\sum \xi_{i}=1$. If $\xi_{i} \neq 0$ then $T_{P_{i}} X \subset T_{P}\left(\operatorname{Sec}^{r} X\right)$.
Proof: There is a morphism $X \times \ldots \times X \rightarrow \operatorname{Sec}^{r} X ;\left(a_{1}, \ldots, a_{r}\right) \mapsto \sum \xi_{i} a_{i}$ with derivative $T_{P_{1}} X \times \ldots \times T_{P_{r}} X \rightarrow T_{P}\left(\operatorname{Sec}^{r} X\right) ;\left(b_{1}, \ldots, b_{r}\right) \mapsto \sum \xi_{i} b_{i}$.

In fact Proposition 9.5 is true without the hypothesis that $P_{1}, \ldots, P_{r}$ are distinct. However, since we do not need this, we omit the details.
9.3. Proof of Proposition 4.3. We must prove the following.

Proposition 9.7. Suppose $n-2 r=2$ and write $\operatorname{Sec}^{r} C=\left\{F_{1}=F_{2}=0\right\}$. Then the variety $X \subset \mathbb{P}^{n-1}$ defined by

$$
\operatorname{rank}\left(\begin{array}{lll}
\frac{\partial F_{1}}{\partial x_{1}} & \cdots & \frac{\partial F_{1}}{\partial x_{n}} \\
\frac{\partial F_{2}}{\partial x_{1}} & \cdots & \frac{\partial F_{2}}{\partial x_{n}}
\end{array}\right) \leq 1
$$

has codimension 3.
If $n=4$ then $C=\left\{F_{1}=F_{2}=0\right\} \subset \mathbb{P}^{3}$ is the intersection of two quadrics. There are 4 singular quadrics in the pencil spanned by $F_{1}$ and $F_{2}$, and each is singular at just one point. Then $X$ is the union of these 4 singular points, and so has codimension 3 .

We now generalise this argument. Let $H$ be the divisor of a hyperplane section. We identify $\mathcal{L}(H)$ with the space of linear forms on $\mathbb{P}^{n-1}$. Let $D_{1}$ and $D_{2}$ be divisors on $C$ of degree $r+1$ with $D_{1}+D_{2}=H$. Let $\Phi\left(D_{1}, D_{2}\right)$ be the $(r+1) \times(r+1)$ matrix of linear forms representing the multiplication map

$$
\mathcal{L}\left(D_{1}\right) \times \mathcal{L}\left(D_{2}\right) \rightarrow \mathcal{L}(H) .
$$

Since $\Phi\left(D_{1}, D_{2}\right)$ has rank at most 1 on $C$, it has rank at most $r$ on $\mathrm{Sec}^{r} C$. Therefore $\operatorname{det} \Phi\left(D_{1}, D_{2}\right)$ is a form of degree $r+1$ vanishing on $\operatorname{Sec}^{r} C$. In particular it belongs to the pencil spanned by $F_{1}$ and $F_{2}$.

Lemma 9.8. Every linear combination of $F_{1}$ and $F_{2}$ arises in this way. Moreover there are exactly 4 forms in the pencil arising as $\operatorname{det} \Phi\left(D_{1}, D_{2}\right)$ with $D_{1} \sim D_{2}$.

Proof: We say that divisor pairs $\left(D_{1}, D_{2}\right)$ and $\left(D_{1}^{\prime}, D_{2}^{\prime}\right)$ are equivalent if $D_{1} \sim D_{1}^{\prime}$ or $D_{1} \sim D_{2}^{\prime}$. It is shown in [F2, Lemma 2.9] that if $\left(D_{1}, D_{2}\right)$ and $\left(D_{1}^{\prime}, D_{2}^{\prime}\right)$ are inequivalent then $\operatorname{Sec}^{r} C=\left\{\operatorname{det} \Phi\left(D_{1}, D_{2}\right)=\operatorname{det} \Phi\left(D_{1}^{\prime}, D_{2}^{\prime}\right)=0\right\} \subset \mathbb{P}^{n-1}$. In particular these two forms are linearly independent.

We claim that the map $\left(D_{1}, D_{2}\right) \mapsto \Phi\left(D_{1}, D_{2}\right)$ is a bijection between the equivalence classes of divisor pairs and the pencil of forms spanned by $F_{1}$ and $F_{2}$. To prove this let $C$ be the image of an elliptic curve $E$ embedded in $\mathbb{P}^{n-1}$ by $\left|n .0_{E}\right|$. Then writing

$$
\operatorname{det} \Phi\left(r .0_{E}+P,(r+2) \cdot 0_{E}-P\right)=s(P) F_{1}+t(P) F_{2}
$$

for $P \in E$, we can see that $s / t$ is a rational function on $E$. It therefore defines a morphism $(s: t): E \rightarrow \mathbb{P}^{1}$. By the previous paragraph, this morphism is nonconstant, and indeed has fibres of the form $\{P,-P\}$. It must therefore be surjective. This proves the claim.

For the final statement we note that $r \cdot 0_{E}+P \sim(r+2) .0_{E}-P$ if and only if $P \in E[2]$.

Lemma 9.9. Let $S$ be the singular locus of $V=\left\{\operatorname{det} \Phi\left(D_{1}, D_{2}\right)=0\right\} \subset \mathbb{P}^{n-1}$. Then $S$ contains $\operatorname{Sec}^{r-1} C$. Moreover
(i) If $D_{1} \nsim D_{2}$ then $S=\operatorname{Sec}^{r-1} C$.
(ii) If $D_{1} \sim D_{2}$ then $S$ has codimension 3 .

Proof: Since $C$ spans $\mathbb{P}^{n-1}$ it is clear that for each $P \in \operatorname{Sec}^{r-1} C$ we have $T_{P} \operatorname{Sec}^{r} C=$ $\mathbb{P}^{n-1}$. Therefore $S$ contains $\operatorname{Sec}^{r-1} C$.
(i) Let $P \in V \backslash \mathrm{Sec}^{r-1} C$ be any point. According to [F2, Theorem 1.3] the $r \times r$ minors of $\Phi\left(D_{1}, D_{2}\right)$ generate $I\left(\operatorname{Sec}^{r-1} C\right)$. Therefore evaluating $\Phi\left(D_{1}, D_{2}\right)$ at $P$ gives a matrix of rank $r$. Moving $P$ to $(1: 0: \ldots: 0)$ and picking suitable bases for $\mathcal{L}\left(D_{1}\right)$ and $\mathcal{L}\left(D_{2}\right)$ we have

$$
\Phi\left(D_{1}, D_{2}\right)=x_{1}\left(\begin{array}{cccc}
0 & 0 & \cdots & 0 \\
0 & 1 & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & 1
\end{array}\right)+\Phi^{\prime}
$$

where $\Phi^{\prime}$ is an $(r+1) \times(r+1)$ matrix of linear forms in $x_{2}, \ldots, x_{n}$. Now the top left entry of $\Phi\left(D_{1}, D_{2}\right)$ is an equation for $T_{P} V$. Since the product of non-zero rational functions on $C$ is again non-zero, the entries of $\Phi\left(D_{1}, D_{2}\right)$ are non-zero. Therefore $P \in V$ is a smooth point.
(ii) Picking suitable bases for $\mathcal{L}\left(D_{1}\right)$ and $\mathcal{L}\left(D_{2}\right)$ we may suppose that $\Phi\left(D_{1}, D_{2}\right)$ is symmetric. Since $\left\{\operatorname{rank} \Phi\left(D_{1}, D_{2}\right) \leq r-1\right\} \subset S$, and the quadratic forms of rank at most $m-2$ have codimension 3 in the space of all quadratic forms in $m$ variables, it follows that $S$ has codimension at most 3. Suppose for a contradiction that $S$ has codimension at most 2 . Then its intersection with $\operatorname{Sec}^{r} C=\left\{F_{1}=F_{2}=0\right\}$ has codimension at most 3. But this intersection is contained in the singular locus of $\operatorname{Sec}^{r} C$, which by Theorem 9.1 has codimension 4 . This is the required contradiction.

To complete the proof of Proposition 9.7, we note that $X$ is the union of the singular loci of the hypersurfaces defined by linear combinations of $F_{1}$ and $F_{2}$. It follows by Lemmas 9.8 and 9.9 that $X$ has codimension 3.

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