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# Actions of the Neumann systems 

## via Picard-Fuchs equations

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July 6, 2000


#### Abstract

The Neumann system describing the motion of a particle on an $n$-dimensional sphere with an anisotropic harmonic potential, has been celebrated as one of the best understood integrable systems of classical mechanics. The present paper adds a detailed discussion and the determination of its action integrals, using differential equations rather than standard integral formulas. We show that the actions of the Neumann system satisfy a Picard-Fuchs equation which in suitable coordinates has a rather simple form for arbitrary $n$. We also present an explicit form of the related Gauß-Manin equations. These formulas are used for the numerical calculation of the actions of the Neumann system.


## 1 Introduction

The classical mechanical problem of a point particle moving on a unit sphere, under the influence of an anisotropic harmonic force, was the first case solved in terms of hyperelliptic theta functions of genus 2 (also called ultraelliptic theta functions). This was done by C. Neumann in his 1856 doctoral dissertation [8]. An amazing relation between the Neumann system and the spectral theory of Schrödinger operators was discovered by J. Moser and E. Trubowitz, see $[7,12]$. The action variables of the Neumann system from this point of view have been discussed by S. P. Novikov and A. P. Veselov in [9].

We will consider the natural $n$-dimensional generalization of the Neumann system, which is a Hamiltonian system in $\mathbb{R}^{2 n+2}$ with the Hamiltonian

$$
\begin{equation*}
H(\boldsymbol{x}, \boldsymbol{y})=\frac{1}{2}\left(|\boldsymbol{y}|^{2}|\boldsymbol{x}|^{2}-(\boldsymbol{y}, \boldsymbol{x})^{2}\right)+\frac{1}{2}(\boldsymbol{x}, \boldsymbol{A} \boldsymbol{x}) \tag{1}
\end{equation*}
$$

where $\boldsymbol{x}=\left(x_{0}, \ldots, x_{n}\right)^{t}, \boldsymbol{y}$ is the conjugate momentum, and $($,$) is the standard Euclidean$ scalar product. $\boldsymbol{A}$ is a positive definite matrix which may be assumed diagonal, with eigenvalues $0<a_{0}<\ldots<a_{n}$. It leaves the unit sphere $S^{n} \subset \mathbb{R}^{n+1}$ invariant: $|\boldsymbol{x}|^{2}$ is a geometric constant which is chosen as 1. $H(\boldsymbol{x}, \boldsymbol{y})$ possesses a set of algebraic integrals $F_{0}, \ldots, F_{n}$ in involution, which, surprisingly, have only recently been found (see [7]):

$$
\begin{equation*}
F_{\nu}(\boldsymbol{x}, \boldsymbol{y})=x_{\nu}^{2}+\sum_{\mu \neq \nu} \frac{\left(x_{\nu} y_{\mu}-x_{\mu} y_{\nu}\right)^{2}}{a_{\nu}-a_{\mu}}, \quad \quad \sum_{\nu=0}^{n} F_{\nu}=1 \tag{2}
\end{equation*}
$$

These integrals are related to another set of integrals, $\eta_{1}, \ldots, \eta_{n}$, defined by the formula

$$
\begin{equation*}
\sum_{\nu=0}^{n} \frac{F_{\nu}}{z-a_{\nu}}=\frac{Q(z)}{A(z)} \tag{3}
\end{equation*}
$$

where

$$
\begin{align*}
& Q(z)=z^{n}+2 \eta_{1} z^{n-1}+2 \eta_{2} z^{n-2}+\ldots+2 \eta_{n}  \tag{4}\\
& A(z)=\prod_{\nu=0}^{n}\left(z-a_{\nu}\right) \tag{5}
\end{align*}
$$

The Liouville tori of this system are in general certain coverings of real parts of Jacobi varieties of the hyperelliptic curve

$$
\begin{equation*}
w^{2}=R(z), \quad R(z)=-A(z) Q(z) \tag{6}
\end{equation*}
$$

The actions of the Neumann system are the periods of the Abelian integral

$$
\begin{equation*}
I_{j}=\frac{1}{2 \pi} \oint \frac{Q(z)}{\sqrt{R(z)}} d z, \quad j=1, \ldots, n \tag{7}
\end{equation*}
$$

over certain real cycles $\gamma_{j}$ (see section 3 for the details).

Similar representations of the actions in terms of Abelian integrals can be given for almost all known integrable systems; for the Kovalevskaya top this has only recently been achieved in [6, 2]. To use the actions e.g. for the computation of semiclassical approximations to quantum mechanical eigenvalues $[6,10,13]$, they have to be evaluated numerically. Of course, (7) is an explicit formula, but from the computational point of view it is not efficient when the integral (7) must be determined for all values of the constants of motion. We suggest another way for calculating the actions, using the Picard-Fuchs equation they satisfy. The existence of such linear ordinary differential equations follows from general de Rham type arguments, but their explicit form can be very complicated (see e. g. [2] where the case of the Kovalevskaya top is discussed). The main goal of this paper is to show that for the Neumann system the Picard-Fuchs (and the related Gauß-Manin) equations in appropriate variables have a relatively simple form, which makes them very convenient for numerical computations.

We should mention that the simplicity of the Picard-Fuchs/Gauß-Manin equations depends very much on the choice of constants of motion. In particular, if we take $F_{\nu}$ or $\eta_{j}$ as independent variables in the differential equation they are quite cumbersome already for $n=2$. However, with the variables $q_{i}$ which are the zeros of the polynomial $Q(z)$,

$$
\begin{equation*}
Q(z)=\prod_{j=1}^{n}\left(z-q_{j}\right) \tag{8}
\end{equation*}
$$

they become much simpler. Although the Picard-Fuchs equations look more natural and are simpler, the Gauß-Manin equations are more suitable from the computational point of view. The reason is that for each $q_{i}$ the corresponding Picard-Fuchs equation has a different set of dependent variables, while the corresponding $n$ Gauß-Manin equations give commuting flows for the same set of dependent variables.

The outline of this paper is as follows. In Section 2 we recall the basic properties of the Neumann system. Separation of variables is achieved in a standard way [7] using elliptical spherical coordinates. In section 3 we discuss the explicit formulas for the actions and introduce reduced actions using the obvious reflection symmetry group $\mathbb{Z}_{2}^{n+1}$. These two actions differ only by some constant coefficients (and therefore satisfy the same linear equations) but the reduced actions are continuous functions of the constants of motion while the usual actions are not. Sections 4 and 5 contain the main results. We find the explicit form of the corresponding Picard-Fuchs and Gauß-Manin equations in the variables $q_{i}$. In the final section we use these results to compute the actions for the Neumann system.

## 2 Separation of variables

In the first part of this section we follow closely the presentation of section 7 of Moser's paper [7].

Consider a particle moving on a unit sphere $S^{n} \subset \mathbb{R}^{n+1}(\boldsymbol{x})$, and let it be attracted towards the center of the sphere by a non-isotropic harmonic potential

$$
\begin{equation*}
V(\boldsymbol{x})=\frac{1}{2}(\boldsymbol{x}, \boldsymbol{A} \boldsymbol{x}) \tag{9}
\end{equation*}
$$

where the matrix $\boldsymbol{A}$ is positive symmetric. Without loss of generality, we may assume $\boldsymbol{A}=\operatorname{diag}\left(a_{\nu}\right)$ with spring constants $0<a_{0}<\ldots<a_{n}$. The constraint $|\boldsymbol{x}|^{2}=1$ may be accounted for in different ways. One possibility is in terms of the Hamiltonian (1) where the kinetic energy has an appropriate dependence on the point $\boldsymbol{x}$. A more familiar starting point is to include a force $\lambda \boldsymbol{x}$ perpendicular to the sphere,

$$
\begin{equation*}
\ddot{\boldsymbol{x}}+\boldsymbol{A} \boldsymbol{x}=\lambda \boldsymbol{x}, \tag{10}
\end{equation*}
$$

and to determine its strength $\lambda$ from the constraint. Differentiating $|\boldsymbol{x}|^{2}=1$ twice with respect to time $t$, and taking the scalar product of (10) with $\boldsymbol{x}$, we find that $\lambda$ is

$$
\begin{equation*}
\lambda=(\boldsymbol{x}, \boldsymbol{A} \boldsymbol{x})-|\dot{\boldsymbol{x}}|^{2} . \tag{11}
\end{equation*}
$$

It may be checked by direct computation that the equations (10) are then equivalent to the canonical equations derived from (1).

In the neighborhood of the stable equilibrium points $\boldsymbol{x}= \pm \boldsymbol{e}_{0}=( \pm 1,0, \ldots, 0)$, we have, to lowest order, $\lambda=a_{0}$ and therefore $\ddot{x}_{j} \approx-\left(a_{j}-a_{0}\right) x_{j}$. Hence the frequencies of small oscillations at energies close to $h=V_{\min }=a_{0} / 2$ are

$$
\begin{equation*}
\omega_{j 0}=\sqrt{a_{j}-a_{0}} \quad(j=1, \ldots, n) . \tag{12}
\end{equation*}
$$

The $n$ eigendirections are tangent to the great circles through $\boldsymbol{e}_{0}$ and the other equilibrium points $\boldsymbol{e}_{j}$. The equilibria at $\boldsymbol{x}= \pm \boldsymbol{e}_{j}$ have energies $a_{j} / 2$ and $n-j$ stable eigendirections (tangent to the great circles through $\boldsymbol{e}_{j}$ and $\boldsymbol{e}_{k}, k=j+1, \ldots, n$ ), with eigenfrequencies $\omega_{k j}=\sqrt{a_{k}-a_{j}}$. The $j$ great circles connecting to $\boldsymbol{e}_{0}, \ldots, \boldsymbol{e}_{j-1}$ give unstable eigendirections, with eigenvalues $\sqrt{a_{j}-a_{i}}, j>i$.

The $n(n+1) / 2$ great circles connecting all pairs of equilibrium points, plus the corresponding momenta, are invariant sets for the equations of motion. This implies the existence of periodic motion along them, of oscillatory or rotational type, depending on the available energy. Parametrizing the circle through $\boldsymbol{e}_{\mu}$ and $\boldsymbol{e}_{\nu}(\mu<\nu)$ with the angle $\varphi$ via $\left(x_{\mu}, x_{\nu}\right)=(\cos \varphi, \sin \varphi)$, it is straightforward to describe this periodic motion with the equation

$$
\begin{equation*}
2 \ddot{\varphi}=-\left(a_{\nu}-a_{\mu}\right) \sin 2 \varphi, \tag{13}
\end{equation*}
$$

which, for each half of the great circle, is nothing but the pendulum equation. It admits solutions for energies $h \geq a_{\mu} / 2$. For $n \geq 2$ and intermediate values of the energy there also exist isolated periodic orbits not related to great circles on $\mathrm{S}^{n}$.

To solve the equations in general, one uses the fact that they separate in Jacobi's "elliptical spherical coordinates" $\left(u_{1}, \ldots, u_{n}\right)$ on $\mathrm{S}^{n}$. These are implicitly given as the zeros $u_{j} \equiv z_{j}(\boldsymbol{x})$ of

$$
\begin{equation*}
f(z):=\sum_{\nu=0}^{n} \frac{x_{\nu}^{2}}{z-a_{\nu}}=0 . \tag{14}
\end{equation*}
$$

(The coordinates $u_{j}$ are not to be confused with the $n+1$ elliptic coordinates in $\mathbb{R}^{n+1}$, also due to Jacobi, but defined as the zeros of $f(z)-1=0$.) The $u_{j}$ are labeled according to the scheme

$$
\begin{equation*}
a_{0} \leq u_{1} \leq a_{1} \leq \ldots \leq a_{n-1} \leq u_{n} \leq a_{n} . \tag{15}
\end{equation*}
$$



Figure 1: Jacobi's elliptical spherical coordinates $u_{1}, u_{2}$ on $\mathrm{S}^{2}$, for $\left(a_{0}, a_{1}, a_{2}\right)=(1,2,4)$. Coordinate lines $u_{1}=$ const are more or less deformed circles with centers on the $x_{0}$-axis; $u_{1}=a_{0}$ is the great circle $x_{0}=0$. Coordinate lines $u_{2}=$ const surround the $x_{2}$-axis; $u_{2}=a_{2}$ is the great circle $x_{2}=0$. The great circle $x_{1}=0$ is composed from pieces with $u_{1}=a_{1}$ or $u_{2}=a_{1}$. One full sheet of coordinates $u_{1}, u_{2}$ is needed for each octant of $S^{2}$; one octant is shaded. The lines are drawn at equidistant values of $u_{1}$ and $u_{2}$.

Note that $2^{n+1}$ copies of the set $\left\{u_{1}, \ldots, u_{n}\right\}$ are needed to cover the entire sphere $\mathrm{S}^{n}$. On the circle $n=1$, the coordinate $u=a_{1} x_{0}^{2}+a_{0} x_{1}^{2}$ runs from $u=a_{0}$ at $\left(x_{0}, x_{1}\right)=(0, \pm 1)$ to $u=a_{1}$ at $\left(x_{0}, x_{1}\right)=( \pm 1,0)$. On the 2 -sphere $S^{2}$, the coordinate lines for $u_{1}$ and $u_{2}$ are shown in Figure 1. Their orthogonality is demonstrated below, see (19). The circle $x_{1}=0$ is a singular line made up of segments $u_{1}=a_{1}$ and $u_{2}=a_{1}$ which meet at the four points $x_{0}^{2}=\left(a_{1}-a_{0}\right) /\left(a_{2}-a_{0}\right)$.

We shall use latin dummies for indices that run from 1 to $n$, and greek for those that run from 0 to $n$. The inverse transformation $\left\{u_{j}\right\} \rightarrow\left\{x_{\nu}\right\}$ is obtained from writing (14) as

$$
\begin{equation*}
\sum_{\nu=0}^{n} \frac{x_{\nu}^{2}}{z-a_{\nu}}=\frac{\prod_{j}\left(z-u_{j}\right)}{\prod_{\nu}\left(z-a_{\nu}\right)}=: \frac{U(z)}{A(z)} \tag{16}
\end{equation*}
$$

This is easily checked by comparing the zeros and poles of both sides. The relation shows that the $x_{\nu}^{2}$ are given as the residues of $f(z) \mathrm{d} z$,

$$
\begin{equation*}
x_{\nu}^{2}=\frac{U\left(a_{\nu}\right)}{A^{\prime}\left(a_{\nu}\right)}=\frac{\prod_{j}\left(a_{\nu}-u_{j}\right)}{\prod_{\mu \neq \nu}\left(a_{\nu}-a_{\mu}\right)} . \tag{17}
\end{equation*}
$$

Let us now prove that the $u_{j}$ are an orthogonal system on $\mathrm{S}^{n}$. To this end we calculate the metric

$$
\begin{equation*}
\mathrm{d} s^{2}=\sum_{\nu=0}^{n} \mathrm{~d} x_{\nu}^{2} \tag{18}
\end{equation*}
$$

in terms of $u_{j}$. Taking the derivative of the logarithm of Eq. (17), we get

$$
\frac{\mathrm{d} x_{\nu}}{x_{\nu}}=\frac{1}{2} \sum_{j} \frac{\mathrm{~d} u_{j}}{u_{j}-a_{\nu}},
$$

and hence

$$
4 \mathrm{~d} s^{2}=\sum_{\nu} x_{\nu}^{2}\left(\sum_{j} \frac{\mathrm{~d} u_{j}}{u_{j}-a_{\nu}}\right)^{2}=\sum_{\nu} \sum_{j} \frac{x_{\nu}^{2} \mathrm{~d} u_{j}^{2}}{\left(u_{j}-a_{\nu}\right)^{2}}+\sum_{\nu} \sum_{j \neq k} \frac{x_{\nu}^{2} \mathrm{~d} u_{j} \mathrm{~d} u_{k}}{\left(u_{j}-a_{\nu}\right)\left(u_{k}-a_{\nu}\right)} .
$$

Rewriting the second term as

$$
\sum_{\nu} \sum_{j \neq k} \frac{x_{\nu}^{2} \mathrm{~d} u_{j} \mathrm{~d} u_{k}}{u_{k}-u_{j}}\left(\frac{1}{u_{j}-a_{\nu}}-\frac{1}{u_{k}-a_{\nu}}\right),
$$

we see that the off-diagonal elements of the metric vanish, and obtain

$$
\begin{equation*}
\mathrm{d} s^{2}=\sum_{j=1}^{n} g_{j}(\boldsymbol{u}) \mathrm{d} u_{j}^{2}, \tag{19}
\end{equation*}
$$

where

$$
\begin{equation*}
g_{j}(\boldsymbol{u})=\frac{1}{4} \sum_{\nu} \frac{x_{\nu}^{2}}{\left(u_{j}-a_{\nu}\right)^{2}}=-\left.\frac{1}{4} \frac{\mathrm{~d}}{\mathrm{~d} z} \sum_{\nu} \frac{x_{\nu}^{2}}{z-a_{\nu}}\right|_{z=u_{j}}=-\frac{1}{4} \frac{U^{\prime}\left(u_{j}\right)}{A\left(u_{j}\right)} . \tag{20}
\end{equation*}
$$

The kinetic and potential energies can now be expressed in terms of the $u_{j}$ :

$$
\begin{align*}
& T=\frac{1}{2}|\dot{\boldsymbol{x}}|^{2}=\frac{1}{2} \sum_{j=1}^{n} g_{j}(\boldsymbol{u}) \dot{u}_{j}^{2},  \tag{21}\\
& V=\frac{1}{2}(\boldsymbol{x}, \boldsymbol{A} \boldsymbol{x})=\frac{1}{2} \sum_{\nu=0}^{n} a_{\nu}-\frac{1}{2} \sum_{j=1}^{n} u_{j} . \tag{22}
\end{align*}
$$

To see this we expand both sides of Eq. (16) in $1 / z$. The left hand side is

$$
\sum_{\nu} \frac{x_{\nu}^{2}}{z\left(1-a_{\nu} / z\right)}=\frac{1}{z} \sum_{\nu} x_{\nu}^{2}\left(1+\frac{a_{\nu}}{z}+\frac{a_{\nu}^{2}}{z^{2}}+\ldots\right)=\frac{1}{z}+\frac{1}{z^{2}} \sum_{\nu}\left(a_{\nu} x_{\nu}^{2}+\frac{a_{\nu}^{2} x_{\nu}^{2}}{z}+\ldots\right),
$$

and the right hand side

$$
\begin{aligned}
\frac{\prod_{j}\left(z-u_{j}\right)}{\prod_{\nu} z\left(1-a_{\nu} / z\right)} & =\frac{1}{z^{n+1}} \prod_{j}\left(z-u_{j}\right) \prod_{\nu}\left(1+\frac{a_{\nu}}{z}+\frac{a_{\nu}^{2}}{z^{2}}+\ldots\right) \\
& =\frac{1}{z^{n+1}}\left(z^{n}-\left(\sum_{j} u_{j}\right) z^{n-1}+\ldots\right)\left(1+\frac{1}{z} \sum_{\nu} a_{\nu}+\ldots\right) .
\end{aligned}
$$

Comparison of the coefficients of $1 / z^{2}$ on both sides leads to the equality (22).

We introduce the canonical momenta

$$
\begin{equation*}
p_{j}=\frac{\partial T}{\partial \dot{u}_{j}}=g_{j}(\boldsymbol{u}) \dot{u}_{j} \tag{23}
\end{equation*}
$$

and obtain the Hamiltonian

$$
\begin{equation*}
H(\boldsymbol{u}, \boldsymbol{p})=T+V=\frac{1}{2} \sum_{j=1}^{n}\left(\frac{p_{j}^{2}}{g_{j}(\boldsymbol{u})}-u_{j}\right)+c \tag{24}
\end{equation*}
$$

where $c=\frac{1}{2} \sum_{\nu} a_{\nu}$ is a constant of no interest. The Hamilton-Jacobi equation $H\left(u, \frac{\partial S}{\partial u}\right)=$ const in this case has the form

$$
\begin{equation*}
\sum_{j}\left(-4 \frac{A\left(u_{j}\right)}{U^{\prime}\left(u_{j}\right)}\left(\frac{\partial S}{\partial u_{j}}\right)^{2}-u_{j}\right)-2 \eta_{1}=0 . \tag{25}
\end{equation*}
$$

We are looking for the solution of this equation $S(u, \eta)$ which depends on $n$ additional parameters $\eta_{1}, \ldots, \eta_{n}$. To separate the variables we need the following two identities found by Jacobi. For arbitrary distinct $u_{j}$, the corresponding polynomial $U(z)=\prod_{j}\left(z-u_{j}\right)$ as in (16), and an arbitrary polynomial $P$ of the form

$$
\begin{equation*}
P(z)=\eta_{1} z^{n-1}+\eta_{2} z^{n-2}+\ldots+\eta_{n}, \tag{26}
\end{equation*}
$$

it holds that

$$
\begin{equation*}
\eta_{1}=\sum_{j=1}^{n} \frac{P\left(u_{j}\right)}{U^{\prime}\left(u_{j}\right)}, \quad \text { and } \quad \sum_{j} u_{j}=\sum_{j} \frac{u_{j}^{n}}{U^{\prime}\left(u_{j}\right)} . \tag{27}
\end{equation*}
$$

To prove these identities, consider the integrals

$$
\frac{1}{2 \pi i} \oint \frac{P(z)}{U(z)} \mathrm{d} z \quad \text { and } \quad \frac{1}{2 \pi i} \oint \frac{z^{n}}{U(z)} \mathrm{d} z
$$

with a contour that encloses all $u_{j}$, and calculate them using the theorem of residues.
Replacing $-2 \eta_{1}$ and $\sum u_{j}$ in (25) by means of Eqs. (27), the Hamilton-Jacobi equation is transformed into

$$
\begin{equation*}
\sum_{j} \frac{1}{U^{\prime}\left(u_{j}\right)}\left(-4 A\left(u_{j}\right)\left(\frac{\partial S}{\partial u_{j}}\right)^{2}-Q\left(u_{j}\right)\right)=0 \tag{28}
\end{equation*}
$$

where $Q(z)=z^{n}+2 P(z)=z^{n}+2 \eta_{1} z^{n-1}+\ldots+2 \eta_{n}$. Setting each term in the sum equal to zero, we obtain

$$
\begin{equation*}
\left(\frac{\partial S}{\partial u_{j}}\right)^{2}=-\frac{Q\left(u_{j}\right)}{4 A\left(u_{j}\right)}, \tag{29}
\end{equation*}
$$

so that

$$
S(u, \eta)=\frac{1}{2} \sum_{j=1}^{n} \int^{u_{j}} \sqrt{\frac{Q(z)}{-A(z)}} \mathrm{d} z
$$

is the solution we are looking for. The parameters $\eta_{1}, \ldots, \eta_{n}$ are integrals of the motion which are related to the involutive algebraic integrals (2) by the formula (3) [7].


Figure 2: The real part of the hyperelliptic curve $w^{2}=R(u)$

## 3 Actions and reduced actions

According to the Liouville-Arnold theorem [1] the generic level set corresponding to fixed values of these integrals must be a union of $n$-dimensional tori. In the neighborhood of these generic tori the actions are defined by the formula

$$
I_{j}=\frac{1}{2 \pi} \int_{\gamma_{j}} \boldsymbol{y} \mathrm{~d} \boldsymbol{x}, \quad j=1, \ldots, n
$$

where $\gamma_{1}, \ldots, \gamma_{n}$ is a basis of cycles on these tori. In the coordinates $u_{1}, \ldots, u_{n}$ the tori are given by fixing $\eta_{1}, \ldots, \eta_{n}$, and the canonical one-form is

$$
\begin{equation*}
\boldsymbol{y} \mathrm{d} \boldsymbol{x}=\sum_{i=1}^{n} p_{i} \mathrm{~d} u_{i}=\frac{1}{2} \sum_{j=1}^{n} \sqrt{\frac{Q\left(u_{j}\right)}{-A\left(u_{j}\right)}} \mathrm{d} u_{j} \tag{30}
\end{equation*}
$$

since

$$
\begin{equation*}
p_{j}=\frac{\partial S}{\partial u_{j}}=\frac{1}{2} \sqrt{\frac{Q\left(u_{j}\right)}{-A\left(u_{j}\right)}} . \tag{31}
\end{equation*}
$$

The last equations can be rewritten in the form

$$
\begin{equation*}
w_{j}^{2}=R\left(u_{j}\right), \quad \text { where } \quad w_{j}=2 A\left(u_{j}\right) p_{j} \quad \text { and } \quad R(z)=-A(z) Q(z) \tag{32}
\end{equation*}
$$

The zeros of $R(z)$ are the constants $a_{0}, \ldots, a_{n}$ plus the $n$ zeros $q_{1}, \ldots, q_{n}$ of the polynomial $Q(z)$ which are known to be real [7]. We label them in increasing order: $z_{0} \leq z_{1} \leq \cdots \leq$ $z_{2 n}$. Equations (32) show that the points $P_{j}=\left(w_{j}, u_{j}\right)$ for any $j$ belong to the real part of the hyperelliptic curve $\Gamma$, of genus $n$,

$$
\begin{equation*}
\Gamma=\left\{(w, u) \in \mathbb{C}^{2}: w^{2}=R(u)\right\} \tag{33}
\end{equation*}
$$

The real part of $\Gamma$ consists of $n$ ovals and one unbounded part; all $P_{i}$ belong to the ovals since by definition the $u_{j}$ are bounded. Let us denote by $c_{j}$ the oval projecting into the segments $\left[z_{2 j-1}, z_{2 j}\right]$ under the map $\pi:(w, u) \rightarrow u$, see Figure 2.

The elliptical spherical coordinates $u_{1}, \ldots, u_{n}$ determine uniquely only $x_{0}^{2}, \ldots, x_{n}^{2}$, see (17), but not the signs of $x_{0}, \ldots, x_{n}$. Now, $x_{\nu}=0$ when one of the coordinates satisfies
$u_{j}=a_{\nu}$. Consequently, when one of the ends of the segment $\left[z_{2 j-1}, z_{2 j}\right]$ is equal to $a_{\nu}$ then the coordinate $x_{\nu}$ changes sign along the cycle $c_{j}$, i.e., it crosses over to another " $2{ }^{n+1}$-tant" of $\mathbb{R}^{n+1}$. This means that the cycle $\gamma_{j}$ on the Liouville torus corresponds to the oval $c_{j}$ taken twice. The same is true when both ends of the segment $\left[z_{2 j-1}, z_{2 j}\right]$ are values $a_{\nu}$ and $a_{\nu+1}$, respectively; then the oval $c_{j}$ involves two cross-overs, at $x_{\nu}=0$ and $x_{\nu+1}=0$, while the cycle $\gamma_{j}$ involves four. This observation leads to the following definition of cycles in the homology group of the curve $\Gamma$ :

$$
\begin{aligned}
\gamma_{j} & =c_{j} \quad \text { if the segment }\left[z_{2 j-1}, z_{2 j}\right] \text { contains none of } a_{0}, \ldots, a_{n}, \\
\gamma_{j} & =2 c_{j} \quad \text { otherwise }
\end{aligned}
$$

We can summarize these results in the following
Theorem 1. The actions of the Neumann system are the complete hyperelliptic integrals

$$
\begin{equation*}
I_{j}=\frac{1}{2 \pi} \oint_{\gamma_{j}} \frac{Q(u)}{2 w} d u \tag{34}
\end{equation*}
$$

over the cycles $\gamma_{j}$ on the curve $\Gamma$ defined above.
This actually means that the Liouville tori for the Neumann system are coverings of the real Jacobian varieties of the curves $\Gamma$, which are isomorphic to the direct product of the cycles $c_{j}$. We now illustrate this result for the simplest case $n=1$. With the natural angular coordinate $\phi$ on the unit circle $S^{1}$ we have the energy

$$
h=\frac{1}{2} \dot{\phi}^{2}+\frac{1}{2}\left(a_{0} \cos ^{2} \phi+a_{1} \sin ^{2} \phi\right) .
$$

The coordinate $u$ is defined by the equation

$$
\frac{\cos ^{2} \phi}{a_{0}-u}+\frac{\sin ^{2} \phi}{a_{1}-u}=0,
$$

which gives $u=a_{1} \cos ^{2} \phi+a_{0} \sin ^{2} \phi$. In the coordinate $\phi$ the action is obviously

$$
I=\frac{1}{2 \pi} \oint \sqrt{2 h-\left(a_{0} \cos ^{2} \phi+a_{1} \sin ^{2} \phi\right)} \mathrm{d} \phi,
$$

which in the coordinate $u$ gives

$$
I=\frac{1}{2 \pi} \oint \sqrt{\frac{u-q}{\left(a_{1}-u\right)\left(u-a_{0}\right)}} \mathrm{d} u, \quad q=a_{0}+a_{1}-2 h
$$

where the integral is taken over the cycles $c_{1}$ in the previous notation. Comparing this formula with (34) we see the factor 2 which is compensated by the choice $\gamma_{1}=2 c_{1}$. The oval $c_{1}$ projects to the interval $\left[q, a_{1}\right]$ in the case of oscillatory motion $\left(q>a_{0}\right)$, and to the interval $\left[a_{0}, a_{1}\right]$ for rotations $\left(q<a_{0}\right)$. Hence, for $n=1$, the possibility $\gamma_{j}=c_{j}$ does not occur.

Another interpretation of the factor 2 derives from the following observation. Writing the energy as

$$
h=\frac{1}{2} \dot{\phi}^{2}+\frac{1}{4}\left(a_{0}+a_{1}\right)+\frac{1}{4}\left(a_{0}-a_{1}\right) \cos 2 \phi,
$$

we see that up to a constant it is the energy of a pendulum in the coordinate $\psi=2 \phi$. Thus the Neumann system for $n=1$ corresponds to two pendulums on the $\phi$-circle.

The first case where it is possible that a $u_{j}$ interval $\left[z_{2 j-1}, z_{2 j}\right]$ contains no $a_{\nu}$ is $n=2$, see section 6. It is obvious that at the transition from $\gamma_{j}=2 c_{j}$ to $\gamma_{j}=c_{j}$ the action $I_{j}$ must be a discontinuous function on phase space. This defect can be removed if we consider the reduced actions $J_{j}$, which are the actions of the Neumann system after reduction by the obvious symmetry group $G=\mathbb{Z}_{2}^{n+1}$ generated by the $n+1$ reflections $\left(x_{i}, y_{i}\right) \rightarrow\left(-x_{i},-y_{i}\right)$. The reduced space is smooth outside the fixed hyperplanes $x_{i}=y_{i}=0$ which correspond to Neumann systems of dimension $n-1$. Since the actions are defined only in the vicinity of a regular Liouville torus we should exclude these motions anyway; the corresponding tori have dimension $n-1$ and therefore are singular.

The coordinates $\left(u_{1}, \ldots, u_{n}\right)$ are well-defined on the reduced space because by definition they depend only on $x_{0}^{2}, \ldots, x_{n}^{2}$. Equation (31) gives the equation of the Liouville tori in the reduced space, which is isomorphic to the product of the ovals (or equivalently to the bounded real components of the Jacobian variety of $\Gamma$ ). Equation (30) gives the canonical 1 -form in the reduced coordinates, so the reduced actions are simply the complete Abelian integrals

$$
\begin{equation*}
J_{j}=\frac{1}{2 \pi} \oint_{c_{j}} \frac{Q(u)}{w} \mathrm{~d} u . \tag{35}
\end{equation*}
$$

Notice that the $J_{j}$ are obviously continuous functions of the constants of motion. The "explanation" of this fact is that the reduced system has only one Liouville torus at each level set obtained by fixing the constants of motion. One can show that the original Neumann system has $N=2^{m+1}$ disjoint Liouville tori, where $m$ is the number of cycles with $\gamma_{j}=c_{j}$, see the appendix. Therefore the volume forms $d I_{1} \ldots d I_{n}$ and $d J_{1} \ldots d J_{n}$ are related by

$$
\begin{equation*}
|G| d J_{1} \ldots d J_{n}=N d I_{1} \ldots d I_{n}, \quad|G|=2^{n+1} . \tag{36}
\end{equation*}
$$

The existence of such a modified version of action variables for integrable systems with several Liouville tori at a given level of the integrals seems to be an interesting phenomenon which deserves further investigation.

In the rest of this paper we will consider the reduced actions $J_{j}$. Since the definitions for $J_{j}$ and $I_{j}$ only differ by a constant the results can be easily reformulated for the usual action $I_{j}$.

## 4 Picard-Fuchs Equations

Using ideas from de Rham theory, the task of calculating the actions can be formulated in terms of ordinary linear differential equations. It is numerically much less expensive to integrate these equations than to compute the integrals. More important, this method
provides insight into the analytic nature of actions and frequencies of the system, in particular their behavior near critical tori. It also provides a way to use analytic continuation in the complex in order to define global actions.

The basic observation is that all actions $J_{j}$ in (35), as well as their derivatives with respect to the integration constants, are hyperelliptic integrals of the second kind, and that the space of Abelian differentials of the second kind, modulo total differentials, has dimension $2 g=2 n$. As a possible basis we can take the action differential

$$
\begin{equation*}
\beta=\frac{Q(z)}{w} \mathrm{~d} z \tag{37}
\end{equation*}
$$

$2 \pi J=\oint \beta$, and its successive derivatives with respect to one root, e.g. $q_{1}$. At this point it is important that the integration paths are actually the cycles $c_{j}$ and not the cycles $\gamma_{j}$ related to the local actions. Denoting the partial derivatives of $\beta$ with respect to $q_{1}$ by $\beta^{\prime}, \beta^{\prime \prime}, \ldots, \beta^{(k)}, \ldots, \beta^{(2 n)}$, and using the abbreviations $X=z-q_{1}, Y=\left(z-q_{2}\right) \ldots\left(z-q_{n}\right)$, we have

$$
\begin{align*}
\beta & =Q(z) \frac{\mathrm{d} z}{w}=X Y \frac{\mathrm{~d} z}{w} \\
\beta^{\prime} & =-\frac{1}{2} Y \frac{\mathrm{~d} z}{w} \\
\beta^{\prime \prime} & =-\frac{1}{4} \frac{Y}{X} \frac{\mathrm{~d} z}{w} \\
& \vdots  \tag{38}\\
\beta^{(k)} & =c_{k} \frac{Y}{X^{k-1}} \frac{\mathrm{~d} z}{w} \\
& \vdots \\
\beta^{(2 n)} & =c_{2 n} \frac{Y}{X^{2 n-1}} \frac{\mathrm{~d} z}{w}
\end{align*}
$$

where

$$
\begin{equation*}
c_{k}=-\frac{1}{2}\left(+\frac{1}{2}\right)\left(+\frac{3}{2}\right) \ldots\left(\frac{2 k-3}{2}\right)=-\frac{(2 k-3)!!}{2^{k}}, \quad c_{0}=1 \tag{39}
\end{equation*}
$$

From de Rham theory it follows that any $2 n+1$ Abelian differentials of the second kind on the curve $\Gamma$ of genus $n$ are linearly dependent modulo total derivatives since the first cohomology of $\Gamma$ has dimension $2 n$. Thus there exist coefficients $b_{0}, \ldots, b_{2 n}$ depending on $q_{j}\left(\right.$ and $\left.a_{\nu}\right)$ such that

$$
\begin{equation*}
\sum_{s=0}^{2 n} b_{s} \beta^{(s)}=\mathrm{d} F \tag{40}
\end{equation*}
$$

for some meromorphic function $F$ on $\Gamma$. Using the fact that the actions are periods of $\beta$ along closed loops $c$ on $\Gamma, 2 \pi J=\oint_{c} \beta$, we integrate this relation along $c$ and obtain a linear differential equation of order $2 n$ for $J$ as a function of $q_{1}$,

$$
\begin{equation*}
\sum_{s=0}^{2 n} b_{s} J^{(s)}=0 \tag{41}
\end{equation*}
$$

Such equations are called Picard-Fuchs equations.
The main problem is now to find some $F$ and coefficients $b_{0}, \ldots, b_{2 n}$ that satisfy (40). The form of the derivatives $\beta^{(k)}$, Eq. (38), suggests the following form for $F$ :

$$
F=\frac{Y w}{X^{2 n-1}}
$$

Indeed, with this choice

$$
\mathrm{d} F=\frac{w \mathrm{~d} Y}{X^{2 n-1}}+\frac{Y \mathrm{~d} w^{2}}{2 w X^{2 n-1}}-(2 n-1) \frac{Y w \mathrm{~d} z}{X^{2 n}}
$$

Rewriting the last term with

$$
\frac{Y w}{X^{2 n}}=\frac{Y w^{2}}{X^{2 n} w}=-\frac{X Y^{2} A(z)}{X^{2 n} w}=-\frac{Y^{2} A(z)}{X^{2 n-1} w}
$$

we have

$$
\mathrm{d} F=Y \frac{-X A(z) Y^{\prime}+\frac{1}{2} R^{\prime}+(2 n-1) Y A(z)}{X^{2 n-1} w} \mathrm{~d} z
$$

where $Y^{\prime}$ and $R^{\prime}$ denote derivatives with respect to $z$. In order to solve (40) we rewrite this as a sum of powers of $X$ similar to those in (38). Thus we consider the polynomial of degree $2 n$ in the numerator,

$$
\begin{equation*}
P(z)=-X A(z) Y^{\prime}+\frac{1}{2} R^{\prime}+(2 n-1) Y A(z) \tag{42}
\end{equation*}
$$

and its Taylor expansion at $z=q_{1}$,

$$
\begin{equation*}
P(z)=\sum_{s=0}^{2 n} p_{s}\left(z-q_{1}\right)^{s}=\sum_{s=0}^{2 n} p_{s} X^{s}=P_{1}(X) \tag{43}
\end{equation*}
$$

where $P_{1}(X):=P\left(X+q_{1}\right)$ is the polynomial $P$ shifted by $q_{1}$. Now we have

$$
\begin{aligned}
\mathrm{d} F & =\frac{Y \sum_{s=0}^{2 n} p_{s} X^{s}}{w X^{2 n-1}} \mathrm{~d} z \\
& =p_{0} \frac{Y \mathrm{~d} z}{w X^{2 n-1}}+p_{1} \frac{Y \mathrm{~d} z}{w X^{2 n-2}}+\cdots+p_{2 n-1} \frac{Y \mathrm{~d} z}{w}+p_{2 n} \frac{X Y \mathrm{~d} z}{w} \\
& =b_{2 n} \beta^{(2 n)}+b_{2 n-1} \beta^{(2 n-1)}+\cdots+b_{1} \beta^{\prime}+b_{0} \beta
\end{aligned}
$$

so that

$$
\begin{equation*}
b_{2 n-s}=\frac{p_{s}}{c_{2 n-s}}=\left.\frac{1}{s!c_{2 n-s}} \frac{\mathrm{~d}^{s} P(z)}{\mathrm{d} z^{s}}\right|_{z=q_{1}} \tag{44}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
b_{2 n}=\frac{p_{0}}{c_{2 n}}=\frac{P_{1}(0)}{c_{2 n}}=\frac{P\left(q_{1}\right)}{c_{2 n}}=\frac{\left(-\frac{1}{2}+(2 n-1)\right) \prod_{s=2}^{n}\left(q_{1}-q_{s}\right) \prod_{\nu=0}^{n}\left(q_{1}-a_{\nu}\right)}{c_{2 n}} \tag{45}
\end{equation*}
$$

Dividing by $b_{2 n}$ we get the normalized form of the Picard-Fuchs equation

$$
\begin{equation*}
J^{(2 n)}+\sum_{s=0}^{2 n-1} \frac{b_{s}}{b_{2 n}} J^{(s)}=0 \tag{46}
\end{equation*}
$$

Notice that the singularities of this equation correspond to the cases when either $q_{1}=q_{j}$ $(j=2, \ldots, n)$ or $q_{1}=a_{\nu}$, i.e. when the curve $\Gamma$ becomes singular. Indeed, as one can see from the previous formulas all the coefficients $p_{0}, \ldots, p_{2 n}$ and therefore $b_{0}, \ldots, b_{2 n}$ are polynomials in $q_{1}, \ldots, q_{n}$. It is obvious what the corresponding equations in the variables $q_{2}, \ldots, q_{n}$ are.

Coefficients for the Picard-Fuchs equation of the differential $\mathrm{d} z / w$ have been given by Schlesinger in his classical book on linear differential equations [11]. Our formula can be considered as the generalization of those results to the differentials which determine the action in the Neumann system.

An important benefit of the linear differential Picard-Fuchs equations is that they give information on the nature of singularities of their solutions, i. e., of the actions and their derivatives. The singular points of the equation in the variable $q_{j}$ are at $q_{j}=a_{\nu}, q_{j}=q_{k}$ $(k \neq j)$ and $q_{j}=\infty$; they are all found to be regular singular points. A linear differential equation with this property is called Fuchsian; from the general theory if follows that Picard-Fuchs equations are always Fuchsian. The indicial equations can be solved for the leading indices of independent fundamental solutions.

To summarize the results of this section we formulate
Theorem 2. The Picard-Fuchs equation for the actions of the Neumann system with respect to the independent variable $q_{i}$ has the form

$$
\sum_{s=0}^{2 n} \frac{b_{s}}{b_{2 n}} J^{(s)}=0
$$

where the coefficients $b_{s}$ are given by formulas (44) with $q_{1}$ replaced by $q_{i}$, (42), and (39), where $X=z-q_{i}$ and $Y=Q / X, R=-Q A$.

## $4.1 \quad \mathrm{n}=1$

Consider the Picard-Fuchs equation for the elliptic case $n=1$. In this case $b_{1}=0$ and $A(q)=\left(q-a_{0}\right)\left(q-a_{1}\right)$ so that

$$
\begin{equation*}
4\left(q-a_{0}\right)\left(q-a_{1}\right) J^{\prime \prime}+J=0 \tag{47}
\end{equation*}
$$

where the prime denotes differentiation with respect to $q$. This second order ode has singularities at the three points $q=s_{i},\left(s_{0}, s_{1}, s_{2}\right)=\left(a_{0}, a_{1}, \infty\right)$. The singularity at $\infty$ becomes obvious with the substitution $q \rightarrow 1 / t$, which turns (47) into

$$
\begin{equation*}
\frac{\mathrm{d}^{2} J}{\mathrm{~d} t^{2}}+\frac{2}{t} \frac{\mathrm{~d} J}{\mathrm{~d} t}+\frac{1}{4 t^{2}\left(1-a_{0} t\right)\left(1-a_{1} t\right)} J=0 \tag{48}
\end{equation*}
$$

This is a particular case of the classical hypergeometric equations, see e.g. [14]. The indicial equation is $\alpha(\alpha-1)=0$ at both finite singular points, and $\left(\alpha+\frac{1}{2}\right)^{2}=0$ at $\infty$. In terms of the Riemann P symbol we thus have

$$
J(q)=P\left\{\begin{array}{cccc}
a_{0} & a_{1} & \infty  \tag{49}\\
0 & 0 & -\frac{1}{2} & q \\
1 & 1 & -\frac{1}{2}
\end{array}\right\}=F\left(-\frac{1}{2},-\frac{1}{2} ; 0 ;\left(q-a_{0}\right) /\left(a_{1}-a_{0}\right)\right)
$$

where $F$ is the hypergeometric function [14]. The indices tell us [4] that one regular solution in the neighborhood of each finite singular point $s_{i}=a_{0}$ or $a_{1}$ is analytic with a linear leading term,

$$
\begin{equation*}
f_{1}\left(q ; s_{i}\right)=q-s_{i}+O\left(\left(q-s_{i}\right)^{2}\right) \tag{50}
\end{equation*}
$$

and that an independent solution can be given in the form

$$
\begin{equation*}
f_{2}\left(q ; s_{i}\right)=1+O\left(q-s_{i}\right)+c_{i} f_{1}\left(q ; s_{i}\right) \ln \left(q-s_{i}\right) \tag{51}
\end{equation*}
$$

where the constant $c_{i}$ in this case turns out to be $\pm 1 / 4$ for $s_{i}=a_{0,1}$.

## $4.2 \mathrm{n}=2$

Evaluating the general formulas for the case $n=2$, we obtain the Picard-Fuchs equation

$$
\begin{align*}
0= & 4 A\left(q_{1}\right) Y\left(q_{1}\right) J^{\prime \prime \prime \prime} \\
& +\left(4 q_{1}^{2}\left(7 q_{1}-6 q_{2}\right)-4 q_{1}\left(5 q_{1}-4 q_{2}\right) T+4\left(3 q_{1}-2 q_{2}\right) S-4 P\right) J^{\prime \prime \prime}  \tag{52}\\
& +\left(9 q_{1}\left(4 q_{1}-3 q_{2}\right)-3\left(5 q_{1}-3 q_{2}\right) T+3 S\right) J^{\prime \prime}+3\left(q_{1}-q_{2}\right) J^{\prime}+\frac{3}{4} J=0
\end{align*}
$$

where the prime denotes differentiation with respect to $q_{1}$, and $T, S, P$ are the symmetric functions of the fixed zeros $a_{\nu}$,

$$
\begin{align*}
T & =a_{0}+a_{1}+a_{2} \\
S & =a_{0} a_{1}+a_{1} a_{2}+a_{2} a_{0}  \tag{53}\\
P & =a_{0} a_{1} a_{2}
\end{align*}
$$

The corresponding equation for the derivatives with respect to $q_{2}$ is obtained by exchanging $q_{1}$ and $q_{2}$.

The coefficient of the second highest derivative in the normalized Picard-Fuchs equation has the following simple partial fraction decomposition

$$
\begin{equation*}
\frac{b_{2 n-1}}{b_{2 n}}=\frac{2}{q_{1}-a_{0}}+\frac{2}{q_{1}-a_{1}}+\frac{2}{q_{1}-a_{2}}+\frac{1}{q_{1}-q_{2}} \tag{54}
\end{equation*}
$$

This can be used to read off the indices of the finite regular singular points as long as $q_{2}$ is not equal to one of the $a_{\nu}$. They are $0,1,1,2$ for $q_{1}=a_{\nu}, 0,1,2,2$ for $q_{1}=q_{2}$, and $-1 / 2,-1 / 2,1 / 2,3 / 2$ at infinity. When $q_{2}$ is confluent with one of the $a_{\nu}$, the indices at $a_{\nu}=q_{2}$ are $0, \frac{1}{2}, 1, \frac{3}{2}$.

### 4.3 General $n$

In the general case it is not obvious whether it is possible (or useful) to obtain more explicit expressions for the $b_{s}$ than (44). $b_{2 n-1} / b_{2 n}$ is interesting because it determines the indices at the regular singular points when there is just one double root in $R(z)$. Using $b_{2 n-1} / b_{2 n}=\left.(2 n-3 / 2)(\log P)^{\prime}\right|_{z=q_{1}}$ a simple calculation shows that

$$
\frac{b_{2 n-1}}{b_{2 n}}=(2 n-3) \sum_{j=2}^{n} \frac{1}{q_{1}-q_{j}}+(2 n-2) \sum_{r=0}^{n} \frac{1}{q_{1}-a_{\nu}}
$$

The fact that the coefficients in the partial fraction decomposition of $b_{2 n-1} / b_{2 n}$ are either $p_{c}=2 n-2$ or $p_{c}=2 n-3$ determines the indices at the corresponding regular singular points as the solutions of the equation $\alpha(\alpha-1) \ldots(\alpha-2 n+2)\left(\alpha-2 n+1+p_{c}\right)$; hence they are $0,1, \ldots, 2 n-2$ and 1 or 2 , respectively. At infinity the indices are $-\frac{1}{2},-\frac{1}{2}, \frac{1}{2}, \frac{3}{2}, \frac{5}{2}, \ldots, \frac{2 n+1}{2}$. One can also verify the following expressions for the coefficients $b_{0}$ and $b_{1}$ :

$$
\begin{align*}
& b_{1}=-2 n q_{1}+2 \sum_{j=1}^{n} q_{j}  \tag{55}\\
& b_{0}=-\frac{1}{2} \tag{56}
\end{align*}
$$

For nondegenerate potentials with $a_{\mu} \neq a_{\nu}, n>3$, the discriminant $b_{2 n}$ can at most have a triple zero for real motions when $a_{j}=q_{j-1}=q_{j}=q_{j+1}$. The indices are then $0,1, \ldots, 2 n-4$ and $-\frac{1}{2}, \frac{1}{2}, 2$. There are two ways in which double zeros can occur in $b_{2 n}$. When $q_{j}=q_{j+1}=a_{\nu}$, the indices become $0,1, \ldots, 2 n-3$ and $\frac{1}{2}, \frac{3}{2}$. In case of a double zero of the type $q_{j-1}=q_{j}=q_{j+1}$, the indices are $0,1, \ldots, 2 n-3$ and $\frac{1}{2}, \frac{5}{2}$.

## 5 Gauß-Manin Equations

Instead of constructing a basis from successive derivatives of the action differential $\beta$, a standard basis of Abelian differentials such as

$$
\begin{equation*}
\left(\alpha_{0}, \alpha_{1}, \ldots, \alpha_{2 n-1}\right)=\left(\frac{\mathrm{d} z}{w}, \frac{z \mathrm{~d} z}{w}, \ldots, \frac{z^{2 n-1} d z}{w}\right) \tag{57}
\end{equation*}
$$

may be taken to express $\beta$ and any of its derivatives, up to total derivatives $\mathrm{d} F$, as linear combinations thereof.

The dependence of the integrals $K_{i}:=\oint \alpha_{i}$ on $\boldsymbol{q}=\left(q_{1}, \ldots, q_{n}\right)^{t}$ is governed by equations of the type

$$
\begin{equation*}
\frac{\mathrm{d} \boldsymbol{K}}{\mathrm{~d} q_{i}}=\boldsymbol{M}^{(i)} \boldsymbol{K}, \quad i=1, \ldots, n \tag{58}
\end{equation*}
$$

where $\boldsymbol{K}=\left(K_{0}, K_{1}, \ldots, K_{2 n-1}\right)^{t}$ and the $2 n \times 2 n$ matrices $\boldsymbol{M}^{(i)}$ are found by decomposition of $\partial \alpha_{k} / \partial q_{i}$ in terms of $\alpha_{l}$ and suitable total derivatives. Such equations are called Gauß-Manin equations. ¿From a knowledge of $K_{i}\left(q_{1}, \ldots, q_{n}\right)$ the action is obtained with
the linear relation $\beta=\alpha_{n}+2 \sum_{i=0}^{n-1} \eta_{i} \alpha_{i}$. The simplicity of this relation is the main advantage of this basis.

Notice one important difference between Picard-Fuchs and Gauß-Manin equations. The latter define $n$ flows with respect to the "times" $q_{i}$ on the $2 n$ dimensional space of Abelian integrals. In the $\alpha$-basis each flow advances the same $2 n$ quantities $K_{i}(\boldsymbol{q})$. The individual Picard-Fuchs equations may of course also be written as first order systems with $2 n$ components, namely, $J$ and its first $2 n-1$ derivatives with respect to $q_{i}$; but these components are different for each flow. This has implications for the numerical evaluation of the equations. If we want to integrate, say, from $\boldsymbol{q}=\boldsymbol{q}^{i}$ to $\boldsymbol{q}^{f}$, the GaußManin equations for the directions $q_{i}$ and $q_{j}$ may be combined to follow a direct route. In the case of Picard-Fuchs equations, on the other hand, integration along the $q_{i}$ direction must be followed by a linear but $\boldsymbol{q}$ dependent transformation between the respective sets of $I$ derivatives before the integration can proceed in the $q_{j}$ direction. We find that this is a severe limitation for their practical use. For example, in order to obtain the actions for an energy surface $\mathrm{d} \eta_{1}=0$ in the case $n=2$, it is by far simpler to use the Gauß-Manin equations, in a combination $\mathrm{d} q_{1}=-\mathrm{d} q_{2}$.

The $n$ flows given by the solution of the linear $\boldsymbol{q}$ dependent systems (58) must commute, otherwise the solution $K(\boldsymbol{q})$ would not exist. Therefore the compatibility relations

$$
\begin{equation*}
\frac{\partial \boldsymbol{M}^{(i)}}{\partial q_{j}}-\frac{\partial \boldsymbol{M}^{(j)}}{\partial q_{i}}+\left[\boldsymbol{M}^{(i)}, \boldsymbol{M}^{(j)}\right]=0 \tag{59}
\end{equation*}
$$

must be satisfied.
The Gauß-Manin equation for Neumann systems of arbitrary $n$ can be given explicitly once the following general result on Abelian integrals of the second kind is established:

Theorem 3. Let $P(z)=\prod_{j=1}^{2 n}\left(z-z_{j}\right)=\sum_{i=0}^{2 n}(-1)^{i} \sigma_{i} z^{2 n-i}$ be a polynomial of degree $2 n$, all zeros $z_{j}$ different, and consider the Riemann surface $w^{2}=(q-z) P(z)$. Let $\left(K_{0}, \ldots, K_{2 n-1}\right)=\boldsymbol{K}=\boldsymbol{K}(q)$ be integrals over the same closed loop, of the standard basis of Abelian differentials of the second kind, $K_{i}=\oint \alpha_{i}=\oint z^{i} d z / w$. Then the following ordinary differential equation holds:

$$
\begin{equation*}
\frac{d \boldsymbol{K}}{d q}=\boldsymbol{M} \boldsymbol{K} \quad \text { with } \quad \boldsymbol{M}=\sum_{k=1}^{2 n} \frac{\boldsymbol{M}_{k}}{q-z_{k}} \quad \text { and } \quad \boldsymbol{M}_{k}=-\frac{\boldsymbol{l}_{k} \otimes \boldsymbol{r}_{k}}{2 D_{k}} \tag{60}
\end{equation*}
$$

where the partial discriminant $D_{k}$ is

$$
\begin{equation*}
D_{k}=\left.\frac{P(z)}{z-z_{k}}\right|_{z=z_{k}}=\prod_{j \neq k}\left(z_{k}-z_{j}\right)=: \sum_{i=0}^{2 n-1}(-1)^{i} \sigma_{i}^{(k)} z_{k}^{2 n-1-i} \tag{61}
\end{equation*}
$$

and the matrix $\boldsymbol{M}_{k}$ derives from the tensor product of the two vectors $\boldsymbol{l}_{k}$ and $\boldsymbol{r}_{k}$, with components $(m=1, \ldots, 2 n)$

$$
\begin{equation*}
\boldsymbol{l}_{k}=\left(1, z_{k}, z_{k}^{2}, \ldots, z_{k}^{2 n-1}\right)^{t} \equiv\left(l_{k 1}, \ldots, l_{k, 2 n}\right)^{t} \tag{62}
\end{equation*}
$$

and

$$
\begin{equation*}
\boldsymbol{r}_{k}=\left(r_{k 1}, \ldots, r_{k, 2 n}\right)^{t}, \quad r_{k m}=(-1)^{m-1}\left[m \sigma_{2 n-1-m}^{(k)} z_{k}+(m-1) \sigma_{2 n-m}^{(k)}\right] . \tag{63}
\end{equation*}
$$

The $\sigma_{i}^{(k)}$ are the symmetric functions of the zeros of $P(z) /\left(z-z_{k}\right)$, with $\sigma_{-1}^{(k)}:=0$. The matrices $\boldsymbol{M}_{k}$ satisfy

$$
\begin{equation*}
M_{k}^{2}=0 . \tag{64}
\end{equation*}
$$

Proof. The task is to find a decomposition of $\mathrm{d} \alpha_{i} / \mathrm{d} q$ into the basic differentials $\alpha_{j}$ and a total derivative $\mathrm{d} F$ (the same for all $i$ ),

$$
\begin{equation*}
\frac{\mathrm{d} \alpha_{i}}{\mathrm{~d} q}=\sum_{j=0}^{2 n-1} b_{i j} \alpha_{j}+b_{i, 2 n} \mathrm{~d} F \quad(i=0, \ldots, 2 n-1) . \tag{65}
\end{equation*}
$$

Once the $b_{i j}$ have been determined, they constitute the elements of the matrix $M$. Now, the decomposition (65) can be achieved with $F=w /(q-z)$ and

$$
\begin{equation*}
\mathrm{d} F=\frac{\mathrm{d} F}{\mathrm{~d} z} \mathrm{~d} z=\left(\frac{1}{2} \frac{w}{(q-z)^{2}}+\frac{P^{\prime}(z)}{2 w}\right) \mathrm{d} z=\frac{P+(q-z) P^{\prime}}{2(q-z) w} \mathrm{~d} z . \tag{66}
\end{equation*}
$$

Differentiating $\alpha_{i}=z^{i} \mathrm{~d} z / w$ with respect to $q$, we obtain $-z^{i} \mathrm{~d} z /(2(q-z) w)$; hence (65) becomes

$$
\begin{equation*}
-z^{i}=\sum_{j=0}^{2 n-1} 2 b_{i j}(q-z) z^{j}+b_{i, 2 n}\left(P+(q-z) P^{\prime}\right) . \tag{67}
\end{equation*}
$$

The coefficients $b_{i j}$ for fixed $i$ are determined by comparison of like powers $z^{j}$ :

$$
\left(\begin{array}{cccccc}
q & 0 & & & \cdots & -\left(q \cdot \sigma_{2 n-1}-1 \cdot \sigma_{2 n}\right)  \tag{68}\\
-1 & q & 0 & & \cdots & +\left(2 q \cdot \sigma_{2 n-2}+0 \cdot \sigma_{2 n-1}\right) \\
0 & -1 & q & 0 & \cdots & -\left(3 q \cdot \sigma_{2 n-3}+1 \cdot \sigma_{2 n-2}\right) \\
\vdots & & & & & \vdots \\
0 & \cdots & & -1 & q & +\left(2 n q \cdot \sigma_{0}+(2 n-2) \cdot \sigma_{1}\right) \\
0 & 0 & \cdots & 0 & -1 & -(2 n-1) \cdot \sigma_{0}
\end{array}\right)\left(\begin{array}{c}
2 b_{i 0} \\
2 b_{i 1} \\
2 b_{i 2} \\
\vdots \\
2 b_{i, 2 n-1} \\
b_{i, 2 n}
\end{array}\right)=-\boldsymbol{e}_{i},
$$

where $\boldsymbol{e}_{i}$ has components $e_{i j}=\delta_{i j}$. Let us denote the $(2 n+1) \times(2 n+1)$ matrix in this equation by $\boldsymbol{N}$; it is the same for all $i$. Multiplying with $\boldsymbol{N}^{-1}$ and taking the transposed equations for $i=0, \ldots, 2 n-1$, we see that the matrix $2 b_{i j}$, hence $2 M$, is the upper left $(2 n \times 2 n)$ minor of $-\boldsymbol{N}^{-1 t}$.

Developing with respect to the last column and sorting terms, the determinant is found to be $\operatorname{det}(\boldsymbol{N})=P(q)=\prod_{j=1}^{2 n}\left(q-z_{j}\right)$. This shows that $\boldsymbol{N}$ is singular at $q=z_{j}$ $(j=1, \ldots, 2 n)$, regular else. The inverse matrix $\boldsymbol{N}^{-1}$ must therefore have corresponding first order poles. To obtain $\boldsymbol{N}^{-1 t}$, we decompose $\boldsymbol{N}$ as

$$
\begin{equation*}
\boldsymbol{N}=\boldsymbol{A}+q \boldsymbol{B}, \tag{69}
\end{equation*}
$$

with $\boldsymbol{A}$ and $\boldsymbol{B}$ independent of $q$. The matrix $\boldsymbol{A}$ has full rank $2 n+1$ whereas $\boldsymbol{B}$ has corank 1 , its last row being identically zero. The matrices $\boldsymbol{N}_{k}:=\boldsymbol{A}+z_{k} \boldsymbol{B}$ are singular, with left kernels $\tilde{\boldsymbol{l}}_{k}=\left(1, z_{k}, z_{k}^{2}, \ldots, z_{k}^{2 n}\right)^{t}$ and right kernels $\tilde{\boldsymbol{r}}_{k}=\left(r_{k 1}, r_{k 2}, \ldots, r_{k, 2 n}, 1\right)^{t}$, where the ${\underset{\sim}{r}}_{k m}$ are given in (63). We also need the left and right null eigenvectors of $\boldsymbol{B}$; they are $\tilde{\boldsymbol{l}}_{2 n+1}=(0, \ldots, 0,1)^{t}$ and $\tilde{\boldsymbol{r}}_{2 n+1}=\left(\sigma_{2 n-1},-2 \sigma_{2 n-2}, 3 \sigma_{2 n-3}, \ldots,(2 n-1) \sigma_{1},-2 n \sigma_{0}, 1\right)^{t}$, respectively.

Let us now form the two $(2 n+1) \times(2 n+1)$ matrices $\boldsymbol{L}$ and $\boldsymbol{R}$ with elements $L_{k j}=\tilde{l}_{j k}$ and $R_{k j}=\tilde{r}_{j k}$. We assert that

$$
\begin{equation*}
\boldsymbol{L}^{t} \boldsymbol{N} \boldsymbol{R}=\boldsymbol{L}^{t}(\boldsymbol{A}+q \boldsymbol{B}) \boldsymbol{R}=\operatorname{diag}\left(\left(q-z_{1}\right) D_{1}, \ldots,\left(q-z_{2 n}\right) D_{2 n}, 1\right) \tag{70}
\end{equation*}
$$

To see this, consider first the elements $\left(\boldsymbol{L}^{t} \boldsymbol{N} \boldsymbol{R}\right)_{j k}=\tilde{\boldsymbol{l}}_{j}^{t} \boldsymbol{N} \tilde{\boldsymbol{r}}_{k}$ for $j, k \leq 2 n$. We decompose

$$
\begin{equation*}
\tilde{\boldsymbol{l}}_{j}^{t} \boldsymbol{N} \tilde{\boldsymbol{r}}_{k}=\tilde{\boldsymbol{l}}_{j}^{t}\left(\boldsymbol{N}-\boldsymbol{N}_{j}\right) \tilde{\boldsymbol{r}}_{k}+\tilde{\boldsymbol{l}}_{j}^{t} \boldsymbol{N}_{j} \tilde{\boldsymbol{r}}_{k}=\left(q-z_{j}\right) \tilde{\boldsymbol{l}}_{j}^{t} \boldsymbol{B} \tilde{\boldsymbol{r}}_{k} \tag{71}
\end{equation*}
$$

where the last equation follows because $\tilde{\boldsymbol{l}}_{j}$ is in the left kernel of $\boldsymbol{N}_{j}$. With a similar argument, we find that $\tilde{\boldsymbol{l}}_{j}^{t} \boldsymbol{N} \tilde{\boldsymbol{r}}_{k}=\left(q-z_{k}\right) \tilde{\boldsymbol{l}}_{j}^{t} \boldsymbol{B} \tilde{\boldsymbol{r}}_{k}$, and since $z_{j} \neq z_{k}$, we conclude that the nondiagonal elements vanish. As for the diagonal elements, it is not difficult to check that

$$
\begin{equation*}
\tilde{\boldsymbol{l}}_{k}^{t} \boldsymbol{B} \tilde{\boldsymbol{r}}_{k}=-\sigma_{2 n-1}+2 \sigma_{2 n-2} z_{k}-3 \sigma_{2 n-3} z_{k}^{2}+\ldots=\left.\frac{\mathrm{d} P}{\mathrm{~d} z}\right|_{z=z_{k}}=D_{k} \tag{72}
\end{equation*}
$$

Similarly, we have $\left(\boldsymbol{L}^{t} \boldsymbol{N} \boldsymbol{R}\right)_{2 n+1, k}=\left(q-z_{k}\right) \tilde{\boldsymbol{l}}_{2 n+1}^{t} \boldsymbol{B} \tilde{\boldsymbol{r}}_{k}=0$ for $k \leq 2 n$, because $\tilde{\boldsymbol{l}}_{2 n+1}^{t} \boldsymbol{B}=0$, and $\left(\boldsymbol{L}^{t} \boldsymbol{N} \boldsymbol{R}\right)_{j, 2 n+1}=0$ for $j \leq 2 n$ with analogous arguments. Finally, $\left(\boldsymbol{L}^{t} \boldsymbol{N} \boldsymbol{R}\right)_{2 n+1,2 n+1}=$ $\tilde{\boldsymbol{l}}_{2 n+1}^{t} \boldsymbol{A} \tilde{\boldsymbol{r}}_{2 n+1}=1$.

Having established the diagonalization (70), we invert it to obtain

$$
\begin{equation*}
\boldsymbol{N}^{-1 t}=\boldsymbol{L} \operatorname{diag}\left(\frac{1}{D_{1}\left(q-z_{1}\right)}, \ldots, \frac{1}{D_{2 n}\left(q-z_{2 n}\right)}, 1\right) \boldsymbol{R}^{t} \tag{73}
\end{equation*}
$$

and using the fact that the last column of $\boldsymbol{L}$ is $(0, \ldots, 0,1)^{t}$, we conclude that neither the $(2 n+1)$ st components of $\tilde{\boldsymbol{l}}_{k}$ and $\tilde{\boldsymbol{r}}_{k}$ nor the vectors $\tilde{\boldsymbol{l}}_{2 n+1}$ and $\tilde{\boldsymbol{r}}_{2 n+1}$ contribute to the upper left $(2 n \times 2 n)$ block of $\boldsymbol{N}^{-1 t}$. This establishes Eq. (60). The properties (64) of $\boldsymbol{M}_{k}$ follow from the orthogonality $\boldsymbol{r}_{k}^{t} \boldsymbol{l}_{k}=0$ which is easily checked with (62) and (63).

Application of this theorem to the Neumann system is straightforward. The $n$ equations (58) are obtained with $w^{2}=-A(z) Q(z)=:\left(q_{m}-z\right) P_{m}(z), m=1, \ldots, n$, where the set $\left\{z_{k}^{(m)}: k=1, \ldots, 2 n\right\}$ of zeros of $P_{m}(z)$ consists of the $a_{\nu}$ and all $q_{i}$ except $q_{m}$. This completes the derivation of the Gauß-Manin equations. One general statement can be made in view of the nature of the matrices $-2 D_{k}^{(m)} \boldsymbol{M}_{k}^{(m)}=\boldsymbol{l}_{k}^{(m)} \otimes \boldsymbol{r}_{k}^{(m)}$. The rank of $\boldsymbol{M}_{k}^{(m)}$ is at most 1 , and the $(2 n-1)$-dimensional orthogonal complement of the vector $\boldsymbol{r}_{k}^{(m)}$ is its eigenspace with eigenvalue 0 , i. e., if initial conditions are chosen from that space, the solutions are regular at $q_{m}=z_{k}^{(m)}$. From $\left(\boldsymbol{l}_{k}^{(m)}, \boldsymbol{r}_{k}^{(m)}\right)=0$ for all $k$ it follows
that the characteristic polynomial of each matrix $\boldsymbol{M}_{k}^{(m)}$ is $\lambda^{2 n}$, i. e., all $2 n$ eigenvalues are zero. Nevertheless, there exist only $2 n-1$ independent eigenvectors; the Jordan normal form contains a non-trivial block which means that one of the solutions has a logarithmic divergence at $q_{m}=z_{k}^{(m)}$.

## $5.1 \mathrm{n}=1$

The Gauß-Manin equation for the differentials $\alpha_{0}=\mathrm{d} z / w$ and $\alpha_{1}=z \mathrm{~d} z / w$ are easily derived with the help of

$$
\binom{\beta^{(0)}}{\beta^{(1)}}=\frac{1}{2}\left(\begin{array}{cc}
-2 q & 2  \tag{74}\\
-1 & 0
\end{array}\right)\binom{\alpha_{0}}{\alpha_{1}}
$$

compare with (38). For $n=1$ this gives no particular advantage over the Picard-Fuchs equation, but as an illustration of theorem 3 it is interesting to note that the equations take on the form

$$
\begin{equation*}
2\left(a_{1}-a_{0}\right) \frac{\mathrm{d}}{\mathrm{~d} q}\binom{\alpha_{0}}{\alpha_{1}}=\left[\frac{1}{q-a_{0}}\binom{1}{a_{0}} \otimes\binom{a_{0}}{-1}-\frac{1}{q-a_{1}}\binom{1}{a_{1}} \otimes\binom{a_{1}}{-1}\right]\binom{\alpha_{0}}{\alpha_{1}} \tag{75}
\end{equation*}
$$

## $5.2 \mathrm{n}=2$

The Gauß-Manin equations are obtained from theorem 3. The change of $K_{i}=\oint \alpha_{i}$ with $q_{1}$ is

$$
\begin{equation*}
\frac{\mathrm{d} K_{i}}{\mathrm{~d} q_{1}}=M_{i j}^{(1)} K_{j} \tag{76}
\end{equation*}
$$

with

$$
\begin{equation*}
\boldsymbol{M}^{(1)}=-\sum_{k=1}^{4} \frac{1}{2 D_{k}^{(1)}} \frac{\boldsymbol{l}_{k}^{(1)} \otimes \boldsymbol{r}_{k}^{(1)}}{q_{1}-z_{k}^{(1)}} \tag{77}
\end{equation*}
$$

The sum is over the $2 n=4$ singular points $z_{k}^{(1)}$ of the equation for $q_{1},\left(z_{1}^{(1)}, z_{2}^{(1)}, z_{3}^{(1)}, z_{4}^{(1)}\right)=$ $\left(a_{0}, a_{1}, a_{2}, q_{2}\right)$, and the constants $D_{k}^{(j)}$ are defined via

$$
\begin{equation*}
\frac{1}{D_{k}^{(j)}}=\operatorname{Res}_{z=z_{k}} \frac{z-q_{j}}{w^{2}} \tag{78}
\end{equation*}
$$

(The constants $D_{k}^{(2)}$ appear in the equation for $q_{2}$.) The vectors $\boldsymbol{l}_{k}^{(1)}$ and $\boldsymbol{r}_{k}^{(1)}$ are

$$
\begin{align*}
\boldsymbol{l}_{k}^{(1) t} & =\left(1, z_{k}^{(1)}, z_{k}^{(1) 2}, z_{k}^{(1) 3}\right)  \tag{79}\\
\boldsymbol{r}_{k}^{(1) t} & =\left(z_{k}^{(1)} \sigma_{2}^{(1 k)},-\left(2 z_{k}^{(1)} \sigma_{1}^{(1 k)}+\sigma_{2}^{(1 k)}\right), 3 z_{k}^{(1)}+2 \sigma_{1}^{(1 k)},-3\right)
\end{align*}
$$

with

$$
\begin{align*}
\sigma_{1}^{(j k)} & =\sum_{l \neq k} z_{l}^{(j)} \\
\sigma_{2}^{(j k)} & =\sum_{l<m} z_{l}^{(j)} z_{m}^{(j)} \quad(l \neq k, m \neq k) \tag{80}
\end{align*}
$$

Performing the sum in (77) explicitly, leads to $M_{i j}^{(1)}=-\left(1 / 2 D_{1}\right) m_{i j}^{(1)}$ with

$$
m_{i j}^{(1)}=\left(\begin{array}{cccc}
q_{1}^{2} q_{12}-q_{1} q_{12} T+q_{1} S & q_{1}^{2}-q_{1} A-B & q_{1}+2 A & -3  \tag{81}\\
q_{1} q_{2} S+q_{12} P & q_{1} m_{12}^{(1)} & q_{1} m_{13}^{(1)} & -3 q_{1} \\
q_{1} m_{21}^{(1)} & m_{21}^{(1)}-2 q_{1}^{2} B & q_{1}^{2} m_{13}^{(1)} & -3 q_{1}^{2} \\
q_{1}^{2} m_{21}^{(1)} & q_{1} m_{32}^{(1)} & m_{21}^{(1)}+3 q_{1}^{3} A-q_{1}^{2} B & -3 q_{1}^{3}
\end{array}\right),
$$

where in addition to (53) we used the abbreviations $q_{12}:=q_{1}-q_{2}, A=q_{2}+T, B=q_{2} T+S$. The corresponding expressions in the equation for $\mathrm{d} K_{i} / \mathrm{d} q_{2}$ are obtained with $q_{1} \leftrightarrow q_{2}$.

Remark We should mention that there is another natural set of differentials

$$
\varphi_{\nu}=\frac{\mathrm{d} z}{\left(z-z_{\nu}\right) \sqrt{R(z)}}, \quad \nu=0,1, \ldots, 2 n
$$

where $z_{\nu}$ are the zeros of $R(z): z_{k}=a_{k}$ for $k=0, \ldots, n, z_{k+n}=q_{k}$ for $k=1, \ldots, n$. Their sum is cohomological to zero,

$$
\sum_{\nu=0}^{2 n} \varphi_{\nu}=\mathrm{d} F, \quad F=-2 R(z)^{-1 / 2}
$$

but any $2 n$ of them form a basis in the first cohomology of the corresponding hyperelliptic curve. These differentials satisfy the following simple system

$$
\partial_{\mu} \varphi_{\nu}=\frac{1}{2} \frac{\varphi_{\nu}-\varphi_{\mu}}{z_{\nu}-z_{\mu}}, \quad \mu \neq \nu
$$

So, the corresponding Gauß-Manin equations for

$$
L_{\sigma}=\oint_{\gamma} \varphi_{\sigma}
$$

have a very simple form (which is actually a particular case of the Knizhnik-Zamolodchikov equation $[5]^{1}$ )

$$
\begin{align*}
& \partial_{\mu} L_{\sigma}=\frac{1}{2} \frac{L_{\sigma}-L_{\mu}}{z_{\sigma}-z_{\mu}}, \quad \mu \neq \sigma  \tag{82}\\
& \partial_{\sigma} L_{\sigma}=\partial_{\sigma}\left(-\sum_{\nu \neq \sigma} L_{\nu}\right)=-\frac{1}{2} \sum_{\nu \neq \sigma} \frac{L_{\sigma}-L_{\nu}}{z_{\sigma}-z_{\nu}} \tag{83}
\end{align*}
$$

Unfortunately in this basis it seems to be impossible to find simple explicit formulae for the action of the Neumann system, so the basis we used in the present paper looks more suitable from this point of view.

[^0]
## 6 Calculation of energy surfaces in action variables

Let us introduce the moment map $F: T^{*} S^{n} \rightarrow \mathbb{R}^{n}$ given by the integrals $F_{1}, \ldots, F_{n}$ of the Neumann system. We call the image of this map in $\mathbb{R}^{n}$ the moment space. The PicardFuchs and Gauß-Manin equations are equations on moment space. For each dimension in moment space there is one independent flow. For the Gauß-Manin equations the dependent variables are all the same for the different commuting flows.

In the Neumann system the moment space has the trivial topology of $\mathbb{R}^{n}$; however, the differential equations have singularities at the critical values of the moment map. From the explicit formula (35) we know that there exists a global continuous solution. To obtain this particular solution from integration of the differential equation we have to 1 ) understand the structure of the singularities, 2) find appropriate initial conditions, and 3) determine how to continue the solution across the singularities.

The bifurcation diagram is the set of critical values of the moment map. We have already seen that the singularities of the differential equations occur when the curve $\Gamma$ becomes singular, i. e., when there is a collision between the roots of $w^{2}$. This set is called the discriminant surface. The bifurcation diagram is contained in the discriminant surface; for the determination of actions we are only interested in that part of the discriminant surface whose preimage contains real tori.

The bifurcation diagram divides the moment space into regions of regularity which we call m-phases. Throughout a given m-phase, all preimages of the moment map are topologically equivalent. In particular, the numbers $(\rho, o, m)$ defined in the appendix are constant within an m-phase: $\rho$ is the number of rotational degrees of freedom; $o$ is the number of oscillatory modes that cross a hyperplane $x_{\nu}=0 ; m$ is the number of oscillatory modes without hyperplane crossing. $N=2^{m+1}$ is the number of disjoint tori in phase space.

The m-phases are most conveniently determined from the nature of the $u_{j}$-intervals [ $z_{2 j-1}, z_{2 j}$ ]. Consider Figure 2 and focus on the forbidden real intervals between the ovals. From each $a_{j}$ there extends exactly one such region, either to the left ( $l$ ) or to the right $(r)$. Moreover, the rightmost forbidden interval always extends to the right, otherwise there could not be one oval in each interval $\left[a_{j-1}, a_{j}\right]$, and hence no real motion. The sequence of forbidden intervals can therefore be encoded as a string of $n+1$ letters $l$ and $r$, starting with the letter for the interval $\left[z_{0}, z_{1}\right]$, and ending with the letter $r$. The total number of such strings is $2^{n}$.

For $n=1$ the only two possibilities are $\operatorname{rr}$ (oscillation) and $\operatorname{lr}$ (rotation). In general there are four different types of motion for the variable $u_{j}$, depending on the letters for the neighboring forbidden intervals:
$l r=: \mathrm{R} u_{j}$ is free to move from $a_{j-1}$ to $a_{j}$; this is a rotation (when $u_{j}$ reaches an endpoint of its interval, it changes to a different leaf of the covering of $S^{n}$ by the coordinate system).
$r r=: \mathrm{O}_{\mathrm{r}} u_{j}$ oscillates in the right part of its interval, between $q_{j}$ and $a_{j}$.
$l l=: \mathrm{O}_{l} u_{j}$ oscillates in the left part of its interval, between $a_{j-1}$ and $q_{j+1}$.
$r l=: \mathrm{O}_{\mathrm{m}} u_{j}$ oscillates in the interior part of its interval, between $q_{j}$ and $q_{j+1}$.
In the case $n=2$ there are four types of tori encoded by $l l r, l r r$, $r l r$, and $r r r$. The projections of Liouville-Arnold tori onto the sphere $\mathrm{S}^{2}$ are shown as dark patches in Figure 3 where again $\left(a_{0}, a_{1}, a_{2}\right)=(1,2,4)$. In detail, the four types of tori are:
$l r r=\mathrm{RO}_{\mathrm{r}} u_{1}$ varies in its full interval $\left[a_{0}, a_{1}\right]$, and with two sweeps back and forth completes one rotation about the $x_{2}$-axis; $u_{2}$ oscillates in $\left[q_{2}, a_{2}\right]$, but two of its periods are needed to complete one period of $x_{2}$ oscillations. The same projection on $\mathrm{S}^{2}$ corresponds to two tori in $\mathrm{T}^{*} \mathrm{~S}^{2}$ which differ in the sense of rotation. $(\rho, o, m)=(1,1,0), N=2$.
$r r r=\mathrm{O}_{\mathrm{r}} \mathrm{O}_{\mathrm{r}}$ Both variables oscillate, $u_{1}$ in $\left[q_{1}, a_{1}\right]$ and $u_{2}$ in $\left[q_{2}, a_{2}\right]$. Each of the two patches in the projection corresponds to one torus in $\mathrm{T}^{*} \mathrm{~S}^{2} .(\rho, o, m)=(2,2,0), N=2$.
$l l r=\mathrm{O}_{1} \mathrm{R}$ The motion type is similar to $\mathrm{RO}_{\mathrm{r}}$, but with the behavior of $u_{1}$ and $u_{2}$ interchanged. The rotation is now about the $x_{0}$-axis, and the oscillation in the $x_{0}$ coordinate. The projection corresponds to two tori in $\mathrm{T}^{*} \mathrm{~S}^{2} .(\rho, o, m)=(1,1,0), N=2$.
$r l r=\mathrm{O}_{\mathrm{m}} \mathrm{R}$ The oscillations of $u_{1}$ are restricted to the interval $\left[q_{1}, q_{2}\right]$ so that there are two patches in the projection, each corresponding to two tori in $\mathrm{T}^{*} \mathrm{~S}^{2}$ which differ in the sense of rotation about the $x_{0}$-axis. $(\rho, o, m)=(1,0,1), N=4$.

Figure 4 contains four different versions of m-phase diagrams which show how the type of motion depends on the values of the constants of motion. The plane of zeros $\left(q_{1}, q_{2}\right)$ gives the simplest picture, see Fig. 4a. The disadvantage of this set of constants of motion is that they are not smooth function on phase space. Only $\eta_{j}$ and $F_{\nu}$ are smooth functions on phase space, and therefore qualify as good coordinates for a bifurcation diagram. The transformation from these smooth constants to $\left(q_{1}, q_{2}\right)$ is singular when $q_{1}=q_{2}$, while the transformation to $\left(J_{1}, J_{2}\right)$ is singular, in addition, when $q_{j}=a_{\nu}$. In spite of these singularities, each set has its advantages. The $q_{j}$ are best suited to determine the PicardFuchs and Gauß-Manin equations; singularities in the real transformation are then less important because the equations are considered in the complex domain.

The Neumann system on the sphere $\mathrm{S}^{2}$ is related to a hyperelliptic curve $\Gamma$ of genus $g=2$ given by the equation $w^{2}=R(z)=-A(z) Q(z)$ with polynomials

$$
\begin{align*}
& Q(z)=\left(z-q_{1}\right)\left(z-q_{2}\right)=z^{2}+2 \eta_{1} z+2 \eta_{2}, \\
& A(z)=\left(z-a_{0}\right)\left(z-a_{1}\right)\left(z-a_{2}\right), \tag{84}
\end{align*}
$$

where

$$
\begin{equation*}
2 \eta_{1}=2 h-\left(a_{0}+a_{1}+a_{2}\right)=-\left(q_{1}+q_{2}\right) \quad \text { and } \quad 2 \eta_{2}=q_{1} q_{2} ; \tag{85}
\end{equation*}
$$

$\eta_{1}$ represents the energy constant $h$. Lines of constant energy $h$ are given as $q_{1}+q_{2}=$ const in Fig. 4a. Fig. 4b displays the same information in the plane of separation constants $\left(\eta_{1}, \eta_{2}\right)$; the correspondence is given by (85). Isoenergy lines are $\eta_{1}=$ const. By the linear transformation $2 \eta_{1}=F_{0}\left(a_{1}+a_{2}\right)+F_{1}\left(a_{2}+a_{0}\right)+F_{2}\left(a_{0}+a_{1}\right), 2 \eta_{2}=F_{0} a_{1} a_{2}+F_{1} a_{2} a_{0}+F_{2} a_{0} a_{1}$ the phase diagram is taken to the plane of the algebraic integrals ( $F_{0}, F_{1}, F_{2}$ ), see Fig. 4c.


Figure 3: Projections of invariant tori (dark) from the phase space $T^{*} S^{2}$ onto the sphere $S^{2}$. The boundaries (caustics) are given by coordinate lines of the elliptical spherical coordinate system. a) type $l r r=\mathrm{RO}_{\mathrm{r}} ; \mathrm{b}$ ) type $r r r=\mathrm{O}_{\mathrm{r}} \mathrm{O}_{\mathrm{r}} ;$ c) type $\left.l l r=\mathrm{O}_{1} \mathrm{R} ; \mathrm{d}\right)$ type $r l r=\mathrm{O}_{\mathrm{m}} \mathrm{R}$.


Figure 4: Phase diagrams with integrals a) $\left(q_{1}, q_{2}\right)$, b) $\left(\eta_{1}, \eta_{2}\right)$, c) $\left(F_{0}, F_{1}, F_{2}\right)$, and d) $\left(J_{1}, J_{2}\right)$. The actions are given for the symmetry reduced system. The correspondence between the four figures can be read off the points $\mathrm{A}, \ldots, \mathrm{F}$. The thin lines are isoenergy lines.

The eight points $\mathrm{A}, \ldots, \mathrm{H}$ help to see the connection between the different phase diagrams. Three of these points represent equilibria: A corresponds to the stable equilibria $\left(x_{0}, x_{1}, x_{2}\right)=( \pm 1,0,0)$ at minimum energy $h(\mathrm{~A})=a_{0} / 2$. B corresponds to the unstable elliptic-hyperbolic equilibria $\left(x_{0}, x_{1}, x_{2}\right)=(0, \pm 1,0)$ at energy $h(\mathrm{~B})=a_{1} / 2$, and includes motion along the circle $x_{2}=0$ which is asymptotic to the equilibrium points. C corresponds to the unstable equilibria $\left(x_{0}, x_{1}, x_{2}\right)=(0,0, \pm 1)$ at energy $h(\mathrm{C})=a_{2} / 2$. The preimage in phase space is a separatrix which projects onto the entire configuration sphere.

Equilibria mark transitions in the topology of isoenergy surfaces in phase space. Increasing $h$ through $h(\mathrm{~A})$, the transition is from an empty set to two 3 -spheres $\mathrm{S}^{3}$, corresponding to the two accessible disks on the configuration sphere $\mathrm{S}^{2}$. At energy $h(\mathrm{~B})$, the accessible region on $\mathrm{S}^{2}$ becomes an annulus, and the two $\mathrm{S}^{3}$ merge to form an isoenergy surface of type $\mathrm{S}^{1} \times \mathrm{S}^{2}$. Finally, at energies above $h(\mathrm{C})$, the entire configuration sphere is accessible, and the isoenergy surface becomes the real projective space $\mathbb{R P}^{3}$.

The phase diagrams show that these transitions at A, B, C involve bifurcations of invariant tori: new types of tori appear as energy is increased. In addition, there occur bifurcations of critical tori at D and E . Increasing $h$ through $h(\mathrm{D})=a_{0}+a_{2}-a_{1}$, oscillation along the circle $x_{1}=0$ about the stable equilibrium points becomes unstable, and two stable rotations $u_{1}=q_{1}=q_{2}$ emerge which are images of each other under time reversal (pitchfork bifurcation). An inverse pitchfork bifurcation takes place when $h$ increases through $h(\mathrm{E})=a_{1}+a_{2}-a_{0}$ : the two stable rotations $u_{1}=q_{1}=q_{2}$ and the unstable rotation along the circle $x_{0}=0$ merge into one stable rotation. For the Picard-Fuchs equation the points D and E are also special because they do not have the generic indices as a result of a coalescence of singularities.

From the point of view of foliation of phase space by isoenergy surfaces and invariant tori, the points A, ..., E mark five topological transitions. The points F, G, H play no particular role in the organization of phase space; they merely serve for orientation in the different phase diagrams.

For arbitrary $n$ the singularity structure of moment space is quite similar. The general phase diagram for the Liouville-Arnold tori of the Neumann system is easily discussed in terms of the constants $q_{j}$. Its outer boundary in $\mathbb{R}^{n}\left(q_{1}, \ldots, q_{n}\right)$ is given by the $3 n-2$ planes $q_{j}=a_{j}(j=1, \ldots, n), q_{j}=a_{j-2},(j=2, \ldots, n)$, and $q_{j}=q_{j+1}(j=1, \ldots, n-1)$. They correspond to critical tori of elliptic type. The inner planes $q_{j}=a_{j-1}(j=1, \ldots, n)$ correspond to hyperbolic critical tori and their associated separatrices. In the following discussion of the change of tori at separatrices we assume that all the $q_{j}$ except the one under consideration are fixed at non-critical values.

1. When $q_{j}=a_{j}$ for any $j$, the variable $u_{j}$ cannot vary, $u_{j}=a_{j}$ and $p_{j}=0$. For $q_{j}$ slightly lower, $u_{j}$ oscillates between $q_{j}$ and $a_{j}$. This shows that the bifurcation involves an elliptic critical torus.
2. Similarly, when $q_{j}=a_{j-2}$ for any $j$, the coordinate $u_{j-1}$ is fixed to $a_{j-2}$, and $p_{j-1}=0$. For $q_{j}$ slightly larger, $u_{j-1}$ oscillates between $q_{j}$ and $a_{j-2}$. Again, the critical torus is elliptic.
3. When $q_{j}=q_{j+1}$ (which may only happen between $a_{j-1}$ and $a_{j}$ ), then $u_{j}=q_{j}$ is fixed, $p_{j}=0$. For $q_{j}$ slightly lower, $u_{j}$ oscillates between $q_{j}$ and $q_{j+1}$. Also here, the critical torus is elliptic.
4. When $q_{j}=a_{j-1}$, the situation is different. For $q_{j}$ slightly above $a_{j-1}$, the variable $u_{j}$ has $q_{j}$ as its left endpoint and varies from there to either $a_{j}$ or $q_{j+1}$, whichever is smaller. For $q_{j}$ slightly below $a_{j-1}$, the variable $u_{j-1}$ has $q_{j}$ as its right endpoint and varies from there to either $a_{j-2}$ or $q_{j-1}$, whichever is larger. Thus there is a transition between different types of tori; at the transition point the tori coalesce and form a separatrix.

The point $\left(q_{1}, q_{2}, \ldots, q_{n}\right)=\left(a_{0}, a_{1}, \ldots, a_{n-1}\right)$ may be viewed as the organizing center of the bifurcation diagram. It is the intersection point of all inner bifurcation planes, hence it is corner point to $2^{n}$ regions with different types of tori. Each region is locally bounded by the $n$ planes $q_{j}=a_{j-1}$, and extends without further internal boundaries to the outer boundary (or to $-\infty$ in the variable $q_{1}$ ). Again this shows that there are exactly $2^{n}$ topologically different types of tori. The codimension $n$ corners of the regions that are not contained in any of the hyperplanes $q_{j}=q_{j+1}$ correspond to stable isolated periodic orbits. All other isolated periodic orbits are unstable.

Let us now comment on the computation of actions based on the Gauß-Manin equations. We start at A, the simplest point in the bifurcation diagram. For arbitrary $n$ it is given by $q_{i}=a_{i}$. The corresponding motion is the stable equilibrium point at the minimum of the potential. Accordingly all actions are zero at A. To find initial conditions for the Gauß-Manin equations we have to find the value of the integrals $K_{i}^{j}$. Since at A the curve has $n$ double roots the genus of the curve drops to 0 , and the integrals can be evaluated using residue calculus. To calculate the action $J_{j}$ we have to chose the cycle $c_{j}$ which can be contracted to a point at A such that

$$
\begin{equation*}
\left.K_{l}^{j}\right|_{\mathrm{A}}=\left.\oint_{c_{j}} \frac{z^{l}}{w}\right|_{\mathrm{A}} \mathrm{~d} z=2 \pi i \operatorname{Res}_{z=a_{j}} \frac{z^{l} \sqrt{z-a_{0}}}{A(z)}=2 \pi \frac{a_{j}^{l}}{\sqrt{a_{j}-a_{0}} \prod_{s}^{\prime}\left(a_{j}-a_{s}\right)}, \tag{86}
\end{equation*}
$$

where $\Pi^{\prime}$ runs from $s=1$ to $n$ excluding $s=j$.
For each $j$ this initial condition corresponds to a solution which is regular at the singular point A . The singularity in $\mathbf{M}$ at A is canceled by $\mathbf{K}^{j}$; yet, the first step of the integration cannot be done numerically. Since it is not obvious that our solution must be analytic we cannot use the differential equation to find the Taylor expansion of the solution at A. However, we can again use the integral solution and calculate its series at A . We choose to step away from A in the direction $\Delta \mathbf{q}=(-1, \ldots,-1)$, i.e. we want to calculate $\mathbf{K}^{j}\left(\mathbf{q}_{\mathrm{A}}-\epsilon \Delta \mathbf{q}\right)$. The first term in the expansion is already known, the second term is given by

$$
\begin{align*}
\left.\sum_{k} \frac{\partial}{\partial q_{k}} K_{l}^{j}\right|_{\mathrm{A}} & =2 \pi i \underset{z=a_{j}}{\operatorname{Res}}-\frac{z^{l}}{2 \sqrt{a_{0}-z} \prod_{s}\left(z-a_{s}\right)} \sum_{k} \frac{1}{z-a_{k}}  \tag{87}\\
& =\frac{\pi}{2} a_{j}^{l-1} \frac{2 l a_{0}-(2 l-1) a_{j}}{{\sqrt{a_{j}-a_{0}}}^{3} \prod_{s}^{\prime}\left(a_{j}-a_{s}\right)}
\end{align*}
$$

Therefore we find to first order in $\epsilon$

$$
\begin{equation*}
K_{l}^{j}\left(\mathbf{q}_{\mathrm{A}}+\epsilon \Delta \mathbf{q}\right)=\left.K_{l}^{j}\right|_{\mathrm{A}}\left(1-\frac{\epsilon}{4}\left(\frac{2 l}{a_{j}}-\frac{1}{a_{j}-a_{0}}\right)\right)+O\left(\epsilon^{2}\right) \tag{88}
\end{equation*}
$$

From there it is straightforward to integrate the Gauß-Manin equations in the entire m-phase $r r \ldots r$. The procedure runs into numerical problems when it meets a bifurcation point whose preimage contains a hyperbolic object. The integrals $K_{i}^{j}$ then typically possess logarithmic singularities.

In principal we could choose a special initial point in every m-phase, compute the integrals $K_{i}^{j}=\oint_{c_{j}} \alpha_{i}$ for that point, and start the integration from there. However, it is possible to cross from one m-phase to a neighboring one by encircling the singularity in the complex. This will give the analytic continuation of the solution chosen at A. We now show that the real part of this analytic continuation is unique and continuous. Therefore it is exactly the solution we are looking for when calculating the reduced actions $J_{j}$.

Consider the curve $w^{2}=R(u)$ in Fig. 2. We have the real cycles $c_{j}$ which map onto the intervals $\left[z_{2 j-1}, z_{2 j}\right]$ under the projection $(w, u) \rightarrow u$. To obtain a (non-canonical) basis of cycles we add the imaginary cycles $d_{j},(j=1, \ldots, n)$ which map onto the intervals $\left[z_{2 j-2}, z_{2 j-1}\right]$. We found that crossing an interior plane in the bifurcation diagram implies $q_{j}=a_{j-1}$. That means we take the two "endpoints" of a $d$-cycle and deform the adjacent cycles so as to exchange the endpoints. The resulting new cycles $\tilde{c}, \tilde{d}$, are then expressed in terms of the old. For example, for $j=2$ we find from elementary topological arguments that

$$
\begin{aligned}
& \tilde{c}_{1}=c_{1}+d_{2} \\
& \tilde{d}_{2}=d_{2} \\
& \tilde{c}_{2}=c_{2}+d_{2} .
\end{aligned}
$$

This shows that the real $c$-cycles are only changed by imaginary cycles. Since the periods of the forms over imaginary cycles are purely imaginary the real part of the analytic continuation around an interior plane of the bifurcation diagram is unique and well defined. In particular this means that when integrating a fixed linear combination of the GaußManin equations, corresponding to a direction $\Delta \mathbf{q}$ in moment space transversely crossing an interior plane, we can use the following procedure. Integrate towards the plane and stop a distance $\epsilon$ before reaching it. The solutions blow up logarithmically with $\epsilon$. Denote the current point by $\mathbf{q}_{s}-\epsilon \Delta \mathbf{q}$, where $\mathbf{q}_{s}$ is the point of intersection with the singular hyperplane. The solution at that point is $\mathbf{K}\left(\mathbf{q}_{s}-\epsilon \Delta \mathbf{q}\right)$. Then restart the integration on the other side of the singularity at $\mathbf{q}_{s}+\epsilon \Delta \mathbf{q}$ with the initial condition $\mathbf{K}\left(\mathbf{q}_{s}-\epsilon \Delta \mathbf{q}\right)$. Of course this introduces an additional error in the numerical solution, but it can be well controlled by $\epsilon$. To obtain a better initial value on the other side one would have to perform analytical continuation.

The energy surfaces in action space shown in Fig. 4d can therefore be obtained by the following procedure. For each action $J_{j}$ calculate the initial conditions $K_{j}^{i}\left(\mathbf{q}_{\mathrm{A}}+\epsilon \Delta \mathbf{q}\right)$ given by $(86,88)$. To find a particular value of the energy $h$ integrate the flow corresponding to $-\sum_{s=1}^{j} \mathbf{M}^{j}$ with this initial condition. When the desired value of $h$ is reached, use
$n-1$ flows that preserve the value of $h$ to generate the energy surface in action space. For $n=2$, only one such flow is needed, corresponding to $\mathbf{M}^{1}-\mathbf{M}^{2}$. The new initial condition is obtained from the previous integration. Whenever a codimension one singularity of the bifurcation diagram is approached, with any flow, use the above procedure to get across it.

## 7 Concluding remarks

Ideally, one would like to consider the dependence of the actions on any set of $n$ physical integration constants $k_{1}, \ldots, k_{n}$, and to have Picard-Fuchs equations with the derivatives $\mathrm{d}^{i} J / \mathrm{d} k_{j}^{i}$, or Gauß-Manin equations for $\mathrm{d} \boldsymbol{K} / \mathrm{d} k_{j}$. In principle this should be possible whenever the actions are complete Abelian integrals, and the derivatives of first or second kind. In practice, however, there tend to be complications when more than $n$ zeros of the polynomial $P(z)$ in theorem 3 depend on the parameters $k_{j}$. The Neumann system is special in that only $n$ branching points $q_{1}, \ldots, q_{n}$ of the Riemann surface $w^{2}=-A(z) Q(z)$ depend on the physical constants $\eta_{j}$, or $F_{\nu}$. The values $q_{j}$ can then be taken as integration constants, and only then do the Picard-Fuchs or Gauss-Manin equations exhibit the simplicity that we found.

However, if the constants $\eta_{j}$ are taken as variables of moment space, instead of the $q_{j}$, it becomes highly cumbersome to write down Gauß-Manin or Picard-Fuchs equations even for the Neumann system. The reason is the nonlinear relationship between the two sets (except for $\eta_{1}$ ). In other integrable systems with Riemann surfaces $w^{2}=P(z)$, it is usually not the case that the $n$ constants of integration are mapped to just $n$ zeros of the polynomial $P(z)$, and then things become much more complicated. The Kovalevskaya system, for instance, has a Riemann surface of genus 2 , but all five zeros change with the two integration constants. The complexity of the Gauß-Manin equations in terms of energy and separation constant as given in [2] is so enormous that their practical usefulness may be questioned. For the Neumann system, however, we have shown that simple and effective formulas for arbitrary $n$ exist.

Finally, we want to stress that there is more to this approach than its effectiveness for numerical computation: the Gauss-Manin equations contain all information about the derivatives of actions, in particular the frequencies. In this connection we refer to an interesting paper by Emil Horozov [3] where he applied Picard-Fuchs equations to the proof of nondegeneracy (for all non-critical tori) of the spherical pendulum. It would be interesting to use our results to investigate the nondegeneracy of the Neumann system.

## Acknowledgements

One of us (APV) is grateful to Vitaly Tarasov for the useful discussions.

## Appendix

Proposition 4. For the Neumann system the number $N$ of disjoint tori in the preimage of the moment map is $N=2^{m+1}$, where $m$ is the number of real ovals of the corresponding curve $w^{2}=R(u)$ which do not contain any $a_{j}$.

Proof. The reduced phase space has one " $2^{n+1}$-tant" of $\mathbb{R}^{n+1}$ as configuration space, with the identification induced by $G$ on its boundaries. The elliptical coordinates are not a smooth coordinate system for the reduced phase space. The transformation to elliptical spherical coordinates smoothly maps each " $2^{n+1}$-tant" of $\mathrm{S}^{n}$ with constant signs of $x_{\nu}$ onto the $n$-cube $I^{n}$ with sides $a_{j-1} \leq u_{j} \leq a_{j}$. The transformation is singular on the faces $u_{j}=a_{\nu}$ of the cube, which are the images of $x_{\nu}=0$. Except for $x_{0}=0$ and $x_{n}=0$ all other sub-spheres are mapped to two faces.

The coordinates $\phi_{j}$ defined by

$$
\int_{0}^{\phi_{j}} d \phi=\int_{a_{j-1}}^{u_{j}} \frac{\mathrm{~d} u}{2 \sqrt{-A(u)}}
$$

remove the corresponding pole singularities in the metric $g_{j}(\boldsymbol{u})$. The $\phi_{j}$ are coordinates on a torus $\mathrm{T}^{n}$ that covers the sphere $\mathrm{S}^{n}$. The codimension two branch points of the covering are given by by $u_{j}=u_{j-1}=a_{j}$, the poles of $g_{j}^{-1}$. Locally the transformation from $u_{j}$ to $\phi_{j}$ is just a reparameterization of coordinate lines. Globally each coordinate line $\phi_{j}$ on $\mathrm{T}^{n}$ covers its image in the cube $I^{n}$ four times: crossing each of the two faces $u_{j}=a_{\nu}$ of the cube in the $u_{j}$-direction restores the original signs of $x_{\nu}$. We can picture the $\mathrm{T}^{n}$ covering the configuration space $\mathrm{S}^{n}$ as $4^{n}$ cubes with periodic boundary conditions, see Figure 5. Since the map from $\mathrm{S}^{n}$ to $I^{n}$ is $2^{n+1}$-fold the covering from $\mathrm{T}^{n}$ to $\mathrm{S}^{n}$ is $4^{n} / 2^{n+1}=2^{n-1}$-fold.

The projection onto configuration space of the set of $n$-dimensional invariant tori in phase space, corresponding to fixed constants of motion, is a union of disjoint rectangular boxes (with periodic boundary conditions) in the $\phi_{j}$ coordinates. The faces of a box correspond to the ovals $z_{j-1} \leq u_{j} \leq z_{j}$. Denote by $\mu_{j}$ the number of endpoints of the $j$ th oval that are not among $a_{0}, \ldots, a_{n}$. The number of fundamental pieces $I^{n}$ needed to cover this box in $\mathrm{T}^{n}$ is given by the product of the corresponding numbers $2^{\mu_{j}}$ for each degree of freedom. The cases $\mu=2$ (rotation) and $\mu=1$ (oscillation) already appear for $n=1$. Another oscillatory type of motion, with $\mu=0$, occurs for $n>1$. Denote by $\rho, o$, $m$, the number of intervals with $\mu_{j}=2,1,0$, respectively; $\rho+o+m=n$. The number of fundamental pieces needed to cover an invariant torus is then $4^{\rho} 2^{o} 1^{m}=2^{2 \rho+o}=4^{n} / 2^{2 m+o}$. As the total number of fundamental cubes is $4^{n}$, there are $2^{2 m+o}$ copies of the torus in $\mathrm{T}^{n}$. Dividing this by the covering number $2^{n-1}$ we find that for this type of torus there are $2^{m+1-\rho}$ disjoint projections onto configuration space $\mathrm{S}^{n}$. But notice that for each rotational degree of freedom there are two different possibilities for the momentum, so that the total number of disjoint tori in phase space is $N=2^{m+1}$.


Figure 5: The $\mathrm{T}^{2}$ covering $\mathrm{S}^{2}$ twice, and the four types of tori (gray) for $n=2$. a) horizontal strip: type $l r r=\mathrm{RO}_{\mathrm{r}} ;$ b) box in the upper left: type $r r r=\mathrm{O}_{\mathrm{r}} \mathrm{O}_{\mathrm{r}} ;$ c) vertical strip at the center: type $l l r=\mathrm{O}_{1} \mathrm{R}$; d) left vertical strip: type $r l r=\mathrm{O}_{\mathrm{m}} \mathrm{R}$. Compare with Fig. 3.

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[^0]:    ${ }^{1}$ We are grateful to Vitaly Tarasov, who explained this to us

