# A lattice version of the KP equation 

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## 1. Introduction

Let $N$ and $M$ be relatively prime integers. Let $V_{N, M}$ be the set of all real valued functions $\psi$ on $\mathbf{Z} \times \mathbf{Z}$ satisfying $\psi(n+N, m)=\psi(n, m+M)=\psi(n, m) . V_{N, M}$ is a vector space of dimension $N M$ over $\mathbf{R}$. Let $A$ and $B$ be functions from an interval $I=(a, b)$ to $V_{N, M} . A(n, m, t)$ will denote the value of $A(t)$ at the point $(n, m) \in \mathbf{Z} \times \mathbf{Z}$. In $\S 3$, we will define two explicit real polynomial maps $f_{N, M}$ and $g_{N, M}$ on $V_{N, M} \times V_{N, M} \times \mathbf{R}^{3}$ to $V_{N, M}$. We will investigate solutions $A(t)$ and $B(t)$ to the following differential-difference equation:

$$
\begin{align*}
& \frac{d A(t)}{d t}=f_{N, M}(A(t), B(t), \alpha, \beta, \gamma)  \tag{1.1}\\
& \frac{d B(t)}{d t}=g_{N, M}(A(t), B(t), \alpha, \beta, \gamma) \tag{1.2}
\end{align*}
$$

for fixed $\alpha, \beta$ and $\gamma$. More intrisically, one may think of $f_{N, M}$ and $g_{N, M}$ as defining a vector field on $V_{N, M} \times V_{N, M}$ depending on parameters $\alpha, \beta$ and $\gamma$. Thus for any given $t$, $f_{N, M}(A(t), B(t), \alpha, \beta, \gamma)$ is a function on $\mathbf{Z} \times \mathbf{Z}$, and this function evaluated at ( $n, m$ ) is a polynomial in $\alpha, \beta$ and $\gamma$ and the numbers $A(i, j, t)$ and $B(i, j, t)$ which will turn out to be of degree 4 , and $g_{N, M}(A(t), B(t), \alpha, \beta, \gamma)$ will turn out to be of degree 5 . Actually these polynomials enjoy certain homogeneity properties explained in $\S 3$.

These equations are derived from a certain algebro-geometric construction, which is in some sense a variant of a construction of Mumford and van Moerbeke (as will be explained in $\S 3$ ). This construction starts with certain algebraic curves $X$ with a distinguished point $P$ (with certain additional structure). Using $X$ and this structure, we
will define a map $\Phi$ from an open subset $U$ of the Jacobian of $X, \operatorname{Jac}(X)$, to $V_{N, M} \times V_{N, M}$. Let $\Psi$ be the canonical map from $X$ to $\operatorname{Jac}(X)$, this map being canonical up to translation. Let $T_{1}, T_{2}$ and $T_{3}$ be a basis of the osculating three space of the curve $\Psi(X)$ at the point $\Psi(P)$. We choose $T_{1}$ to be a tangent vector to the curve at $\Psi(P)$ and $T_{2}$ to be in the osculating plane of the curve at $\Psi(P)$. These vectors can be translated over the whole of $\operatorname{Jac}(X)$ to obtain vector fields again denoted by $T_{i}$. We will define $f_{N, M}$ and $g_{N, M}$ in such a way that the following holds: Suppose we start with a line bundle $\mathscr{L} \in U$ and allow it to flow along the vector field $T_{i}$ for time $t$ to a line bundle $\mathscr{L}_{t}$. Then there are $\alpha, \beta$ and $\gamma$ (depending on $i$ ) so that if we let $Y=\Phi(\mathscr{L})$ flow along the vector field on $V_{N, M} \times V_{N, M}$ defined above by (1.1) and (1.2) to $Y_{t}$, then $\Phi\left(\mathscr{L}_{l}\right)=Y_{t}$. Thus the complicated flow defined by the non-linear equations (1.1) and (1.2) can be 'linearized' to a straight line flow on a Jacobian. Conjecturally, the generic $A$ and $B$ come from such a curve and line bundle. We also write down explicit conserved quantities of these flows.

Having derived these equations, we can ask about the behavior of solutions of these equations, especially with $N$ and $M$ large. One way of analyzing the behavior of such solutions is to try to construct a continuous model for these equations. Another way is to look graphically at numerical solutions of our equations. For an interesting account of these two ways, see [Z]. We have not attempted such an analysis in this paper. Instead, we will exhibit some solutions to our equations which do have interesting continuous models. Thus our results indicate that such an analysis would be interesting. One way of precisely defining these rather vague comments is the following rather ad hoc definition:

Definition 1.1. $\mathscr{C}$ is the class of all functions $f$ on $\mathbf{R}^{3}$ satisfying the following properties:
(i) $f(x+1, y, t)=f(x, y+1, t)=f(x, y, t)$ for all $(x, y, t) \in \mathbf{R}^{3}$.
(ii) Given $\varepsilon>0$, there are $(\alpha, \beta, \gamma) \in \mathbf{R}^{3}$, an integer $N$, a constant $C$ and functions $A(t)$ and $B(t)$ from $\mathbf{R}$ to $V_{N, N^{2}+1}$ so that

$$
\left|f\left(\frac{n}{N}, \frac{m}{N^{2}+1}, t\right)-N A(n, m, t)-C\right|<\varepsilon
$$

and so that $A$ and $B$ satisfy the equations (1.1) and (1.2).
The main result of this paper (Theorem 2.5) is that $\mathscr{C}$ contains many of the solutions of the KP equation arising from algebraic geometry [D]. The definition of the class $\mathscr{C}$ is quite restrictive in that we require the discrete $\operatorname{NA}(n, m, t)$ to be close to $f(x, y, t)$ for all $t$. It would be interesting to know what further conditions on the class $\mathscr{C}$
would imply that an $f \in \mathscr{C}$ satisfying these further conditions satisfies the KP equation. Another question unanswered here is whether a solution to the KP hierarchy (in three variables) belongs to some variant of the class $\mathscr{C}$. The definition of the class $\mathscr{C}$ is (to be frank) based on what we can prove.
§ 2 reviews the theory of curves and their Jacobians defined over R. § 2 concludes with a precise statement of our main theorem. $\S 3$ derives the expressions for $f_{N, M}$ and $g_{N, M}$. We conclude $\S 3$ with a few observations on the relation of this work to the work of Mumford and van Moerbeke [MM] on spectral curves. In §4, we give a proof of the following theorem:

Theorem 1.1. If $C$ is non-hyperelliptic and $V$ is a generic three dimensional subspace of $H^{0}(C, \Omega)$, then the map from $V \otimes H^{0}(C, \Omega) \rightarrow H^{0}\left(C, \Omega^{\otimes 2}\right)$ is surjective.

The proof of Theorem 1.1 was supplied by Lazarsfeld based on the ideas of [GL]. Green also supplied a proof, and Eisenbud provided a simpler proof by direct computation for trigonal curves, and so for a generic curve. § 5 develops some Kodaira-Spencer type deformation theory. In $\S 6$ we show that a certain class of 'good' curves exists using a monodromy argument as well as our Kodaira-Spencer theory and Theorem 1.1. § 7 gives the proof of our main theorem.

The work in this paper was motivated by a hope of Trubowitz that understanding the spectral theory of lattice models of the KP equation might yield some insight into the transcendental spectral theory of the KP equation.

## 2. Curves defined over $\mathbf{R}$

Let $C$ be a non-singular curve defined over $C$, i.e. a compact Riemann surface. Thus we can find a holomorphic embedding of $C$ into $\mathbf{P}^{n}$ so that $C$ is the locus of zeros of homogeneous polynomials with complex coefficients. We say that $C$ is defined over $\mathbf{R}$ if we can choose the embedding so that the polynomials all have real coefficients. Note that $\mathbf{P}^{n}$ has a natural antiholomorphic involution

$$
\iota\left(z_{0}, \ldots, z_{n}\right)=\left(\bar{z}_{0}, \ldots, \bar{z}_{n}\right)
$$

and that $\iota$ leaves $C$ invariant when $C$ is defined over $\mathbf{R}$. We denote the restriction of $\iota$ to $C$ by $\iota$ again. A function $f$ on an $\iota$ invariant open set of $C$ is said to be defined over $\mathbf{R}$ if $f\left(z^{l}\right)=\overline{f(z)}$. A point or divisor is defined over $\mathbf{R}$ if it is invariant under $\iota$. A holomorphic one form $\omega$ is defined over $\mathbf{R}$ if locally $\omega=d f$, where $f$ is defined over $\mathbf{R}$. If $\omega$ is defined over $\mathbf{R}$ and $\gamma$ is a path on the surface, then

$$
\int_{\gamma^{\prime}} \omega=\overline{\int_{\gamma^{\prime}}} \omega
$$

as we can readily see by dividing the path $\gamma$ into subpaths on which $\omega$ is exact.
Note that $\iota$ acts on the cohomology $H^{1}(C, \mathbf{Z})$. Let $\Lambda^{+}(C)$ be the elements of $H^{1}(C, \mathbf{Z})$ fixed by $\iota$ and let $\Lambda^{-}$be the elements $\gamma$ of $H^{1}(C, \mathbf{Z})$ with $\gamma^{\iota}=-\gamma$. Note there is a natural integration map:

$$
\int: H^{\mathrm{i}}(C, \mathbf{Z}) \rightarrow H^{0}(C, \Omega)^{*}
$$

where $H^{0}(C, \Omega)^{*}$ is the dual pace of the holomorphic one forms of $C$. Let $H^{0}(C, \Omega)^{*}(\mathbf{R})$ be the set of real points of $H^{0}(C, \Omega)^{*}$. Then $\Lambda^{+}$maps to $H^{0}(C, \Omega)^{*}(\mathbf{R})$ and $\Lambda^{-}$maps to $i H^{0}(C, \Omega)^{*}(\mathbf{R}) . H^{0}(C, \Omega)^{*}(\mathbf{R})$ is a real vector space of dimension $g$, and the complex span of the vectors in $H^{0}(C, \Omega)^{*}(\mathbf{R})$ is just $H^{0}(C, \Omega)^{*} . H^{1}(C, \mathbf{Z})$ maps to a lattice in $H^{0}(C, \Omega)^{*}$. Thus we see that $\Lambda^{+}$maps to a lattice in $H^{0}(C, \Omega)^{*}(\mathbf{R})$. The Jacobian of $C$ has a natural real structure and the quotient of $H^{0}(C, \Omega)^{*}(\mathbf{R})$ by the image of $\Lambda^{+}$is the component of the real points of the Jacobian of $C$ which contains the identity of the Jacobian.

We next discuss theta functions following [M]. We regard $H^{1}(C, Z)$ as a subgroup of $H^{0}(C, \Omega)^{*}$. The involution $\iota$ extends to an antiholomorphic involution on $H^{0}(C, \Omega)^{*}$ again denoted by $\iota$. Let $\mathbf{C}_{1}^{*}$ be the set of all complex numbers of absolute value one. Choose a map

$$
\alpha: H^{1}(C, \mathbf{Z}) \rightarrow \mathbf{C}_{1}^{*}
$$

so that

$$
\begin{equation*}
\alpha\left(u^{\imath}\right)=\overline{\alpha(u)} \tag{2.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\alpha\left(u_{1}+u_{2}\right)}{\alpha\left(u_{1}\right) \alpha\left(u_{2}\right)}=e^{i \pi\left(u_{1}, u_{2}\right\rangle} \tag{2.2}
\end{equation*}
$$

There is a unique Hermitian form $H$ on $H^{0}(C, \Omega)^{*}$ so that

$$
\operatorname{Im} H(x, y)=\langle x, y\rangle
$$

We see that since

$$
\left\langle x^{l}, y^{l}\right\rangle=-\langle x, y\rangle,
$$

we have that

$$
\overline{H\left(x^{l}, y^{\prime}\right)}=H(x, y) .
$$

Let $\vartheta$ defined on $H^{0}(C, \Omega)^{*}$ be the function satisfying the functional equation

$$
\vartheta(z+u)=\alpha(u) e^{\pi H(z, u)+\pi H(u, u) / 2} \vartheta(z)
$$

for $z \in H^{0}(C, \Omega)^{*}$ and $u \in H^{1}(C, \mathbf{Z})$. The function $\vartheta^{\prime}$ defined by

$$
\vartheta^{\prime}(z)=\overline{\vartheta\left(z^{l}\right)}
$$

satisfies the same functional equation as $\vartheta$. Since $\vartheta$ is defined up to a constant multiple by its functional equation, we see that we can choose $\vartheta$ to be real on the fixed set of $\iota$, which is $H^{0}(C, \Omega)(\mathbf{R})^{*}$. Consider $K_{1} \in H^{0}(C, \Omega)^{*}$ and suppose that $K_{1}^{\prime} \in K_{1}+H^{1}(C, \mathbf{Z})$. Choose a point $P \in C$ and a parameter $z$ around $P$. We can define linear functionals in $H^{0}(C, \Omega)(\mathbf{R})^{*}$ by the formulas:

$$
v_{i}(\omega)=\left(\frac{d^{i-1}}{d z^{i-1}} \frac{\omega}{d z}\right)_{z=0} .
$$

The $v_{i}$ form a Frenet frame for the natural map $\phi$ of a neighborhood of $P$ in $C$ to $H^{0}(C, \Omega)^{*}$ defined by the formula

$$
\phi(Q)(\omega)=\int_{P}^{Q} \omega .
$$

We call the span of $v_{1}, v_{2}$ and $v_{3}$ the osculating three space at $P$.
Define a meromorphic function $f$ on $\mathbf{R}^{3}$ by the formula

$$
\begin{equation*}
f(x, y, t)=\frac{\partial^{2}}{\partial x^{2}} \log \vartheta\left(x v_{1}+y v_{2}+t v_{3}+K_{1}\right) . \tag{2.3}
\end{equation*}
$$

Then

$$
\begin{aligned}
\overline{f(x, y, t)} & =\frac{\partial^{2}}{\partial x^{2}} \log \vartheta\left(x v_{1}+y v_{2}+t v_{3}+K_{1}^{t}\right) \\
& =\frac{\partial^{2}}{\partial x^{2}} \log \vartheta\left(x v_{1}+y v_{2}+t v_{3}+K_{1}+u\right) \\
& =\frac{\partial^{2}}{\partial x^{2}} \log \vartheta\left(x v_{1}+y v_{2}+t v_{3}+K_{1}\right) \\
& =f(x, y, t)
\end{aligned}
$$

for some $u \in H^{1}(C, \mathbf{Z})$. We have used that the second logarithmic derivative of the $\vartheta$ function is periodic, as follows from differentiating the functional equation.

Let $\mathcal{M}_{g, 1}$ be the set of all $(C, P)$, where $C$ is a curve of genus $g$ and $P$ is a point on $C$. $\mathcal{M}_{g, 1}$ has the natural structure of an analytic space. Let $\mathcal{M}_{g, 1}(\mathbf{R})$ be the subset of $\mathcal{M}_{g, 1}$ consisting of all curves $C$ defined over $\mathbf{R}$ and $P$ a point of $C$ defined over $\mathbf{R}$. Let $(C, P) \in \mathcal{M}_{g, 1}$.

Definition 2.1. $v_{i} \in H^{0}(C, \Omega((i+1) P))$ for $i$ from 1 to 3 are called adapted if the $v_{i}$ map to linearly independent elements of $H^{0}(C, \Omega)^{*}$.

Define

$$
\phi: H^{0}(C, \Omega(-P)) \oplus H^{0}(C, \Omega(-2 P)) \oplus H^{0}(C, \Omega(-3 P)) \rightarrow H^{0}\left(C, \Omega^{\otimes 2}(P)\right)
$$

by

$$
\phi\left(\omega_{1}, \omega_{2}, \omega_{3}\right)=v_{1} \omega_{1}+v_{2} \omega_{2}+v_{3} \omega_{3},
$$

where $H^{0}\left(C, \Omega^{\otimes 2}(P)\right)$ is the set of quadratic differentials which have a pole at $P$.
Definition 2.2. The $v_{i}$ are acceptable if the map $\phi$ is injective.
Let $V_{j}$ be the annihilator of $H^{0}(C, \Omega(-j P))$ in $H^{0}(C, \Omega)^{*}$, where $H^{0}(C, \Omega(-j P))$ is the set of all one forms which vanish $j$ times at $P$. We assume that $V_{3}$ has dimension 3, which is the same as assuming that there is not a non-trivial function having a pole of order 3 or less at $P$. Let $\Lambda^{+}(\mathbf{R})$ be the real span of the vectors in $\Lambda^{+}$. Note that both $\Lambda^{+}(\mathbf{R})$ and $H^{0}(C, \Omega)(\mathbf{R})$ are in $H^{1}(C, \mathbf{C})$ and that cup product on $H^{1}(C, \mathbf{C})$ induces a perfect pairing between these two real vector spaces and hence that we have a natural isomorphism from $\Lambda^{+}(\mathbf{R})$ to $H^{0}(C, \Omega)^{*}(\mathbf{R})$. Choose $v_{1}, v_{2}$ and $v_{3}$ in $\Lambda^{+}(\mathbf{R})$ so that $v_{i} \in V_{i}$. $H^{0}(C, \Omega((j+1) P))$ is also included in $H^{1}(C, C)$, since all the differentials in $H^{0}(C, \Omega((j+1) P))$ are of the second kind. Further $H^{0}(C, \Omega((j+1) P))$ is the annihilator of $H^{0}(C, \Omega(-j P))$ under cup product. It follows that $v_{j} \in H^{0}(C, \Omega((j+1) P))$.

Definition 2.3. ( $C, P) \in \mathcal{M}_{g, 1}(\mathbf{R})$ is good if:
(i) $\Lambda^{+} \cap H^{0}(C, \Omega(2 P))$ has rank one in $\Lambda^{+}$.
(ii) $\Lambda^{+} \cap H^{0}(C, \Omega(3 P))$ has rank two in $\Lambda^{+}$.
(iii) There are adapted $v_{i} \in \Lambda^{+}$which are acceptable.
(iv) The dimension of $V_{3}$ is three.

Proposition 2.4. The good points of $\mathcal{M}_{g, 1}(\mathbf{R})$ are dense (in the classical topology) if $g>2$.

We next state our main theorem.
Theorem 2.5. Suppose that $K_{1}-K_{1}^{t}$ is in the image of $H^{1}(C, Z)$ and that the pair $(C, P)$ is good. Let the $v_{i}$ be an adapted acceptable set in $\Lambda^{+}$and let $\bar{v}_{i}$ be the image of $v_{i}$ in $H^{0}(C, \Omega)^{*}$. Suppose that the function

$$
\vartheta\left(x \bar{v}_{1}+y \bar{v}_{2}+t \bar{v}_{3}+K_{1}\right)
$$

does not vanish on $\mathbf{R}^{3}$. Let

$$
f(x, y, t)=\frac{\partial^{2}}{\partial x^{2}} \log \vartheta\left(x \bar{v}_{1}+y \bar{v}_{2}+t \bar{v}_{3}+K_{1}\right)
$$

Then $f$ is in the class $\mathscr{C}$.
Under the hypotheses of Theorem 2.5 , it is well known that there is a constant $K$ so that $f+K$ satisfies the KP equation.

## 3. Equations of motion

Before deriving the formulas for $f_{N, M}$ and $g_{N, M}$, we define the homogeneity properties of these polynomials mentioned in the Introduction. Define an $\mathbf{R}^{*}$ action on the space of all polynomials on $V_{N, M} \times V_{N, M} \times \mathbf{R}^{3}$ by

$$
P^{\lambda(s)}(A, B, \alpha, \beta, \gamma)=P\left(s A, s^{2} B, s \alpha, s^{2} \beta, s^{3} \gamma\right)
$$

We say $P$ has weight $r$ if

$$
P^{\lambda(s)}=s^{r} P
$$

Thus the polynomial $P_{i, j}$ defined by $P_{i, j}(A, B, \alpha, \beta, \gamma)=A(i, j)$ has weight 1 . By abuse of notation, we will denote $P_{i, j}$ by $A(i, j)$. Similarly, $B(i, j)$ will denote the analogous polynomial of weight 2 , while $\alpha, \beta$ and $\gamma$ will denote the analogous polynomial of degree 1,2 , and 3 respectively. We will show that our expressions for $f_{N, M}$ and $g_{N, M}$ will have weights four and five respectively.

Let $X$ be a smooth curve defined over $\mathbf{R}$ of genus $g$, and let $P$ and $Q$ be real points of $X$. Suppose that $N(P-Q)$ is linearly equivalent to 0 . Thus there is a function $\alpha$ on $X$ having a pole of order $N$ at $P$ and a zero of order $N$ at $Q$, and having no other poles or zeros. Let $R_{i}$ and $S_{i}$ be points of $X$ for $i=1, \ldots, M$ so that $R_{i}+S_{i}$ is defined over $\mathbf{R}$.

Suppose there is a function $\beta$ having divisor

$$
M(P+Q)-\sum_{i=1}^{M}\left(R_{i}+S_{i}\right) .
$$

For general $i$, we define $R_{i}$ and $S_{i}$ by periodicity, $R_{i}=R_{[i]}$ and $S_{i}=S_{[i]}$, where $[i]$ is the positive residue of $i \bmod m$. Let $\mathscr{L}$ be a line bundle of degree $g$ on $X$. Let

$$
\mathscr{L}_{n, m}=\mathscr{L}\left((n+m) P+(m-n) Q+D_{m}\right)
$$

Here $D_{0}=0$ and $D_{m+1}=D_{m}-R_{m+1}-S_{m+1}$.
Definition 3.1. $\mathscr{L}$ is non-degenerate if $H^{0}\left(X, \mathscr{L}_{n, m}(-P)\right)=0$.
The Riemann-Roch theorem then implies that $h^{0}\left(X, \mathscr{L}_{n, m}\right)=1$ if $\mathscr{L}$ is non-degenerate. Assume $\mathscr{L}$ is nondegenerate. Let $z$ be a parameter defined at $P$ and choose a section $s_{0,0}$ of $\mathscr{L}$. There is a nonzero section $s_{n, m}$ of $H^{0}\left(X, \mathscr{L}_{n, m}\right)$, which is defined up to a constant, and $s_{n, m}$ considered as a meromorphic section of $\mathscr{L}$ has a pole of order exactly $n+m$ at $P$. Let

$$
f_{n, m}=\frac{s_{n, m}}{s_{0,0}}
$$

We normalize $s_{n, m}$ so that

$$
f_{n, m} z^{n+m}(P)=1
$$

i.e. the leading term in the Laurent expansion of $f$ in terms of $z$ is one. We can also normalize $\alpha$ and $\beta$ so that $\left(\alpha z^{N}\right)(P)=1$ and $\left(\beta z^{M}\right)(P)=1$

Given a non-degenerate line bundle $\mathscr{L}$ and the parameter $z$, we can form several functions on $\mathbf{Z} \times \mathbf{Z}$, namely

$$
d_{i}(n, m)=\left(\frac{d^{i} f_{n, m} z^{n+m}}{i!d z^{i}}\right)
$$

In this paper, we will be mostly considering $d_{1}, d_{2}$ and $d_{3}$. They are the coefficients of the Laurent expansion of $f_{n, m}$.

First, let's notice that we can write down a linear relation between $s_{n+1, m}, s_{n, m+1}$, $s_{n, m}$ and $s_{n-1, m}$ in terms of the functions $d_{1}$ and $d_{2}$. Such a relation must exist, since $s_{n+1, m}, s_{n, m+1}, s_{n, m}$ and $s_{n-1, m}$ are all in $H^{0}\left(\mathscr{L}_{n, m}(P+Q)\right)$ and $h^{0}\left(\mathscr{L}_{n, m}(P+Q)\right)=3$ by Rie-mann-Roch and our assumptions on non-degeneracy. Fixing $n$ and $m$ for the moment,
we can write

$$
a f_{n+1, m}+b f_{n, m+1}=c f_{n, m}+d f_{n-1, m} .
$$

Both $f_{n+1, m}$ and $f_{n, m+1}$ have a pole of order $n+m+1$ at $P$ and $f_{n+1, m} / f_{n, m+1}$ has value 1 there. Both $f_{n, m}$ and $f_{n-1, m}$ have poles of order less than $n+m+1$ at $P$. Thus we have $a=b \neq 0$. So we may choose $a=-b=1$. Thus we may write

$$
s_{n+1, m}-s_{n, m+1}=A(\mathscr{L}, n, m) s_{n, m}+B(\mathscr{L}, n, m) s_{n-1, m} .
$$

$A$ and $B$ are uniquely determined. We denote $A(\mathscr{L}, n, m)$ by $A(n, m)$ when $\mathscr{L}$ is understood. By comparing the Laurent expansions of the above equations around $z=0$, we see that we have the following recursion relations:

$$
\begin{gather*}
d_{1}(n+1, m)-d_{1}(n, m+1)=A(n, m)  \tag{3.1}\\
d_{2}(n+1, m)-d_{2}(n, m+1)=A(n, m) d_{1}(n, m)+B(n, m)  \tag{3.2}\\
d_{3}(n+1, m)-d_{3}(n, m+1)=A(n, m) d_{2}(n, m)+B(n, m) d_{1}(n, m) . \tag{3.3}
\end{gather*}
$$

The $d_{i}$ also have periodicity properties with respect to translation by $(N, 0)$ and $(0, M)$. Specifically, let

$$
\alpha=z^{-N}+a_{1} z^{-N+1}+\ldots
$$

and

$$
\beta=z^{-M}+b_{1} z^{-M+1}+\ldots
$$

be the Laurent expansions of $\alpha$ and $\beta$. Note that $\alpha s_{n, m} \in H^{0}\left(\mathscr{L}_{n+N, m}\right)$. So $\alpha s_{n, m}$ is a constant multiple of $s_{n+N, m}$. By our normalization, this constant must be 1 so

$$
\begin{equation*}
\alpha s_{n, m}=s_{n+N, m} . \tag{3.4}
\end{equation*}
$$

Similarly,

$$
\begin{equation*}
\beta s_{n, m}=s_{n, m+M} . \tag{3.5}
\end{equation*}
$$

Consider the Laurent expansions of the equation (3.4). Comparing coefficients we see that

$$
\begin{gathered}
d_{1}(n+N, m)=a_{1}+d_{1}(n, m) \\
d_{2}(n+N, m)=a_{2}+a_{1} d_{1}(n, m)+d_{2}(n, m)
\end{gathered}
$$

and

$$
d_{3}(n+N, m)=a_{3}+a_{2} d_{1}(n, m)+a_{1} d_{2}(n, m)+d_{3}(n, m)
$$

We have similar formulas for $d_{i}(n, m+M)$. Note that $A(n+N, m)=A(n, m)=A(n, m+M)$ and that $B(n, m+M)=B(n, m)=B(n+N, m)$.

The key observation here is that given the $a_{i}$ for $i$ from 1 to 3 , we can compute the $d_{i}$ and the $b_{i}$ for $i$ from 1 to 3 in terms of $A$ and $B$ by universal polynomials which depend only on $N$ and $M$. Since $f_{0,0}=1$, we have $d_{i}(0,0)=0$. Hence we can use the recurrence relation (3.1) to solve for $d_{1}(l,-l)$ directly. Since $N$ and $M$ are relatively prime, by the Chinese remainder theorem for any $n$ and $m$, we can find $a$ and $b$ so that

$$
n=l+a N
$$

and

$$
m=-l+b M .
$$

So

$$
d_{1}(n, m)=d_{1}(l,-l)+a_{1} a+b_{1} b .
$$

We have

$$
d_{1}(N M,-N M)=M a_{1}-N b_{1} .
$$

But $d_{1}(N M,-N M)$ is expressed directly in terms of the $A$ 's. So we can determine $b_{1}$ in terms of the $A$ 's and $a_{1}$. We see $b_{1}$ and $d_{1}$ have weight 1 . Having determined $d_{1}$, we can now determine $d_{2}(n, m)$ from the recurrence relation (3.2) and the Chinese remainder theorem. We can similarly find an expression for $b_{2}$ in terms of the $A$ 's, the $B$ 's, and $a_{1}$ and $a_{2}$. Finally, $d_{3}$ and $b_{3}$ are determined in the same way. Note that these formulas only involve the $a_{i}$, and $A$ and $B$, and not $X$ or $\mathscr{L}$. Further, the $d_{i}$ and $b_{i}$ have weight $i$.

Actually, one can continue this process and find that all the $b_{i}$ can be expressed in terms of the $A$ 's, the $B$ 's and the $a_{i}$. If we allow the $\mathscr{L}$ to evolve while fixing the curve, the points $P, Q$, and the $R_{i}$ and $S_{i}$ as well as the parameter $z$, the $b_{i}$ will of course remain constant. This means that the $b_{i}$ are conserved quantities of such an evolution.

Let $D_{0}$ be a fixed divisor of degree $g$ on $X$, and let $J_{g}$ be the Jacobian of all line bundles of degree $g$ on $X$. If $z$ is a point of $X$ close to $P$, we can define a linear functional $\phi(z)$ on $H^{0}(X, \Omega)$ by the formula:

$$
\phi(z)(\omega)=\int_{P}^{z} \omega .
$$

There is a natural analytic homomorphism

$$
\Phi: H^{0}(X, \Omega)^{*} \rightarrow J_{g}
$$

so that $\Phi(0)=\mathcal{O}\left(D_{0}\right)$, and so that $\Phi(\phi(z)+\alpha)=\Phi(\alpha)(z-P)$. The kernel of $\Phi$ is just $H^{1}(X, Z)$. On the other hand, we have the previously introduced linear functionals defined on $H^{0}(X, \Omega)$ :

$$
v_{i}(\omega)=\left(\frac{d^{i-1}}{d z^{i-1}} \frac{\omega}{d z}\right)_{z=0}
$$

Let

$$
A_{n, m}^{\prime}(f)=A(\Phi(f), n, m)
$$

for $f \in H^{0}(X, \Omega)^{*}$. Our aim is to compute

$$
\nabla A_{n, m}^{\prime} \cdot v_{3}
$$

which is the directional derivative of $A^{\prime}$ in the direction $v_{3}$, in terms of polynomials in the $A$ 's and $B$ 's. For $t$ small, let $D(t)$ be the divisor $D_{1}+D_{2}+D_{3}$ on $X$, where the $z\left(D_{i}\right)$ are the three cube roots of $t$. We first show that we have the formula:

$$
\nabla A_{n, m}^{\prime}(f) \cdot v_{3}=2 \frac{d}{d z} A(\Phi(f)(D(z)-3 P), n, m)_{z=0}
$$

Let $w(x)$ be the point of $X$ so that $z(w(x))=x$ for $x$ small. Let $\psi(x)=f+\Phi^{-1}(D(x)-3 P)$ so that $\Phi(\psi(x))=\Phi(f)(D(x)-3 P)$. Then

$$
\begin{equation*}
\psi\left(x^{3}\right)=\phi(x)+\phi(\zeta x)+\phi\left(\zeta^{2} x\right)+f \tag{3.6}
\end{equation*}
$$

where $\zeta$ is a primitive cube root of 1 . We have

$$
\frac{d}{d z} A(\Phi(f)(D(z)-3 P), n, m)_{z=0}=\nabla A^{\prime}(f)_{n, m} \cdot \psi^{\prime}(0)
$$

and

$$
\begin{aligned}
& v_{1}=\phi^{\prime}(0) \\
& v_{2}=\phi^{\prime \prime}(0)
\end{aligned}
$$

and

$$
v_{3}=\phi^{\prime \prime \prime}(0)
$$

So our claim will follow from

$$
\phi^{\prime \prime \prime}(0)=2 \psi^{\prime}(0)
$$

However, this follows by differentiating the identity (3.6) three times and setting $x=0$.
Note that $s_{n, m}, s_{n+1, m}, s_{n+2, m}$ and $s_{n+3, m}$ are a basis of $H^{0}\left(\mathscr{L}_{n, m}(3 P)\right)$. For $t$ small, let $\mathscr{L}_{n, m, t}$ be the line bundle $\mathscr{L}_{n, m}\left(3 P-D_{1}-D_{2}-D_{3}\right)$, where the $z\left(D_{i}\right)$ are the three cube roots of $t$. Note that $h^{0}\left(\mathscr{L}_{n, m, t}(-P)\right)=0$, for $t$ sufficiently small, as $\mathscr{L}_{n, m, t}(-P) \rightarrow \mathscr{L}_{n, m}(-P)$ as $t \rightarrow 0$. Let $s_{n, m, t}$ be a non-zero section of $\mathscr{L}_{n, m, t}$ varying holomorphically with $t$. We can write

$$
\begin{equation*}
s_{n, m, t}=\sum_{i=0}^{3} a_{i}(n, m, t) s_{n+i, m} \tag{3.7}
\end{equation*}
$$

where the $a_{i}$ are all holomorphic in $t$ and one is non-zero at $t=0$. In fact, since $s_{n, m, t} \rightarrow s_{n, m}$ as $t \rightarrow 0$ modulo multiplication by constants, we see that $a_{0}$ is non-zero.

Let

$$
\begin{equation*}
a_{i}(n, m, t)=\sum_{j} a_{i j}(n, m) t^{j} \tag{3.8}
\end{equation*}
$$

be the Taylor expansion of $a_{i}(n, m, t)$, where we may assume that $a_{0,0}(n, m)=1$. We have the identity

$$
\begin{equation*}
0=s_{n, m, t^{3}}(t), \tag{3.9}
\end{equation*}
$$

since $s_{n, m, t^{3}}$ vanishes at all the cube roots of $t^{3}$, including $t$. On the other hand, we have

$$
\begin{equation*}
s_{n+i, m}(z)=\sum_{j=0}^{\infty} d_{j}(n+i, m) z^{-n-i-m+j} s_{0,0} . \tag{3.10}
\end{equation*}
$$

Now substitute (3.10) and (3.8) into (3.7) and (3.7) into (3.9) and compute the first few nonzero coefficients of $t$. The coefficient of $t^{-n-m}$ in $s_{n, m, t^{3}}(t)$ is just $a_{3,1}(n, m)+$ $a_{0,0}(n, m)$, since

$$
a_{3,0}(n, m)=a_{2,0}(n, m)=a_{1,0}(n, m)=0
$$

Thus $a_{3,1}(n, m)=-1$, since $a_{0,0}=1$. We can replace $s_{n, m, t}$ by $s_{n, m, t} t / a_{3}(n, m, t)$ and assume that $a_{3}(n, m, t)=-t$ for all $n$ and $m$. With our new choice, we have

$$
\frac{s_{n, m, t} z^{n+m}}{s_{0,0, t}}(P)=1 .
$$

So we can write

$$
\begin{equation*}
s_{n+1, m, t}-s_{n, m+1, t}=A\left(\mathscr{L}_{t}, n, m\right) s_{n, m, t}+B\left(\mathscr{L}_{t}, n, m\right) s_{n-1, m, t} \tag{3.11}
\end{equation*}
$$

The coefficient of $t^{-n-m+1}$ in $s_{n, m, t^{3}}(t)$ is just $a_{2,1}(n, m)-d_{1}(n+3, m)+d_{1}(n, m)$, so

$$
\begin{equation*}
a_{2,1}(n, m)=d_{1}(n+3, m)-d_{1}(n, m) \tag{3.12}
\end{equation*}
$$

So $a_{2,1}$ has weight 1 . The coefficient of $t^{-n-m+2}$ in $s_{n, m, t^{3}}(t)$ is just

$$
a_{1,1}(n, m)+a_{2,1}(n, m) d_{1}(n+2, m)-d_{2}(n+3, m)+d_{2}(n, m)
$$

so

$$
\begin{equation*}
a_{1,1}(n, m)=-d_{1}(n+2, m) a_{2,1}(n, m)+d_{2}(n+3, m)-d_{2}(n, m) . \tag{3.13}
\end{equation*}
$$

So $a_{1,1}(n, m)$ has weight 2 . The coefficient of $t^{-n-m+3}$ in $f_{n, m, t^{3}}(t)$ is just

$$
a_{0,1}(n, m)+a_{1,1}(n, m) d_{1}(n+1, m)+a_{2,1}(n, m) d_{2}(n+2, m)-d_{3}(n+3, m)+d_{3}(n, m)
$$

so
$a_{0,1}(n, m)=-a_{1,1}(n, m) d_{1}(n+1, m)-a_{2,1}(n, m) d_{2}(n+2, m)+d_{3}(n+3, m)-d_{3}(n, m)$.
So $a_{0,1}(n, m)$ has weight 3 .
Let's look at the expression

$$
\Psi=s_{n+1, m, t}-s_{n, m+1, t}-A\left(\mathscr{L}_{t}, n, m\right) s_{n, m, t}-B\left(\mathscr{L}_{t}, n, m\right) s_{n-1, m, t}
$$

Equation (3.7) allows us to express the $s_{a, b, t}$ in terms of the $a_{i}(a, b, t)$ and in terms of the $s_{c, d}$. Further,

$$
s_{n+i, m+1}=s_{n+i+1, m}-A(n+i, m) s_{n+i, m}-B(n+i, m) s_{n+i-1, m} .
$$

We can therefore express $\Psi$ as a linear combination of the $s_{n+i, m}$. But the vectors $s_{n+i, m}$ for fixed $m$ are all linearly independent. Since $\Psi=0$, all the coefficients of this expresson for $\Psi$ must be zero. In particular, since the coefficients are power series in $t$, the individual terms in this power series are zero. If $Q(t)=\Sigma_{i} Q_{i}(t) s_{n+i, m}$, let $p(Q)$ be the coefficient of $t$ in the power series expansion of the coefficient of $Q_{0}$ and let $q(Q)$ be the coefficient of $t$ in the power series expansion of $Q_{-1}$. Note that $p\left(s_{n+1, m, t}\right)=0$, since $s_{n+1, m, t}$ does not involve $s_{n, m}$. Next, let's compute $p\left(s_{n, m+1, t}\right)$. We have

$$
s_{n, m+1, t}=-t s_{n+3, m+1}+a_{2}(n, m+1, t) s_{n+2, m+1}+a_{1}(n, m+1, t) s_{n+1, m+1}+a_{0}(n, m+1, t) s_{n, m+1} .
$$

The coefficient of $s_{n, m}$ in $s_{n, m+1}$ is $-A(\mathscr{L}, n, m)$ and that the coefficient of $s_{n, m}$ in $s_{n+1, m+1}$ is $-B(\mathscr{L}, n+1, m)$ so

$$
-p\left(s_{n, m+1, t}\right)=a_{1,1}(n, m+1) B(\mathscr{L}, n+1, m)+a_{0,1}(n, m+1) A(\mathscr{L}, n, m)
$$

Let $\dot{A}(\mathscr{L}, n, m)$ denote

$$
\frac{d}{d t} A\left(\mathscr{L}_{t}, n, m\right)_{t=0}
$$

Next we compute

$$
p\left(A\left(\mathscr{L}_{t}, n, m\right) s_{n, m, t}\right)=\dot{A}(\mathscr{L}, n, m)+a_{0,1}(n, m) A(\mathscr{L}, n, m)
$$

We further compute

$$
p\left(B\left(\mathscr{L}_{t}, n, m\right) s_{n-1, m, t}\right)=B(\mathscr{L}, n, m) a_{1,1}(n-1, m)
$$

So we obtain from $p(\Psi)=0$,

$$
\begin{align*}
\dot{A}(\mathscr{L}, n, m)=a_{1,1}(n, m+1) B(\mathscr{L}, n+1, m) & +a_{0,1}(n, m+1) A(\mathscr{L}, n, m) \\
& -a_{0,1}(n, m) A(\mathscr{L}, n, m)  \tag{3.15}\\
& -B(\mathscr{L}, n, m) a_{1,1}(n-1, m)
\end{align*}
$$

thus $\dot{A}(n, m)$ has weight four. Similarly,

$$
q\left(s_{n, m+1, t}\right)=-B(\mathscr{L}, n, m) a_{0,1}(n, m+1)
$$

and

$$
q\left(B\left(\mathscr{L}_{t}, n, m\right) s_{n-1, m, t}\right)=\dot{B}(n, m)+a_{0,1}(n-1, m) B(\mathscr{L}, n, m)
$$

so

$$
\dot{B}(n, m)=-a_{0,1}(n-1, m) B(\mathscr{L}, n, m)+B(\mathscr{L}, n, m) a_{0,1}(n, m+1)
$$

Taking into account the formulas (3.12), (3.13) and (3.14) for $a_{i, j}$, we have formulas for $\dot{A}(\mathscr{L}, n, m)$ and $\dot{B}(\mathscr{L}, n, m)$. So we have formulas for

$$
\nabla A_{n, m}^{\prime}(f) \cdot v_{3}=-2 \dot{A}(n, m)
$$

and

$$
\nabla B_{n, m}^{\prime}(f) \cdot v_{3}=-2 \dot{B}(n, m)
$$

These formulas are the $f_{N, M}$ and $g_{N, M}$ referred to in the Introduction. Note that $\dot{B}$ has weight five.

We can express the functions $A$ and $B$ in terms of $\vartheta$ functions. For simplicity, let us assume that all the $R_{i}$ are the same point $R$ and all the $S_{i}$ are the same point $S$, and $Q, S$ and $R$ are all close to $P$. The formula

$$
\phi(z)(\omega)=\int_{P}^{z} \omega
$$

always gives well defined element of $J_{g}$. If $\{\vartheta=0\}+K$ does not contain the image of $\phi$, then we get a well defined divisor $\mathscr{D}_{K}$ on $X$ by pulling back this divisor by $\phi$ locally. This is well defined, since a choice of a different path from $P$ to $z$ would yield the same divisor. There is a constant $K_{\mathscr{L}} \in H^{0}(X, \Omega)^{*}$ so that

$$
\mathscr{D}_{K_{\mathscr{P}}}
$$

is the divisor of a non-zero section of $\mathscr{L}$. Fix a divisor $D_{0}$ of degree $g-1$. For $x$ near $P$ let

$$
C_{x}=K_{O\left(x+D_{0}\right)}
$$

Notice that

$$
(-n-m) C_{P}+(n-m) C_{Q}+m\left(C_{R}+C_{S}\right)+K_{\mathscr{L}}=K_{\mathscr{L}_{n, m}} .
$$

Consider the following meromorphic function on $H^{0}(X, \Omega)^{*}$ :

$$
g_{n, m}(Z)=\vartheta^{-n-m}\left(Z+C_{P}\right) \vartheta^{n-m}\left(Z+C_{Q}\right) \vartheta^{m}\left(Z+C_{R}\right) \vartheta^{m}\left(Z+C_{S}\right) \vartheta\left(Z+K_{\mathscr{L}_{n, m}}\right)\left(\vartheta\left(Z+K_{\mathscr{L}}\right)\right)^{-1}
$$

This function is periodic on $H^{0}(X, \Omega)^{*}$ and so $h_{n, m}=g_{n, m} \circ \phi$ is a well defined rational function on $X$. The divisor of $h_{n, m}$ is just the same as the divisor of $f_{n, m}$. So to compute $d_{1}(n, m)$ all we have to do is to compute the logarithmic derivative of $h_{n, m}$ at $P$ with respect to $z$. We have that the derivative of $f \circ \phi$ is just $\nabla f \cdot v_{1}$. So we obtain

$$
d_{1}(n, m)=C^{\prime}+C_{1} n+C_{2} m+v_{1} \nabla \log \vartheta\left(K_{\mathscr{L}_{n, m}}\right)
$$

for suitable constants $C_{1}$ and $C_{2}$ independent of $\mathscr{L}$ and a constant $C^{\prime}$ dependent on $\mathscr{L}$. Thus we have

$$
A(n, m)=C+v_{1} \nabla \log \vartheta\left(K_{\mathscr{L}_{n+1, m}}\right)-v_{1} \nabla \log \vartheta\left(K_{\mathscr{L}_{n, m+1}}\right)
$$

for a suitable $C$ independent of $\mathscr{L}$. There is a similar formula for $B$.

Given $A$ and $B$ in $V_{N, M}$, here is a conjectural construction of a curve $X$, points

$$
P, Q, R_{1}, \ldots, R_{M}, S_{1}, \ldots, S_{M} \in X
$$

and a line bundle on $X$ which will give back $A$ and $B$ when we apply the construction of this section, at least for generic $A$ and $B$. We can consider the following difference operator

$$
L(\psi)=\psi(n+1, m)-\psi(n, m+1)-A(n, m) \psi(n, m)-B(n, m) \psi(n-1, m)
$$

on the space of all complex functions on $\mathbf{Z}^{2}$. If $\alpha, \beta \in \mathbf{C}^{*}$, let $\mathcal{M}_{\alpha, \beta}$ be the set of all $\psi$ so that $L(\psi)=0$ and $\psi(n+N, m)=\alpha \psi(n, m)$ and $\psi(n, m+M)=\beta \psi(n, m)$. Let $\mathscr{B}$ be the set of ( $\alpha, \beta$ ) so that the dimension of $\mathcal{M}_{\alpha, \beta}$ is positive. If there is a curve $X$ as in this section having associated $A$ and $B$, let $X^{\prime}$ be $X-\left\{P, Q, R_{1}, \ldots, R_{M}, S_{1}, \ldots, S_{M}\right\}$. Then there is a natural map $\pi: X^{\prime} \rightarrow \mathscr{B}$ by sending $x \in X^{\prime}$ to $(\alpha(x), \beta(x))$. To see that the image of $\pi$ is in $\mathscr{B}$, we choose an isomorphism of the fiber $\mathscr{L}_{x}$ with C and let $\psi(n, m)=s_{n, m}(x) \in \mathscr{L}_{x}$. Further, there is even a line bundle on the subset $\mathscr{B}_{1}$ of $\mathscr{B}$ on which the dimension of $\mathcal{M}_{\alpha, \beta}$ is exactly one. This suggests that for generic $A$ and $B$, that $\mathscr{B}_{1}=\mathscr{B}$ and that $\mathscr{B}$ can be compactified to a curve $X$ by adding points $\left\{P, Q, R_{1}, \ldots, R_{M}, S_{1}, \ldots, S_{M}\right\}$. Further extending the line bundle to a line bundle of $X$ and applying the construction of this section will give back the original generic $A$ and $B$. But we have not worked out here this conjectural correspondence.

The construction of this section is very close to that of Mumford and van Moerbeke [MM], although we have not worked out the exact relation here. Start with a complex function $\phi$ on $\mathbf{Z}$. Define $\psi$ on $\mathbf{Z}^{2}$ by the following inductive procedure on $m$ :

$$
\psi(n, 0)=\phi(n)
$$

and

$$
\psi(n, m+1)=-A(n, m) \psi(n, m)-B(n, m) \psi(n-1, m)+\psi(n+1, m)
$$

Define $L_{1}(\phi)(n)=\psi(n, M)$. Then $\mathscr{B}$ is the set of $(\alpha, \beta)$ so that there is a nonzero $\phi$ with $\phi(n+N)=\alpha \phi(n)$ and $L_{1}(\phi)=\beta \phi$. Thus $\mathscr{B}$ is the spectral curve associated by Mumford and van Moerbeke to the operator $L_{1}$. Note that in Mumford and van Moerbeke's theory, the operator $L_{1}$ can be reconstructed from the curve, the line bundle, and the points added to compactify, while the construction here depends on a choice of a decomposition of the zeros of $\beta$ into $M$ divisors of degree 2 . Note that the case $M=1$ is the classical Toda lattice case.

## 4. Proof of Theorem 1.1

We will prove Theorem 1.1 following Lazarsfeld. We will show that there is a rank two vector bundle $E$ on $C$ with $\operatorname{det} E=\Omega, h^{0}(E)=3, h^{0}\left(E^{*}\right)=0$, and $E$ is generated by global sections. Suppose that we are given such an $E$. Then setting $H=H^{0}(E)$, there is a canonical exact sequence

$$
0 \rightarrow \Omega^{-1} \rightarrow H \otimes \mathcal{O} \rightarrow E \rightarrow 0
$$

Next, set $V=H^{*}$, and dualize this to get:

$$
\begin{equation*}
0 \rightarrow E^{*} \rightarrow V \otimes O \rightarrow \Omega \rightarrow 0 . \tag{4.1}
\end{equation*}
$$

Twisting by $\Omega$ and take cohomology.

$$
H^{0}(\Omega) \otimes V \rightarrow H^{0}\left(\Omega^{\otimes 2}\right) \rightarrow H^{1}\left(E^{*} \otimes \Omega\right) \rightarrow V \otimes H^{1}(\Omega) \rightarrow 0 .
$$

Since $h^{0}(E)=h^{1}\left(E^{*} \otimes \Omega\right)=3$, we see that the map $H^{1}\left(E^{*} \otimes \Omega\right) \rightarrow V \otimes H^{1}(\Omega)$ is an isomorphism. On the other hand, (4.1) lets us view $V$ as a subspace of $H^{0}(\Omega)$. So the theorem will follow from the existence of such an $E$.

To construct such an $E$, we fix a line bundle $A$ on $C$ so that the degree of $A$ is $g$, $h^{0}(A)=2$, and $A$ is generated by global sections. Indeed, let $A=\Omega\left(-P_{1}-\ldots-P_{g-2}\right)$, where the $P_{i}$ are chosen generically. Since the map of $C$ to projective space via the canonical map is an embedding, $A$ has the required properties. Note that there is a unique section of $\Omega \otimes A^{*}$ which vanishes at the $P_{i}$. Consider the kernel $K$ of the natural map

$$
\alpha: \operatorname{Ext}^{1}\left(A, \Omega \otimes A^{*}\right) \rightarrow \operatorname{Hom}\left(H^{0}(A), H^{1}\left(\Omega \otimes A^{*}\right)\right),
$$

which takes an extension

$$
0 \rightarrow \Omega \otimes A^{*} \rightarrow E \rightarrow A \rightarrow 0,
$$

to the connecting homomorphism it determines. $\operatorname{Ext}^{1}\left(A, \Omega \otimes A^{*}\right)$ is dual to $H^{0}\left(A^{\otimes 2}\right)$ and

$$
\operatorname{Hom}\left(H^{0}(A), H^{1}\left(\Omega \otimes A^{*}\right)\right)
$$

is dual to $H^{0}(A)^{*} \otimes H^{0}(A)^{*}$, and $\alpha$ is dual to the multiplication map

$$
H^{0}(A) \otimes H^{0}(A) \rightarrow H^{0}\left(A^{\otimes 2}\right)
$$

The base point free pencil trick shows that the cokernel of multiplication has dimension
$h^{0}\left(A^{\otimes 2}\right)-3=g-2$. So $K$ is a vector space of dimension $g-2$. On the other hand, any non-trivial extension in $K$ gives a vector bundle $E$ satisfying the desired properties, except that $E$ might not be generated by its global sections.

We will show that if we choose a generic element of $K$, then the resulting $E$ will be generated by global sections. Suppose $E$ is not generated by global sections. The three sections of $E$ do generate a subsheaf $E^{\prime}$ of $E$, which sits in a diagram of extension as follows:


Since $E$ is generated by global sections away from $P_{1} \ldots P_{g-2}$, we have that $D \subset \cup P_{i}$. So if $E$ comes from an element $e \in K$ that fails to be generated by global sections, then there is a point $P$ among the $P_{i}$ and an extension

$$
0 \rightarrow \Omega \otimes A^{*}(-P) \rightarrow E^{\prime \prime} \rightarrow A \rightarrow 0
$$

so that $e$ is induced from this extension. Note that then such extensions are necessarily surjective on global sections. But such extensions are classified by elements in

$$
\operatorname{ker}\left(H^{0}\left(A^{\otimes 2}(P)\right)\right)^{*} \rightarrow H^{0}(A)^{*} \otimes H^{0}(A(P))^{*}
$$

Noting that $h^{0}(A(P))=3$ (since $P \in P_{1}+\ldots+P_{g-2}$ ), the base point free pencil trick shows that the cokernel of

$$
H^{0}(A) \otimes H^{0}(A(P)) \rightarrow H^{0}\left(A^{\otimes 2}(P)\right)
$$

has dimension $g-3$. Hence the extensions in $K$ which fail to be generated by global sections have codimension at least 1 and so an extension with all the required properties exists.

## 5. Kodaira-Spencer theory

Suppose that $U$ is a simply connected neighborhood of 0 in $\mathbf{C}^{n}$. We will denote the coordinates on $\mathbf{C}^{n}$ by $z_{1}, \ldots, z_{n}$. Let $\pi: \mathscr{X} \rightarrow U$ be a proper smooth map from an $n+1$ dimensional $\mathscr{X}$ so the fibers $\mathscr{X}_{s}$ of $\pi$ are smooth curves of genus $g$. Let $Q$ be a section of $\pi$. Let $\gamma_{1}, \ldots, \gamma_{g}$ be elements of $H^{1}(\mathscr{X}, \mathbf{Z})$. Let $\omega_{1}, \ldots, \omega_{g}$ be global sections of $\Omega_{\mathscr{H} / U}$, the relative one forms on $\mathscr{X} \rightarrow U$. Let $\omega_{i, s}$ denote the restriction of $\omega_{i}$ to $H^{0}\left(\mathscr{X}_{s}, \Omega\right)$. We
assume that the $\omega_{i, s}$ form a basis of $H^{0}\left(\mathscr{X}_{s}, \Omega\right)$ which is dual to the restrictions of the $\gamma_{j}$ to $\boldsymbol{H}^{1}\left(\mathscr{X}_{s}, \mathrm{C}\right)$.

Let $W_{s}$ be the complex span of the $\gamma_{i}$ in $H^{1}\left(\mathscr{X}_{s}, \mathbf{C}\right)$. Let us further assume that there are functions $A_{j}$ on $U$ for $j=2, \ldots, g, B_{j}$ on $U$ for $j=3, \ldots, g$ and $C_{j}$ on $U$ for $j=4, \ldots, g$ so that if we set

$$
\begin{aligned}
& \delta_{1}(s)=\gamma_{1}+\sum_{j=2}^{g} A_{j}(s) \gamma_{j} \\
& \delta_{2}(s)=\gamma_{2}+\sum_{j=3}^{g} B_{j}(s) \gamma_{j} \\
& \delta_{3}(s)=\gamma_{3}+\sum_{j=4}^{g} C_{j}(s) \gamma_{j},
\end{aligned}
$$

then for each $s \in U, \delta_{1}(s)$ is a basis for the annihilator of $H^{0}\left(\mathscr{O}_{s}, \Omega(-Q(s))\right)$ in $W_{s}$, so that $\delta_{1}(s)$ and $\delta_{2}(s)$ are a basis for the annihilator of $H^{0}\left(\mathscr{X}_{s}, \Omega(-2 Q(s))\right)$ in $W_{s}$ and so that $\delta_{1}(s)$ and $\delta_{2}(s)$ and $\delta_{3}(s)$ are a basis for the annihilator of $H^{0}\left(\mathscr{X}_{s}, \Omega(-3 Q(s))\right)$ in $W_{s}$. Our aim is to compute the partials of the $A_{j}, B_{j}$ and $C_{j}$ in terms of Kodaira-Spencer theory. In particular, we wish to know when the map $\Phi$ from $U$ to $\mathbf{C}^{38-6}$ defined by sending $s$ to the vector

$$
\left(A_{2}(s), \ldots, A_{g}(s), B_{3}(s), \ldots, B_{g}(s), C_{4}(s), \ldots, C_{g}(s)\right)
$$

has maximal rank.
To compute these partials, we introduce the following functions $a_{i}, b_{i}$ and $c_{i}$ so that

$$
\begin{gathered}
\omega_{2, s}^{\prime}=\omega_{2, s}-a_{2}(s) \omega_{1, s} \in H^{0}\left(\mathscr{X}_{s}, \Omega(-Q(s))\right) \\
\omega_{3, s}^{\prime}=\omega_{3, s}-a_{3}(s) \omega_{1, s}-b_{3}(s) \omega_{2, s} \in H^{0}\left(\mathscr{X}_{s}, \Omega(-2 Q(s))\right)
\end{gathered}
$$

and

$$
\omega_{j, s}^{\prime}=\omega_{j, s}-a_{j}(s) \omega_{1, s}-b_{j}(s) \omega_{2, s}-c_{j}(s) \omega_{3, s} \in H^{0}\left(\mathscr{X}_{s}, \Omega(-3 Q(s))\right)
$$

for $j>3$. By evaluating the identities

$$
\begin{aligned}
& \left\langle\delta_{1}(s), \omega_{j, s}^{\prime}\right\rangle=0, \\
& \left\langle\delta_{2}(s), \omega_{j, s}^{\prime}\right\rangle=0
\end{aligned}
$$

for $j>2$ and

$$
\left\langle\delta_{3}(s), \omega_{j, s}^{\prime}\right\rangle=0
$$

for $j>3$, we see that the $a_{i}, b_{i}$ and $c_{i}$ can be expressed in term of the $A_{i}, B_{i}$ and $C_{i}$. For instance, $c_{i}=C_{i}$. On the other hand, the functions $a_{i}^{\prime}=\left\langle\delta_{1}(0), \omega_{j}^{\prime}\right\rangle, b_{i}^{\prime}=\left\langle\delta_{2}(0), \omega_{j}^{\prime}\right\rangle$, $c_{i}^{\prime}=\left\langle\delta_{3}(0), \omega_{j}^{\prime}\right\rangle$ can all be expressed in terms of the functions $a_{i}, b_{i}$ and $c_{i}$ by expanding their definitions. So it suffices to determine when the map $\phi$ from $U$ to $\mathbf{C}^{3 g-6}$ defined by sending $u$ to the vector

$$
\left(a_{2}^{\prime}(s), \ldots, a_{g}^{\prime}(s), b_{3}^{\prime}(s), \ldots, b_{8}^{\prime}(s), c_{4}^{\prime}(s), \ldots, c_{g}^{\prime}(s)\right)
$$

has maximal rank.
Let us assume for the moment that $n=\operatorname{dim} U$ is one and that $z=z_{1}$ is the coordinate. By shrinking $U$, we can find a cover $\left\{U_{\alpha}\right\}$ of $\mathscr{X}_{0}$ and holomorphic embeddings $h_{a}$ : $U_{a} \times U \rightarrow \mathscr{X}$ so that $\pi \circ h_{a}$ is just the projection of $U_{a} \times U$ to $U$ and so that if $Q(0) \in U_{a}$, then $Q(s)=h_{\alpha}(Q(0), s)$. Let $\Theta$ be the sheaf of holomorphic derivations on $\mathscr{O}_{0}$, i.e. the dual of the sheaf of holomorphic one forms. If $f$ is a function on an open subset $V$ of $U_{a}$, define $T_{a}(f)$ to be the function on $h_{a}(V \times U)$ defined by $T_{a}(f)\left(h_{a}(v, s)\right)=f(v)$. There is a cocycle $D_{\alpha, \beta} \in H^{0}\left(U_{a} \cap U_{\beta}, \Theta\right)$ so that if $f$ is a function on a non-empty open $V$ of $\mathscr{R}_{0}$ then

$$
\lim _{z \rightarrow 0}\left(\frac{T_{\beta}(f)-T_{a}(f)}{z}\right)=D_{\alpha, \beta}(f)
$$

where both sides are defined. Thus we get a Kodaira-Spencer class

$$
K S \in H^{1}\left(\mathscr{X}_{0}, \Theta(-Q(0))\right)
$$

which is easily seen to be independent of the choices of covers we have made.
Let $\omega$ be a meromorphic form on $\mathscr{P}_{0}$ which is of the second kind, i.e. it is locally exact. By choosing the $U_{\alpha}$ simply connected, we can write $\omega=d f_{a}$, where $f_{\alpha}$ is meromorphic on $U_{\alpha}$. Then $c_{\alpha, \beta}=\left\{f_{\alpha}-f_{\beta}\right\}$ defines a cocycle with values in $\mathbf{C}$ and gives a well defined element $L_{\omega} \in H^{1}\left(\mathscr{O}_{0}, \mathbf{C}\right)=H^{1}\left(\mathscr{X}_{s}, \mathbf{C}\right)$. Let $\omega^{\prime}$ be a section of $\Omega_{\mathscr{Y} \mid D}((-k+1) Q(U))$. So for each $s \in U, \omega_{s}^{\prime}$ is a holomorphic one form on $\mathscr{X}_{s}$ vanishing $k-1$ times at $Q(s)$. Suppose that $\omega \in H^{0}\left(\mathscr{O}_{0}, \Omega(k Q(0))\right)$. We have a function defined on $U$ by the following process: The product of $L_{\omega}$ and $\omega_{s}$ is a well defined element denoted $\left\langle\omega, \omega^{\prime}\right\rangle_{s}$ of $H^{1}\left(\mathscr{X}_{s}, \Omega\right)=\mathbf{C}$. Note that $\left\langle\boldsymbol{w}, w^{\prime}\right\rangle_{0}=0$, since $L_{\omega}$ maps to zero in $H^{1}(\mathscr{O}(k-1)(Q(0)))\left(L_{\omega}\right.$ is the coboundary of the $\left.\left\{f_{a}\right\}\right)$. We have the following formula:

$$
\left(\frac{d}{d z}\left\langle\omega, \omega^{\prime}\right\rangle_{z}\right)_{z=0}=K S(\omega)\left(\omega^{\prime}\right)
$$

We have used the multiplication maps

$$
H^{1}(\Theta(-Q(0))) \times H^{0}(\Omega(k(Q(0)))) \rightarrow H^{1}(\mathscr{O}(k-1)(Q(0)))
$$

to evaluate $K S(\omega)$ and

$$
H^{1}(\mathscr{O}((k-1)(Q(0)))) \times H^{0}(\Omega((1-k)(Q(0)))) \rightarrow H^{1}(\Omega)=\mathbf{C}
$$

to evaluate $K S(\omega)\left(\omega^{\prime}\right)$.
This formula is easily proved. For let $\lambda_{\alpha, \beta}=T_{\alpha}\left(f_{\alpha}\right)-T_{\beta}\left(f_{\beta}\right)-c_{\alpha, \beta}$ define a cohomology class $\Lambda$ in $H^{1}(\mathcal{O}((k-1) Q(U)))$. Note that the $\lambda_{\alpha, \beta}$ all vanish when $z=0$. So

$$
\lim _{z \rightarrow 0} \frac{\lambda_{\alpha, b}}{z}=D_{\alpha, \beta}\left(f_{\alpha}\right),
$$

since

$$
T_{\beta}\left(f_{\alpha}\right)=T_{\beta}\left(f_{\beta}\right)-c_{\alpha, \beta}
$$

Consequently,

$$
\lim _{z \rightarrow 0}\left\langle\frac{\Lambda}{z}, \omega^{\prime}\right\rangle_{z}=K S(\omega)\left(\omega^{\prime}\right)
$$

On the other hand, for $z \neq 0$, then $\Lambda$ is the image of $L_{\omega}$ in $H^{1}\left(\mathscr{X}_{s}, \mathscr{O}((k-1) Q(s))\right)$. Thus $\left\langle\Lambda, \omega^{\prime}\right\rangle_{z}=\left\langle\omega, \omega^{\prime}\right\rangle_{z}$, and so our formula is established.

Let us return to the case of $U$ of dimension $n$. We can apply the analysis of the preceding paragraph to the curve $C_{i}$ defined by be setting all the $z_{j}=0$ for $j \neq i$. This will give an element $K S\left(\partial / \partial z_{i}\right) \in H^{1}\left(\mathscr{X}_{0}, \Theta(-Q(0))\right)$. We have

$$
\left(\frac{\partial a_{j}^{\prime}}{\partial z_{i}}\right)_{z=0}=K S\left(\frac{\partial}{\partial z_{i}}\right)\left(\delta_{1}(0)\right)\left(\omega_{j}^{\prime}\right)
$$

for $j>1$

$$
\left(\frac{\partial b_{j}^{\prime}}{\partial z_{i}}\right)_{z=0}=K S\left(\frac{\partial}{\partial z_{i}}\right)\left(\delta_{2}(0)\right)\left(\omega_{j}^{\prime}\right)
$$

for $j>2$

$$
\left(\frac{\partial c_{j}^{\prime}}{\partial z_{i}}\right)_{z=0}=K S\left(\frac{\partial}{\partial z_{i}}\right)\left(\delta_{3}(0)\right)\left(\omega_{j}^{\prime}\right)
$$

for $j>3$.

THEOREM 5.1. Suppose that the $K S\left(\partial / \partial z_{i}\right)$ actually span $H^{1}\left(\mathscr{R}_{0}, \Theta(-Q(0))\right)$ and that the $\delta_{i}$ are acceptable. Then the map $\Phi$ defined above has maximal rank.

Proof. By duality, the map from $H^{1}\left(\mathscr{X}_{0}, \Theta(-Q(0))\right)$ to

$$
H^{1}\left(\mathscr{X}_{0}, \mathscr{O}(Q(0))\right) \oplus H^{1}\left(\mathscr{X}_{0}, \mathscr{O}(2 Q(0))\right) \oplus H^{1}\left(\mathscr{X}_{0}, \mathscr{O}(3 Q(0))\right)
$$

defined by $\theta \rightarrow \Sigma \delta_{i} \theta$ is surjective.

## 6. A monodromy argument

There are smooth analytic manifolds $U_{1}$ and $\mathscr{X}$ and a proper smooth morphism $\pi: \mathscr{X} \rightarrow U_{1}$ and a section $Q$ of $\pi$ so that the dimension of $U_{1}$ is $3 g-2$ and the dimension of $\mathscr{X}$ is $3 g-1$ and so that the induced map $G: U_{1} \rightarrow M_{g, 1}$ is surjective, where $G$ is defined by $G(u)=\left(\pi^{-1}(u), Q(u)\right)$. Further we may choose $U_{1}$ and $\pi$ so that they are defined over $\mathbf{R}$ and so that for any point $u \in U_{1}$, there is a coordinate system $z_{i}$ so that the KodairaSpencer classes $K S\left(\partial / \partial z_{i}\right)$ generate $H^{1}\left(\mathscr{X}_{u}, \Theta(-Q(u))\right)$. In such a situation there is a monodromy map $T: \pi_{1}\left(U_{1}, s\right) \rightarrow \operatorname{Sp}\left(H^{1}\left(\mathscr{X}_{s}, Z\right)\right)$, where $\operatorname{Sp}\left(H^{1}\left(\mathscr{X}_{s}, \mathbf{Z}\right)\right)$ is the group of symplectic automorphisms of $H^{1}\left(\mathscr{X}_{s}, \mathbf{Z}\right)$. We can assume that the image of $T$ is a subgroup of finite index in $\operatorname{Sp}\left(H^{1}\left(\mathscr{X}_{s}, Z\right)\right)$. We will establish the following later in this section:

Proposition 6.1. There is a dense set of points $u \in U_{1}(\mathbf{R})$ so that if $v_{1}, v_{2}$ and $v_{3}$ are adapted and in $\Lambda^{+}\left(\mathscr{X}_{u}\right)(\mathbf{R})$, then the $v_{i}$ are acceptable.

Let $U_{2}$ be the set of all $(u, \Lambda)$ so that $u \in U_{1}$ and $\Lambda$ is a complex subspace of $H^{1}\left(\mathscr{X}_{u}, \mathbf{C}\right)$ of dimension $g$. Note that $U_{2}$ inherits the natural structure of a complex manifold of dimension $3 g-2$ so that the projection map $P_{1}$ from $U_{2}$ to $U_{1}$ is a covering map. Indeed, let $W$ be simply connected neighborhood of $u \in U_{1}$. Then $H^{1}\left(\mathscr{X}_{w}, \mathrm{C}\right)$ for $w \in W$ form a local system of vector spaces on $W$, which is trivial, since $W$ is simply connected. Thus we get an identification $\phi_{w}$ of $H^{1}\left(\mathscr{X}_{u}, \mathbf{C}\right)$ with $H^{1}\left(\mathscr{X}_{w}, \mathbf{C}\right)$. The map $\psi: w \mapsto\left(w, \phi_{w}(\Lambda)\right)$ is a section of $P_{1}$ and defines a chart for the holomorphic structure of $U_{2}$, by definition. Note that $U_{2}$ also inherits a real structure. Indeed if $(u, \Lambda) \in U_{2}$, then the antiholomorphic involution on $\mathscr{X}$ restricts to an antiholomorphic map from $\mathscr{X}_{u}$ to $\mathscr{X}_{u^{\prime}}$. Let $\Lambda^{\prime}$ be the image of $\Lambda$ under this map. We define $\iota$ on $U_{2}$ by $(u, \Lambda)^{\iota}=\left(u^{\iota}, \Lambda^{\prime}\right)$. In particular, if $(u, \Lambda) \in U_{2}$ is a real point, then $\Lambda$ is invariant under $\iota$ and so $\Lambda=\Lambda^{+}\left(\mathscr{X}_{u}\right) \cdot \mathbf{C}$.

If $(u, \Lambda) \in U_{2}$, there is a natural map from $\Lambda$ to $H^{0}\left(\mathscr{X}_{u}, \Omega\right)^{*}$ induced by cup product. Let $U_{3} \subset U_{2}$ be the set of all $(u, \Lambda)$ so that this map is an isomorphism. Note that the real
points of $U_{2}$ are all in $U_{3}$. Let $U_{4}$ be the set of all $(u, \Lambda) \in U_{3}$ so that there are adapted $v_{i} \in \Lambda$ which are acceptable. $U_{4}$ is an open subset of $U_{3}$ whose complement is defined by analytic equations.

We wish to show
Proposition 6.2. $\bar{U}_{4}=\bar{U}_{3}$.
Proof. Both $U_{4}$ and $U_{3}$ are defined locally by the non-vanishing of analytic equations. If the proposition were false, there is a whole component $U_{5}$ of $\bar{U}_{3}$ contained in the complement of $U_{4} . \pi_{1}\left(U_{1}, s\right)$ acts on $U_{5}$, since the projection from $U_{5}$ to $U_{1}$ is a covering map. If $\gamma \in \pi_{1}\left(U_{1}, s\right)$, the image of $(s, \Lambda)$ under $\gamma$, which is just $(s, T(\gamma)(\Lambda))$, would be in $U_{5}$. Thus if $\left(s, \Lambda^{\prime}\right)$ is any point in $U_{5}$, then $\left(s, T(\gamma)\left(\Lambda^{\prime}\right)\right) \in U_{5}$. In particular, if $s$ is a generic point of $U_{5}$, we would have that $A(\Lambda)$ would not contain an acceptable adapted set for all $A$ in some subgroup $H$ of finite index in $\operatorname{Sp}\left(H^{\text {l }}\left(\mathscr{X}_{s}, \mathbf{Z}\right)\right)$. But $H$ is Zariski dense in $\operatorname{Sp}\left(H^{1}\left(\mathscr{X}_{s}, \mathrm{C}\right)\right)$. It follows that for all $A \in \operatorname{Sp}\left(H^{1}\left(\mathscr{X}_{s}, \mathrm{C}\right)\right)$, we would have that $A(\Lambda)$ would not contain an acceptable adapted set.

Consider the transvection

$$
T_{w}(v)=v+\langle v, w\rangle w,
$$

where $w$ is a holomorphic one form on $\mathscr{X}_{s}$. If $v_{i}^{\prime}$ for $i$ from one to three form a basis of $H^{0}\left(\mathscr{X}_{s}, \Omega(4 Q(s))\right) \cap \Lambda^{\prime}$, the $T_{w}\left(v_{i}^{\prime}\right)$ form a basis of $H^{0}\left(\mathscr{X}_{s}, \Omega(4 Q(s))\right) \cap T_{w}\left(\Lambda^{\prime}\right)$. We say that $v_{i}^{\prime}$ for $i$ from one to three satisfy a nontrivial relation if there are nontrivial $\omega_{i}$ in $H^{0}\left(\mathscr{X}_{s}, \Omega(-i Q(s))\right)$ so that $v_{1}^{\prime} \omega_{1}+v_{2}^{\prime} \omega_{2}+v_{3}^{\prime} \omega_{3}=0$. Let's suppose that if $w_{i} \in H^{0}\left(\mathscr{X}_{s}, \Omega\right)$ are chosen generically, then the $w_{i}$ do not satisfy any relation and that $(s, \Lambda) \in U_{5}$. Suppose the $v_{i}$ are an adapted set in $\Lambda$. We may choose the $v_{i}$ and the $w_{i}$ so that $\left\langle v_{1}, w_{2}\right\rangle=0$, $\left\langle v_{1}, w_{3}\right\rangle=0$, and $\left\langle v_{2}, w_{3}\right\rangle=0$, but that $\left\langle v_{1}, w_{1}\right\rangle \neq 0,\left\langle v_{2}, w_{2}\right\rangle \neq 0$, and $\left\langle v_{3}, w_{3}\right\rangle \neq 0$. Let $S_{t}=t^{2}\left(T_{t^{-1} w_{3}} \circ T_{t^{-1} w_{2}} \circ T_{t^{-1} w_{1}}\right)$. Then

$$
\begin{aligned}
& S_{t}\left(v_{1}\right)=t^{2} v_{1}+\left\langle v_{1}, w_{1}\right\rangle w_{1} \\
& S_{t}\left(v_{2}\right)=t^{2} v_{2}+\left\langle v_{2}, w_{1}\right\rangle w_{1}+\left\langle v_{2}, w_{2}\right\rangle w_{2} \\
& S_{t}\left(v_{3}\right)=t^{2} v_{3}+\left\langle v_{3}, w_{1}\right\rangle w_{1}+\left\langle v_{3}, w_{2}\right\rangle w_{2}+\left\langle v_{3}, w_{3}\right\rangle w_{3} .
\end{aligned}
$$

The $S_{t}\left(v_{i}\right)$ are adapted and $\left(s, S_{t}(\Lambda)\right) \in U_{5}$, so $S_{t}\left(v_{i}\right)$ satisfy a nontrivial relation for all $t$. By taking the limit as $t \rightarrow 0$, we see that the $w_{i}$ would satisfy a non-trivial relation. Thus if $U_{5}$ is nonempty, $w_{i}$ chosen generically would satisfy a nontrivial relation.

We know that the map $\psi: H^{0}\left(\mathscr{X}_{s}, \Omega\right)^{3} \rightarrow H^{0}\left(\mathscr{X}_{s}, \Omega^{\otimes 2}\right)$ defined by $\left(\omega_{1}, \omega_{2}, \omega_{3}\right) \rightarrow$
$w_{1} \omega_{1}+w_{2} \omega_{2}+w_{3} \omega_{3}$ is surjective for generic $s$ and $w_{i}$. We may assume that $w_{1}$ does not vanish at $Q(s)$, that $w_{2}$ vanishes exactly once at $Q(s)$, and that $w_{3}$ vanishes exactly twice at $Q(s)$. Since the vectors $\left(w_{2},-w_{1}, 0\right),\left(w_{3}, 0,-w_{1}\right)$ and $\left(0, w_{3},-w_{2}\right)$ are in the kernel of $\psi$, these vectors generate the kernel of $\psi$. So if $w_{1} \omega_{1}+w_{2} \omega_{2}+w_{3} \omega_{3}=0$ is a nontrivial relation, we must have

$$
\begin{aligned}
& \omega_{1}=a w_{2}+b w_{3} \\
& \omega_{2}=-a w_{1}+c w_{3} \\
& \omega_{3}=-b w_{1}-c w_{2}
\end{aligned}
$$

But $\omega_{3}$ vanishes two times at $Q(s)$, so $b=c=0$. So $a=0$. So $U_{5}$ is empty.
Proof of Proposition 6.1. Note that Proposition 6.2 implies that $U_{4}(\mathbf{R})$ is dense in $U_{2}(\mathbf{R})$. Indeed, $U_{4}$ is dense in $U_{2}$ and the complement of $U_{4}$ is defined locally by analytic equations. If $U_{4}(\mathbf{R})$ were not dense in $U_{2}(\mathbf{R})$, then the equations defining the complement of $U_{4}(\mathbf{R})$ would vanish of an open subset of $U_{2}(\mathbf{R})$ and hence on an open subset of $U_{2}$. Further, $U_{2}(\mathbf{R})$ actually maps onto $\mathcal{M}_{g, 1}(\mathbf{R})$, and if $(u, \Lambda) \in U_{2}(\mathbf{R})$, then $\Lambda$ is the lattice fixed by the antiholomorphic involution of $\mathscr{X}_{u}$. It follows that for a dense set of points $(C, P)$ in $\mathscr{M}_{g, 1}(\mathbf{R})$ so that if the $v_{i} \in \Lambda^{+}$adapted, then the $v_{i}$ are acceptable.

Proof of Proposition 2.4. Let $u \in U_{4}(\mathbf{R})$. We can choose a basis $\gamma_{i}$ of $\Lambda^{+}\left(\mathscr{X}_{u}\right)$ so that for suitable choice of $a_{i}, b_{i}$ and $c_{i}$ we have that

$$
\begin{aligned}
& V_{1}=\gamma_{1}+\sum_{j=2}^{g} a_{j} \gamma_{j} \\
& V_{2}=\gamma_{2}+\sum_{j=3}^{g} b_{j} \gamma_{j} \\
& V_{3}=\gamma_{3}+\sum_{j=4}^{g} c_{j} \gamma_{j}
\end{aligned}
$$

are an adapted and therefore acceptable set $v_{i}$. We can therefore find $A_{j}(s), B_{j}(s), C_{j}(s)$ locally as in section 3 so that $V_{i}=\delta_{i}(u)$, and so that the $\delta_{i}(s)$ satisfy the condition in the first paragraph of section 5 . The map $\Phi$ restricts locally to a map of maximal rank from $U_{4}(\mathbf{R})$ to $\mathbf{R}^{3 g-6}$. In particular, we can find points $s$ near to $u$ so $A_{j}(s)$ and $B_{j}(s)$ are rational. Thus we see that ( $\mathscr{X}_{s}, Q(s)$ ) is good.

## 7. Proof of Theorem 2.5

We use the notation of $\S 4$. Assume that $\gamma_{1}, \gamma_{2}$ and $\gamma_{3}$ are adapted at 0 , i.e. that $\bar{\gamma}_{i} \in H^{0}\left(\mathscr{R}_{0}, \Omega(i+1) Q(0)\right)$ and that they form a basis of $H^{0}\left(\mathscr{X}_{0}, \Omega(4(Q(0)))\right)$. It follows that $\omega_{2,0} \in H^{0}(\Omega(-Q(0)))$ and that $\omega_{3,0} \in H^{0}(\Omega(-2 Q(0)))$. Assume the rest of the $\gamma_{i}$ defined over $\mathbf{R}$ and that $\left(\mathscr{X}_{0}, Q(0)\right)$ is good. Then we can assume that $\gamma_{1}$ and $\gamma_{2}$ are both in $H^{1}\left(\mathscr{O}_{0}, \mathbf{Z}\right)$. By replacing $U$ by a smaller neighborhood of 0 , we can find a neighborhood of $Q(u)$ and a function $z$ defined over $\mathbf{R}$ so that $z=0$ is the defining equation for $Q(U)$ and so that $d z=\omega_{1}$ as relative forms. If $s \in U$ and $R \in \mathscr{X}$ so that $z(R)=\zeta$ and $\omega_{s}$ is a holomorphic one form on $\mathscr{X}_{s}$, we denote

$$
\int_{Q(s)}^{R} \omega_{s}
$$

by

$$
\int_{0}^{\zeta} \omega_{s}
$$

where the integral is to be taken on a path connecting $Q(s)$ and $R$ lying close to $Q(s)$.
If $e \in \mathbf{C}$, there are $e^{+}(z)$ and $e^{-}(z)$ so that $e^{+}(z)+e^{-}(z)=z$ and $z^{2}-\left(e^{+}(z)\right)^{2}-\left(e^{-}(z)\right)^{2}=$ $2 z^{2} e$ for $z$ small. Consider the following functions:

$$
A_{j}^{\prime}(s, z, e)=\frac{\int_{0}^{z} \omega_{j}}{z}
$$

where $j=2, \ldots, g$,

$$
B_{j}^{\prime}(s, z, e)=\left(1+\frac{1}{z^{2}}\right)\left(\int_{0}^{z} \omega_{j}-\int_{0}^{e^{+}(z)} \omega_{j}-\int_{0}^{e^{-}(z)} \omega_{j}\right)
$$

and

$$
C_{j}^{\prime}=\left(\frac{d^{2}}{d z^{2}}\left(\frac{\omega_{j}}{d z}\right) \frac{d}{d z}\left(\frac{\omega_{2}}{d z}\right)-\frac{d^{2}}{d z^{2}}\left(\frac{\omega_{2}}{d z}\right) \frac{d}{d z}\left(\frac{\omega_{j}}{d z}\right)\right)_{z=0}
$$

These functions are defined on $U \times(D-\{0\}) \times D$, where $D$ is some neighborhood of $0 \in \mathbf{C}$ so that all the integrals above are defined. The significance of these functions is the following: Suppose we have a point $(s, 1 / N, e)$ with $N \in Z$ with $N \gg 0$ and that $A_{j}^{\prime}(s, 1 / N, e)=0$ for $j>1, B_{j}^{\prime}(s, 1 / N, e)=0$ for $j>2, B_{2}^{\prime}(s, 1 / N, e)=1$ and $C_{j}^{\prime}(s, 1 / N, e)=0$ for $j>3$. Then there are points $P, Q, S$ and $R$ all in $\mathscr{X}_{s}$ so that $z(P)=0, z(Q)=1 / N$,
$z(S)=e^{+}(1 / N)$, and $z(R)=e^{-}(1 / N)$. Further, the divisor $P-Q$ is a point of order $N$ in the Jacobian of $\mathscr{X}_{s}$, and the divisor $P+Q-R-S$ is a point of order $N^{2}+1$, since $B_{1}^{\prime}(s, z, e)=1$. Finally, the linear functional $L$ on $H^{0}(\Omega)$ :

$$
\omega \rightarrow\left(\frac{d^{2}}{d z^{2}}\left(\frac{\omega}{d z}\right) \frac{d}{d z}\left(\frac{\omega_{2}}{d z}\right)-\frac{d^{2}}{d z^{2}}\left(\frac{\omega_{2}}{d z}\right) \frac{d}{d z}\left(\frac{\omega}{d z}\right)\right)_{z=0}
$$

vanishes on the span of $\omega_{1}, \omega_{2}, \omega_{4} \ldots$ and so $L$ must to a multiple of

$$
\omega \rightarrow\left\langle\omega, \gamma_{3}\right\rangle .
$$

Note that $L$ is in the osculating three space of the curve $\mathscr{X}_{s}$ at $Q(s)$.
We claim that the $A_{j}^{\prime}$ and $B_{j}^{\prime}$ can be extended as holomorphic functions to $U \times D \times D$, and that $A_{j}^{\prime}(s, 0, e)$ and $B_{j}^{\prime}(s, 0, e)$ can be computed in terms of $e$ and the $a_{j}$ and $b_{j}$ of $\S 3$. We can write near $Q(s)$

$$
\omega_{i, s}=\alpha_{i}(s) d z+\beta_{i}(s) z d z+\frac{1}{2} \varepsilon_{i}(s) z^{2} d z+D_{i} z^{3} d z
$$

where the $\alpha_{i}$ and $\beta_{i}$ are functions on $U$ and the $D_{i}$ are functions on $\mathscr{X}$ defined locally. Note that

$$
a_{1}=1
$$

Similarly, we see that

$$
\beta_{1}=\varepsilon_{1}=0
$$

and that

$$
\alpha_{2}(0)=\alpha_{3}(0)=\beta_{3}(0)=0 .
$$

Note that $\beta_{2}$ is nonzero near $s=0$, since the form $\omega_{2,0}$ is in $H^{0}\left(\mathscr{O}_{0}, \Omega(-Q(0))\right)$, but not in $H^{0}\left(\mathscr{X}_{0}, \Omega(-2 Q(0))\right)$. Similarly, $\varepsilon_{3}(0) \neq 0$. Since $\omega_{2, s}-a_{2}(s) \omega_{1, s}$ vanishes at $Q(s)$, we see that $\alpha_{2}(s)=a_{2}(s)$. Similarly, $\omega_{3, s}-a_{3}(s) \omega_{1, s}-b_{3}(s) \omega_{2, s}$ vanishes twice at $Q(s)$, so

$$
\alpha_{3}(s)-a_{3}(s)-b_{3}(s) \alpha_{2}(s)=0
$$

and

$$
\beta_{3}(s)-b_{3}(s) \beta_{2}(s)=0
$$

For $i>3$ we have

$$
\alpha_{i}(s)-a_{i}(s)-b_{i}(s) \alpha_{2}(s)-c_{i}(s) \alpha_{3}(s)=0
$$

and

$$
\beta_{i}(s)-b_{i}(s) \beta_{2}(s)-c_{i}(s) \beta_{3}(s)=0,
$$

and

$$
\varepsilon_{i}(s)-b_{i}(s) \varepsilon_{2}(s)-c_{i}(s) \varepsilon_{3}(s)=0 .
$$

Further, we have

$$
\lim _{z \rightarrow 0} A_{j}=\alpha_{j}
$$

and

$$
\lim _{z \rightarrow 0} B_{j}=e \beta_{j}
$$

These can be seen by the formulas:

$$
\int_{0}^{\zeta} \omega_{i, s}=\zeta \alpha_{i, s}+\frac{\zeta^{2}}{2} \beta_{i, s}+\ldots
$$

Note that

$$
C_{j}^{\prime}=\varepsilon_{j} \beta_{2}-\varepsilon_{2} \beta_{j}
$$

By shrinking $U$, we may assume that $\beta_{2}$ and $\varepsilon_{3}$ never vanish on $U$. Examining these equations, we see that the subvariety of $U$ defined by the vanishing of $\alpha_{i}$ for $i>1$, the $\beta_{i}$ for $i>2$ and the $C_{i}^{\prime}>3$ is contained in the subvariety of $U$ defined by the vanishing of $a_{i}$ for $i>1$, the $b_{i}$ for $i>2$ and the $c_{i}$ for $i>3$. Using the fact that the $a_{j}$ for $j>1$, the $b_{j}$ for $j>2$ and the $c_{j}$ for $j>3$ all have independent gradients, we see that the $\alpha_{j}$ for $j>1$, the $\beta_{j}$ for $j>2$ and the $C_{j}^{\prime}$ for $j>3$ all have independent gradients. Thus, $A_{j}^{\prime}$ for $j>1$, the $B_{j}^{\prime}$ for $j>2$ and the $C_{j}^{\prime}$ for $j>3$ have independent gradients when restricted to the set $z=0$ near $s=0$. So if

$$
R=\left(0,0, \frac{1}{\beta_{2}(0)}\right)
$$

then the equations $A_{j}^{\prime}=0$ for $j>1, B_{j}^{\prime}=0$ for $j>2, B_{2}^{\prime}=1$ and $C_{j}^{\prime}=0$ for $j>0$ defines a
smooth manifold $W \subset U \times D \times C$ in a neighborhood of $R$ and that the function $z$ defines a smooth map from $W$ to $D$ near $R$. We have established:

Proposition 7.1. If $(C, P)$ is a good pair and $v_{i}$ for $i$ from 1 to 3 is an adapted set with $v_{1}$ and $v_{2}$ in $\Lambda^{+}(C)$ and $v_{3} \in \Lambda^{+}(C)(\mathbf{R})$, then there is a family $\pi: \mathscr{Y}_{s} \rightarrow D$, which is proper and smooth over $D$ and sections $P_{i}: D \rightarrow \mathscr{Y}$ such that if $\mathscr{Y}_{z}$ denotes the fiber of $\mathscr{Y}$ over $z$, then $(C, P)=\left(\mathscr{Y}_{0}, P_{i}(0)\right)$. Further, if we denote by $v_{i}$ the element of $H^{1}\left(\mathscr{Y}_{z}, \mathrm{C}\right)$ obtained by transport of $v_{i}$, then for $z \neq 0, \bar{v}_{3}$ is in the osculating three space of $\mathscr{y}_{z}$, and for all $\omega \in H^{0}\left(\mathscr{Y}_{z}, \Omega\right)$, we have

$$
\int_{P_{1}(z)}^{P_{2}(z)} \omega=z\left\langle\omega, v_{1}\right\rangle
$$

and

$$
\int_{P_{1}(z)}^{P_{2}(z)} \omega-\int_{P_{1}(z)}^{P_{3}(z)} \omega-\int_{P_{1}(z)}^{P_{4}(z)} \omega=\frac{z^{2}\left\langle\omega, v_{2}\right\rangle}{1+z^{2}}
$$

Further, if $z$ is real, then $\mathscr{Y}_{z}$ is defined over $\mathbf{R}$ and the points $P_{1}(z), P_{2}(z)$ and the divisor $P_{3}(z)+P_{4}(z)$ are defined over $\mathbf{R}$.

Choose $\gamma_{1}, \ldots, \gamma_{g}$ so that $\gamma_{i}=v_{i}$ for $i$ from 1 to 3 and $\gamma_{i} \in \Lambda^{+}(\mathbf{R})$ and choose

$$
\alpha: H^{1}(\mathscr{Y}, \mathbf{Z}) \rightarrow \mathbf{C}_{\mathbf{1}}^{*}
$$

satisfying 2.1 and 2.2. Let $\vartheta_{z}$ be the theta function on $H^{0}\left(\mathscr{Y}_{z}, \Omega\right)^{*}$ attached to $\alpha$. Let

$$
f\left(z, x_{1}, \ldots, x_{g} ; K_{1}\right)=\vartheta_{z}\left(\sum x_{i} \bar{\gamma}_{i}+K_{1}\right)
$$

where $\bar{\gamma}_{i}$ is the image of $\gamma_{i}$ in $H^{0}(\Omega)^{*}$. We assume that we have chosen $K_{1}$ so that $K_{1}-K_{1}^{\prime} \in H^{1}(C, Z)$ and that $f\left(0 ; x_{1}, \ldots, x_{g} ; K_{1}\right)$ never vanishes for $\left(x_{1}, \ldots, x_{g}\right) \in \mathbf{R}^{g}$. Define

$$
\begin{aligned}
H\left(z ; x_{1}, \ldots, x_{g} ; K_{1}\right)= & \frac{\partial \log f\left(z ; x_{1}+z, x_{2}, \ldots, x_{g} ; K_{1}\right)}{z \partial x_{1}} \\
& -\frac{\partial \log f\left(z ; x_{1}, x_{2}+z^{2} /\left(1+z^{2}\right), x_{3}, \ldots, x_{g} ; K_{1}\right)}{z \partial x_{1}} \\
& -\frac{\partial^{2} \log f\left(0 ; x_{1}, x_{2}, \ldots, x_{g} ; K_{1}\right)}{\partial x_{1}^{2}}
\end{aligned}
$$

For fixed $K_{1}$, note that $H\left(z ; x_{1}, \ldots, x_{g} ; K_{1}\right)$ is periodic on $\mathbf{R}^{g}$ for each $z$ with respect
to a lattice in $\mathbf{R}^{8}$ independent of $z$. Further, $H\left(z ; x_{1}, \ldots, x_{g} ; K_{1}\right)$ can be extended to an analytic function on $D \times \mathbf{R}^{g}$, which vanishes on $\{0\} \times \mathbf{R}^{g}$. Thus given $\varepsilon$, there is a $\delta$ so that if $|z|<\delta$, then $\left|H\left(z ; x_{1}, \ldots, x_{g} ; K_{1}\right)\right|<\varepsilon$. Fix $N \in \mathbf{Z}$ so that $1 / N<\delta$. Let $\mathscr{L}$ be the line bundle on $C=\mathscr{Y}_{1 / N}$ associated to $K_{1}$. Let $P=P_{1}(1 / N), Q=P_{2}(1 / N), R=P_{3}(1 / N)$ and $S=P_{4}(1 / N)$. Choose a parameter $z$ around $P$ so that the Frenet frame associated to $z$ is the image of the $v_{i}$. Let $A(n, m, \mathscr{L})$ be the functions introduced in section 3 . The line bundle $\mathscr{L}_{n, m}$ is associated to

$$
K_{1}+\frac{n}{N} \gamma_{1}+\frac{m}{N^{2}+1} \gamma_{2}=K_{\mathscr{L}_{n, m}}
$$

and the equations say that $\mathscr{O}(Q-P)$ is associated to $\gamma_{1} / N$ and $\mathcal{O}(P+Q-R-S)$ is associated to

$$
\frac{\gamma_{2}}{N^{2}+1}
$$

Furthermore,

$$
\frac{\partial \log f\left(1 / N ; x_{1} \ldots, x_{g} ; K_{1}\right)}{\partial x_{1}}=v_{1} \nabla \log \vartheta_{1 / N}\left(\sum x_{i} \gamma_{i}+K_{1}\right) .
$$

Let

$$
h_{1}\left(x, y ; K_{1}\right)=\frac{\partial^{2} \log f\left(0 ; x, y, 0, \ldots, 0 ; K_{1}\right)}{\partial x^{2}}
$$

Then we have the following formula:

$$
H\left(\frac{1}{N} ; \frac{n}{N}, \frac{m}{N^{2}+1}, 0, \ldots ; K_{1}\right)=N A(n, m, \mathscr{L})+C-h_{1}\left(\frac{n}{N}, \frac{m}{N^{2}+1} ; K_{1}\right) .
$$

In particular,

$$
\left|N A(n, m, \mathscr{L})+C-h_{1}\left(\frac{n}{N}, \frac{m}{N^{2}+1} ; K_{1}\right)\right|<\varepsilon .
$$

Fix $K_{1}$ and apply the preceding discussion to $K_{1}+t \gamma_{3}$, noting that

$$
\left|H\left(z, x_{1}, \ldots, x_{g} ; K_{1}+t \gamma_{3}\right)\right|<\varepsilon
$$

independent of $t$. Let

$$
h(x, y, t)=h_{1}\left(x, y ; K_{1}+t \gamma_{3}\right)=\frac{\partial^{2} \log \left(0 ; x, y, t, \ldots, 0 ; K_{1}\right)}{\partial x^{2}}
$$

Let $\mathscr{L}_{t}$ be the line bundle associated to $K_{1}+t \gamma_{3}$. Then we have

$$
\left|N A\left(n, m, \mathscr{L}_{t}\right)+C-h\left(\frac{n}{N}, \frac{m}{N^{2}+1}, t\right)\right|<\varepsilon .
$$

Thus Theorem 2.5 is established.

## References

[D] Dobrovin, B. A., Theta functions and non-linear equations. Russian Math. Surveys, 36 (1981), 11-92.
[GL] Green, M. \& Lazarsfeld, R., On the projective normality of complete linear series on an algebraic curve. Invent. Math., 83 (1986), 73-90.
[MM] van Moerbeke, P. \& Mumford, D., The spectrum of difference operators and algebraic curves. Acta Math., 143 (1979), 93-154.
[M] Mumpord, D., Abelian Varieties. Tata Institute, Oxford University Press, 1970.
[Z] Zabusky, N. J., Computational synergetics and mathematical innovation. J. Comput. Phys., 43 (1981), 195-249.

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