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# MEAN FIELD EQUATIONS, HYPERELLIPTIC CURVES AND MODULAR FORMS: II 

by Chang-Shou Lin a Chin-Lung Wang

Abstract. - A pre-modular form $Z_{n}(\sigma ; \tau)$ of weight $\frac{1}{2} n(n+1)$ is introduced for each $n \in \mathbb{N}$, where $(\sigma, \tau) \in \mathbb{C} \times \mathbb{H}$, such that for $E_{\tau}=\mathbb{C} /(\mathbb{Z}+\mathbb{Z} \tau)$, every non-trivial zero of $Z_{n}(\sigma ; \tau)$, i.e. $\sigma$ is not a 2 -torsion of $E_{\tau}$, corresponds to a (scaling family of) solution to the equation

$$
\Delta u+e^{u}=\rho \delta_{0},
$$

on the flat torus $E_{\tau}$ with singular strength $\rho=8 \pi n$.
In Part I [1], a hyperelliptic curve $\bar{X}_{n}(\tau) \subset \operatorname{Sym}^{n} E_{\tau}$, the Lamé curve, associated to the MFE was constructed. Our construction of $Z_{n}(\sigma ; \tau)$ relies on a detailed study of the correspondence $\mathbb{P}^{1}(\mathbb{C}) \leftarrow \bar{X}_{n}(\tau) \rightarrow E_{\tau}$ induced from the hyperelliptic projection and the addition map.

As an application of the explicit form of the weight 10 pre-modular form $Z_{4}(\sigma ; \tau)$, a counting formula for Lamé equations of degree $n=4$ with finite monodromy is given in the appendix (by Y.-C. Chou)

Résumé (Équations de champ moyen, courbes hyperelliptiques et formes modulaires: II)
Nous introduisons une forme pré-modulaire $Z_{n}(\sigma ; \tau)$ de poids $\frac{1}{2} n(n+1)$ pour chaque $n \in \mathbb{N}$, $\operatorname{avec}(\sigma, \tau) \in \mathbb{C} \times \mathbb{H}$, de sorte que pour $E_{\tau}=\mathbb{C} /(\mathbb{Z}+\mathbb{Z} \tau)$, tout zéro non trivial de $Z_{n}(\sigma ; \tau)$, c'est-à-dire que $\sigma$ n'est pas de 2 -torsion dans $E_{\tau}$, correspond à une (famille de) solution de l'équation
(MFE) $\triangle u+e^{u}=\rho \delta_{0}$,
sur le tore plat $E_{\tau}$ avec $\rho=8 \pi n$.
Dans la partie I [1], nous avons construit une courbe hyperelliptique $\bar{X}_{n}(\tau) \subset \operatorname{Sym}^{n} E_{\tau}$, la courbe de Lamé, associée à l'équation (MFE). Notre construction de $Z_{n}(\sigma ; \tau)$ s'appuie sur une étude détaillée de la correspondance $\mathbb{P}^{1}(\mathbb{C}) \leftarrow \bar{X}_{n}(\tau) \rightarrow E_{\tau}$ induite par la projection hyperelliptique et l'application d'addition.

Dans l'appendice, Y.-C. Chou donne, comme application de l'expression explicite de la forme $Z_{4}(\sigma ; \tau)$ pré-modulaire de poids 10 , une formule de comptage pour les équations de Lamé de degré $n=4$ avec monodromie finie.

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## 0. Introduction

Let $E=E_{\tau}=\mathbb{C} / \Lambda_{\tau}$ be a flat torus, where $\tau \in \mathbb{H}=\{\tau \in \mathbb{C} \mid \operatorname{Im} \tau>0\}$ and $\Lambda=\Lambda_{\tau}=\mathbb{Z} \omega_{1}+\mathbb{Z} \omega_{2}$ with $\omega_{1}=1$ and $\omega_{2}=\tau$. Also $\omega_{3}:=\omega_{1}+\omega_{2}$.

Convention. - For $z \in \mathbb{C}$ we denote $[z]:=z(\bmod \Lambda) \in E$. For a point $[z]$ in $E$ we often write $z$ instead of $[z]$ to simplify notations when no confusion should arise. For $N \in \mathbb{N}, E[N]:=\{[z] \in E \mid N z \in \Lambda\}$ is the group of $N$-torsion points in $E$. Also $E^{\times}:=E \backslash\{[0]\}$.

We will use the Weierstrass elliptic function $\wp(z)=\wp(z ; \Lambda)$ and its associated functions $\zeta(z ; \Lambda)$ and $\sigma(z ; \Lambda)$ extensively. We often write $\tau$ instead of $\Lambda$ and even omit it in the notation when no confusion should arise. We take [14] as our general reference on elliptic functions.

In this paper, we continue our study, initiated in [9] and developed in Part I [1], on the singular Liouville (mean field) equation:

$$
\begin{equation*}
\triangle u+e^{u}=8 \pi n \delta_{0} \quad \text { on } E \tag{0.1}
\end{equation*}
$$

under the flat metric. Here, $\triangle=\partial^{2} / \partial x^{2}+\partial^{2} / \partial y^{2}$ is the Laplace operator on $E$ induced from $\mathbb{C}, n \in \mathbb{N}$, and $\delta_{0}$ is the Dirac measure at $[0] \in E$.

The solvability of equation (0.1) depends on the moduli $\tau$ in a sophisticated manner. For $n=1$, this was only recently settled in $[9,11,2]$. The aim of this paper is to develop a theory via modular forms to investigate such a dependence for all $n \in \mathbb{N}$ and to lay the foundation towards a complete resolution to equation (0.1).

We review briefly in Section 0.1 what had been done in earlier works (mainly in Part I) to reformulate the problem using Lamé equations and Green functions. More technical statements will be recalled in later sections whenever needed. In Section 0.2 we describe new results proved in this paper.

### 0.1. Reduction to a Green function equation over the Lamé curve

0.1.1. The Liousille curve. - It was shown in [1, Th. 0.3 \& Th. 0.6 ] that if there is a solution $u$ to equation (0.1) then it lies in a scaling family of solutions $u_{\lambda}$ through the Liouville formula:

$$
\begin{equation*}
u_{\lambda}(z)=\log \frac{8 e^{2 \lambda}\left|f^{\prime}(z)\right|^{2}}{\left(1+e^{2 \lambda}|f(z)|^{2}\right)^{2}}, \quad \lambda \in \mathbb{R} \tag{0.2}
\end{equation*}
$$

where $f$ is a meromorphic function on $\mathbb{C}$, known as a developing map, which can be normalized to satisfy the type II constraints:

$$
\begin{equation*}
f\left(z+\omega_{j}\right)=e^{2 i \theta_{j}} f(z), \quad \theta_{j} \in \mathbb{R}, \quad j=1,2 \tag{0.3}
\end{equation*}
$$

(i) There is a unique $\lambda$ so that $u_{\lambda}$ is even. Moreover, the normalized developing map $f$ has precisely $n$ simple zeros $\left[a_{1}\right], \ldots,\left[a_{n}\right]$ in $E^{\times}:=E \backslash\{[0]\}$ and $n$ simple poles $\left[-a_{1}\right], \ldots,\left[-a_{n}\right]$. They are characterized by
(ii) The non-degenerate constraints: $\left[a_{i}\right] \notin E[2]$ for all $i,\left[a_{i}\right] \neq \pm\left[a_{j}\right]$ for $i \neq j$.
(iii) The following $n-1$ algebraic equations:

$$
\begin{equation*}
\sum_{i=1}^{n} \wp^{\prime}\left(a_{i}\right) \wp^{r}\left(a_{i}\right)=0, \quad r=0, \ldots, n-2 . \tag{0.4}
\end{equation*}
$$

(iv) The transcendental equation on the Green function: ${ }^{(1)}$

$$
\begin{equation*}
\sum_{i=1}^{n} \nabla G\left(a_{i}\right)=0 \tag{0.5}
\end{equation*}
$$

The affine algebraic curve $X_{n} \subset \operatorname{Sym}^{n} E$ defined by equations (0.4) and the nondegenerate constraints is called the ( $n$-th) Liouville curve.
0.1.2. The Lamé curve. - The Liouville curve $X_{n}$ has the important hyperelliptic structure arising from its connection with the integral Lamé equations on $E$ :

$$
\begin{equation*}
w^{\prime \prime}=(n(n+1) \wp+B) w, \tag{0.6}
\end{equation*}
$$

where $B \in \mathbb{C}$ is usually known as the auxiliary or spectral parameter. For $a=$ $\left(a_{1}, \ldots, a_{n}\right) \in \mathbb{C}^{n}$, let $w_{a}(z)$ be the classical Hermite-Halphen ansatz:

$$
\begin{equation*}
w_{a}(z):=e^{z \sum \zeta\left(a_{i} ; \tau\right)} \prod_{i=1}^{n} \frac{\sigma\left(z-a_{i} ; \tau\right)}{\sigma(z ; \tau)} . \tag{0.7}
\end{equation*}
$$

Then the following statements were proved in [1, Th. 0.7].
(i) The point $[a]:=a(\bmod \Lambda)$ lies in $X_{n}$ if and only if $w_{a}$ and $w_{-a}$ are independent solutions to equation (0.6). In that case, the parameter $B$ equals

$$
\begin{equation*}
B_{a}:=(2 n-1) \sum_{i=1}^{n} \wp\left(a_{i}\right) \text {. } \tag{0.8}
\end{equation*}
$$

(ii) The compactified curve

$$
\bar{X}_{n} \subset \operatorname{Sym}^{n} E
$$

is a hyperelliptic curve, known as the Lamé curve, with the added points $\bar{X}_{n} \backslash X_{n}$ being the branch points of the hyperelliptic projection

$$
B: \bar{X}_{n} \longrightarrow \mathbb{P}^{1}(\mathbb{C})
$$

[^0](iii) A point $[a] \in \bar{X}_{n}$ is a branched point if and only if $[-a]=[a]$. In fact
$$
\left\{\left[-a_{i}\right]\right\} \cap\left\{\left[a_{i}\right]\right\} \neq \varnothing \Longrightarrow[-a]=[a]
$$

Also $[0] \in\left\{\left[a_{i}\right]\right\} \Rightarrow[a]=0^{n}$.
(iv) The unique point at infinity $[0]^{n} \in \bar{X}_{n}$ is a non-singular point.
(v) The finite branch points satisfy $[a] \in\left(E^{\times}\right)^{n},\left[a_{i}\right] \neq\left[a_{j}\right]$ for $i \neq j$, and $[a]=[-a]$; $w_{a}=w_{-a}$ is still a solution to equation (0.6) with $B=B_{a}$. These solutions are known as the Lamé functions.
(vi) Let $Y_{n}=B^{-1}(\mathbb{C})$ be the finite part of $\bar{X}_{n}$. Then $Y_{n}$ is parametrized by

$$
Y_{n} \cong\left\{(B, C) \mid C^{2}=\ell_{n}(B)\right\}
$$

where $\ell_{n}(B)$ is the Lamé polynomial in $B$ of degree $2 n+1$, and $\bar{X}_{n}$ coincides with the projective hyperelliptic model of $Y_{n}$. In particular, the Lamé curve $\bar{X}_{n}$ is irreducible and is smooth if and only if $\ell_{n}(B)$ has no multiple roots.

Under this description, for $[a] \in X_{n}$, the ratio

$$
f=f_{a}:=\frac{w_{a}}{w_{-a}}=e^{2 \sum_{i=1}^{n} \zeta\left(a_{i}\right) z} \prod_{i=1}^{n} \frac{\sigma\left(z-a_{i} ; \tau\right)}{\sigma\left(z+a_{i} ; \tau\right)}
$$

gives the candidate of a developing map in (0.2). The original singular Liouville equation (0.1) is then equivalent to the Green function equation (0.5) over the unramified (Liouville) loci $X_{n}$ of the map $B: \bar{X}_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})$.
0.2. Main results: a theory of pre-modular forms. - Recall the Weierstrass equation $\wp^{\prime 2}=4 \wp^{3}-g_{2} \wp-g_{3}=\prod_{i=1}^{3}\left(\wp-e_{i}\right)$, where $e_{i}(\tau)=\wp\left(\frac{1}{2} \omega_{i} ; \tau\right), i=1,2,3$. We will also use the quasi-periods $\eta_{i}(\tau):=\zeta\left(z+\omega_{i} ; \tau\right)-\zeta(z ; \tau)=2 \zeta\left(\frac{1}{2} \omega_{i} ; \tau\right), i=1,2$, extensively.
0.2.1. The Hecke function and pre-modular forms. - For $z=x+i y=r \omega_{1}+s \omega_{2}$, $r, s \in \mathbb{R}$, it was shown in $[9$, Lem. 2.3, Lem. 7.1] that

$$
\begin{equation*}
-4 \pi \frac{\partial G}{\partial z}(z ; \tau)=\zeta(z ; \tau)-r \eta_{1}(\tau)-s \eta_{2}(\tau) \tag{0.9}
\end{equation*}
$$

For $[z] \in E_{\tau}[N] \backslash\{[0]\}$, the right-hand side of equation (0.9) first appeared in [8], where Hecke showed that it is a modular form of weight one with respect to $\Gamma(N)=$ $\left\{A \in \operatorname{SL}(2, \mathbb{Z}) \mid A \equiv I_{2}(\bmod N)\right\}$. Thus we call the following function

$$
\begin{equation*}
Z(z ; \tau)=Z_{r, s}(\tau):=\zeta\left(r \omega_{1}+s \omega_{2} ; \tau\right)-r \eta_{1}(\tau)-s \eta_{2}(\tau), \tag{0.10}
\end{equation*}
$$

$(z, \tau) \in \mathbb{C} \times \mathbb{H}$, the Hecke function. Notice that it is holomorphic only in $\tau$, and for fixed $\tau$ it depends only on $[z]=z\left(\bmod \Lambda_{\tau}\right) \in E_{\tau}$. In this paper, functions of this sort are called pre-modular forms.

Definition 0.1. - An analytic function $h$ in $(z, \tau) \in \mathbb{C} \times \mathbb{H}$ is pre-modular of weight $k \in \mathbb{N}$ if it satisfies
(1) For any fixed $\tau$, the function $h(z)$ is analytic in $z$ and $\bar{z}$ and it depends only on $z\left(\bmod \Lambda_{\tau}\right) \in E_{\tau}$.
(2) For any fixed torsion type $z\left(\bmod \Lambda_{\tau}\right) \in E_{\tau}[N]$, the function $h(\tau)$ is modular of weight $k$ with respect to $\Gamma(N)$.
We write $h(z ; \tau)$ for a pre-modular form $h$.
By writing $z=r+s \tau$ with $r, s \in \mathbb{R}$, it is easy to see that condition (1) is equivalent to saying that $h(z)$ is analytic and periodic in $r, s$ with period 1. A torsion type in condition (2) is simply a choice of $r, s, \in\left(\frac{1}{N} \mathbb{Z}\right) / \mathbb{Z}$. In particular, the Hecke function $Z$ is pre-modular of weight one.

We may regard pre-modular forms as the restriction of holomorphic functions in three complex variables $(r, s, \tau) \in \mathbb{C}^{2} \times \mathbb{H}$ to the $\mathbb{R}$-linear slice $L$ defined by $r, s \in \mathbb{R}$. Although $L \cong \mathbb{C} \times \mathbb{H}$, the embedding is not $\mathbb{C}$-linear.

The notion of pre-modular forms allows us to study deformations in $z$ to relate different modular forms corresponding to different torsion points.

Recently this idea was applied in [2] to achieve a complete solution to equation (0.1) for $n=1$ and for all $\tau .{ }^{(2)}$ In that case equations (0.4) are vacuous and the problem is equivalent to solving non-trivial zeros of $Z(z ; \tau)$, i.e., $z \notin E_{\tau}[2]$. Thus, a key step towards the general cases is to generalize the pre-modular form $Z_{1}=Z$ to certain " $Z_{n}$ " for all $n \geqslant 2$.
0.2.2. The main constructions. - By the anti-symmetry of $\nabla G$, equation (0.5) holds automatically on the branch points of $Y_{n}$, hence they are referred to as trivial solutions. Nevertheless further investigations on the local structures of the branch points are indispensable. This is done in Section 1.

We proceed to construct a pre-modular form $Z_{n}(\sigma ; \tau)$, with $\sigma \in E_{\tau}$, associated to the family of Lamé curves $\bar{X}_{n}(\tau), \tau \in \mathbb{H} .{ }^{(3)}$ It should have the property that every non-trivial solution $[a]=\left\{\left[a_{1}\right], \ldots,\left[a_{n}\right]\right\} \in X_{n}(\tau)$ to equation (0.5) comes from a zero of $Z_{n}(\sigma ; \tau)$ with $\sigma=\sum_{i=1}^{n}\left[a_{i}\right] \notin E_{\tau}[2]$, and vice versa. The construction is stated in (0.15). Its justification consists of three steps corresponding to Theorems 0.2, 0.3 and 0.4 in the following.

Consider the meromorphic function

$$
\begin{equation*}
\mathbf{z}_{n}(a):=\zeta\left(\sum_{i=1}^{n} a_{i}\right)-\sum_{i=1}^{n} \zeta\left(a_{i}\right) \tag{0.11}
\end{equation*}
$$

on $E^{n}$. Write $a_{i}=r_{i} \omega_{1}+s_{i} \omega_{2}$. If $\sum_{i=1}^{n}\left[a_{i}\right] \neq 0$ then from equation (0.9)

$$
-4 \pi \sum_{i=1}^{n} \frac{\partial G}{\partial z}\left(a_{i}\right)=\sum_{i=1}^{n}\left(\zeta\left(r_{i} \omega_{1}+s_{i} \omega_{2}\right)-r_{i} \eta_{1}-s_{i} \eta_{2}\right)=Z\left(\sum_{i=1}^{n} a_{i}\right)-\mathbf{z}_{n}(a) .
$$

[^1]Hence the Green function equation (0.5) is equivalent to

$$
\begin{equation*}
\mathbf{z}_{n}(a)=Z\left(\sum_{i=1}^{n} a_{i}\right) \tag{0.12}
\end{equation*}
$$

This motivates us to study the map

$$
\begin{equation*}
\sigma_{n}: \bar{X}_{n} \longrightarrow E, \quad[a] \longmapsto \sigma_{n}([a]):=\sum_{i=1}^{n}\left[a_{i}\right] \tag{0.13}
\end{equation*}
$$

induced from the addition map $E^{n} \rightarrow E$. Since the algebraic curve $\bar{X}_{n}$ is irreducible, $\sigma_{n}$ is a finite morphism and $\operatorname{deg} \sigma_{n}$ is defined.

Theorem $0.2\left(=\right.$ Theorem 2.4). - The map $\sigma_{n}: \bar{X}_{n} \rightarrow E$ has degree $\frac{1}{2} n(n+1)$.
From Theorem 0.2, there is a polynomial

$$
W_{n}(\mathbf{z}) \in \mathbb{Q}\left[g_{2}, g_{3}, \wp(\sigma), \wp^{\prime}(\sigma)\right][\mathbf{z}]
$$

of degree $\frac{1}{2} n(n+1)$ in $\mathbf{z}$ which defines the (branched) covering map $\sigma_{n}$.
The next task is to find a natural primitive element of this covering map, namely a rational function on $\bar{X}_{n}$ which has $W_{n}$ as its minimal polynomial. This is achieved by the following fundamental theorem.

Theorem 0.3 (= Theorem 3.2). - The rational function $\mathbf{z}_{n} \in K\left(\bar{X}_{n}\right)$ is a primitive generator for the field extension $K\left(\bar{X}_{n}\right)$ over $K(E)$ which is integral over the affine curve $E^{\times}$.

This means that $W_{n}\left(\mathbf{z}_{n}\right)=0$, and conversely for generic choices of $\sigma=\sigma_{0} \in E_{\tau}$, the roots of $W_{n}(\mathbf{z})\left(\sigma_{0} ; \tau\right)=0$ are precisely those $\frac{1}{2} n(n+1)$ values $\mathbf{z}=\mathbf{z}_{n}(a)$ with $\sigma_{n}(a)=\sigma_{0}$. The proof is contained in Section 3. Here we give a brief sketch of the idea used in the proof.

A major tool used is the tensor product of two Lamé equations $w^{\prime \prime}=I_{1} w$ and $w^{\prime \prime}=I_{2} w$, where $I=n(n+1) \wp(z), I_{1}=I+B_{a}$ and $I_{2}=I+B_{b}$.

For a general point $\sigma_{0} \in E$, we need to show that the $\frac{1}{2} n(n+1)$ points on the fiber of $\bar{X}_{n} \rightarrow E$ above $\sigma_{0}$ has distinct $\mathbf{z}_{n}$ values. From (0.11), it suffices to show that for $\sigma_{n}(a)=\sigma_{n}(b)=\sigma_{0}$,

$$
\sum_{i=1}^{n} \zeta\left(a_{i}\right)=\sum_{i=1}^{n} \zeta\left(b_{i}\right) \Longrightarrow B_{a}=B_{b}
$$

Indeed, then we conclude $[a]=[b]$ if $\sigma_{0} \notin E[2]$.
If $w_{1}^{\prime \prime}=I_{1} w_{1}$ and $w_{2}^{\prime \prime}=I_{2} w_{2}$, then the product $q=w_{1} w_{2}$ satisfies the fourth order ODE

$$
\begin{equation*}
q^{\prime \prime \prime \prime}-2\left(I_{1}+I_{2}\right) q^{\prime \prime}-6 I^{\prime} q^{\prime}+\left(\left(B_{a}-B_{b}\right)^{2}-2 I^{\prime \prime}\right) q=0 . \tag{0.14}
\end{equation*}
$$

We remark that if $B_{a}=B_{b}$, then $I_{1}=I_{2}$ and $q$ actually satisfies a third order ODE as the second symmetric product of a Lamé equation, which is a useful tool used in Part I in the study of the Lamé curve.

If however $B_{a} \neq B_{b}$, by the definition of $w_{a}$ in (0.7) and the addition law,

$$
q=w_{a} w_{-b}+w_{-a} w_{b}
$$

is an even elliptic function solution to equation (0.14), hence a polynomial in $x=\wp(z)$. This leads to strong constraints on equation (0.14) in the variable $x$ and eventually leads to a contradiction for generic choices of $\sigma_{0}$.

Now we set

$$
\begin{equation*}
Z_{n}(\sigma ; \tau):=W_{n}(Z)(\sigma ; \tau) \tag{0.15}
\end{equation*}
$$

Then $Z_{n}(\sigma ; \tau)$ is pre-modular of weight $\frac{1}{2} n(n+1)$. From the construction and equation (0.12) it is readily seen that $Z_{n}(\sigma ; \tau)$ is the generalization of the Hecke function we are looking for. In fact, for $n \geqslant 1$, we have

Theorem 0.4. - To every scaling family $\left\{u_{\lambda}\right\}$ of solutions to the singular Liouville equation (0.1) on $E_{\tau}$, the zero set $a \in X_{n}$ of its normalized developing map $f$ satisfies $Z_{n}\left(\sigma_{n}(a) ; \tau\right)=0$ with $\sigma_{n}(a) \notin E_{\tau}[2]$. Conversely, given $\sigma_{0} \in E_{\tau} \backslash E_{\tau}[2]$ with $Z_{n}\left(\sigma_{0} ; \tau\right)=0$, there is a unique $a \in X_{n}$ with $\sigma_{n}(a)=\sigma_{0}$ and it determines a developing map $f=w_{a} / w_{-a}$ of a scaling family of solutions to equation (0.1).

The proof is given in Section 4, where we also present a version of it in terms of monodromy groups of Lamé equations in Theorem 4.5.

For $\sigma \in E_{\tau}[N]$, the $N$-torsion points, the modular form $Z_{2}(\sigma ; \tau)$ and $Z_{3}(\sigma ; \tau)$ were first constructed by Dahmen [3] in his study on integral Lamé equations (0.6) with algebraic solutions (i.e., with finite monodromy group). For $n \geqslant 4$, the existence of a modular form $Z_{n}(\sigma ; \tau)$ of weight $\frac{1}{2} n(n+1)$ was also conjectured in [3]. This is now settled by our results.
0.2.3. Relation with finite gap theory. - It remains to find effective and explicit constructions of $Z_{n}$. Since $\sigma_{n}$ is defined by the addition map, which is purely algebraic, in principle this allows us to compute the polynomial $W_{n}(\mathbf{z})$ for any $n \in \mathbb{N}$ by eliminating variables $B$ and $C$. In practice the needed calculations are very demanding and time consuming.

In a different direction, the Lamé curve had also been studied extensively in the finite band integration theory. In the complex case this theory concerns the eigenvalue problem on a second order ODE

$$
L w:=w^{\prime \prime}-I w=B w
$$

with eigenvalue $B$. The potential $I=I(z)$ is called a finite-gap (band) potential if the ODE has only logarithmic free solutions except for finitely many $B \in \mathbb{C}$. The integral Lamé equations (with $I(z)=n(n+1) \wp(z))$ provide the first non-trivial examples of them. Using this theory, Maier [12] had recently written down an explicit map $\pi_{n}: \bar{X}_{n} \rightarrow E$ in terms of the coordinates $(B, C)$ on $\bar{X}_{n}$ (cf. Theorem 5.3). It turns out we can prove

Theorem 0.5. - The map $\pi_{n}$ agrees with $\sigma_{n}: \bar{X}_{n} \rightarrow E$.

This is part of Theorem 5.6 where another presentation of $\mathbf{z}_{n}$ in this context is also given. The main idea in the proof is to compare the Hermite-Halphen ansatz (0.7) with the Hermite-Krichever ansatz (given in (5.1)) of solutions to the Lamé equations (0.6).

This provides an alternative way to compute $W_{n}(\mathbf{z})$ by eliminating $B, C$. In particular the weight 10 pre-modular form $Z_{4}(\sigma ; \tau)$ is explicitly written down in Example 5.10. The existence and effective construction of $Z_{n}(\sigma ; \tau)$ opens the door to extend our complete results on equation (0.1) for $n=1$ to general $n \in \mathbb{N}$.

As a related application, the explicit expression of $Z_{4}$ is used to solve Dahmen's conjecture on a counting formula for Lamé equations (0.6) with finite monodromy in the $n=4$ case. The method works for general $n$ once $Z_{n}$ is shown to have expected asymptotic behavior at cusps. The details are given in the appendix, written by Y.-C. Chou.

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$$
\text { 1. GEOMETRY OF } B: \bar{X}_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})
$$

In this section we investigate the local structure of the branch points of the hyperelliptic projection $B: \bar{X}_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})$.
1.1. Some useful formulas from Part I. - We give quantitative descriptions on those results recalled in Section 0.1 which will be used in this paper.

Let $f$ be a normalized developing map of a solution $u$ to equation (0.1) with simple zeros $\left\{\left[a_{1}\right], \ldots,\left[a_{n}\right]\right\}$ and simple poles $\left\{\left[b_{1}\right], \ldots,\left[b_{n}\right]\right\}$ in $E$. One of the crucial properties proved in Part I (cf. Section 0.1.1 (i)) is the equality

$$
\left\{\left[b_{1}\right], \ldots,\left[b_{n}\right]\right\}=\left\{\left[-a_{1}\right], \ldots,\left[-a_{n}\right]\right\} .
$$

With this being established, the logarithmic derivative $g:=(\log f)^{\prime}=f^{\prime} / f$ is readily seen to be an even elliptic function on $E$ of the form

$$
\begin{equation*}
g(z)=\sum_{i=1}^{n} \frac{\wp^{\prime}\left(a_{i}\right)}{\wp(z)-\wp\left(a_{i}\right)} . \tag{1.1}
\end{equation*}
$$

It has simple poles at $\pm\left[a_{i}\right]$ and only one zero at $[0]$. Hence $\operatorname{ord}_{z=0} g(z)=2 n$. It leads to the $n-1$ equations in $\left[a_{1}\right], \ldots,\left[a_{n}\right]$ given in equations (0.4): under the algebraic coordinates $\left(w, x_{i}, y_{i}\right)=\left(\wp(z), \wp\left(a_{i}\right), \wp^{\prime}\left(a_{i}\right)\right)$,

$$
\begin{aligned}
g(z) & =\sum_{i=1}^{n} \frac{1}{w} \frac{y_{i}}{1-x_{i} / w} \\
& =\sum_{i=1}^{n} \frac{y_{i}}{w}+\sum_{i=1}^{n} \frac{y_{i} x_{i}}{w^{2}}+\cdots+\sum_{i=1}^{n} \frac{y_{i} x_{i}^{r}}{w^{r+1}}+\cdots .
\end{aligned}
$$

Since $\operatorname{ord}_{z=0} g(z)=2 n$ and $1 / w$ has a zero at $z=0$ of order two, we get $x_{i} \neq x_{j}$ for $i \neq j$ and

$$
\begin{equation*}
\sum_{i=1}^{n} y_{i} x_{i}^{r}=0, \quad r=0, \ldots, n-2 \tag{1.2}
\end{equation*}
$$

Equations (1.2), together with the Weierstrass equation $y_{i}^{2}=4 x_{i}^{3}-g_{2} x_{i}-g_{3}$ for all $i=1, \ldots, n$, give the algebraic form of equations (0.4).

The Green function equation (0.5) is equivalent to the type II constraints (0.3) ([10, Lem. 2.4]). Indeed, by the addition law,

$$
\begin{align*}
f=\exp \int g d z & =\exp \int \sum_{i=1}^{n}\left(2 \zeta\left(a_{i}\right)-\zeta\left(a_{i}-z\right)-\zeta\left(a_{i}+z\right)\right) d z \\
& =e^{2 \sum_{i=1}^{n} \zeta\left(a_{i}\right) z} \prod_{i=1}^{n} \frac{\sigma\left(z-a_{i}\right)}{\sigma\left(z+a_{i}\right)} . \tag{1.3}
\end{align*}
$$

The monodromy effect on $f$ is then calculated from the standard formula

$$
\begin{equation*}
\sigma\left(z+\omega_{j}\right)=-e^{\eta_{j}\left(z+\frac{1}{2} \omega_{j}\right)} \sigma(z), \quad j=1,2 . \tag{1.4}
\end{equation*}
$$

Let $a_{i}=r_{i} \omega_{1}+s_{i} \omega_{2}$ for $i=1, \ldots, n$. By way of the Legendre relation $\eta_{1} \omega_{2}-\eta_{2} \omega_{1}=2 \pi i$ we compute easily that

$$
\begin{align*}
& f\left(z+\omega_{1}\right)=e^{-4 \pi i \sum_{i} s_{i}+2 \omega_{1}\left(\sum \zeta\left(a_{i}\right)-\sum r_{i} \eta_{1}-\sum s_{i} \eta_{2}\right)} f(z), \\
& f\left(z+\omega_{2}\right)=e^{4 \pi i \sum_{i} r_{i}+2 \omega_{2}\left(\sum \zeta\left(a_{i}\right)-\sum r_{i} \eta_{1}-\sum s_{i} \eta_{2}\right)} f(z) . \tag{1.5}
\end{align*}
$$

By equation (0.9) and the linear independence of $\omega_{1}$ and $\omega_{2}$, the equivalence of equation (0.5) and (0.3) follows immediately.

In Section 0.1 .2 we have reviewed the hyperelliptic structure $B: \bar{X}_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})$ on the Lamé curve induced by $a \mapsto B_{a}=(2 n-1) \sum \wp\left(a_{i}\right)$. Also $\bar{X}_{n}$ contains the Liouville curve $X_{n}$ as the unramified loci. By way of Lamé equation (0.6) with $B=B_{a}$ and by setting $f=w_{a} / w_{-a}$, where $w_{a}$ is the ansatz solution ( 0.7 ), we see that solving equation (0.1) is equivalent to solving the integral Lamé equation (0.6) with unitary projective monodromy groups.

The finite part $Y_{n}$ of $\bar{X}_{n}$ is defined by equation $C^{2}=\ell_{n}(B)$ where the Lamé polynomial $\ell_{n}(B)$ is of degree $2 n+1$ and can be effectively computed (cf. [1, Th. 7.4]). Later we will discuss factorization properties of $\ell_{n}(B)$ in Proposition 2.2. Here, we focus on formulas which lead to parametrization of the Lamé curve near branched points, i.e., $a \in \bar{X}_{n}$ with $a=-a$.

Proposition 1.1 ([1, (7.5.3) \& Prop. 7.5])
(1) Let $a \in Y_{n}$, then $(B, C)$ can be parameterized by $B(a)=B_{a}$ and

$$
\begin{equation*}
C(a)=\wp^{\prime}\left(a_{i}\right) \prod_{j \neq i}\left(\wp\left(a_{i}\right)-\wp\left(a_{j}\right)\right) . \tag{1.6}
\end{equation*}
$$

Here, formula (1.6) is valid for any $i \in\{1, \ldots, n\}$.
(2) The limiting equations of equations (1.2) at $a=0^{n} \in \bar{X}_{n}$ are given by

$$
\begin{equation*}
\sum_{i=1}^{n} t_{i}^{2 r+1}=0, \quad r=1, \ldots, n-1 \tag{1.7}
\end{equation*}
$$

subject to the non-degenerate constraints $t_{i} \neq 0, t_{i} \neq-t_{j}$ for $i \neq j$.
Equations (1.7) have a unique non-degenerate solution in $\mathbb{P}^{n-1}(\mathbb{C})$ up to permutations. It gives rise to the unique tangent direction

$$
[t]=\left[t_{1}: \cdots: t_{n}\right] \in \mathbb{P}\left(T_{0^{n}}\left(\bar{X}_{n}\right)\right) \subset \mathbb{P}\left(T_{0^{n}}\left(\operatorname{Sym}^{n} E\right)\right)
$$

Remark 1.2. - Notice that (1.6) arises from (1.1) and $\operatorname{ord}_{z=0} g_{a}(z)=2 n$ in

$$
g_{a}(z):=\sum_{i=1}^{n} \frac{\wp^{\prime}\left(a_{i}\right)}{\wp(z)-\wp\left(a_{i}\right)}=\frac{\sum_{i=1}^{n} \wp^{\prime}\left(a_{i}\right) \prod_{j \neq i}\left(\wp(z)-\wp\left(a_{j}\right)\right)}{\prod_{i=1}^{n}\left(\wp(z)-\wp\left(a_{i}\right)\right)},
$$

where the numerator reduces to the constant $C(a)$.
1.2. Local structures on finite and infinite branch points. - By working on formula (1.6), we may relate various local parameters of $Y_{n}$ as follows:

Lemma 1.3. - Let $a=\left\{a_{1}, \ldots, a_{n}\right\} \in Y_{n} \backslash X_{n}$ be a finite branched point and $b=$ $\left\{b_{1}, \ldots, b_{n}\right\} \in X_{n}$ be a point near $a$.
(i) Let $i$ be an index with $a_{i} \in E[2]$, then $C$ can be used as a parameter for $b_{i}-a_{i}$ with $b_{i}^{\prime}(0) \neq 0, \infty$.
(ii) Let $i$ be an index with $a_{i} \notin E[2]$, and let $i^{\prime}$ be the corresponding index with $a_{i^{\prime}}=-a_{i}$. Then $C$ can be used as a parameter for $b_{i}+b_{i^{\prime}}$ with $\left(b_{i}+b_{i^{\prime}}\right)^{\prime}(0) \neq 0, \infty$.

## Proof

(i) For $a_{i}=-a_{i}\left(2\right.$-torsion) in $E$, we have $\wp^{\prime}\left(a_{i}\right)=0$ and $\wp^{\prime}\left(b_{i}\right)=\wp^{\prime \prime}\left(a_{i}\right)\left(b_{i}-a_{i}\right)+$ $o\left(\left|b_{i}-a_{i}\right|\right)$. Formula (1.6) then implies that

$$
C(b)=\left[\wp^{\prime \prime}\left(a_{i}\right) \prod_{j \neq i}\left(\wp\left(a_{i}\right)-\wp\left(a_{j}\right)\right)\right]\left(b_{i}-a_{i}\right)+o\left(\left|b_{i}-a_{i}\right|\right) .
$$

We get the inverse map $C \mapsto b_{i}(C)-a_{i}$ since $\wp^{\prime \prime}\left(a_{i}\right) \neq 0$ for $a_{i} \in E[2]$.
(ii) Similarly, for $a_{i} \notin E[2]$ with $a_{i^{\prime}}=-a_{i}$, we have

$$
\wp\left(b_{i}\right)-\wp\left(b_{i^{\prime}}\right)=\wp^{\prime}\left(-b_{i^{\prime}}\right)\left(b_{i}+b_{i^{\prime}}\right)+o\left(\left|b_{i}+b_{i^{\prime}}\right|\right) .
$$

Since $-b_{i^{\prime}}$ is close to $-a_{i^{\prime}}=a_{i}$, we get

$$
C(b)=\left[\wp^{\prime}\left(a_{i}\right)^{2} \prod_{j \neq i, i^{\prime}}\left(\wp\left(a_{i}\right)-\wp\left(a_{j}\right)\right)\right]\left(b_{i}+b_{i^{\prime}}\right)+o\left(\left|b_{i}+b_{i^{\prime}}\right|\right) .
$$

Then we get the inverse map $C \mapsto\left(b_{i}+b_{i^{\prime}}\right)(C)$ since $\wp^{\prime}\left(a_{i}\right) \neq 0$.
Remark 1.4. - From formula (1.6) (with $a$ being substituted by $b$ ) we have $b_{i^{\prime}}(C)=$ $-b_{i}(-C)$, hence $b_{i^{\prime}}^{\prime}(C)=b_{i}^{\prime}(-C)$. If $b_{i}^{\prime}(0)$ and $b_{i^{\prime}}^{\prime}(0)$ are finite then they are equal and non-vanishing. If this holds for all $i$ then $C \mapsto b(C)$ is a holomorphic map in each local branch of $Y_{n}$ at $a$ and we conclude that $a \in Y_{n}$ is either a smooth point or a nodal
singularity $\left(y^{2}=x^{2}\right)$. We will see in Remark 5.14 that this is indeed the case. At this point we conclude only the finiteness of $\left(b_{i}+b_{i^{\prime}}\right)^{\prime}(0)$ as stated in Lemma 1.3 (ii).

Now we give a precise description of the unique tangent direction at $0^{n} \in \bar{X}_{n}$. Denote by $[t]=\left[t_{1}, \ldots, t_{n}\right]$ the homogeneous coordinates satisfying the limiting equations (1.7) with non-degenerate constraints.

Let $p_{j}=\sum_{i=1}^{n} t_{i}^{j}$ be the $j$-th Newton symmetric polynomial and $\lambda_{j}$ be the $j$-th elementary symmetric polynomial of $t_{1}, \ldots, t_{n}$. We use the convention that $p_{0}=0$, $\lambda_{0}=1$ and $\lambda_{j}=0$ for $j>n$. By a Vandermonde determinant argument we have $\lambda_{1}=p_{1} \neq 0$.

Proposition 1.5. - The point $[t] \in \mathbb{P}\left(T_{0^{n}}\left(\bar{X}_{n}\right)\right)$ is characterized by the recursions

$$
\begin{equation*}
\lambda_{k+1}=2 \frac{(k-n)}{(k-2 n)(k+1)} \lambda_{k} \lambda_{1}, \quad 1 \leqslant k \leqslant n-1 . \tag{1.8}
\end{equation*}
$$

Proof. - By Proposition 1.1 (2), namely the uniqueness of non-degenerate solutions to (1.7), we only need to verify that that the point defined by (1.8) satisfies $p_{3}=p_{5}=$ $\cdots=p_{2 n-1}=0$.

Since $\lambda_{1} \neq 0$, without loss of generality we assume that $\lambda_{1}=1$. The recursions in (1.8), with $\lambda_{1}=1$, are equivalent to saying that the polynomial

$$
Q(t):=\prod_{i=1}^{n}\left(1+t_{i} t\right)=\sum_{k=0}^{n} \lambda_{k} t^{k}
$$

coincides with the hypergeometric function $F(-n ;-2 n \mid 2 t)$. That is, $f(t):=Q\left(\frac{1}{2} t\right)$ satisfies the hypergeometric equation

$$
\begin{equation*}
(\delta(\delta-(2 n+1))-t(\delta-n)) f=t\left(t f^{\prime \prime}-(t+2 n) f^{\prime}+n f\right)=0 \tag{1.9}
\end{equation*}
$$

where $\delta=t d / d t$. Then $g:=(\log f)^{\prime}=f^{\prime} / f$ satisfies

$$
\begin{equation*}
g^{\prime}+g^{2}-\left(1+\frac{2 n}{t}\right) g+\frac{n}{t}=0 . \tag{1.10}
\end{equation*}
$$

Write $g=\sum_{k=0}^{\infty} g_{k} t^{k}$. From

$$
g=\frac{f^{\prime}}{f}=\sum_{i=1}^{n} \frac{\frac{1}{2} t_{i}}{1+\frac{1}{2} t_{i} t}=\sum_{k=0}^{\infty} \frac{p_{k+1}}{2^{k+1}} t^{k},
$$

we have $g_{k}=p_{k+1} / 2^{k+1}$ and it suffices to show $g_{2}=g_{4}=\cdots=g_{2 n-2}=0$.
From the series expansions

$$
\begin{aligned}
g^{\prime} & =g_{1}+\sum_{k=1}^{\infty}(k+1) g_{k+1} t^{k}, \\
g^{2} & =\frac{1}{4}+\sum_{k=1}^{\infty}\left(\sum_{j=0}^{k} g_{j} g_{k-j}\right) t^{k}, \\
\left(1+\frac{2 n}{t}\right) g & =\frac{n}{t}+\sum_{k=0}^{\infty}\left(g_{k}+2 n g_{k+1}\right) t^{k},
\end{aligned}
$$

we get by equation (1.10) that

$$
\begin{aligned}
-\frac{n}{t} & +\frac{n}{t}+\left(g_{1}+\frac{1}{4}-\left(\frac{1}{2}+2 n g_{1}\right)\right) \\
& +\sum_{k=1}^{\infty}\left((k+1) g_{k+1}+\sum_{j=0}^{k} g_{j} g_{k-j}-\left(g_{k}+2 n g_{k+1}\right)\right) t^{k}=0
\end{aligned}
$$

Hence $g_{1}=-1 /(4(2 n-1))$ and for all $k \geqslant 1$ we have recursions

$$
\begin{equation*}
(2 n-(k+1)) g_{k+1}=\sum_{j=1}^{k-1} g_{j} g_{k-j} . \tag{1.11}
\end{equation*}
$$

(We have used the fact $g_{0}=\frac{1}{2}$ to remove one $g_{k}$.)
For $k=1$, the sum is empty and we get $g_{2}=0$. Now we conclude the proof by induction. Suppose that $g_{2}=g_{4}=\cdots=g_{2 m}=0$ with $m<n-1$. Then for $k=2 m+1$ we have $2 n-(k+1)=2(n-(m+1))>0$, and the recursions (1.11) show that $g_{2(m+1)}$ is a sum of $g_{j} g_{k-j}$ with either $j$ or $k-j$ being an even number no bigger than $2 m$. Hence $g_{2(m+1)}=0$, and this completes the induction.

## 2. Geometry of $\sigma_{n}: \bar{X}_{n} \rightarrow E$

The aim of this section is to prove Theorem 0.2.

### 2.1. Lamé functions [cf. Section 0.1.2 (v)]

Definition 2.1. - The type of a point $a \in Y_{n} \backslash X_{n}$ is defined to be the number of half periods contained in $a=\left\{a_{i}\right\}$. Hence there are four types O, I, II, III. For $n=2 k$, $a$ must be of type O or II. For $n=2 k+1, a$ must be of type I or III.

The type of a Lame functions $w_{a}$ (with $[a]=[-a] \in Y_{n} \backslash X_{n}$ ) is defined to be the type of its zero set $a$ accordingly.

There are factorizations of the polynomial $\ell_{n}(B)$ according to the types:
Proposition $2.2([5,14])$. - We may decompose $\ell_{n}\left(B ; g_{2}, g_{3}\right)$ as

$$
\ell_{n}\left(B ; g_{2}, g_{3}\right)=c_{n}^{2} l_{0}(B) l_{1}(B) l_{2}(B) l_{3}(B)
$$

where $c_{n} \in \mathbb{Q}_{>0}$ is a constant, $l_{i}(B)$ 's are monic polynomials in $B$ whose coefficients are polynomials in $e_{1}, e_{2}, e_{3}$ such that
(1) For $n=2 k, l_{0}(B)=\prod\left(B-B_{a}\right)$ with a being of type $O$, and $\operatorname{deg} l_{0}(B)=$ $\frac{1}{2} n+1=k+1$. For $i=1,2,3, l_{i}(B)=\prod\left(B-B_{a}\right)$ with a being of type II which does not contain $\frac{1}{2} \omega_{i}$, also $\operatorname{deg} l_{i}(B)=\frac{1}{2} n=k$.
(2) For $n=2 k+1, l_{0}(B)=\Pi\left(B-B_{a}\right)$ with a being of type III, and $\operatorname{deg} l_{0}(B)=$ $\frac{1}{2}(n-1)=k$. For $i=1,2,3, l_{i}(B)=\prod\left(B-B_{a}\right)$ with a being of type $I$ which contains $\frac{1}{2} \omega_{i}$, also $\operatorname{deg} l_{i}(B)=\frac{1}{2}(n+1)=k+1$.

We remark that Proposition 2.2, together with Proposition 1.1 and Lemma 1.3, will be used in the proof of Theorem 0.2. Here are some examples to illustrate Proposition 2.2:

Example 2.3 (Decompositions of $\ell_{n}(B)$ for $1 \leqslant n \leqslant 5$ )
(1) $n=1, k=0, \bar{X}_{1} \cong E$,

$$
C^{2}=\ell_{1}(B)=4 B^{3}-g_{2} B-g_{3}=4 \prod_{i=1}^{3}\left(B-e_{i}\right)
$$

(2) $n=2, k=1$, (notice that $\left.e_{1}+e_{2}+e_{3}=0\right)$

$$
\begin{aligned}
C^{2}=\ell_{2}(B) & =\frac{4}{81} B^{5}-\frac{7}{27} g_{2} B^{3}+\frac{1}{3} g_{3} B^{2}+\frac{1}{3} g_{2}^{2} B-g_{2} g_{3} \\
& =\frac{2^{2}}{3^{4}}\left(B^{2}-3 g_{2}\right) \prod_{i=1}^{3}\left(B+3 e_{i}\right) .
\end{aligned}
$$

(3) $n=3, k=1, \operatorname{deg} l_{i}(B)=2$ for $i=1,2,3$,

$$
\begin{aligned}
C^{2}=\ell_{3}(B)= & \frac{1}{2^{2} 3^{4} 5^{4}} B\left(16 B^{6}-504 g_{2} B^{4}+2376 g_{3} B^{3}\right. \\
& \left.+4185 g_{2}^{2} B^{2}-36450 g_{2} g_{3} B+91125 g_{3}^{2}-3375 g_{2}^{3}\right) \\
= & \frac{2^{2}}{3^{4} 5^{4}} B \prod_{i=1}^{3}\left(B^{2}-6 e_{i} B+15\left(3 e_{i}^{2}-g_{2}\right)\right)
\end{aligned}
$$

(4) $n=4, k=2, \operatorname{deg} l_{0}(B)=3$,

$$
C^{2}=\ell_{4}(B)=\frac{1}{3^{8} 5^{4} 7^{4}}\left(B^{3}-52 g_{2} B+560 g_{3}\right) \prod_{i=1}^{3}\left(B^{2}+10 e_{i} B-7\left(5 e_{i}^{2}+g_{2}\right)\right)
$$

(5) $n=5, k=2, \operatorname{deg} l_{i}(B)=3$ for $i=1,2,3$,

$$
\begin{aligned}
C^{2}=\ell_{5}(B)= & \frac{1}{3^{12} 5^{4} 7^{4} 11^{2}}\left(B^{2}-27 g_{2}\right) \\
& \times \prod_{i=1}^{3}\left(B^{3}-15 e_{i} B^{2}+\left(315 e_{i}^{2}-132 g_{2}\right) B+e_{i}\left(2835 e_{i}^{2}-540 g_{2}\right)\right)
\end{aligned}
$$

2.2. The degree of the addition map $\sigma_{n}$. - We are now ready to study the addition $\operatorname{map} \sigma_{n}: \bar{X}_{n} \rightarrow E, a \mapsto \sigma_{n}(a)=\sum_{i=1}^{n} a_{i}$ defined in (0.13) and determine its degree $\operatorname{deg} \sigma_{n}$.

For a finite morphism of irreducible curves $f: X \rightarrow Y$, the function field $K(X)$ is a finite extension of $K(Y)$ and $\operatorname{deg} f=[K(X): K(Y)]$. Geometrically, $\operatorname{deg} f$ is the number of points for a general fiber $f^{-1}(p), p \in Y$. If the image curve $Y$ is smooth, the degree is equal to the length of the scheme-theoretic fiber $f^{-1}(p)$ for any $p \in Y$. A standard reference is [6].

Theorem 2.4 ( $=$ Theorem 0.2). - The map $\sigma_{n}: \bar{X}_{n} \rightarrow E$ has degree $\frac{1}{2} n(n+1)$.
Proof. - The idea is to apply Theorem of the Cube [13, p. 58, Cor. 2] for morphisms from an arbitrary variety $V$ (not necessarily smooth) into abelian varieties (here the torus $E$ ). For any three morphisms $f, g, h: V \rightarrow E$ and a line bundle $L \in \operatorname{Pic} E$, we have
(2.1) $(f+g+h)^{*} L \cong(f+g)^{*} L \otimes(g+h)^{*} L \otimes(h+f)^{*} L \otimes f^{*} L^{-1} \otimes g^{*} L^{-1} \otimes h^{*} L^{-1}$.

We will apply it to the algebraic curve $V=V_{n} \subset E^{n}$ which consists of the ordered $n$-tuples a's so that $V_{n} / S_{n}=\bar{X}_{n}$. (Here, $S_{n}$ is the permutation group on $n$ letters.)

For any line bundle $L$ and any finite morphism $f: V \rightarrow E$, we have $\operatorname{deg} f^{*} L=$ $\operatorname{deg} f \operatorname{deg} L$. In the following we fix an $L$ with $\operatorname{deg} L=1$.

We prove inductively that for $j=1, \ldots, n$ the morphism $f_{j}: V_{n} \rightarrow E$ defined by

$$
f_{j}(a):=a_{1}+\cdots+a_{j}
$$

has $\operatorname{deg} f_{j}^{*} L=\frac{1}{2} j(j+1) n$ !. The case $j=n$ then gives the result since $f_{n}$ is a finite morphism which descends to $\sigma_{n}$ under the $S_{n}$ action. (Notice that for $j<n$ the map $f_{j}$ does not descend to a map on $\bar{X}_{n}$.)

Assuming first that it has been proved for $j=1,2$. To go from $j$ to $j+1$, we let $f(a)=f_{j-1}(a), g(a)=a_{j}$, and $h(a)=a_{j+1}$. Then by $(2.1), f_{j+1}^{*} L$ has degree $n!$ times

$$
\frac{1}{2} j(j+1)+3+\frac{1}{2} j(j+1)-\frac{1}{2}(j-1) j-1-1=\frac{1}{2}(j+1)(j+2)
$$

as expected. It remains to investigate the case $j=1$ and $j=2$.
For $j=1$, by Section 0.1 .2 (iii)-(iv), the inverse image of $0 \in E$ under $f_{1}: V_{n} \rightarrow E$ consists of a single point $0^{n}$. By Proposition 1.1 (2), the limiting system of equations (1.7) has a unique non-degenerate solution in $\mathbb{P}^{n-1}(\mathbb{C})$ up to permutations. From this, we conclude that there are precisely $n!$ branches of $V_{n} \rightarrow E$ near $0^{n}$. For a point $b \in E^{\times}$close to 0 , each branch will contribute a point $a$ with $a_{1}=b$. In particular, $f_{1}$ is a finite morphism and $\operatorname{deg} f_{1}^{*} L=\operatorname{deg} f_{1}=n$ !.

For $j=2$, we consider the (scheme-theoretic) inverse image of $0 \in E$ under $f_{2}: V_{n} \rightarrow E$. Namely $V_{n} \ni a \mapsto a_{1}+a_{2}=0$.

The point $a=0$ again contributes degree $n!$ by a similar branch argument. Indeed, over each branch near $0^{n}$ we may represent $a=\left(a_{i}(z)\right)$ by an analytic curve in $z$. Then condition $t_{i}+t_{j} \neq 0$ in Proposition 1.1 (2) implies that $z \mapsto a_{1}(z)+a_{2}(z) \in E$ is still locally biholomorphic for $z$ close to 0 . As a byproduct, since every irreducible component contains a branch near $0^{n}, f_{2}$ is necessarily a finite morphism and $\operatorname{deg} f_{2}^{*} L=\operatorname{deg} f_{2}$.

For those points $a \neq 0$ with $f_{2}(a)=0$, we have $a_{1}=-a_{2}$ and thus $a=-a$ by Section 0.1.2 (iii), i.e., $a$ corresponds to a branch point for the hyperelliptic projection $Y_{n} \rightarrow \mathbb{C}$. Let $b=\left(b_{1}, \ldots, b_{n}\right) \in V_{n}$ be a point near $a$. By Lemma 1.3, we see that $C \mapsto\left(b_{1}+b_{2}\right)(C)$ is bi-holomorphic near $C=0$. If $a \in V_{n}$ is a non-singular point then $C$ is a local parameter and $f_{2}$ is unramified at $a$. In that case the degree contribution at $a$ is one. We first treat the case that all branch points are non-singular points:

If $n=2 k$, by Proposition 2.2 (1) the degree contribution from type O points $a=$ $\left\{ \pm a_{1}, \ldots, \pm a_{k}\right\}$ is given by

$$
(k+1) \times(k \times 2 \times(n-2)!),
$$

while the degree from the type II points $\left\{ \pm a_{1}, \ldots, \pm a_{k-1}, \frac{1}{2} \omega_{i}, \frac{1}{2} \omega_{j}\right\}$ is

$$
3 \times k \times((k-1) \times 2 \times(n-2)!) .
$$

The sum is $2\left(4 k^{2}-2 k\right)(n-2)!=2 n!$.

If $n=2 k+1$, by Proposition $2.2(2)$, the degree contribution from type III points $\left\{ \pm a_{1}, \ldots, \pm a_{k-1}, \frac{1}{2} \omega_{1}, \frac{1}{2} \omega_{2}, \frac{1}{2} \omega_{3}\right\}$ is

$$
k \times((k-1) \times 2 \times(n-2)!),
$$

while the type I points $\left\{ \pm a_{1}, \ldots, \pm a_{k}, \frac{1}{2} \omega_{i}\right\}$ contribute

$$
3 \times(k+1) \times(k \times 2 \times(n-2)!)
$$

The sum is again $2\left(4 k^{2}+2 k\right)(n-2)!=2 n!$. Thus in both cases we get the total degree $n!+2 n!=3 n!$ as expected.

If $Y_{n}=Y_{n}\left(\tau_{0}\right)$ is singular, let $a \in Y_{n} \backslash X_{n}$ be a singular (branch) point with $C^{2}=$ $h(B)\left(B-B_{a}\right)^{m}, m \geqslant 2$ and $h\left(B_{a}\right) \neq 0$. The curve $Y_{n}$ arises from flat degenerations of smooth curves $Y_{n}(\tau)$ where $m$ linear factors become the same. By Lemma 1.3, $f_{2}$ leads to an analytic equivalence between $C$ and $b_{1}+b_{2}$ near $a$. That is, the equation $b_{1}+b_{2}=0$ is simply $C=0$. Let $\widetilde{B}=B-B_{a}$. Then the analytic structure sheaf of $f_{2}^{-1}(0)$ at $a$ is given by

$$
\mathbb{C} \llbracket \widetilde{B}, C \rrbracket /\left\langle C^{2}-h\left(B_{a}+\widetilde{B}\right) \widetilde{B}^{m}, C\right\rangle \cong \mathbb{C} \llbracket \widetilde{B} \rrbracket /\left\langle\widetilde{B}^{m}\right\rangle,
$$

which also has length $m$. This shows that $f_{2}$ is compatible with the degeneration and the degree counting then follows from the smooth case.

## 3. The primitive generator $\mathbf{z}_{n}$

We prove Theorem 0.3, in a more precise form, in Theorem 3.2.

### 3.1. Setup of the proof

Definition 3.1 (The fundamental rational function $\mathbf{z}_{n}$ on $\bar{X}_{n}$ )
The function

$$
\mathbf{z}_{n}\left(a_{1}, \ldots, a_{n}\right):=\zeta\left(\sum_{i=1}^{n} a_{i}\right)-\sum_{i=1}^{n} \zeta\left(a_{i}\right), \quad a_{i} \in \mathbb{C}
$$

is meromorphic and periodic in each $a_{i}$, hence it defines a rational function on $E^{n}$. By symmetry, it descends to a rational function on $\operatorname{Sym}^{n} E$. We denote the restriction $\left.\mathbf{z}_{n}\right|_{\bar{X}_{n}}$ also by $\mathbf{z}_{n}$, which is a rational function on $\bar{X}_{n}$ with poles along the fiber $\sigma_{n}^{-1}(0)$.

The importance of $\mathbf{z}_{n}$ is readily seen from investigation on the Green function equation (0.5): Let $a_{i}=r_{i} \omega_{1}+s_{i} \omega_{2}$. Then by (0.9),

$$
\begin{align*}
-4 \pi \sum \frac{\partial G}{\partial z}\left(a_{i}\right) & =\sum\left(\zeta\left(r_{i} \omega_{1}+s_{i} \omega_{2}\right)-r_{i} \eta_{1}-s_{i} \eta_{2}\right)  \tag{3.1}\\
& =Z\left(\sum a_{i}\right)-\mathbf{z}_{n}(a)
\end{align*}
$$

Hence $\sum_{i=1}^{n} \nabla G\left(a_{i}\right)=0 \Longleftrightarrow \mathbf{z}_{n}(a)=Z\left(\sigma_{n}(a)\right)$. This links $\sigma_{n}(a)$ with $\mathbf{z}_{n}$.
Theorem 3.2 ( $=$ Theorem 0.3). - There is a weighted homogeneous polynomial

$$
W_{n}(\mathbf{z}) \in \mathbb{Q}\left[g_{2}, g_{3}, \wp(\sigma), \wp^{\prime}(\sigma)\right][\mathbf{z}]
$$

of $\mathbf{z}$-degree $d_{n}=\operatorname{deg} \sigma_{n}$ such that for $\sigma=\sigma_{n}(a)=\sum a_{i}$, we have

$$
W_{n}\left(\mathbf{z}_{n}\right)(a)=0
$$

Here, the weights of $\mathbf{z}, \wp(\sigma), \wp^{\prime}(\sigma), g_{2}, g_{3}$ are $1,2,3,4,6$ respectively.
Indeed, $\mathbf{z}_{n}(a)$ is a primitive generator of the finite extension of rational function field $K\left(\bar{X}_{n}\right)$ over $K(E)$ with $W_{n}(\mathbf{z})$ being its minimal polynomial. ${ }^{(4)}$

Moreover, the extension is integral over the affine curve $E^{\times}$.
Proof. - There is nothing to prove for $n=1$, so we assume that $n \geqslant 2$.
Since $\mathbf{z}_{n} \in K\left(\bar{X}_{n}\right)$, which is algebraic over $K(E)$ with degree $d_{n}$, its minimal polynomial $W_{n}(\mathbf{z}) \in K(E)[\mathbf{z}]$ exists with $d:=\operatorname{deg} W_{n} \mid d_{n}$.

Notice that for $\sigma_{0} \in E$ being outside the branch loci of $\sigma_{n}: \bar{X}_{n} \rightarrow E$, there are precisely $d_{n}$ different points $a=\left\{a_{1}, \ldots, a_{n}\right\} \in \bar{X}_{n}$ with $\sigma_{n}(a)=\sum a_{i}=\sigma_{0}$. Thus for the rational function $\mathbf{z}_{n}=\zeta\left(\sum a_{i}\right)-\sum \zeta\left(a_{i}\right) \in K\left(\bar{X}_{n}\right)$ to be a primitive generator, it is sufficient to show that $\mathbf{z}_{n}$ has exactly $d_{n}$ branches over $K(E)$. That is, $\sum \zeta\left(a_{i}\right)$ gives different values for different choices of those $a$ above $\sigma_{0}$. Indeed, for any given $\sigma=\sigma_{0}$, the polynomial $W_{n}(\mathbf{z})=0$ has at most $d$ roots. But now $\mathbf{z}_{n}(a)$ with $\sigma_{n}(a)=\sigma_{0}$ gives $d_{n}$ distinct roots of $W_{n}(\mathbf{z})$, hence we must conclude $d=d_{n}$ and $\mathbf{z}_{n}$ is a primitive generator.

Hence it is sufficient to show the following more precise result:
Theorem 3.3. - Let $a, b \in Y_{n}$ and $\left(a_{1}, \ldots, a_{n}\right),\left(b_{1}, \ldots, b_{n}\right) \in \mathbb{C}^{n}$ be representatives of $a, b$ such that

$$
\begin{equation*}
\sum_{i=1}^{n} a_{i}=\sum_{i=1}^{n} b_{i}, \quad \sum_{i=1}^{n} \zeta\left(a_{i}\right)=\sum_{i=1}^{n} \zeta\left(b_{i}\right) . \tag{3.2}
\end{equation*}
$$

Suppose that $\sum \wp\left(a_{i}\right) \neq \sum \wp\left(b_{i}\right)$. Then $a, b$ are branch points of $Y_{n} \rightarrow \mathbb{C}$ which contains the same number of half periods. Equivalently, the Lamé functions $w_{a}$ and $w_{b}$ are of the same type.

We emphasize that $\bar{X}_{n}$ is not required to be smooth. Theorem 3.2 follows immediately by choosing $\sigma_{0}$ outside the branch loci of $\bar{X}_{n} \rightarrow E$ and $\sigma_{0} \notin E[2]$. Indeed, let $a, b \in Y_{n}$ with $\sigma_{n}(a)=\sigma_{n}(b)=\sigma_{0}$ and $\mathbf{z}_{n}(a)=\mathbf{z}_{n}(b)$, or more precisely with conditions in (3.2) satisfied. By Theorem 3.3 we are left with the case $\sum \wp\left(a_{i}\right)=\sum \wp\left(b_{i}\right)$ but $a \neq b$. Then $a=-b$ by Proposition $1.1(1)$, and in particular $\sigma_{n}(a)=-\sigma_{n}(b)$. Together with $\sigma_{n}(a)=\sigma_{n}(b)$ we conclude that $\sigma_{0}=\sigma_{n}(a)=\sigma_{n}(b) \in E[2]$. This contradicts to the assumption $\sigma_{0} \notin E[2]$. Hence we must have $a=b$.

Since $\mathbf{z}_{n}$ has no poles over $E^{\times}$, it is indeed integral over the affine Weierstrass model of $E^{\times}$with coordinate ring (let $x_{0}=\wp(\sigma), y_{0}=\wp^{\prime}(\sigma)$ )

$$
R\left(E^{\times}\right)=\mathbb{C}\left[x_{0}, y_{0}\right] /\left(y_{0}^{2}-4 x_{0}^{3}-g_{2} x_{0}-g_{3}\right) .
$$

The homogeneity of $W_{n}(\mathbf{z})$ also follows from this.

[^2]Remark 3.4. - By Theorem 0.2 we have $d_{n}=\frac{1}{2} n(n+1)$. We do not use this result in the formulation nor in the proof of Theorem 3.2.

We will give two proofs to Theorem 3.3 in Section 3.2 and Section 3.3. The first proof is longer but contains more information. Both proofs are based on the following basic lemma.

Lemma 3.5 (Tensor product). - Let $I=n(n+1) \wp(z), I_{1}=I+B_{a}$ and $I_{2}=I+B_{b}$. Suppose that $w_{1}^{\prime \prime}=I_{1} w_{1}$ and $w_{2}^{\prime \prime}=I_{2} w_{2}$. Then the product $q:=w_{1} w_{2}$ satisfies the following fourth order ODE:

$$
\begin{equation*}
q^{\prime \prime \prime \prime}-2\left(I_{1}+I_{2}\right) q^{\prime \prime}-6 I^{\prime} q^{\prime}+\left(\left(B_{a}-B_{b}\right)^{2}-2 I^{\prime \prime}\right) q=0 . \tag{3.3}
\end{equation*}
$$

Proof. - This follows from a straightforward computation. Indeed,

$$
\begin{aligned}
q^{\prime} & =w_{1}^{\prime} w_{2}+w_{1} w_{2}^{\prime} \\
q^{\prime \prime} & =\left(I_{1}+I_{2}\right) q+2 w_{1}^{\prime} w_{2}^{\prime} \\
q^{\prime \prime \prime} & =2 I^{\prime} q+\left(I_{1}+I_{2}\right) q^{\prime}+2\left(I_{1} w_{1} w_{2}^{\prime}+I_{2} w_{1}^{\prime} w_{2}\right)
\end{aligned}
$$

Notice that if $a=b$ (or just $B_{a}=B_{b}$ ) then $I_{1}=I_{2}$ and we stop here to get the third order ODE as the symmetric product of the Lamé equation.

In general, we take one more differentiation to get

$$
\begin{aligned}
q^{\prime \prime \prime \prime} & =2 I^{\prime \prime} q+4 I^{\prime} q^{\prime}+\left(I_{1}+I_{2}\right) q^{\prime \prime}+2 I^{\prime} q^{\prime}+2\left(I_{1}+I_{2}\right) w_{1}^{\prime} w_{2}^{\prime}+4 I_{1} I_{2} q \\
& =2\left(I_{1}+I_{2}\right) q^{\prime \prime}+6 I^{\prime} q^{\prime}+\left(2 I^{\prime \prime}-\left(I_{1}-I_{2}\right)^{2}\right) q .
\end{aligned}
$$

This proves the lemma.
Recall from the Hermite-Halphen ansatz in (0.7) that

$$
w_{ \pm a}(z)=e^{ \pm z \sum \zeta\left(a_{i}\right)} \prod_{i=1}^{n} \frac{\sigma\left(z \mp a_{i}\right)}{\sigma(z)}
$$

are solutions to $w^{\prime \prime}=\left(n(n+1) \wp(z)+B_{a}\right) w=: I_{1} w$, and

$$
w_{ \pm b}(z)=e^{ \pm z \sum \zeta\left(b_{i}\right)} \prod_{i=1}^{n} \frac{\sigma\left(z \mp b_{i}\right)}{\sigma(z)}
$$

are solutions to $w^{\prime \prime}=\left(n(n+1) \wp(z)+B_{b}\right) w=: I_{2} w$. Then $q_{a,-b}:=w_{a} w_{-b}$ and $q_{-a, b}:=w_{-a} w_{b}$ are solutions to equation (3.3).
3.2. First proof of Theorem 3.3. - By assumption we have $\sum a_{i}=\sum b_{i}$, hence

$$
\begin{equation*}
q_{a,-b}(z)=\prod_{i=1}^{n} \frac{\sigma\left(z-a_{i}\right) \sigma\left(z+b_{i}\right)}{\sigma^{2}(z)} \tag{3.4}
\end{equation*}
$$

is an elliptic function. Similarly $q_{-a, b}(z)=q_{a,-b}(-z)$ is elliptic. In particular there exists an even elliptic function solution to equation (3.3), namely

$$
\begin{equation*}
Q:=\frac{1}{2}\left(q_{a,-b}+q_{-a, b}\right)=(-1)^{n} \frac{\prod_{i=1}^{n} \sigma\left(a_{i}\right) \sigma\left(b_{i}\right)}{z^{2 n}}(1+O(|z|)) . \tag{3.5}
\end{equation*}
$$

Let $q$ be an even elliptic solution to equation (3.3). Then we may investigate it in variable $x=\wp(z)$. To avoid confusion, we denote

$$
\dot{f}=\partial f / \partial x \quad \text { and } \quad f^{\prime}=\partial f / \partial z
$$

Let $y^{2}=p(x)=4 x^{3}-g_{2} x-g_{3}$. Then $\wp^{\prime}=y, \wp^{\prime \prime}=6 \wp^{2}-\frac{1}{2} g_{2}=\frac{1}{2} \dot{p}(x) \cdot \wp^{\prime \prime \prime}=12 \wp \wp^{\prime}=$ $12 x y, \wp^{\prime \prime \prime \prime}=12 \wp^{\prime 2}+12 \wp \wp^{\prime \prime}=12 p(x)+6 x \dot{p}(x)$. Also

$$
\begin{aligned}
q^{\prime} & =\dot{q} \wp^{\prime}=y \dot{q} \\
q^{\prime \prime} & =\ddot{q} \wp^{\prime 2}+\dot{q} \wp^{\prime \prime}=p(x) \ddot{q}+\frac{1}{2} \dot{p}(x) \dot{q}, \\
q^{\prime \prime \prime} & =\dddot{q} \wp^{\prime 3}+3 \ddot{q} \wp^{\prime} \wp^{\prime \prime}+\dot{q} \wp^{\prime \prime \prime}, \\
q^{\prime \prime \prime \prime} & =\dddot{q} \wp^{\prime 4}+6 \dddot{q} \wp^{\prime 2} \wp^{\prime \prime}+3 \ddot{q}\left(\wp^{\prime \prime}\right)^{2}+4 \ddot{q} \wp^{\prime} \wp^{\prime \prime \prime}+\dot{q} \wp^{\prime \prime \prime \prime} \\
& =p(x)^{2} \dddot{q} \ddot{q}+3 p(x) \dot{p}(x) \dddot{q}+\left(\frac{3}{4} \dot{p}(x)^{2}+48 x p(x)\right) \ddot{q}+(12 p(x)+6 x \dot{p}(x)) \dot{q} .
\end{aligned}
$$

By substituting these into equation (3.3) we get the ODE in $x$ :

$$
\begin{align*}
L_{4} q:= & p^{2} \dddot{q} \\
& +3 p \dot{p} \ddot{q}+\left(\frac{3}{4} \dot{p}^{2}-2\left(2\left(n^{2}+n-12\right) x+\beta\right) p\right) \ddot{q}  \tag{3.6}\\
& -\left(\left(2\left(n^{2}+n-3\right) x+\beta\right) \dot{p}+6\left(n^{2}+n-2\right) p\right) \dot{q} \\
& +\left(\alpha^{2}-n(n+1) \dot{p}\right) q=0,
\end{align*}
$$

where

$$
\begin{equation*}
\alpha:=B_{a}-B_{b} \quad \text { and } \quad \beta:=B_{a}+B_{b} . \tag{3.7}
\end{equation*}
$$

We would like to find constraints for equation $L_{4} q=0$ with $\alpha \neq 0$ to have a polynomial solution $q(x)$. Here, $g_{2}$ and $g_{3}$ could be arbitrary, not necessarily satisfying the non-degeneracy condition $g_{2}^{3}-27 g_{3}^{2} \neq 0$.

Suppose that $q(x)$ is a polynomial in $x$ of degree $m \geqslant 1$ :

$$
\begin{equation*}
q(x)=x^{m}-s_{1} x^{m-1}+s_{2} x^{m-2}-\cdots+(-1)^{m} s_{m} \tag{3.8}
\end{equation*}
$$

which satisfies

$$
\begin{equation*}
\operatorname{deg}_{x} L_{4} q(x) \leqslant 1 \tag{3.9}
\end{equation*}
$$

Then we can solve $s_{j}$ recursively in terms of $\alpha^{2}, \beta$ and $g_{2}, g_{3}$.
Indeed, the top degree term $x^{m+2}$ in equation (3.6) has coefficient

$$
\begin{aligned}
& 16 m(m-1)(m-2)(m-3)+144 m(m-1)(m-2)+108 m(m-1) \\
& \quad-16\left(n^{2}+n-12\right) m(m-1)-24\left(n^{2}+n-3\right) m \\
& \quad-24\left(n^{2}+n-2\right) m-12 n(n+1) \\
& =(m-n)\left(4 m^{3}+(4 n+68) m^{2}+(8 n-101) m+3(n+1)\right)
\end{aligned}
$$

which vanishes precisely when $m=n$. Thus we may assume $m=n$.
The next order term $x^{n+1}$ without the $s_{1}$ factor has coefficient

$$
-8 n(n-1) \beta-12 n \beta=-4 n(2 n+1) \beta
$$

and the coefficient of $-s_{1} x^{n+1}$ is given by

$$
\begin{aligned}
& 16(n-1)(n-2)(n-3)(n-4)+144(n-1)(n-2)(n-3) \\
& \quad+108(n-1)(n-2)-16\left(n^{2}+n-12\right)(n-1)(n-2) \\
& \quad-24\left(n^{2}+n-3\right)(n-1)-24\left(n^{2}+n-2\right)(n-1)-12 n(n+1) \\
& =-8 n(2 n-1)(2 n+1) .
\end{aligned}
$$

Hence

$$
\begin{equation*}
s_{1}=\frac{\beta}{2(2 n-1)} \tag{3.10}
\end{equation*}
$$

Inductively the coefficients of $x^{n+2-i}$ in equation (3.6) for $i=1, \ldots, n$ give rise to recursions to solve $s_{i}$ in terms of $\beta, \alpha^{2}$ and $g_{2}, g_{3}$. Hence we get

Lemma 3.6. - For $i=1, \ldots, n$, there is a polynomial expression

$$
s_{i}=s_{i}\left(\alpha^{2}, \beta, g_{2}, g_{3}\right)=C_{i} \beta^{i}+\cdots
$$

which is homogeneous of degree $i$ with $\operatorname{deg} \alpha=\operatorname{deg} \beta=1$ and $\operatorname{deg} g_{2}=2$, $\operatorname{deg} g_{3}=3$. ${ }^{(5)}$ Moreover, $C_{i}$ is a non-zero rational number.

There are still two remaining terms in (3.9). That is,

$$
\begin{equation*}
L_{4} q=F_{1}\left(\alpha, \beta, g_{2}, g_{3}\right) x+F_{0}\left(\alpha, \beta, g_{2}, g_{3}\right) \tag{3.11}
\end{equation*}
$$

The basic structure on the consistency equations is described by the following two lemmas:

Lemma 3.7. - We have

$$
\begin{aligned}
& F_{1}\left(\alpha, \beta, g_{2}, g_{3}\right)=\alpha^{2} G_{1}\left(\alpha, \beta, g_{2}, g_{3}\right)=\alpha^{2}\left((-1)^{n-1} s_{n-1}\left(\alpha^{2}, \beta, g_{2}, g_{3}\right)+\cdots\right), \\
& F_{0}\left(\alpha, \beta, g_{2}, g_{3}\right)=\alpha^{2} G_{0}\left(\alpha, \beta, g_{2}, g_{3}\right)=\alpha^{2}\left((-1)^{n} s_{n}\left(\alpha^{2}, \beta, g_{2}, g_{3}\right)+\cdots\right) .
\end{aligned}
$$

For the remaining terms, each term of them has either $g_{2}$ or $g_{3}$ as a factor, hence it has lower degree in $\alpha, \beta$.

Proof. - Equation (3.11) gives

$$
\begin{aligned}
& F_{1}\left(\alpha, \beta, g_{2}, g_{3}\right)=(-1)^{n-1} \alpha^{2} s_{n-1}+\text { terms in } s_{1}, \ldots, s_{n-2} \\
& F_{0}\left(\alpha, \beta, g_{2}, g_{3}\right)=(-1)^{n} \alpha^{2} s_{n}+\text { terms in } s_{1}, \ldots, s_{n-1}
\end{aligned}
$$

If $\alpha=0$ then for any $\beta \in \mathbb{C}$ there is an $a \in Y_{n}$ with $B_{a}=\beta / 2$ and a polynomial solution $q_{a}(x)=w_{a} w_{-a}=\prod_{i=1}^{n}\left(x-\wp\left(a_{i}\right)\right)$ to the symmetric product of the Lamé equation, hence a polynomial solution to $L_{4}(q)=0$.

Thus $F_{1}\left(0, \beta, g_{2}, g_{3}\right)=0=F_{0}\left(0, \beta, g_{2}, g_{3}\right)$. Since $F_{i}$ depends on $\alpha^{2}$, we have $F_{i}=$ $\alpha^{2} G_{i}, i=0,1$, for some homogeneous polynomials $G_{0}, G_{1}$ in $\alpha^{2}, \beta, g_{2}, g_{3}$ of degree $n$ and $n-1$ respectively, and $G_{i}$ 's can be written as

$$
\begin{aligned}
& G_{1}=(-1)^{n-1} s_{n-1}+\cdots, \\
& G_{0}=(-1)^{n} s_{n}+\cdots .
\end{aligned}
$$

[^3]To see the dependence of the remaining terms on $g_{2}$ and $g_{3}$, we let $g_{2}=0=g_{3}$, and then $L_{4}(q) \equiv \alpha^{2}\left((-1)^{n-1} s_{n-1} x+(-1)^{n} s_{n}\right)\left(\bmod x^{2}\right)$ because both $p(x)=4 x^{3}$ and $\dot{p}(x)=12 x^{2}$ vanish modulo $x^{2}$. Thus we have $F_{1}(\alpha, \beta)=(-1)^{n-1} \alpha^{2} s_{n-1}$ and $F_{0}(\alpha, \beta)=(-1)^{n} \alpha^{2} s_{n}$ whenever $g_{2}=0=g_{3}$. This proves the lemma.

Lemma 3.8. - The polynomials $G_{1}$ and $G_{0}$ have no common factors for any $g_{2}, g_{3}$.
Proof. - We consider first the special case $g_{2}=g_{3}=0$. Then (3.9) becomes

$$
\begin{align*}
& 16 x^{6} \dddot{q}+144 x^{5} \dddot{q}+\left(108 x^{4}-8 x^{3}\left(2\left(n^{2}+n-12\right) x+\beta\right)\right) \ddot{q} \\
&-\left(12 x^{2}\left(2\left(n^{2}+n-3\right) x+\beta\right)+24 x^{3}\left(n^{2}+n-2\right)\right) \dot{q}  \tag{3.12}\\
&+\left(\alpha^{2}-12 n(n+1) x^{2}\right) q \equiv 0 \quad(\bmod \mathbb{C} \oplus \mathbb{C} x) .
\end{align*}
$$

The coefficients of $x^{n-k}, k=0, \ldots, n-2$, lead to recursive equations

$$
\begin{equation*}
(-1)^{k}\left(m_{k} s_{k+2}+n_{k} \beta s_{k+1}+\alpha^{2} s_{k}\right)=0, \tag{3.13}
\end{equation*}
$$

where the constants $m_{k}$ and $n_{k}$ are given by

$$
\begin{aligned}
m_{k}= & 16(n-(k+2))(n-(k+3))(n-(k+4))(n-(k+5)) \\
& +144(n-(k+2))(n-(k+3))(n-(k+4)) \\
& +\left(108-16\left(n^{2}+n-12\right)\right)(n-(k+2))(n-(k+3)) \\
& \quad-24\left(2 n^{2}+2 n-5\right)(n-(k+2))-12 n(n+1) \\
= & -4(k+2)(2 n-(k+1))(2 n-(2 k+1))(2 n-(2 k+3)), \\
n_{k}= & (8(n-(k+1))(n-(k+2))+12(n-(k+1))) \\
= & 4(n-(k-1))(n-(k+1)) .
\end{aligned}
$$

Since $k \leqslant n-2$, we have $m_{k} \neq 0$ and $n_{k} \neq 0$.
Let $\gamma$ be a non-trivial common factor of $G_{1}$ and $G_{0}$. Under the assumption $g_{2}=$ $g_{3}=0$ we have $G_{1}=(-1)^{n-1} s_{n-1}$ and $G_{0}=(-1)^{n} s_{n}$. Then $\gamma$ and $\alpha$ are co-prime, because if $\alpha=0$ then $s_{n-1}(0, \beta)=c_{n-1} \beta^{n-1}$ and $s_{n}(0, \beta)=c_{n} \beta^{n}$ for some nonzero constants $c_{n-1}$ and $c_{n}$. By the recursive equation (3.13) for $k=n-2$, we have $\gamma \mid s_{n-2}\left(\alpha^{2}, \beta, 0,0\right)$ too. By induction on $k$ for $k=n-3, \ldots, 0$ in decreasing order we conclude that $\gamma \mid s_{0}=1$, which leads to a contradiction.

For $g_{2}, g_{3} \in \mathbb{C}$, we see by Lemma 3.7 that the leading terms of $G_{1}, G_{0}$, as polynomials in $\alpha$ and $\beta$, are $(-1)^{n-1} s_{n-1}\left(\alpha^{2}, \beta, 0,0\right)$ and $(-1)^{n} s_{n}\left(\alpha^{2}, \beta, 0,0\right)$ respectively. Since $s_{n-1}\left(\alpha^{2}, \beta, 0,0\right)$ and $s_{n}\left(\alpha^{2}, \beta, 0,0\right)$ are co-prime, as we have just seen, we conclude that $G_{1}\left(\alpha, \beta, g_{2}, g_{3}\right)$ and $G_{0}\left(\alpha, \beta, g_{2}, g_{3}\right)$ are also co-prime. The proof is complete.

Proposition 3.9. - The common zeros of $G_{1}=0$ and $G_{0}=0$ consist of pairs of branch points $(a, b)$ corresponding to Lame functions of the same type. If $\bar{X}_{n}$ is nonsingular, there are exactly $n(n-1)$ such ordered pairs $(a, b)$ 's.

Proof. - It suffices to prove the (generic) case that $\bar{X}_{n}$ is non-singular, namely the case that all the Lamé functions are distinct. The general case follows from the nonsingular case by a limiting argument.

For any two Lamé functions $w_{a}, w_{b}$ of the same type (cf. Section 2.1), it is easy to see that we may arrange the representatives of $a$ and $b$ so that equations (3.2) holds. It follows that $q:=q_{a,-b}=q_{-a, b}$ (cf. (3.4)) is an even elliptic function solution to equation (3.3), or equivalently $q(x)$ is a polynomial solution to $L_{4} q(x)=0$ in variable $x=\wp(z)$.

From the above discussion, $(\alpha, \beta)$ must be a common root of $G_{1}$ and $G_{0}$ (where $\alpha=B_{a}-B_{b}, \beta=B_{a}+B_{b}$ ). By Lemma 3.6 and 3.7, we have $\operatorname{deg} G_{1}=n-1$ and $\operatorname{deg} G_{0}=n$ and $G_{1}, G_{0}$ are co-prime to each other by Lemma 3.8. Hence by Bézout theorem there are at most $n(n-1)$ common zeros.

On the other hand, the number of such ordered pairs can be determined by Proposition 2.2. Indeed, if $n=2 k$ then we have

$$
(k+1) k+3 k(k-1)=4 k^{2}-2 k=n(n-1)
$$

such pairs. If $n=2 k+1$, the number of pairs is given by

$$
k(k-1)+3(k+1) k=4 k^{2}+2 k=n(n-1) .
$$

Hence in all cases the number of ordered pairs coming from the Lamé functions of the same type agrees with the Bézout degree of the polynomial system defined by $G_{1}=0=G_{0}$. Thus these $n(n-1)$ pairs form the zero locus as expected (and there is no contribution from infinity).

The above discussions from Lemma 3.5 to Proposition 3.9 constitute a complete proof of Theorem 3.3. Here is a summary: we already know that $Q$ in (3.5) is an even elliptic function with singularity only at $0 \in E$. Thus

$$
Q(x)=c \prod_{i=1}^{n}\left(\wp(z)-\wp\left(c_{i}\right)\right)=: c \prod_{i=1}^{n}\left(x-x_{i}\right)
$$

is a polynomial solution to the $\operatorname{ODE}(3.6)$ with $\alpha=B_{a}-B_{b}, \beta=B_{a}+B_{b}$.
Since $\alpha=B_{a}-B_{b} \neq 0$, by Lemma $3.7(\alpha, \beta)$ must be a common root of $G_{1}(\alpha, \beta)=$ $0=G_{0}(\alpha, \beta)$. Then Proposition 3.9 says that $(\alpha, \beta)$ is pair of Lamé functions of the same type. This proves Theorem 3.3.

For future reference, we summarize the results into the following statement on a fourth order ODE which arises from the tensor product of two different Lamé equations with the same parameter $n \in \mathbb{N}$.

Theorem 3.10. - Let $I(z)=n(n+1)_{\wp}(z)$. The fourth order ODE

$$
\begin{equation*}
q^{\prime \prime \prime \prime}(z)-2(I+\beta) q^{\prime \prime}(z)-6 I^{\prime} q^{\prime}(z)+\left(\alpha^{2}-2 I^{\prime \prime}\right) q(z)=0 \tag{3.14}
\end{equation*}
$$

with $\alpha \neq 0$ has an elliptic function solution $q$ if and only if $(\alpha, \beta)$ is common zero to $G_{0}(\alpha, \beta)=0$ and $G_{1}(\alpha, \beta)=0$. Moreover, this solution $q$ must be even.

Example 3.11. - For $n=2, \beta=B_{a}+B_{b}, \alpha=B_{a}-B_{b}$, we have

$$
s_{1}=\frac{1}{6} \beta, \quad s_{2}=\frac{1}{36} \beta^{2}+\frac{1}{72} \alpha^{2}-\frac{1}{4} g_{2} .
$$

The first compatibility equation from $x^{1}$ is

$$
s_{1}\left(\alpha^{2}+36 g_{2}\right)-6 \beta g_{2}=0
$$

After substituting $s_{1}$ we get

$$
\begin{equation*}
\frac{1}{6} \alpha^{2} \beta=0 . \tag{3.15}
\end{equation*}
$$

The second compatibility equation from $x^{0}$ is

$$
s_{2}\left(\alpha^{2}+6 g_{2}\right)-s_{1}\left(\beta g_{2}+24 g_{3}\right)+4 \beta g_{3}+\frac{3}{2} g_{2}^{2}=0 .
$$

By substituting $s_{1}, s_{2}$ and noticing the (expected) cancellations we get

$$
\begin{equation*}
\alpha^{2}\left(\frac{1}{36} \beta^{2}+\frac{1}{72} \alpha^{2}-\frac{1}{6} g_{2}\right)=0 . \tag{3.16}
\end{equation*}
$$

If $B_{a} \neq B_{b}$ then (3.15) implies that $B_{b}=-B_{a}$ and then (3.16) leads to

$$
B_{a}^{2}=3 g_{2} \Longrightarrow \wp\left(a_{1}\right)+\wp\left(a_{2}\right)= \pm \sqrt{g_{2} / 3} .
$$

By Example $2.3(2)$, such $a \in \bar{X}_{2}$ lies in the branch loci of the Lamé curve. In particular, $a, b \in \sigma_{2}^{-1}(0)$. Denote by $\wp\left( \pm q_{ \pm}\right)= \pm \sqrt{g_{2} / 12}$. Then $a:=\left\{q_{+},-q_{+}\right\} \neq b:=$ $\left\{q_{-},-q_{-}\right\}$unless $g_{2}=0$. When $g_{2} \neq 0, \mathbf{z}_{2}$ fails to distinguish the two points $a$ and $b$. When $g_{2}=0$ (equivalently $\tau=e^{\pi i / 3}$ ), $a=b$ becomes a singular branch point for $\sigma_{2}: \bar{X}_{2} \rightarrow E_{\tau}$.

Example 3.12. - For $n=3, \beta=B_{a}+B_{b}, \alpha=B_{a}-B_{b}$. Then

$$
\begin{aligned}
& s_{1}=\frac{1}{10} \beta \\
& s_{2}=\frac{1}{600}\left(4 \beta^{2}+\alpha^{2}-150 g_{2}\right), \\
& s_{3}=\frac{1}{3600}\left(2 \beta^{3}+3 \alpha^{2} \beta-120 \beta g_{2}+900 g_{3}\right) .
\end{aligned}
$$

The two compatibility equations from $x^{1}$ and $x^{0}$ are

$$
\begin{aligned}
& 0=\frac{1}{600} \alpha^{2}\left(4 \beta^{2}+\alpha^{2}+60 g_{2}\right) \\
& 0=\frac{1}{3600} \alpha^{2}\left(2 \beta^{3}+3 \alpha^{2} \beta-90 \beta g_{2}+540 g_{3}\right) .
\end{aligned}
$$

If $\alpha \neq 0$ then $\alpha^{2}=-4 \beta^{2}-60 g_{2}$ and the second equation becomes

$$
\beta^{3}+27 g_{2} \beta-54 g_{3}=0 .
$$

It is clear that there are only finite solutions $\left(B_{a}, B_{b}\right)$ 's to this, though it may not be so straightforward to see that these 6 solution pairs (for generic tori) come from the branch loci as proved in Proposition 3.9.

### 3.3. Second proof of Theorem 3.3

Proof. - Following the definition of $q_{a,-b}(z)$ in (3.4), we now consider the odd elliptic solution to equation (3.3) ( $=$ equation (3.14)) instead:

$$
q(z):=\frac{1}{2}\left(q_{a,-b}(z)-q_{-a, b}(z)\right) .
$$

The function $q(z)$ has a pole of order $3+2 l$ at $0 \in E$ with $l \leqslant n-2$. Thus $q(z) / \wp^{\prime}(z)$ is an even elliptic function with the only pole at 0 since $q\left(\frac{1}{2} \omega_{i}\right)=0$ for $1 \leqslant i \leqslant 3$. If $q(z)$ does not vanish completely, then

$$
q(z)=c \wp^{\prime}(z) \prod_{i=1}^{l}\left(\wp(z)-\wp\left(c_{i}\right)\right)=: c \wp^{\prime}(z) f(\wp(z)),
$$

where $f(x)=\prod_{i=1}^{l}\left(x-\wp\left(c_{i}\right)\right)=x^{l}-s_{1} x^{l-1}+\cdots+(-1)^{l} s_{l}$.
Equation (3.14) now reads as

$$
\begin{align*}
& q^{\prime \prime \prime \prime}(z)-2(\beta+2 n(n+1) \wp(z)) q^{\prime \prime}(z)  \tag{3.17}\\
& \\
& \quad-6 n(n+1) \wp^{\prime}(z) q^{\prime}(z)+\left(\alpha^{2}-2 n(n+1) \wp^{\prime \prime}(z)\right) q(z)=0 .
\end{align*}
$$

By straightforward calculations, we get derivatives of $q$ in terms of derivatives of $\wp(z)$ and $f^{\prime}(x)$. For example,

$$
\begin{aligned}
q^{\prime}(z) & =\wp^{\prime \prime}(z) f(x)+\wp^{\prime}(z)^{2} f^{\prime}(x), \\
q^{\prime \prime}(z) & =\wp^{\prime \prime \prime}(z) f(x)+3 \wp^{\prime \prime}(z) \wp^{\prime}(z) f^{\prime}(x)+\wp^{\prime}(z)^{3} f^{\prime \prime}(x), \quad \text { etc. }
\end{aligned}
$$

Then (3.17) is equivalent to

$$
\begin{aligned}
& \quad f(x)\left((360-96 n(n+1)) x^{2}-24 \beta x+(4 n(n+1)-18) g_{2}+\alpha^{2}\right) \\
& +f^{\prime}(x)\left((1320-96 n(n+1)) x^{3}-36 \beta x^{2}\right. \\
& \left.\quad+(12 n(n+1)-150) g_{2} x+(6 n(n+1)-60) g_{3}+3 \beta g_{2}\right) \\
& +f^{\prime \prime}(x)\left((1020-16 n(n+1)) x^{4}-8 \beta x^{3}+(4 n(n+1)-210) g_{2} x^{2}\right. \\
& \left.\quad+\left(2 \beta g_{2}+(4 n(n+1)-120) g_{3}\right) x+2 \beta g_{3}+\frac{15}{4} g_{2}^{2}\right) \\
& + \\
& +f^{\prime \prime \prime}(x)\left(60 x^{2}-30 g_{2}\right)\left(4 x^{3}-g_{2} x-g_{3}\right) \\
& + \\
& f^{\prime \prime \prime \prime}(x)\left(4 x^{3}-g_{2} x-g_{3}\right)^{2}=0 .
\end{aligned}
$$

By comparing the coefficients of $x^{l+2}$, we obtain

$$
\begin{aligned}
(360-96 n(n+1))+l(1320- & 96 n(n+1))+l(l-1)(1020-16 n(n+1)) \\
& +240 l(l-1)(l-2)+16 l(l-1)(l-2)(l-3)=0 .
\end{aligned}
$$

After simplification, this is reduced to

$$
4 n(n+1)=(2 l+3)(2 l+5)
$$

which obviously leads to a contradiction since the number in the right-hand side is odd while the number in the left-hand side is even. Therefore we must have $q \equiv 0$ from the beginning. That is, $\left\{a_{i},-b_{i}\right\}=\left\{-a_{i}, b_{i}\right\}$.

If one of $a, b$ does not correspond to a Lamé function, say $a \in X_{n}$, then $\left\{a_{i}\right\} \cap\left\{-a_{i}\right\}=\varnothing$ by Section 0.1.2 (iii) and we conclude that $\left\{a_{i}\right\}=\left\{b_{i}\right\}$. Otherwise $a$ and $b$ correspond to Lamé functions of the same type.

### 3.4. The degree of the rational function $\mathbf{z}_{n}$

Theorem 3.13. - The structure of the map $\mathbf{z}_{n}: \bar{X}_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})$ over $\infty \in \mathbb{P}^{1}(\mathbb{C})$ is analytically equivalent to $\sigma_{n}: \bar{X}_{n} \rightarrow E$ over 0 . In particular it has the same degree as the one for $\sigma_{n}$, namely $\operatorname{deg} \mathbf{z}_{n}=\operatorname{deg} \sigma_{n}=\frac{1}{2} n(n+1)$.

Proof. - By definition, $\mathbf{z}_{n}^{-1}(\infty)=\sigma_{n}^{-1}(0)$ as sets. So the crucial point is to compare the ramification structures of $\bar{X}_{n} \rightarrow E$ at $0 \in E$ and $\bar{X}_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})$ at $\infty \in \mathbb{P}^{1}(\mathbb{C})$. Let $a \in \bar{X}_{n}$ with $\sigma_{n}(a)=0$. Then for $b=\left\{b_{i}\right\}_{i=1}^{n} \in \bar{X}_{n}$ in a small analytic neighborhood of $a$ we have $b_{i} \neq 0$ all $i$.

If $a \neq 0^{n}$, then every $b_{i}$ is away from 0 and

$$
\left(\mathbf{z}_{n}(b)\right)^{-1}=\left(\zeta\left(\sigma_{n}(b)\right)-\sum_{i=1}^{n} \zeta\left(b_{i}\right)\right)^{-1}=\sigma_{n}(b)+o\left(\sigma_{n}(b)\right)
$$

In terms of the coordinate of $\mathbb{P}^{1}(\mathbb{C})$ at $\infty$, the map $\mathbf{z}_{n}$ near $a \neq 0^{n}$ is seen to be analytically equivalent to $\sigma_{n}$.

At $a=0^{n}$, we compute the expansion of $\left(\mathbf{z}_{n}(b)\right)^{-1}$ as

$$
\frac{\sigma_{n}(b)}{\sigma_{n}(b) \zeta\left(\sigma_{n}(b)\right)-\sigma_{n}(b) \sum_{i=1}^{n} \zeta\left(b_{i}\right)}=\sigma_{n}(b)\left(1+\sigma_{n}(b) \sum_{i=1}^{n} \zeta\left(b_{i}\right)+O\left(\sigma_{n}^{4}(b)\right)\right)
$$

The tangent direction $\left(t_{i}\right)$ at $0^{n}$ is related to $\left(b_{i}\right)$ through the asymptotic

$$
t_{i}|B|^{1 / 2} \sim-1 / b_{i}
$$

(cf. [1, Prop. 7.5] and the proof therein). Hence

$$
\lim _{b \rightarrow 0} \sigma_{n}(b) \sum_{i=1}^{n} \zeta\left(b_{i}\right)=\sum_{i=1}^{n} t_{i}^{-1} \sum_{i=1}^{n} t_{i}=: \Lambda_{n}
$$

The precise value of $\Lambda_{n}$ follows from Proposition 1.5:

$$
\Lambda_{n}=\frac{e_{n-1} e_{1}}{e_{n}}=\frac{1}{2} n(n+1) \neq-1 .
$$

Hence $\left(\mathbf{z}_{n}(b)\right)^{-1}=\left(1+\Lambda_{n}\right) \sigma_{n}(b)+o\left(\sigma_{n}(b)\right)$ and we again have the analytic equivalence (up to a constant multiple).

In particular, $\operatorname{deg} \mathbf{z}_{n}=\operatorname{deg} \sigma_{n}=\frac{1}{2} n(n+1)$ by Theorem 0.2.

## 4. Pre-modular forms $Z_{n}(\sigma ; \tau)$

Pre-modular forms are defined in Definition 0.1. Since the Hecke function is premodular of weight one, Theorem 3.2 then implies

Corollary 4.1. - $Z_{n}(\sigma ; \tau):=W_{n}(Z)(\sigma ; \tau)$ is pre-modular of weight $\frac{1}{2} n(n+1)$, with $Z, \wp(\sigma), \wp^{\prime}(\sigma), g_{2}, g_{3}$ being of weight $1,2,3,4,6$ respectively.
4.1. Completion of the proof of Theorem 0.4. - We call the $2 n+1$ branch points $a \in Y_{n} \backslash X_{n}$ trivial critical points since $a=-a$ and the Green equation (0.5) holds trivially. They satisfy a nice compatibility condition with the case $n=1$ under the addition map:

Lemma 4.2. - Let $a=\left\{a_{1}, \ldots, a_{n}\right\} \in Y_{n}$ be a solution to the Green equation $\sum_{i=1}^{n} \nabla G\left(a_{i}\right)=0$. Then $a$ is trivial, i.e., $a=-a$, if and only if $\sigma_{n}(a) \in E[2]$.

Proof. - If $a$ is trivial, then $\sigma_{n}(a) \in E[2]$ clearly. If $a$ is non-trivial, i.e., $a \in X_{n}$, by equations (1.5), it gives rise to a type II developing map $f$ with

$$
f\left(z+\omega_{1}\right)=e^{-4 \pi i \sum_{i} s_{i}} f(z), \quad f\left(z+\omega_{2}\right)=e^{4 \pi i \sum_{i} r_{i}} f(z)
$$

Here, $a_{i}=r_{i} \omega_{1}+s_{i} \omega_{2}$ for $i=1, \ldots, n$.
If $\sigma_{n}(a) \in E[2]$, then both exponential factors reduce to one and we conclude that $f(z)$ is an elliptic function on $E$. Notice that the only zero of $f^{\prime}(z)$ is at $z=0$ which has order $2 n$, and the only poles of $f^{\prime}(z)$ are at $-a_{i}$ of order $2, i=1, \ldots, n$. This forces that $\sigma_{n}(a)=0$ and

$$
f^{\prime}(z)=\sum_{j=1}^{n} E_{j} \wp\left(z+a_{j}\right)+C_{1}
$$

for some constants $E_{1}, \ldots, E_{n}$ and $C_{1}$, since $f^{\prime}$ is residue free. Then

$$
f(z)=-\sum_{j=1}^{n} E_{j} \zeta\left(z+a_{i}\right)+C_{1} z+C_{2}
$$

for some constant $C_{2}$. But $f(z)$ is elliptic, which implies that $C_{1}=0$ and $\sum_{j=1}^{n} E_{j}=0$. Now $f^{2 k-1}(0)=0$ for $k=1, \ldots, n$ leads to a system of linear equations in $E_{j}$ 's (cf. [1, Lem. 2.3.1]):

$$
\sum_{j=1}^{n} \wp^{k}\left(a_{j}\right) E_{j}=0, \quad k=1, \ldots, n
$$

Then $\wp\left(a_{i}\right) \neq \wp\left(a_{j}\right)$ for $i \neq j$ forces that $E_{j}=0$ for all $j$. This is a contradiction and we conclude that $\sigma_{n}(a) \notin E[2]$.

The following theorem completes the proof of Theorem 0.4:
Theorem 4.3 (Extra critical points versus non-trivial zeros of $Z_{n}(\sigma ; \tau)$ )
(i) Given $\sigma_{0} \in E_{\tau} \backslash E_{\tau}[2]$ with $Z_{n}\left(\sigma_{0} ; \tau\right)=0$, there is a unique $a \in X_{n}$ such that $\sigma_{n}(a)=\sigma_{0}$ and $\mathbf{z}_{n}(a)=Z\left(\sigma_{0}\right)$.
(ii) Conversely, if $a \in X_{n}$ and $\mathbf{z}_{n}(a)=Z(\sigma(a))$ then $Z_{n}(\sigma(a) ; \tau)=0$ and $\sigma_{n}(a) \notin$ $E_{\tau}[2]$.

## Proof

(i) For any given $\sigma_{0}$, by substituting $\sigma$ with $\sigma_{0}$ in $W_{n}(\mathbf{z})$, we get a polynomial $W_{n, \sigma_{0}}(\mathbf{z})$ of degree $\frac{1}{2} n(n+1)$. Since $W_{n}(\mathbf{z})$ is the minimal polynomial of the rational function $\mathbf{z}_{n} \in K\left(\bar{X}_{n}\right)$ over $K(E)$, those $\mathbf{z}_{n}(a)$ with $a \in \bar{X}_{n}$ and $\sigma_{n}(a)=\sigma_{0}$ give rise to all the roots, counted with multiplicities, of $W_{n, \sigma_{0}}(a)=0$.

Now $Z\left(\sigma_{0}\right)$ is a root of $W_{n, \sigma_{0}}(\mathbf{z})=0$ with $\sigma_{0} \notin E[2]$, hence there is a point $a \in X_{n}$ which corresponds to it. That is, $Z\left(\sigma_{0}\right)=\mathbf{z}_{n}(a)$ with $\sigma_{n}(a)=\sigma_{0}$, and $a$ is unique by

Theorem 3.3. Notice that if $a \in \bar{X}_{n} \backslash X_{n}$ then $a=-a$ and then $\sigma_{n}(a) \in E[2]$. So we must have $a \in X_{n}$.
(ii) It is clear that $Z_{n}(\sigma(a))=W_{n}\left(Z(\sigma(a))=W_{n}\left(\mathbf{z}_{n}(a)\right)=0\right.$. Since $a \in X_{n}$, by equation (3.1) we have $\sum_{i=1}^{n} \nabla G\left(a_{i}\right)=0$. But since $a$ is non-trivial, Lemma 4.2 implies that $\sigma_{n}(a) \notin E[2]$.
4.2. Monodromy aspects. - We present below an extended version of Theorem 0.4 in terms of monodromy groups of Lamé equations. The original case of mean field equations corresponds to the case with unitary monodromy.

Let $a=\left\{a_{1}, \ldots, a_{n}\right\} \in X_{n}, B_{a}=(2 n-1) \sum_{i=1}^{n} \wp\left(a_{i}\right)$ and $w_{a}, w_{-a}$ be the independent ansatz solutions (0.7) to $w^{\prime \prime}=\left(n(n+1) \wp(z)+B_{a}\right) w$. From equations (1.4), one calculates that the monodromy matrices are given by

$$
\begin{align*}
\binom{w_{a}}{w_{-a}}\left(z+\omega_{1}\right) & =\left(\begin{array}{cc}
e^{-2 \pi i s} & 0 \\
0 & e^{2 \pi i s}
\end{array}\right)\binom{w_{a}}{w_{-a}}(z), \\
\binom{w_{a}}{w_{-a}}\left(z+\omega_{2}\right) & =\left(\begin{array}{cc}
e^{2 \pi i r} & 0 \\
0 & e^{-2 \pi i r}
\end{array}\right)\binom{w_{a}}{w_{-a}}(z), \tag{4.1}
\end{align*}
$$

where the two numbers $r, s \in \mathbb{C}(\bmod \mathbb{Z})$ are uniquely determined by

$$
\begin{equation*}
r \omega_{1}+s \omega_{2}=\sigma(a)=\sum_{i=1}^{n} a_{i}, \quad r \eta_{1}+s \eta_{2}=\sum_{i=1}^{n} \zeta\left(a_{i}\right) \tag{4.2}
\end{equation*}
$$

The system is non-singular by the Legendre relation $\omega_{1} \eta_{2}-\omega_{2} \eta_{1}=-2 \pi i$.
The next lemma extends Lemma 4.2:
Lemma 4.4. - Let $a \in X_{n}$ with ( $r, s$ ) given by (4.2). Then ( $\left.r, s\right) \notin \frac{1}{2} \mathbb{Z}^{2}$.
Proof. - If $(r, s) \in \frac{1}{2} \mathbb{Z}^{2}$ then $f:=w_{a} / w_{-a}$ is elliptic by equations (4.1). Since

$$
f^{\prime}=\frac{w_{a}^{\prime} w_{-a}-w_{a} w_{-a}^{\prime}}{w_{-a}^{2}}=\frac{C}{w_{-a}^{2}}
$$

we find that $z=0$ is the only zero of $f^{\prime}(z)$ in $E$, which has order $2 n$. The proof of Lemma 4.2 for this $f$ goes through and leads to a contradiction.

Now we consider $Z_{r, s}(\tau)$ in (0.10) but with $r, s, \in \mathbb{C}$, and define

$$
\begin{equation*}
Z_{n ; r, s}(\tau):=W_{n}\left(Z_{r, s}\right)(r+s \tau ; \tau), \quad r, s \in \mathbb{C} \tag{4.3}
\end{equation*}
$$

It reduces to $Z_{n}(\sigma ; \tau)$ for $\sigma=r+s \tau$ when $r, s \in \mathbb{R}$ (see [2] for its role in the isomonodromy problems and Painlevé VI equations).

By substituting $Z_{n}(\sigma ; \tau)$ with $Z_{n ;, r, s}(\tau)$ and using Lemma 4.4 in place of Lemma 4.2, the proof of Theorem 4.3 also leads to:

Theorem 4.5. - Let $r, s \in \mathbb{C}$. Then any non-trivial solution $\tau$ to $Z_{n ; r, s}(\tau)=0$, i.e., with $r+s \tau\left(\bmod \Lambda_{\tau}\right) \notin E_{\tau}[2]$, corresponds to an $a=\left(a_{1}, \ldots, a_{n}\right) \in \mathbb{C}^{n}$ such
that $a\left(\bmod \Lambda_{\tau}\right) \in X_{n}(\tau)$ and

$$
\sum_{i=1}^{n} a_{i}=r+s \tau, \quad \sum_{i=1}^{n} \zeta\left(a_{i} ; \tau\right)=r \eta_{1}(\tau)+s \eta_{2}(\tau)
$$

Equivalently, by equations (4.2), the Lame equation $w^{\prime \prime}=\left(n(n+1) \wp\left(z ; \Lambda_{\tau}\right)+B_{a}\right) w$ has its monodromy representation given by equations (4.1).

We leave the straightforward justifications to the interested readers.

## 5. An explicit determination of $Z_{n}$

From the equations of $\bar{X}_{n} \subset \operatorname{Sym}^{n} E$ (cf. (1.2)) and the recursively defined algebraic formula for the addition map $E^{n} \rightarrow E$, in principle it is possible to compute $W_{n}$ and hence $Z_{n}$ by elimination theory (cf. [7]). However we shall present a more direct approach on this to reveal more structures inside it.
5.1. Comparisons with the Hermite-Krichever ansatz. - Besides the HermiteHalphen ansatz (0.7), there is another ansatz, the Hermite-Krichever ansatz, which can also be used to construct solutions to the integral Lamé equation (0.6). It takes the form

$$
\begin{equation*}
\psi(z):=\left(U(\wp(z))+V(\wp(z)) \frac{\wp^{\prime}(z)+\wp^{\prime}\left(a_{0}\right)}{\wp(z)-\wp\left(a_{0}\right)}\right) \frac{\sigma\left(z-a_{0}\right)}{\sigma(z)} e^{\left(\zeta\left(a_{0}\right)+\kappa\right) z}, \tag{5.1}
\end{equation*}
$$

where $U(x)$ and $V(x)$ are polynomials in $x$, the point $a_{0} \in \mathbb{C} \backslash \Lambda$, and $\kappa \in \mathbb{C}$ is a constant. As usual, we set $(x, y)=\left(\wp(z), \wp^{\prime}(z)\right)$ and $\left(x_{0}, y_{0}\right)=\left(\wp\left(a_{0}\right), \wp^{\prime}\left(a_{0}\right)\right)$ to be the corresponding algebraic coordinates.

The ansatz (5.1) makes sense since $\psi$ only has poles at $z=0$. The two poles at $z= \pm a_{0}$ from $\left(\wp(z)-\wp\left(a_{0}\right)\right)^{-1}$ cancel with the zeros from $\sigma\left(z-a_{0}\right)$ and $\wp^{\prime}(z)+\wp^{\prime}\left(a_{0}\right)$. In order for $\operatorname{ord}_{z=0} \psi(z)=-n$, we must have

Lemma 5.1 (Degree constraints)
(i) If $n=2 m$ with $m \in \mathbb{N}$ then $\operatorname{deg} U \leqslant m-1$ and $\operatorname{deg} V=m-1$.
(ii) If $n=2 m+1$ with $m \in \mathbb{N} \cup\{0\}$ then $\operatorname{deg} U=m$ and $\operatorname{deg} V \leqslant m-1$.

By an obvious normalization, in case (i) we may assume that

$$
U(x)=\sum_{i=0}^{m-1} u_{i} x^{i}, \quad V(x)=\sum_{i=0}^{m-1} v_{i} x^{i} \quad \text { with } v_{m-1}=1
$$

and in case (ii) we assume that $U(x)=\sum_{i=0}^{m} u_{i} x^{i}$ with $u_{m}=1$ and $V(x)=$ $\sum_{i=0}^{m-1} v_{i} x^{i}$. In both cases, the requirement that $\psi(z)$ satisfies equation (0.6) leads to recursions on $u_{i}$ 's and $v_{i}$ 's. It turns out to be convenient to work with coordinates $\left(B, \kappa, x_{0}, y_{0}\right)$ to parametrize $u_{i}$ 's and $v_{i}$ 's, and this was carried out by Maier in [12, §4]. The following is a summary:

In case (i) the recursion determines $v_{i}\left(v_{m-1}=1\right)$ and then $u_{i}$ for $i=m-1$, $m-2, \ldots$ in descending order. In case (ii) it starts with $u_{m}=1$ and determines $v_{i}$ and then $u_{i}$ for $i=m-1, m-2, \ldots$ also in descending order.

There are two compatibility equations coming from

$$
u_{-1}\left(B, \kappa, x_{0}, y_{0}\right)=0 \quad \text { and } \quad v_{-1}\left(B, \kappa, x_{0}, y_{0}\right)=0
$$

The two parameters $x_{0}, y_{0}$ satisfy $y_{0}^{2}=4 x_{0}^{3}-g_{2} x_{0}-g_{3}$. Hence there are four variables $\left(B, \kappa, x_{0}, y_{0}\right) \in \mathbb{C}^{4}$ which are subject to three polynomial equations. By taking into account the limiting cases with $\left(x_{0}, y_{0}\right)=(\infty, \infty)$, this recovers the Lame curve $\bar{Y}_{n}$ (which is denoted by $\Gamma_{\ell}$ in [12] with $\ell=n$ ).

There are four natural coordinate projections (rational functions) $\bar{Y}_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})$, namely $B, \kappa, x_{0}$ and $y_{0}$ respectively. The first one $B: \bar{Y}_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})$ is simply the hyperelliptic structure map. The main result in [12] is an explicit description of the other three maps in terms of the coordinates $(B, C)$ on $\bar{Y}_{n}$. To state it we need to first recall some variants of Lamé polynomials.

Definition 5.2 ([12, Def. 3.2, 3.4, 3.6])
(1) The twisted Lamé polynomials $l t_{j}(B), j=0,1,2,3$ are monic polynomials whose zeros correspond to solutions to (0.6) given by the Hermite-Krichever ansatz with $\kappa \neq 0$ and $a_{0}=0, \frac{1}{2} \omega_{1}, \frac{1}{2} \omega_{2}, \frac{1}{2} \omega_{3}$ respectively, i.e., $\left(x_{0}, y_{0}\right)=(\infty, \infty),\left(e_{1}, 0\right)$, $\left(e_{2}, 0\right),\left(e_{3}, 0\right)$ respectively.
(2) The theta-twisted polynomial $l_{\theta}(B)$ is the monic polynomial whose roots correspond to the case $\kappa=0$ and $a_{0} \notin E[2]$. (For $\kappa=0$ and $a_{0} \in E[2]$ they correspond to the ordinary Lamé polynomials $l_{i}(B)$ 's.)

Theorem 5.3 ([12, Th. 4.1]). - For all $n \in \mathbb{N}$ and $i \in\{1,2,3\}$,

$$
\begin{align*}
x_{0}(B) & =e_{i}+\frac{4}{n^{2}(n+1)^{2}} \frac{l_{i}(B) l t_{i}(B)^{2}}{l_{0}(B) l t_{0}(B)^{2}} \\
y_{0}(B, C) & =\frac{16}{n^{3}(n+1)^{3}} \frac{C}{c_{n}} \frac{l t_{1}(B) l t_{2}(B) l t_{3}(B)}{l_{0}(B)^{2} l t_{0}(B)^{3}}  \tag{5.2}\\
\kappa(B, C) & =-\frac{(n-1)(n+2)}{n(n+1)} \frac{C}{c_{n}} \frac{l_{\theta}(B)}{l_{0}(B) l t_{0}(B)} .
\end{align*}
$$

The formula for $x_{0}(B)$ is valid for all three choices of $i$.
All the factors lie in $\mathbb{Q}\left[e_{1}, e_{2}, e_{3}, g_{2}, g_{3}, B\right]$ and are monic in $B$. They are homogeneous with degrees of $B, e_{i}, g_{2}, g_{3}$ being $1,1,2,3$ respectively.

As a simple consistency check, we have $C^{2}=\ell_{n}(B)$ by Proposition 2.2.
Remark 5.4. - In [12] $\nu=C / c_{n}$ is used instead.
The polynomials $l_{0}(B), l_{i}(B)(i=1,2,3)$ are written there as $L_{\ell}^{I}\left(B ; g_{2}, g_{3}\right)$, $L_{\ell}^{I I}\left(B ; e_{i}, g_{2}, g_{3}\right)$, called the Lamé spectral polynomials, where $\ell=n$.

The polynomials $l t_{0}(B), l t_{i}(B)(i=1,2,3)$ are written there as $L t_{\ell}^{I}\left(B ; g_{2}, g_{3}\right)$, $L t_{\ell}^{I I}\left(B ; e_{i}, g_{2}, g_{3}\right)$, called the twisted Lamé polynomials.

Also $l_{\theta}(B)$ is written there as $L \theta_{\ell}\left(B ; g_{2}, g_{3}\right)$.
The compatibility equations from the recursive formulas for these special cases give rise to explicit formulas for $l t_{j}(B)$ 's and $l_{\theta}(B)$ 's. Tables for $l t_{0}(B), l_{\theta}(B)$ up to $n=8$, and for $l t_{i}(B)$ up to $n=6$, are given in [12, Tables 5,6$]$.

Example 5.5. - We recall Maier's formulas for $l t_{j}(B)$ and $l_{\theta}(B)$ for $n \leqslant 4$.
(1) First of all, $l_{\theta}(B)=1$ for $n \leqslant 3$. For $n=4$,

$$
l_{\theta}(B)=B^{2}-\frac{193}{3} g_{2}
$$

Also for $n=1, l t_{j}(B)=1$ for all $j$.
(2) $n=2: l t_{0}(B)=1, l t_{i}(B)=B-6 e_{i}$ for $i=1,2,3$.
(3) $n=3: l t_{0}(B)=B^{2}-\frac{75}{4} g_{2}$, and for $i=1,2,3$,

$$
l t_{i}(B)=B^{2}-15 e_{i} B+\frac{75}{4} g_{2}-225 e_{i}^{2} .
$$

(4) $n=4: l t_{0}(B)=B^{3}-\frac{343}{4} g_{2} B-\frac{1715}{2} g_{3}$. For $i=1,2,3$,

$$
\begin{aligned}
l t_{i}(B)=B^{4}-55 e_{i} B^{3} & +\left(\frac{539}{4} g_{2}-945 e_{i}^{2}\right) B^{2} \\
& +\left(1960 e_{i} g_{2}+2450 g_{3}\right) B+61740 e_{i}^{2} g_{2}-68600 e_{i} g_{3}-9261 g_{2}^{2} .
\end{aligned}
$$

To apply Theorem 5.2, we need to compare the projection map

$$
\begin{equation*}
\pi_{n}: \bar{Y}_{n} \longrightarrow E, \quad a \longmapsto \pi_{n}(a):=a_{0} . \tag{5.3}
\end{equation*}
$$

with the addition map $\sigma_{n}: \bar{Y}_{n} \rightarrow E$. They turn out to be the same!
Theorem 5.6. $-\pi_{n}(a)=\sigma_{n}(a)$. Moreover, $\kappa(a)=-\mathbf{z}_{n}(a)$.
Proof. - During the proof we view $a_{i} \in \mathbb{C}$ instead of its image $\left[a_{i}\right] \in E$.
Let $a \in Y_{n}$. The two expressions (0.7) and (5.1), which correspond to the same solution to the Lamé equation (0.6), must be proportional to each other by a constant. Hence we get

$$
\kappa(a)=\sum_{i=1}^{n} \zeta\left(a_{i}\right)-\zeta\left(a_{0}\right) .
$$

Recall that $\mathbf{z}_{n}(a)=\zeta\left(\sigma_{n}(a)\right)-\sum_{i=1}^{n} \zeta\left(a_{i}\right)$. Then

$$
\begin{equation*}
\mathbf{z}_{n}(a)+\kappa(a)=\zeta\left(\sigma_{n}(a)\right)-\zeta\left(a_{0}\right) . \tag{5.4}
\end{equation*}
$$

As a well-defined meromorphic function on $\bar{Y}_{n}$, we conclude that

$$
a_{0}(a)=\sigma_{n}(a)+c
$$

for some constant $c \in \mathbb{C}$. Consider a point $a \in Y_{n} \backslash X_{n}$ with $\sigma_{n}(a)=\frac{1}{2} \omega_{1}$, i.e., $l_{1}\left(B_{a}\right)=0$. Such an $a$ exists by Proposition 2.2. Then $\mathbf{z}_{n}(a)=0$ trivially. We also have $\kappa(a)=0$ by Theorem 5.2 since

$$
C_{a}^{2}=c_{n}^{2} l_{0}\left(B_{a}\right) l_{1}\left(B_{a}\right) l_{2}\left(B_{a}\right) l_{3}\left(B_{a}\right)=0
$$

(again by Proposition 2.2). So equation (5.4) implies $0=\frac{1}{2} \eta_{1}-\zeta\left(\frac{1}{2} \omega_{1}+c\right.$ ), and hence $c=0$. This proves $\sigma_{n}(a)=a_{0}$, which represents $\pi_{n}(a)$ in $E$, and also $\kappa(a)=-\mathbf{z}_{n}(a)$. The proof is complete.
5.2. Effective construction and explicit formulas for $Z_{n}, n \leqslant 4$. - Now we describe an explicit construction, based on Theorem 5.3, of the polynomial $W_{n}(\mathbf{z})$ in Theorem 3.2. It is an application of elimination theory using resultants.

By Theorem 5.3 and Theorem 5.6, we may eliminate $C$ to get

$$
\begin{equation*}
\frac{y_{0}}{\mathbf{z}_{n}}=\frac{16}{n^{2}(n+1)^{2}(n-1)(n+2)} \frac{l t_{1}(B) l t_{2}(B) l t_{3}(B)}{l_{0}(B) l t_{0}(B)^{2} l_{\theta}(B)} \tag{5.5}
\end{equation*}
$$

which leads to a polynomial equation $g=0$ for

$$
\begin{equation*}
g:=\mathbf{z} \prod_{i=1}^{3} l t_{i}(B)-y_{0} \frac{n^{2}(n+1)^{2}(n-1)(n+2)}{16} l_{0}(B) l t_{0}(B)^{2} l_{\theta}(B) . \tag{5.6}
\end{equation*}
$$

On the other hand, the three rational expressions of $x_{0}$ lead to $f=0$ for

$$
\begin{align*}
f & :=\frac{1}{3} \sum_{i=1}^{3} l_{i}(B) l t_{i}(B)^{2}-x_{0} \frac{n^{2}(n+1)^{2}}{4} l_{0}(B) l t_{0}(B)^{2}  \tag{5.7}\\
& =l_{i}(B) l t_{i}(B)^{2}-\left(x_{0}-e_{i}\right) \frac{n^{2}(n+1)^{2}}{4} l_{0}(B) l t_{0}(B)^{2}, \quad i \in\{1,2,3\} .
\end{align*}
$$

Notice that $f, g$ are polynomials in $g_{2}, g_{3}$ (and $B, x_{0}, y_{0}$ ) instead of in $e_{i}$ 's.
Let $R(f, g ; B)$ be the resultant of the two polynomials $f$ and $g$ arising from the elimination of the variable $B$. Standard elimination theory (see e.g [7, Chap. 5]) implies that $R(f, g ; B)$ gives the equation defining the branched covering map $\sigma_{n}: \bar{Y}_{n} \rightarrow E$ outside the loci $C=0$ :

Proposition 5.7. - $R(f, g ; B)(\mathbf{z})=\lambda_{n} W_{n}(\mathbf{z}) \in \mathbb{Q}\left[g_{2}, g_{3}, x_{0}, y_{0}\right][\mathbf{z}]$, where $\lambda_{n}=$ $\lambda_{n}\left(g_{2}, g_{3}, x_{0}, y_{0}\right)$ is independent of $\mathbf{z}$.

In particular, the pre-modular form $Z_{n}(\sigma ; \tau)=W_{n}(Z)(\sigma ; \tau)$ can be explicitly computed for any $n \in \mathbb{N}$ by way of $R(f, g ; B)$.

In practice, such a computation is time consuming even using computer. In the following, we apply it to the initial cases up to $n=4$. As before we denote $x_{0}=$ $\wp(\sigma)=: \wp$ and $y_{0}=\wp^{\prime}(\sigma)=: \wp^{\prime}$.

Example 5.8. - For $n=2$, it is easy to see that

$$
\begin{aligned}
& f=B^{3}-9 \wp B^{2}+27\left(g_{2} \wp+g_{3}\right) \\
& g=\mathbf{z} B^{3}-9 \wp^{\prime} B^{2}-9 \mathbf{z} g_{2} B+27\left(g_{2} \wp^{\prime}-2 \mathbf{z} g_{3}\right)
\end{aligned}
$$

The resultant $R(f, g ; B)$ is calculated by the $6 \times 6$ Sylvester determinant:

$$
\left|\begin{array}{cccccc}
1 & -9 \wp & 0 & 27\left(g_{2} \wp+g_{3}\right) & 0 & 0 \\
0 & 1 & -9 \wp & 0 & 27\left(g_{2} \wp+g_{3}\right) & 0 \\
0 & 0 & 1 & -9 \wp & 0 & 27\left(g_{2} \wp+g_{3}\right) \\
\mathbf{z} & -9 \wp^{\prime} & -9 \mathbf{z} g_{2} & 27\left(g_{2} \wp^{\prime}-2 \mathbf{z} g_{3}\right) & 0 & 0 \\
0 & \mathbf{z} & -9 \wp^{\prime} & -9 \mathbf{z} g_{2} & 27\left(g_{2} \wp^{\prime}-2 \mathbf{z} g_{3}\right) & 0 \\
0 & 0 & \mathbf{z} & -9 \wp^{\prime} & -9 \mathbf{z} g_{2} & 27\left(g_{2} \wp^{\prime}-2 \mathbf{z} g_{3}\right)
\end{array}\right| .
$$

A direct evaluation gives

$$
R(f, g ; B)(\mathbf{z})=-3^{9} \Delta\left(\wp^{\prime}\right)^{2}\left(\mathbf{z}^{3}-3 \wp \mathbf{z}-\wp^{\prime}\right) .
$$

Here, $\Delta=g_{2}^{3}-27 g_{3}^{2}$ is the discriminant. This gives $W_{2}(\mathbf{z})=\mathbf{z}^{3}-3 \wp \mathbf{z}-\wp^{\prime}$ and $Z_{2}(\sigma ; \tau)=W_{2}(Z)=Z^{3}-3 \wp Z-\wp^{\prime}$.

Example 5.9. - For $n=3$, we have

$$
\begin{aligned}
f=16 & B^{6}-576 B^{5} \wp+360 B^{4} g_{2}+5400 B^{3}\left(5 g_{3}+4 g_{2} \wp\right) \\
& -3375 B^{2} g_{2}^{2}-84375 \Delta-101250 B g_{2}\left(3 g_{3}+2 g_{2} \wp\right), \\
g= & 16 B^{6} \mathbf{z}-1440 B^{5} \wp^{\prime}-1800 B^{4} g_{2} \mathbf{z}+54000 B^{3}\left(g_{2} \wp^{\prime}-g_{3} \mathbf{z}\right) \\
& -16875 B^{2} g_{2}^{2} \mathbf{z}-506250 B g_{2}^{2} \wp^{\prime}+421875 \Delta \mathbf{z} .
\end{aligned}
$$

It takes a couple seconds to evaluate the corresponding $12 \times 12$ Sylvester determinant (in Mathematica) to get

$$
R(f, g ; B)(\mathbf{z})=2^{36} 3^{27} 5^{30} \Delta^{5}\left(\wp^{\prime}\right)^{4} W_{3}(\mathbf{z})
$$

where $W_{3}(\mathbf{z})$ is given by

$$
W_{3}(\mathbf{z})=\mathbf{z}^{6}-15 \wp \mathbf{z}^{4}-20 \wp^{\prime} \mathbf{z}^{3}+\left(\frac{27}{4} g_{2}-45 \wp^{2}\right) \mathbf{z}^{2}-12 \wp \wp^{\prime} \mathbf{z}-\frac{5}{4} \wp^{\prime 2} .
$$

It seems impractical to compute this resultant by hand.
Both $Z_{2}$ and $Z_{3}$ are known to Dahmen [3]. Here is a new example:
Example 5.10. - For $n=4$, the expansions of the polynomials $f$ and $g$, as given in (5.7) and (5.6) by a direct substitution, are already too long to put here. Nevertheless, a couple hours Mathematica calculation gives

$$
R(f, g ; B)(\mathbf{z})=-2^{80} 3^{63} 5^{60} 7^{63} \Delta^{18}\left(\wp^{\prime}\right)^{8} W_{4}(\mathbf{z})
$$

where $W_{4}(\mathbf{z})$ is the degree 10 polynomial:

$$
\begin{align*}
W_{4}(\mathbf{z})= & \mathbf{z}^{10}-45 \wp \mathbf{z}^{8}-120 \wp^{\prime} \mathbf{z}^{7}+\left(\frac{399}{4} g_{2}-630 \wp^{2}\right) \mathbf{z}^{6}-504 \wp \wp^{\prime} \mathbf{z}^{5} \\
& -\frac{15}{4}\left(280 \wp^{3}-49 g_{2} \wp-115 g_{3}\right) \mathbf{z}^{4}+15\left(11 g_{2}-24 \wp^{2}\right) \wp^{\prime} \mathbf{z}^{3} \\
& -\frac{9}{4}\left(140 \wp^{4}-245 g_{2} \wp^{2}+190 g_{3} \wp+21 g_{2}^{2}\right) \mathbf{z}^{2}  \tag{5.8}\\
& -\left(40 \wp^{3}-163 g_{2 \wp}+125 g_{3}\right) \wp^{\prime} \mathbf{z}+\frac{3}{4}\left(25 g_{2}-3 \wp^{2}\right)\left(\wp^{\prime}\right)^{2} .
\end{align*}
$$

The weight 10 pre-modular form $Z_{4}(\sigma ; \tau)$ is then obtained.
5.3. Remarks on rationality and singularities of the Lamé curve. - We have constructed two affine curves from $\bar{X}_{n}$. One is the hyperelliptic model

$$
Y_{n}=\left\{(B, C) \mid C^{2}=\ell_{n}(B)\right\}
$$

another one is

$$
Y_{n}^{\prime}:=\left\{\left(x_{0}, y_{0}, \mathbf{z}\right) \mid y_{0}^{2}=4 x_{0}^{2}-g_{2} x_{0}-g_{3}, W_{n}\left(x_{0}, y_{0} ; \mathbf{z}\right)=0\right\}
$$

which is a degree $\frac{1}{2} n(n+1)$ branched cover of the original curve

$$
E=\left\{\left(x_{0}, y_{0}\right) \mid y_{0}^{2}=4 x_{0}^{3}-g_{2} x_{0}-g_{3}\right\}
$$

under the projection $\sigma_{n}^{\prime}: Y_{n}^{\prime} \rightarrow E$ with defining equation $W_{n}(\mathbf{z})=0$.
The curve $Y_{n}$ is birational to $Y_{n}^{\prime}$ over $E$, namely the addition map $\sigma_{n}: Y_{n} \rightarrow E$ is compatible with $\sigma_{n}^{\prime}: Y_{n}^{\prime} \rightarrow E$. Notice that both $\ell_{n}$ and $W_{n}$ have coefficients
in $\mathbb{Q}\left[g_{2}, g_{3}\right]$. The explicit birational map $\phi:(B, C) \rightarrow\left(x_{0}, y_{0}, \mathbf{z}\right)$ (given in Theorems 5.3 and 5.6 via $\left.\mathbf{z}_{n}=-\kappa\right)$ also has coefficients in $\mathbb{Q}\left[g_{2}, g_{3}\right]$. This implies that $\phi$ is defined over $\mathbb{Q}$. Moreover $\phi$ extends to a birational morphism

by identifying $\sigma_{n}^{-1}\left(0_{E}\right)$ with $\mathbf{z}_{n}^{-1}(\infty)$. The morphism $\phi$ is an isomorphism outside those branch points for $Y_{n} \rightarrow \mathbb{P}^{1}(\mathbb{C})$ (i.e., $C=0$ ). In particular, the non-isomorphic loci lie in $\mathbf{z}_{n}=0$ by formulas (5.2) and Theorem 5.6.

Remark 5.11. - In contrast to the smoothness of $Y_{n}(\tau)$ for general $\tau$, for all $n \geqslant 3$ the model $Y_{n}^{\prime}(\tau)$ is singular at points $\mathbf{z}=0=y_{0}$ (and hence $x_{0}=e_{i}$ for some $i$ ). Indeed from (5.2) this is equivalent to $C=0$ and $l_{i}(B) l t_{i}(B)^{2}=0$ for some $1 \leqslant i \leqslant 3$. For $n=2$, there is only one solution $B$ for each fixed $i$ (cf. Example 5.5). However, for $n \geqslant 3$ there are more than one such solutions $B$. These points $(B, 0) \in Y_{n}$ are collapsed to the same point $\left(x_{0}, y_{0}, \mathbf{z}\right)=\left(e_{i}, 0,0\right) \in Y_{n}^{\prime}$ under $\phi$, thus $\left(e_{i}, 0,0\right)$ is a singular point of $Y_{n}^{\prime}$.

For $n=3,4$ this is easily seen from the equation $W_{n}(\mathbf{z})=0$ given above since it contains a quadratic polynomial in $\left(\mathbf{z}, \wp^{\prime}\right)$ as its lowest degree terms.

In particular, the birational map $\phi^{-1}$ is also represented by rational functions $B=B\left(x_{0}, y_{0}, \mathbf{z}\right)$ and $C=C\left(x_{0}, y_{0}, \mathbf{z}\right)$ with coefficients in $\mathbb{Q}\left[g_{2}, g_{3}\right]$ and with at most poles along $\mathbf{z}=0$. In principle such an explicit inverse can be obtained by a Groebner basis calculation associated to the ideal of the graph $\Gamma_{\phi}$. The following statement is clear from the above description:

Proposition 5.12. - Let $E$ be defined over $\mathbb{Q}$, i.e., $g_{2}, g_{3} \in \mathbb{Q}$. Then the Lamé curve $\bar{Y}_{n}$ is also defined over $\mathbb{Q}$ for all $n \in \mathbb{N}$. Moreover, $\bar{Y}_{n}^{\prime}$ and all the morphisms $\sigma_{n}, \sigma_{n}^{\prime}, \phi$ are also defined over $\mathbb{Q}$.

A rational point $(B, C) \in \bar{Y}_{n}$ is mapped to a rational point $\left(x_{0}, y_{0}, \mathbf{z}\right) \in \bar{Y}_{n}^{\prime}$ by $\phi$. For the converse, given $\left(x_{0}, y_{0}\right) \in E(\mathbb{Q})$, a point $\left(x_{0}, y_{0}, \mathbf{z}\right)$ in the $\sigma_{n}^{\prime}$-fiber gives a unique $(B, C) \in \bar{Y}_{n}(\mathbb{Q})$ if $\mathbf{z} \in \mathbb{Q}$ and $\left(x_{0}, y_{0}, \mathbf{z}\right) \neq\left(e_{i}, 0,0\right)$ for any $i$.

Remark 5.13. - It is well-known that there are only few (i.e., at most finite) rational points on a non-elliptic hyperelliptic curve. This phenomenon is consistent with the irreducibility of the polynomial $W_{n}(\mathbf{z})$ over $K(E)$ in light of Hilbert's irreducibility theorem that there is an infinite (Zariski dense) set of $\left(g_{2}, g_{3}, x_{0}, y_{0}\right) \in \mathbb{Q}^{4}$ so that the specialization of $W_{n}(\mathbf{z})$ over there are all irreducible. Nevertheless, it might be interesting to see if $\mathbf{z}_{n}$ plays any role in the study of rational points.

Remark 5.14. - It was proved in [12, Prop. 3.2] that a Lamé curve is either smooth or nodal, and there is at most one node. The proof relies on the degree formula
$\operatorname{deg} \pi_{n}=\frac{1}{2} n(n+1)=\operatorname{deg} \kappa$ which was quoted as a significant formula from finiteband integration theory without explicit references in [12, p.1139]. While this might be well-known to experts in this field, we want to point out that it also follows from Theorem 5.6 and our degree formula $\operatorname{deg} \sigma_{n}=\frac{1}{2} n(n+1)=\operatorname{deg} \mathbf{z}_{n}$ in Theorem 2.4 and Theorem 3.13.

## Appendix. A counting formula for Lamé equations by You-Cheng Chou $\left.{ }^{( }\right)$

Using the pre-modular forms constructed in Sections 4 and 5, we verify the $n=4$ case of Dahmen's conjectural counting formula [3, Conj. 73] for integral Lamé equations with finite monodromy. It is known that the finite monodromy group is necessarily a dihedral group.

For $N \in \mathbb{N}$ we denote by $\phi(N):=\#\{k \in \mathbb{Z} \mid \operatorname{gcd}(k, N)=1,0 \leqslant k<N\}$ the Euler function and we set $\phi(N)=0$ if $N \notin \mathbb{N}$. Similarly we define

$$
\Psi(N):=\#\left\{\left(k_{1}, k_{2}\right) \mid \operatorname{gcd}\left(k_{1}, k_{2}, N\right)=1,0 \leqslant k_{i}<N\right\} .
$$

A.1. Dahmen's conjecture. - Let $L_{n}(N)$ be the number of Lamé equations $w^{\prime \prime}=$ $(n(n+1) \wp(z)+B) w$, up to linear equivalence, which has finite monodromy group isomorphic to the dihedral group $D_{N}$. Using the Hermite-Halphen ansatz (0.7) and the theory in Section 4, the problem is reduced to the zero counting of the following $\mathrm{SL}(2, \mathbb{Z})$ modular form

$$
M_{n}(N):=\prod_{\substack{0 \leq k_{1}, k_{2}<N \\ \operatorname{gcd}\left(k_{1}, k_{2}, N\right)=1}} Z_{n}\left(\frac{k_{1}+k_{2} \tau}{N} ; \tau\right) .
$$

Using this, by repeating Dahmen's argument in [3, Lem. 65, 74], we get
Proposition A.1. - Suppose that for all $N \in \mathbb{Z}_{\geqslant 3}$ and $n \in \mathbb{N}$ we have that

$$
\nu_{\infty}\left(M_{n}(N)\right)=a_{n} \phi(N)+b_{n} \phi(N / 2),
$$

where $a_{2 m}=a_{2 m+1}=m(m+1) / 2, b_{2 m}=b_{2 m-1}=m^{2}$. Then

$$
L_{n}(N)=\frac{1}{2}\left(\frac{n(n+1) \Psi(N)}{24}-\left(a_{n} \phi(N)+b_{n} \phi(N / 2)\right)\right)+\frac{2}{3} \varepsilon_{n}(N),
$$

where $\varepsilon_{n}(N)=1$ if $N=3$ and $n \equiv 1(\bmod 3)$, and $\varepsilon_{n}(N)=0$ otherwise.
Furthermore, $Z_{n}(\sigma ; \tau)$ with $\sigma$ a torsion point has only simple zeros in $\tau \in \mathbb{H}$.
Proof. - Recall the formula for $\operatorname{SL}(2, \mathbb{Z})$ modular forms of weight $k$ :

$$
\sum_{P \neq \infty, i, \rho} \nu_{P}(f)+\nu_{\infty}(f)+\frac{\nu_{i}(f)}{2}+\frac{\nu_{\rho}(f)}{3}=\frac{k}{12} .
$$

[^4]The modular form $f=M_{n}(N)$ has weight $k=\frac{1}{2} n(n+1) \Psi(N)$. Notice that the counting is always doubled under the symmetry $\left(k_{1}, k_{2}\right) \rightarrow\left(N-k_{1}, N-k_{2}\right)$, thus by [3, Lem. 65] an upper bound for $L_{n}(N)$ is given by

$$
U_{n}(N):=\frac{1}{2}\left(\frac{n(n+1) \Psi(N)}{24}-\left(a_{n} \phi(N)+b_{n} \phi(N / 2)\right)\right)+\frac{2}{3} \varepsilon_{n}(N)
$$

That is, $L_{n}(N) \leqslant U_{n}(N)$. Moreover, the equality holds if and only if each factor $Z_{n}\left(\left(k_{1}+k_{2} \tau\right) / N ; \tau\right)$ has only simple zeros.

We will show the equality holds by comparing it with the counting formula for the projective monodromy group $P L_{n}(N)$ (cf. [3, Lem. 74]).

We recall the relation between $L_{n}(N)$ and $P L_{n}(N)$ :

$$
P L_{n}(N)= \begin{cases}L_{n}(N)+L_{n}(2 N) & \text { if } N \text { is odd } \\ L_{n}(2 N) & \text { if } N \text { is even }\end{cases}
$$

If $n$ is even and $N$ is odd, we have

$$
\begin{aligned}
& P L_{n}(N)=L_{n}(N)+L_{n}(2 N) \\
& \leqslant \\
& \frac{1}{2}\left(\frac{n(n+1) \Psi(N)}{24}-\left(\frac{\frac{n}{2}\left(\frac{n}{2}+1\right)}{2} \phi(N)+\frac{n^{2}}{4} \phi(N / 2)\right)\right)+\frac{2}{3} \varepsilon_{n}(N) \\
& \quad+\frac{1}{2}\left(\frac{n(n+1) \Psi(2 N)}{24}-\left(\frac{\frac{n}{2}\left(\frac{n}{2}+1\right)}{2} \phi(2 N)+\frac{n^{2}}{4} \phi(N)\right)\right)+\frac{2}{3} \varepsilon_{n}(2 N) \\
& =\frac{n(n+1)}{12}(\Psi(N)-3 \phi(N))+\frac{2}{3} \varepsilon_{n}(N)
\end{aligned}
$$

For the last equality, we use $\varepsilon_{n}(2 N)=0, \Psi(2 N)=3 \Psi(N)$ and $\phi(2 N)=\phi(N)$. (If $N$ is even, the relations are $\varepsilon_{n}(N)=0, \Psi(2 N)=4 \Psi(N)$ and $\phi(2 N)=\phi(N)$.) For the other three cases with $(n, N)$ being (even, even), (odd, odd) or (odd, even), the computations are similar. They lead to

$$
P L_{n}(N) \leqslant \frac{n(n+1)}{12}(\Psi(N)-3 \phi(N))+\frac{2}{3} \varepsilon_{n}(N)
$$

On the other hand, using the method of dessins d'enfants, Dahmen showed directly that the equality holds [4]. Thus all the intermediate inequalities are indeed equalities, and in particular $L_{n}(N)=U_{n}(N)$ holds.
A.2. $q$-expansions for some modular forms. - Recall that

$$
\begin{aligned}
\sum_{m \in \mathbb{Z}} \frac{1}{(m+z)^{k}} & =\frac{1}{(k-1)!}(-2 \pi i)^{k} \sum_{n=1}^{\infty} n^{k-1} e^{2 \pi i n z} \\
\sum_{n \in \mathbb{Z}} \frac{1}{(x+n)^{2}} & =\pi^{2} \cot ^{2}(\pi x)+\pi^{2} \\
\sum_{n \in \mathbb{Z}} \frac{1}{(x+n)^{3}} & =\pi^{3} \cot ^{3}(\pi x)+\pi^{3} \cot (\pi x)
\end{aligned}
$$

We compute the $q$-expansions for $g_{2}, g_{3}, \wp, \wp^{\prime}, Z$, where $q=e^{2 \pi i \tau}$ :

$$
g_{2}=60 \sum_{(n, m) \neq(0,0)} \frac{1}{(n+m \tau)^{4}}=60\left(2 \zeta(4)+2 \frac{(-2 \pi i)^{4}}{3!} \sum_{n=1}^{\infty} \sigma_{3}(n) q^{n}\right)
$$

where $\sigma_{k}(n):=\sum_{d \mid n} d^{k}$. Similarly,

$$
g_{3}=140 \sum_{(n, m) \neq(0,0)} \frac{1}{(n+m \tau)^{6}}=140\left(2 \zeta(6)+2 \frac{(-2 \pi i)^{6}}{5!} \sum_{n=1}^{\infty} \sigma_{5}(n) q^{n}\right) .
$$

Let $z=r+s \tau$. For $s=0$, we have

$$
\begin{aligned}
& \wp^{\prime}(r ; \tau)=-2 \sum_{n, m \in \mathbb{Z}} \frac{1}{(r+n+m \tau)^{3}} \\
& =-2 \sum_{n \in \mathbb{Z}} \frac{1}{(r+n)^{3}}-2 \sum_{m=1}^{\infty} \sum_{n \in \mathbb{Z}}\left(\frac{1}{(m \tau+n+r)^{3}}-\frac{1}{(m \tau+n-r)^{3}}\right) \\
& =-2 \sum_{n \in \mathbb{Z}} \frac{1}{(r+n)^{3}}-2 \sum_{m=1}^{\infty} \frac{(-2 \pi i)^{3}}{2!} \sum_{n=1}^{\infty} n^{2}\left(e^{2 \pi i n(m \tau+r)}-e^{2 \pi i n(m \tau-r)}\right) \\
& =-2 \pi^{3} \cot (\pi r)-2 \pi^{3} \cot ^{3}(\pi r)+16 \pi^{3} \sum_{n, m=1}^{\infty} n^{2} \sin (2 \pi n r) q^{n m} \text {. } \\
& \wp(r ; \tau)=\frac{1}{r^{2}}+\sum_{(n, m) \neq(0,0)}\left(\frac{1}{(r+n+m \tau)^{2}}-\frac{1}{(n+m \tau)^{2}}\right) \\
& =\sum_{n \in \mathbb{Z}} \frac{1}{(r+n)^{2}}-\sum_{n=1}^{\infty} \frac{2}{n^{2}}+\sum_{m=1}^{\infty} \sum_{n \in \mathbb{Z}}\left(\frac{1}{(m \tau+r+n)^{2}}+\frac{1}{(m \tau-r+n)^{2}}-\frac{2}{(m \tau+n)^{2}}\right) \\
& =\pi^{2} \cot ^{2}(\pi r)+\frac{2}{3} \pi^{2}+\sum_{m=1}^{\infty}(-2 \pi i)^{2} \sum_{n=1}^{\infty}\left(e^{2 \pi i n(m \tau+r)}+e^{2 \pi i n(m \tau-r)}-2 e^{2 \pi i n m \tau}\right) \\
& =\pi^{2} \cot ^{2}(\pi r)+\frac{2}{3} \pi^{2}+8 \pi^{2} \sum_{n, m=1}^{\infty}(1-\cos 2 n \pi r) q^{n m} .
\end{aligned}
$$

Also, for the Hecke function $Z$ (cf. (0.10)) we have

$$
Z(r ; \tau)=\pi \cot (\pi r)+4 \pi \sum_{n, m=1}^{\infty}(\sin 2 n \pi r) q^{n m}
$$

For $s=\frac{1}{2}$, we have

$$
\begin{aligned}
\wp^{\prime}\left(r+\frac{1}{2} \tau ; \tau\right) & =-2 \sum_{(n, m) \neq(0,0)} \frac{1}{\left(r+n+\left(\frac{1}{2}+m\right) \tau\right)^{3}} \\
& =-2 \sum_{m=1}^{\infty}\left(\sum_{m \in \mathbb{Z}} \frac{1}{\left(n+r+\left(m-\frac{1}{2}\right) \tau\right)^{3}}-\sum_{n \in \mathbb{Z}} \frac{1}{\left(n-r+\left(m-\frac{1}{2}\right) \tau\right)^{3}}\right) \\
& =-2 \frac{(-2 \pi i)^{3}}{2!} \sum_{n, m=1}^{\infty} n^{2}\left(e^{2 \pi i n(r+(m-1 / 2) \tau)}-e^{2 \pi i n(-r)+(m-1 / 2) \tau}\right) \\
& =16 \pi^{3} \sum_{n, m=1}^{\infty} n^{2}(\sin 2 \pi n r) q^{n(m-1 / 2)} .
\end{aligned}
$$

Similarly,

$$
\wp\left(r+\frac{1}{2} \tau ; \tau\right)=-\frac{1}{3} \pi^{2}+8 \pi^{2} \sum_{n, m=1}^{\infty} n q^{n m}-8 \pi^{2} \sum_{n, m=1}^{\infty} n(\cos 2 \pi n r) q^{n(m-1 / 2)}
$$

and $Z\left(r+\frac{1}{2} \tau ; \tau\right)=4 \pi \sum_{n, m=1}^{\infty}(\sin 2 \pi n r) q^{n(m-1 / 2)}$.
A.3. The counting formula for $n=4$. - Now we give the computations for $n=4$ and prove the formula $L_{4}(N)=U_{4}(N)$ from Proposition A.1.

Theorem A.2. - For $n=4$ and $N \in \mathbb{Z}_{\geqslant 3}$, we have

$$
L_{4}(N)=\frac{1}{2}\left(\frac{5}{6} \Psi(N)-(3 \phi(N)+4 \phi(N / 2))\right) .
$$

Moreover, $Z_{4}(\sigma ; \tau)$ with $\sigma \in E_{\tau}[N]$ has only simple zeros in $\tau \in \mathbb{H}$.
Proof. - For $n=4$, the pre-modular form $Z_{4}=W_{4}(Z)$ is given in (5.8):

$$
\begin{aligned}
& W_{4}(Z)=Z^{10}- 45 \wp Z^{8}-120 \wp^{\prime} Z^{7}+\left(\frac{399}{4} g_{2}-630 \wp^{2}\right) Z^{6}-\left(504 \wp \wp^{\prime}\right) Z^{5} \\
&-\frac{15}{4}\left(280 \wp^{3}-49 g_{2} \wp-115 g_{3}\right) Z^{4}+15\left(11 g_{2}-24 \wp^{2}\right) \wp^{\prime} Z^{3} \\
&-\frac{9}{4}\left(140 \wp^{4}-245 g_{2} \wp^{2}+190 g_{3} \wp+21 g_{2}^{2}\right) Z^{2} \\
& \quad\left(40 \wp^{3}-163 g_{2} \wp+125 g_{3}\right) \wp Z+\frac{3}{4}\left(25 g_{2}-3 \wp^{2}\right) \wp^{\prime 2},
\end{aligned}
$$

where $Z$ is the Hecke function. We compute the asymptotic behavior of $W_{4}(Z)$ when $\tau \rightarrow \infty$. Let $z=r+s \tau$. We divide the problem into two cases
(1) $s \equiv 0(\bmod 1):$ according to the $q$-expansion given in Section A.2, we have

$$
\begin{gathered}
g_{2} \rightarrow \frac{3}{4} \pi^{4}, \quad g_{3} \rightarrow \frac{8}{27} \pi^{6}, \quad Z(z) \rightarrow \pi \cot (\pi r), \\
\wp^{\prime}(z) \rightarrow-2 \pi^{3} \cot (\pi r)-2 \pi^{3} \cot ^{3}(\pi r), \\
\wp(z) \rightarrow \pi^{2} \cot ^{2}(\pi r)+\frac{2}{3} \pi^{2} .
\end{gathered}
$$

A direct computation shows that $W_{4}(Z)$ has a zero at $\infty$ when $s=0$.
By replacing all the modular forms $g_{2}, g_{3}, \wp, \wp^{\prime}$ and $Z$ appeared in $W_{4}(Z)$ with their $q$-expansions, we have (e.g. using Mathematica)

$$
W_{4}(Z)=2^{14} 3^{3} 5^{2} 7 \pi^{10} \cos ^{2}(\pi r) \sin ^{2}(\pi r) q^{3}+O\left(q^{4}\right)
$$

(2) $s \not \equiv 0(\bmod 1):$ in this case we have

$$
Z \rightarrow 2 \pi i\left(s-\frac{1}{2}\right), \quad \wp(z) \rightarrow-\frac{1}{3} \pi^{2}, \quad \wp^{\prime}(z) \rightarrow 0, \quad g_{2} \rightarrow \frac{4}{3} \pi^{4}, \quad g_{3} \rightarrow \frac{8}{27} \pi^{6}
$$

Hence the constant term of $W_{4}(Z)$ is given by

$$
W_{4}(z)=-64 \pi^{10}(-2+s)(-1+s)^{2} s^{2}(1+s)(-3+2 s)(-1+2 s)^{2}(1+2 s)+O(q)
$$

If $s \not \equiv 0(\bmod 1)$ then $W_{4}(Z)$ has a zero at $\tau=\infty \Longleftrightarrow s \equiv \frac{1}{2}(\bmod 1)$.
Now we fix $s=\frac{1}{2}$ and replace the modular forms $g_{2}, g_{3}, \wp, \wp^{\prime}$ and $Z$ appeared in $W_{4}(Z)$ with their $q$-expansions. We get

$$
W_{4}(Z)=2^{10} 3^{3} 5^{2} 7 \pi^{10} \cos ^{2}(\pi r) \sin ^{2}(\pi r) q^{2}+O\left(q^{3}\right)
$$

These computations for the $q$-expansions imply that

$$
\begin{aligned}
\nu_{\infty}\left(M_{4}(N)\right)= & 3 \#\left\{1 \leqslant k_{1} \leqslant N \mid \operatorname{gcd}\left(N, k_{1}\right)=1\right\} \\
& +2 \#\left\{0 \leqslant k_{1} \leqslant N \mid \operatorname{gcd}\left(N / 2, k_{1}\right)=1\right\} \\
= & 3 \phi(N)+4 \phi(N / 2) .
\end{aligned}
$$

Since the value of $\nu_{\infty}\left(M_{4}(N)\right)$ coincides with the assumption in Proposition A. 1 for $n=4$, the theorem follows from it accordingly

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[^0]:    ${ }^{(1)}$ The Green function $G(z)$ on $E$ is defined by $-\triangle G=\delta_{0}-1 /|E|$ and $\int_{E} G=0$, where $|E|$ is the area of $E$. Also $G(z, w)=G(z-w, 0)=G(z-w)$ by the translation invariance.

[^1]:    ${ }^{(2)}$ It was stated in Part I [1, p. 141-142] that such a complete solution for $n=1$ will appear in Part II (this paper), and it was included in the first arXiv version arXiv:1502.03295v1. Later on we found its deep connection with Painlevé VI equations. Therefore we extracted that part and published it separately in [2].
    ${ }^{(3)}$ In this paper, we often use $\sigma$ as the coordinate on $E$ whenever the map $\sigma_{n}: X_{n} \rightarrow E$ defined in (0.13) is involved. This should not be confused with the Weierstrass $\sigma$ function.

[^2]:    ${ }^{(4)}$ The fact that the coefficients lie in $\mathbb{Q}$, instead of just in $\mathbb{C}$, follows from standard elimination theory and two facts (i) The equations of $\bar{X}_{n}$ are defined over $\mathbb{Q}\left[g_{2}, g_{3}\right]$ (cf. equations (1.2)), and (ii) the addition map $E^{n} \rightarrow E$ is defined over $\mathbb{Q}$. In Section 5 , we carry out the elimination procedure using the resultant for another explicit presentation $\pi_{n}$ of $\sigma_{n}$.

[^3]:    ${ }^{(5)}$ Notice that the weight, as assigned in Theorem 3.2, is twice the degree.

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