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# Higher index focus-focus singularities in the Jayne-Cummings-Gaudin model : symplectic invariants and monodromy. 

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#### Abstract

We study the symplectic geometry of the Jaynes-Cummings-Gaudin model with $n=2 m-1$ spins. We show that there are focus-focus singularities of maximal Williamson type $(0,0, m)$. We construct the linearized normal flows in the vicinity of such a point and show that soliton type solutions extend them globally on the critical torus. This allows us to compute the leading term in the Taylor expansion of the symplectic invariants and the monodromy associated to this singularity.


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## 1 Introduction

The theory of integrable systems started with the work of Liouville [1] where he defined the general concept of integrable systems and their integration by "quadratures". Let $\mathcal{M}$ be a symplectic space of dimension $2 n$, the system is Liouville integrable if we can define $n$ Hamiltonians $H_{i}$ in involution such that "generically" $d H_{1} \wedge d H_{2} \wedge \cdots \wedge d H_{n}$ has maximal rank $r=n$. The remarkable Arnold-Liouville theorem states that the space $\mathcal{M}$ is then fibered by invariant $n$-dimensional Lagrangian tori. This is a semi global result as it gives a global information on the fiber, which is a torus, and it is sufficient to assert that the motion is quasi periodic almost everywhere. This fibration however contains singular fibers i.e. fibers containing points where the rank $r$ is not maximal, which play a very important role for the global properties of the system, see e.g. [11]. In this work, we will concentrate on special singularities where the rank $r=0$. They correspond to equilibrium points, which can be stable or unstable, and have also a very rich physical content.

An equilibrium point $x^{(0)}$ is a simultaneous critical point for all the $H_{i}$, in other words, it is a point at which the differential of the moment map $\mu$ vanishes. Expanding the $H_{i}$ around such a point, we get $n$ Poisson commuting quadratic forms $Q_{i}$.

We shall consider the case of a purely focus-focus equilibrium point of Williamson type $\left(m_{e}=m_{h}=0\right.$ and $\left.m_{f f}=m\right)$. In this case, the dimension of $\mathcal{M}$ is equal to $4 m$, and there exist canonical coordinates $\left(p_{i}, q_{i}\right)$ such that the quadratic forms $Q_{i}$ are linear combinations of the quadratic normal forms

$$
\begin{aligned}
K_{j} & =p_{2 j} q_{2 j}+p_{2 j+1} q_{2 j+1}, \quad j=1, \ldots, m \\
L_{j} & =-p_{2 j} q_{2 j+1}+p_{2 j+1} q_{2 j}, \quad j=1, \ldots, m
\end{aligned}
$$

By Eliasson theorem [8], this description extends to a local neighborhood of the point $x_{0}$. There exists a symplectic diffeomorphism $\Phi: U \subset \mathbb{R}^{2 n} \longrightarrow \mathcal{V} \subset \mathcal{M}$ mapping the neighborhood $U$ of the origin to the neighborhood $\mathcal{V}$ of $x^{(0)}$, and a local diffeomorphism $\psi: \mathbb{R}^{n} \longrightarrow \mathbb{R}^{n}$ in a neighborhood of the origin such that:

$$
\mu \circ \Phi=\psi \circ \mu^{(0)}
$$

where $\mu^{(0)}$ is the quadratic moment map $(p, q) \rightarrow(K, L)$. An important consequence of this theorem is that it allows to extend the flows associated to the quadratic generators $K_{j}$ and $L_{j}$ from a neighborhood of $x^{(0)}$ (local level) to a family of fibers containing the level set of $x^{(0)}$ (semi global level). These flows will be called normal flows and will play an important role in this paper.

Using the complex variables

$$
\begin{equation*}
w_{j}=p_{2 j}+i p_{2 j+1}, \quad z_{j}=q_{2 j}+i q_{2 j+1}, \quad j=1, \ldots, m \tag{1}
\end{equation*}
$$

we have

$$
w_{j} \bar{z}_{j}=K_{j}+i L_{j}
$$

hence, the level set $\mathcal{S}$ of the equilibrium point is the image, by the diffeomorphism $\Phi$, of the product of $m$ two dimensional components, themselves the union of two complex planes intersecting transversely

$$
\mathcal{C}_{j}: w_{j}=0 \text { or } z_{j}=0, \quad j=1, \ldots, m
$$

It is a $2 m$-dimensional cone $\mathcal{C}$

$$
\mathcal{C}=\prod_{j=1}^{m} \mathcal{C}_{j}
$$

We see that in a neighborhood of the equilibrium point, the phase-space fibration induced by the moment map is symplectically equivalent to the fibration of a direct product of $m$ independent integrable four dimensional dynamical systems with a focus-focus singularity. An important question is to see whether such a simple description holds also at the semi-global level.

In the Jaynes-Cummings-Gaudin model that we will consider in this paper, the level set of the equilibrium point is a compact pinched torus of dimension $2 m$. A result of Tien Zung [2] asserts that, at the topological level, this pinched torus is equivalent to a product of $m$ two dimensional pinched tori. But this is not true at the symplectic level, and it is the purpose of this work to compute symplectic invariants defined in 3, 4] preventing this simple product decomposition of the singular fiber.

To achieve this goal, we will combine two informations. The first one is a very simple description of the normal modes around the singularity, which provide a complete local description of the system. The second one, of a semi global nature, is provided by the explicit solitonic solutions of the equations of motion on the singular fibre. It turns out that the two objects, normal modes and solitons, match perfectly : solitons are just global extensions of normal modes to the full singular fiber. The $2 m$ real parameters describing the initial conditions of the solitonic solution may be seen as coordinates on the singular fiber which precisely realize its decomposition into a product of $m$ two dimensional singular tori, thereby extending to the whole fiber the decomposition provided by the normal coordinates in the vicinity of the critical point.


Figure 1: The singular torus seen as a product of two dimensional pinched tori whose coordinates are the parameters $X_{k}$ entering the solitonic solution of the equations of motion. A big motion $\left(K_{1}\right)$ on one of the components of the torus induces a motion on the other components of the torus as well.

The plan of the paper is as follows. In Section 2 we recall the definition of the Jaynes-Cummings-Gaudin model and basic facts about its integrability. In Section 3 we study the vicinity of an unstable point of Williamson type $\left(m_{e}=m_{h}=0, m_{f f}=m\right)$. In particular we give a simple description of normal coordinates around that point. In Section 4 we relate these normal modes to solitonic solutions on the pinched torus, the level set of the unstable point. In Section 5we construct closed periodic trajectories on Liouville tori close to the critical one. This requires to take into account both diagonal [4] and non diagonal [5] components of the motion which signal the obstruction to Tien Zung decomposition at the symplectic level. The action variables associated to these periodic motion provide the symplectic invariants in the sense of San Vũ Ngoc. We use them to determine Duistermaat's monodromy [9. A number of technical details are provided in the Appendices for the sake of completeness.

Acknowlegment. We wish to thank San Vũ Ngoc for very helpful discussions.

## 2 The Jaynes-Cummings-Gaudin model

### 2.1 Basic definitions and integrability

This model describes a collection of $n$ spins coupled to a single harmonic oscillator. It derives from the Hamiltonian:

$$
\begin{equation*}
H=\sum_{j=1}^{n}\left(2 \epsilon_{j}+\omega\right) s_{j}^{z}+\omega \bar{b} b+\sum_{j=1}^{n}\left(\bar{b} s_{j}^{-}+b s_{j}^{+}\right) \tag{2}
\end{equation*}
$$

The $\vec{s}_{j}$ are spin variables, and $b, \bar{b}$ is a harmonic oscillator. The Poisson brackets read

$$
\begin{equation*}
\left\{s_{j}^{a}, s_{j}^{b}\right\}=-\epsilon_{a b c} s_{j}^{c}, \quad\{b, \bar{b}\}=i \tag{3}
\end{equation*}
$$

The $\vec{s}_{j}$ brackets are degenerate. We fix the value of the Casimir functions

$$
\left\|\vec{s}_{j}\right\|^{2}=\sum_{a=1}^{3} s_{j}^{a} s_{j}^{a}=s^{2}
$$

Phase space has dimension 2(n+1). In the Hamiltonian we have used $s_{j}^{ \pm}=s_{j}^{1} \pm i s_{j}^{2}$ which have Poisson brackets $\left\{s_{j}^{z}, s_{j}^{ \pm}\right\}= \pm i s_{j}^{ \pm}, \quad\left\{s_{j}^{+}, s_{j}^{-}\right\}=2 i s_{j}^{z}$. The equations of motion read

$$
\begin{align*}
\dot{b} & =-i \omega b-i \sum_{j=1}^{n} s_{j}^{-}  \tag{4}\\
\dot{s_{j}^{z}} & =i\left(\bar{b} s_{j}^{-}-b s_{j}^{+}\right)  \tag{5}\\
\dot{s_{j}^{+}} & =i\left(2 \epsilon_{j}+\omega\right) s_{j}^{+}-2 i \bar{b} s_{j}^{z}  \tag{6}\\
\dot{s_{j}^{-}} & =-i\left(2 \epsilon_{j}+\omega\right) s_{j}^{-}+2 i b s_{j}^{z} \tag{7}
\end{align*}
$$

Introducing the Lax matrices

$$
\begin{align*}
L(\lambda) & =2 \lambda \sigma^{z}+2\left(b \sigma^{+}+\bar{b} \sigma^{-}\right)+\sum_{j=1}^{n} \frac{\vec{s}_{j} \cdot \vec{\sigma}}{\lambda-\epsilon_{j}}  \tag{8}\\
M(\lambda) & =-i \lambda \sigma^{z}-i \frac{\omega}{2} \sigma^{z}-i\left(b \sigma^{+}+\bar{b} \sigma^{-}\right) \tag{9}
\end{align*}
$$

where $\sigma^{a}$ are the Pauli matrices, $\sigma^{ \pm}=\frac{1}{2}\left(\sigma^{x} \pm i \sigma^{y}\right)$,

$$
\sigma^{x}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right), \quad \sigma^{y}=\left(\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right), \quad \sigma^{z}=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right),
$$

it is not difficult to check that the equations of motion are equivalent to the Lax equation

$$
\begin{equation*}
\dot{L}(\lambda)=[M(\lambda), L(\lambda)] \tag{10}
\end{equation*}
$$

Letting

$$
L(\lambda)=\left(\begin{array}{cc}
A(\lambda) & B(\lambda) \\
C(\lambda) & -A(\lambda)
\end{array}\right)
$$

we have

$$
\begin{align*}
& A(\lambda)=2 \lambda+\sum_{j=1}^{n} \frac{s_{j}^{z}}{\lambda-\epsilon_{j}}  \tag{11}\\
& B(\lambda)=2 b+\sum_{j=1}^{n} \frac{s_{j}^{-}}{\lambda-\epsilon_{j}}  \tag{12}\\
& C(\lambda)=2 \bar{b}+\sum_{j=1}^{n} \frac{s_{j}^{+}}{\lambda-\epsilon_{j}} \tag{13}
\end{align*}
$$

The non vanishing Poisson brackets of these functions are simple :

$$
\begin{align*}
\{A(\lambda), B(\mu)\} & =\frac{i}{\lambda-\mu}(B(\lambda)-B(\mu))  \tag{14}\\
\{A(\lambda), C(\mu)\} & =-\frac{i}{\lambda-\mu}(C(\lambda)-C(\mu))  \tag{15}\\
\{B(\lambda), C(\mu)\} & =\frac{2 i}{\lambda-\mu}(A(\lambda)-A(\mu)) \tag{16}
\end{align*}
$$

It follows immediately that $\operatorname{Tr}\left(\mathrm{L}^{2}(\lambda)\right)=2 \mathrm{~A}^{2}(\lambda)+2 \mathrm{~B}(\lambda) \mathrm{C}(\lambda)$ Poisson commute for different values of the spectral parameter:

$$
\left\{\operatorname{Tr}\left(\mathrm{L}^{2}\left(\lambda_{1}\right)\right), \operatorname{Tr}\left(\mathrm{L}^{2}\left(\lambda_{2}\right)\right)\right\}=0
$$

Hence $\Lambda(\lambda) \equiv \frac{1}{2} \operatorname{Tr}\left(\mathrm{~L}^{2}(\lambda)\right)$ generates Poisson commuting quantities. One has

$$
\begin{equation*}
\Lambda(\lambda)=\frac{Q_{2 n+2}(\lambda)}{\prod_{j}\left(\lambda-\epsilon_{j}\right)^{2}}=4 \lambda^{2}+4 H_{n+1}+2 \sum_{j=1}^{n} \frac{H_{j}}{\lambda-\epsilon_{j}}+\sum_{j=1}^{n} \frac{s^{2}}{\left(\lambda-\epsilon_{j}\right)^{2}} \tag{17}
\end{equation*}
$$

where the $(n+1)$ commuting Hamiltonians $H_{j}, j=1, \cdots, n+1$ read

$$
\begin{equation*}
H_{n+1}=b \bar{b}+\sum_{j} s_{j}^{z} \tag{18}
\end{equation*}
$$

and

$$
\begin{equation*}
H_{j}=2 \epsilon_{j} s_{j}^{z}+\left(b s_{j}^{+}+\bar{b} s_{j}^{-}\right)+\sum_{k \neq j} \frac{s_{j} \cdot s_{k}}{\epsilon_{j}-\epsilon_{k}}, \quad j=1, \cdots, n \tag{19}
\end{equation*}
$$

The physically interesting Hamiltonian eq.(22) is

$$
\begin{equation*}
H=\omega H_{n+1}+\sum_{j=1}^{n} H_{j} \tag{20}
\end{equation*}
$$

### 2.2 Critical points

The critical points are equilibrium points for all the Hamiltonians $H_{j}, j=1, \cdots, n+1$. At such points the derivatives with respect of all coordinates on phase space vanish. In particular, since

$$
\frac{\partial H_{n+1}}{\partial \bar{b}}=b, \quad \frac{\partial H_{j}}{\partial \bar{b}}=s_{j}^{-}
$$

we see that the critical points must be located at

$$
\begin{equation*}
b=\bar{b}=0, \quad s_{j}^{ \pm}=0, \quad s_{j}^{z}=e_{j} s, \quad e_{j}= \pm 1 \tag{21}
\end{equation*}
$$

When we expand around a configuration eq.(21), all the quantities $\left(b, \bar{b}, s_{j}^{+}, s_{j}^{-}\right)$are first order, but $s_{j}^{z}$ is second order because

$$
s_{j}^{z}=e_{j} \sqrt{s^{2}-s_{j}^{+} s_{j}^{-}}=s e_{j}-\frac{e_{j}}{2 s} s_{j}^{+} s_{j}^{-}+\cdots, \quad e_{j}= \pm 1
$$

It is then simple to see that all first order terms in the expansions of the Hamiltonians $H_{j}$ vanish. Hence we have found $2^{n}$ critical points.

## 3 The vicinity of an unstable fixed point

### 3.1 Normal Forms

In the following discussion, we are considering the system with an odd number of spins $n=2 m-1$ with $m \geq 1$, and we are assuming that the equilibrium point defined by eq. (21) has the Williamson type $m_{e}=m_{h}=0$ and $m_{f f}=m$, i.e. it consists locally of a direct product of $m$ elementary focus-focus singularities. It has been shown before that such points can be obtained by choosing the Zeeman energies $\epsilon_{j}$ in a suitable region in parameter space [14, 13]. The quadratic normal form around an equilibrium point can be easily obtained from the Lax representation [10, 14 and this is recalled in Section 6 below. The above assumption on the Williamson type is equivalent to the statement that the classical Bethe equation (52) has no real root (hence the choice $n$ odd). Since the equation is real, it has $m$ pairs of complex conjugated roots $E_{j}, \bar{E}_{j}$, for $1 \leq j \leq m$. To first order in small deviations around the equilibrium point, the normal coordinates are $B\left(E_{j}\right), B\left(\bar{E}_{j}\right), C\left(E_{j}\right)$, and $C\left(\bar{E}_{j}\right)$ where the functions $B(\lambda), C(\lambda)$ are those defined in eqs.(12|13). The precise relation with the canonical coordinates eq.(11) is:

$$
\begin{array}{lll}
B\left(E_{j}\right)=-i a^{\prime}\left(E_{j}\right) z_{j}, & & B\left(\bar{E}_{j}\right)=w_{j} \\
C\left(E_{j}\right)=\bar{w}_{j}, & C\left(\bar{E}_{j}\right)=i a^{\prime}\left(\bar{E}_{j}\right) \bar{z}_{j} \tag{23}
\end{array}
$$

These expressions for $z_{j}$ and $w_{j}$ as functions of the dynamical variables of the Gaudin model should be understood as giving the differential of the function $\Phi^{-1}$ at the equilibrium point. An important consequence of Eliasson's theorem is that the knowledge of this differential is sufficient to identify the generators of the normal flows, as long as we are interested in their restriction to the singular torus.

The symplectic form in terms of normal coordinates reads:

$$
\omega^{(0)}=\frac{1}{2} \sum_{j} d z_{j} \wedge d \bar{w}_{j}+d \bar{z}_{j} \wedge d w_{j}
$$

As explained in the Introduction, the components of the moment map $\mu^{(0)}$ are the quadratic Hamiltonians:

$$
L_{j}=\frac{1}{2 i}\left(\bar{z}_{j} w_{j}-z_{j} \bar{w}_{j}\right), \quad K_{j}=\frac{1}{2}\left(\bar{z}_{j} w_{j}+z_{j} \bar{w}_{j}\right)
$$

The Hamilton equations of motion associated to $\omega^{(0)}$ are:

$$
\dot{z}_{j}=2 \frac{\partial H}{\partial \bar{w}_{j}}, \quad \dot{w}_{j}=-2 \frac{\partial H}{\partial \bar{z}_{j}}, \quad \dot{\bar{z}}_{j}=2 \frac{\partial H}{\partial w_{j}}, \quad \dot{\bar{w}}_{j}=-2 \frac{\partial H}{\partial z_{j}}
$$

From them, we see that the flows associated to $K_{j}, L_{j}$ are

$$
K_{j}:\left(z_{j} \rightarrow e^{a_{j}} z_{j}, w_{j} \rightarrow e^{-a_{j}} w_{j}\right), \quad L_{j}:\left(z_{j} \rightarrow e^{i \theta_{j}} z_{j}, w_{j} \rightarrow e^{i \theta_{j}} w_{j}\right)
$$

### 3.2 Generators of normal flows

We begin by expanding the function $\mu^{(0)} \circ \Phi^{-1}$ to second order in small deviations from the equilibrium point, recalling that $B\left(E_{j}\right), B\left(\bar{E}_{j}\right), C\left(E_{j}\right)$, and $C\left(\bar{E}_{j}\right)$ are first order in these deviations. This gives:

$$
\begin{align*}
K_{j} & =\frac{i}{2}\left(\frac{B\left(E_{j}\right) C\left(E_{j}\right)}{a^{\prime}\left(E_{j}\right)}-\frac{B\left(\bar{E}_{j}\right) C\left(\bar{E}_{j}\right)}{a^{\prime}\left(\bar{E}_{j}\right)}\right)  \tag{24}\\
L_{j} & =-\frac{1}{2}\left(\frac{B\left(E_{j}\right) C\left(E_{j}\right)}{a^{\prime}\left(E_{j}\right)}+\frac{B\left(\bar{E}_{j}\right) C\left(\bar{E}_{j}\right)}{a^{\prime}\left(\bar{E}_{j}\right)}\right) \tag{25}
\end{align*}
$$

But Eliasson's theorem implies that $\mu^{(0)} \circ \Phi^{-1}=\psi^{-1} \circ \mu$, so we should be able to express these quantities in terms of the conserved Hamiltonians $H_{1}, \ldots, H_{n+1}$. In Section 6 below, we show that:

$$
\begin{equation*}
B\left(E_{j}\right) C\left(E_{j}\right)=4 E_{j}^{2}+4 H_{n+1}+\sum_{i=1}^{n} \frac{2 H_{i}}{E_{j}-\epsilon_{i}}+\sum_{k=1}^{n} \frac{s^{2}}{\left(E_{j}-\epsilon_{i}\right)^{2}}+\ldots \tag{26}
\end{equation*}
$$

where the neglected terms are of order 3 and higher in small deviations from the equilibrium point. This formula enables to construct the Taylor expansion of the function $\psi^{-1}$ up to first order. We get:

$$
\begin{aligned}
K_{j} & =i\left(\frac{2}{a^{\prime}\left(E_{j}\right)}-\frac{2}{a^{\prime}\left(\bar{E}_{j}\right)}\right) \delta H_{n+1}+i \sum_{i=1}^{n}\left(\frac{1}{a^{\prime}\left(E_{j}\right)\left(E_{j}-\epsilon_{i}\right)}-\frac{1}{a^{\prime}\left(\bar{E}_{j}\right)\left(\bar{E}_{j}-\epsilon_{i}\right)}\right) \delta H_{i}+(.27) \\
L_{j} & =-\left(\frac{2}{a^{\prime}\left(E_{j}\right)}+\frac{2}{a^{\prime}\left(\bar{E}_{j}\right)}\right) \delta H_{n+1}-\sum_{i=1}^{n}\left(\frac{1}{a^{\prime}\left(E_{j}\right)\left(E_{j}-\epsilon_{i}\right)}+\frac{1}{a^{\prime}\left(\bar{E}_{j}\right)\left(\bar{E}_{j}-\epsilon_{i}\right)}\right) \delta H_{i}+(28)
\end{aligned}
$$

Here $\delta H_{i}=H_{i}-H_{i, \text { eq }}$, where $H_{i, \text { eq }}$ denotes the value of $H_{i}$ at the equilibrium point. The neglected term are of order at least two in the $\delta H_{i}$ 's. When we consider fibers of the moment
map closer and closer to the fiber containg the equilibrium point, the flows generated by the linear terms in the $\delta H_{i}$ 's in equations (27) and (28) tend towards the normal flows defined in the Introduction. In particular, they coincide with these normal flows on the fiber containing the equilibrium point.

## 4 Normal flows on a singular torus

### 4.1 Normal flows and soliton solutions

From the previous discussion, the normal flows on the singular torus are given by the Hamiltonian flows associated to the generators $K_{j}$ and $L_{j}$ defined in terms of the conserved Hamiltonians in eqs. (27) and (28). On the singular torus, these flows take the form of solitonic trajectories. For the reader's convenience a complete derivation of the soliton formulae, which were constructed in [13], is given in Section 9 below. Here, we summarize the main features of these solitonic solutions.

We recall that $n=2 m-1$, and we assume that $\mathcal{E}^{0}=\emptyset$ in the construction of Section 9 We define column vectors $\left(E^{j}\right)_{l}=E_{l}^{j}$ and $\left(X E^{j}\right)_{l}=X_{l} E_{l}^{j}, l=1 \cdots 2 m$, and two polynomials

$$
\mathcal{P}^{-}(\lambda)=\frac{1}{D_{n+1}} \operatorname{det}\left(\begin{array}{cccccccc}
1 & \lambda & \cdots & \lambda^{n_{-}} & 0 & 0 & \cdots & 0 \\
1 & E & \cdots & E^{n_{-}} & X & X E & \cdots & X E^{n_{+}}
\end{array}\right)
$$

and

$$
\mathcal{P}^{+}(\lambda)=\frac{1}{D_{0}} \operatorname{det}\left(\begin{array}{cccccccc}
0 & 0 & \cdots & 0 & 1 & \lambda & \cdots & \lambda^{n_{+}} \\
1 & E & \cdots & E^{n_{-}} & X & X E & \cdots & X E^{n_{+}}
\end{array}\right)
$$

where $D_{0}$ and $D_{n+1}$ are the determinants:

$$
\begin{align*}
D_{0} & =\operatorname{det}\left(\begin{array}{llllllll}
1 & E & \cdots & E^{n_{-}} & X & X E & \cdots & X E^{n_{+}-1}
\end{array}\right)  \tag{29}\\
D_{n+1} & =\operatorname{det}\left(\begin{array}{llllllll}
1 & E & \cdots & E^{n_{-}-1} & X & X E & \cdots & X E^{n_{+}}
\end{array}\right) \tag{30}
\end{align*}
$$

The indices $n_{ \pm}$are defined as

$$
n_{+}=\operatorname{degree}\left(\mathcal{P}^{+}(\lambda)\right)=\frac{1}{2}(n-1)=m-1, \quad n_{-}=\operatorname{degree}\left(\mathcal{P}^{-}(\lambda)\right)=\frac{1}{2}(n+1)=m
$$

The $X_{l}$ are given by:

$$
\begin{equation*}
X_{l}(\{t\})=X_{l}(0) e^{\sqrt{-1}}\left(\sum_{i} \frac{s e_{i}}{E_{l}-\epsilon_{i}} t_{i}-t_{n+1}\right), \tag{31}
\end{equation*}
$$

where $t_{i}$ are the times flows associated to the Hamiltonians $H_{i}$ and $e_{i}= \pm 1$ characterize the critical point $s_{i}^{z}=s e_{i}$. The constants $X_{l}(0)$ are subjected to the reality conditions:

$$
\begin{equation*}
\overline{X_{l}(\{t\})} X_{\bar{l}}(\{t\})=\overline{X_{l}(0)} X_{\bar{l}}(0)=-\frac{1}{4} \tag{32}
\end{equation*}
$$

where $X_{\bar{l}}$ is associated to $E_{\bar{l}} \equiv \overline{E_{l}}$. The function $C(\lambda)$ is given by :

$$
\begin{gather*}
C(\lambda)=  \tag{33}\\
2 \xrightarrow{\operatorname{det}\left(\begin{array}{ccccccc}
1 & \lambda & \cdots & \lambda^{n_{-}} & 0 & 0 & \cdots \\
1 & E & \cdots & E^{n_{-}} & X & X E & \cdots \\
1 & X E^{n_{+}}
\end{array}\right) \operatorname{det}\left(\begin{array}{ccccccccc}
0 & 0 & \cdots & 0 & 1 & \lambda & \cdots & \lambda^{n_{+}} \\
1 & E & \cdots & E^{n_{-}} & X & X E & \cdots & X E^{n_{+}}
\end{array}\right)} \\
\prod_{j}\left(\lambda-\epsilon_{j}\right) \operatorname{det}\left(\begin{array}{lllllll}
1 & E & \cdots & E^{n_{-}-1} & X & X E & \cdots \\
& \left.X E^{n_{+}}\right)^{2}
\end{array}\right.
\end{gather*}
$$

Let us now specialize the evolution equation (31) to the normal flows defined through their generators (27) and (28). Let us denote by $a_{l}$ (resp. $\theta_{l}$ ) the time variable of the normal flow generated by $K_{l}\left(\right.$ resp. $\left.L_{l}\right)$. For an observable $X\left(t_{1}, \cdots, t_{n+1}\right)$, we have

$$
\begin{aligned}
& X\left(a_{l}\right)=X\left(t_{i}=i\left(\frac{1}{a^{\prime}\left(E_{l}\right)\left(E_{l}-\epsilon_{i}\right)}-\frac{1}{a^{\prime}\left(\bar{E}_{l}\right)\left(\bar{E}_{l}-\epsilon_{i}\right)}\right) a_{l}, \quad t_{n+1}=i\left(\frac{2}{a^{\prime}\left(E_{l}\right)}-\frac{2}{a^{\prime}\left(\bar{E}_{l}\right)}\right) a_{l}\right) \\
& X\left(\theta_{l}\right)=X\left(t_{i}=-\left(\frac{1}{a^{\prime}\left(E_{l}\right)\left(E_{l}-\epsilon_{i}\right)}+\frac{1}{a^{\prime}\left(\bar{E}_{l}\right)\left(\bar{E}_{l}-\epsilon_{i}\right)}\right) \theta_{l}, \quad t_{n+1}=-\left(\frac{2}{a^{\prime}\left(E_{l}\right)}+\frac{2}{a^{\prime}\left(\bar{E}_{l}\right)}\right) \theta_{l}\right)
\end{aligned}
$$

Let us apply this formula to $X_{k}(\{t\})$ in eq.(31). We have to evaluate:

$$
\frac{1}{a^{\prime}\left(E_{l}\right)} \sum_{i=1}^{n} \frac{s e_{i}}{\left(E_{k}-\epsilon_{i}\right)\left(E_{l}-\epsilon_{i}\right)}
$$

Suppose first $k \neq l$. We have:

$$
\begin{aligned}
\frac{1}{a^{\prime}\left(E_{l}\right)} \sum_{i=1}^{n} \frac{s e_{i}}{\left(E_{k}-\epsilon_{i}\right)\left(E_{l}-\epsilon_{i}\right)} & =\frac{1}{a^{\prime}\left(E_{l}\right)} \sum_{i=1}^{n} \frac{1}{E_{k}-E_{l}}\left(\frac{s e_{i}}{E_{l}-\epsilon_{i}}-\frac{s e_{i}}{E_{k}-\epsilon_{i}}\right) \\
& =\frac{1}{a^{\prime}\left(E_{l}\right)} \sum_{i=1}^{n} \frac{1}{E_{k}-E_{l}}\left(-2 E_{l}+2 E_{k}\right)=\frac{2}{a^{\prime}\left(E_{l}\right)}
\end{aligned}
$$

where we have used the classical Bethe equation:

$$
a(E)=2 E+\sum_{i} \frac{s e_{i}}{E-\epsilon_{i}}=0
$$

If $k=l$, we have directly:

$$
\frac{1}{a^{\prime}\left(E_{l}\right)} \sum_{i=1}^{n} \frac{s e_{i}}{\left(E_{l}-\epsilon_{i}\right)^{2}}=\frac{1}{a^{\prime}\left(E_{l}\right)}\left(2-a^{\prime}\left(E_{l}\right)\right)
$$

Using these sums in eq.(31), we arrive at the very simple result:

$$
\begin{equation*}
X_{k}\left(a_{l}\right)=X_{k}(0) e^{a_{l}\left(\delta_{k l}-\delta_{k \bar{l}}\right)}, \quad X_{k}\left(\theta_{l}\right)=X_{k}(0) e^{i \theta_{l}\left(\delta_{k l}+\delta_{k \bar{l}}\right)} \tag{34}
\end{equation*}
$$

It is clear that the reality condition (32) is preserved by these flows. The meaning of this formula is that the solitons can be viewed as a non-linear global extension of the normal modes on the critical torus. The $a_{l}$ trajectories give $m$ non compact cycles, while the $\theta_{l}$ trajectories give the remaining $m$ compact cycles. In view of eq.(31) we see that the soliton solution depends on $2 m$ complex parameters $X_{l}(0)$ submitted to the $m$ complex conditions eq.(32). Hence there remains $m$ complex (or $2 m$ real) parameters which may be viewed as coordinates on the singular torus. These solitonic coordinates can be interpreted as generalizing at the semi global level the decomposition of the singular fiber into a product of two dimensional singular tori that is performed by normal coordinates in the vicinity of the critical point.

### 4.2 Solitonic picture of the tangent cone

We would like now to precise the connection between these solitonic formulae and the linearized flows in the vicinity of the equilibrium point. As discussed in the Introduction, the tangent cone to the singular torus at the equilibrium point is the product of $m$ two dimensional components $\mathcal{C}_{k}$. There is an important relation between the components $\mathcal{C}_{k}$ of the cone $\mathcal{C}$ and the flows associated to the Hamiltonians $K_{j}$ on the critical manifold. In normal coordinates the orbits of generators $L_{j}$ are closed circles, which become closed orbits when mapped by $\Phi$ on phase-space $\mathcal{M}$. The $m$ other generators $K_{j}$ induce flows which become unbounded when $a_{j}$ goes to $\pm \infty$ and the trajectories eventually leave the open neighborhood $U$ in which the diffeomorphism $\Phi$ is defined. Notice however that when $a_{j} \rightarrow+\infty, z_{j}=0$ is a contracting manifold for the flow $K_{j}$ and $w_{j}=0$ is an expanding manifold. When $a_{j} \rightarrow-\infty$ the situation is opposite. Hence a pattern $\mathcal{S}$ of $m$ signs $a_{j} \rightarrow \pm \infty$ identifies a unique contracting manifold $\mathcal{C}_{\mathcal{S}}^{\text {contracting }}$.

Because there is only one critical point on the critical manifold, when we send all the $a_{j}$ to $\pm \infty$ we end up into a neighborhood of the critical point. The contracting manifold $\mathcal{C}_{\mathcal{S}}^{\text {contracting }}$ is therefore the asymptotic manifold when all the $a_{j}$ 's are sent to infinity with the associated sign pattern.

From eq. (34), the large $a_{j}$ limit corresponds to:

$$
\begin{aligned}
& \left|X_{j}\right| \rightarrow \infty \text { and }\left|X_{\bar{j}}\right| \rightarrow 0 \text { if } a_{j} \rightarrow+\infty \\
& \left|X_{\bar{j}}\right| \rightarrow \infty \text { and }\left|X_{j}\right| \rightarrow 0 \text { if } a_{j} \rightarrow-\infty
\end{aligned}
$$

To connect the soliton solutions with the definition of the tangent cone, we have to take this limit in the soliton formulae. Let $\mathcal{I}$ the set of $m$ indices $\left(i_{1}, i_{2}, \cdots i_{m}\right) \subset\{1 \cdots m, \overline{1} \cdots \bar{m}\}$ such that $X_{i_{j}} \rightarrow \infty$, specifying a contracting component $\mathcal{C}_{\mathcal{S}(\mathcal{I})}^{\text {contracting }}$ of the tangent cone. If $j \in \mathcal{I}$ we still denote by $\bar{j}$ the label of the complex conjugate root $\overline{E_{j}} \equiv E_{\bar{j}}$ (so that $\overline{\bar{j}}=j$ ). Then because of the reality condition eq.(32), $\bar{j} \notin \mathcal{I}$ and the indices $\bar{j}$ belong to the complementary set $\overline{\mathcal{I}}$ of $\mathcal{I}$ in $\{1 \cdots m, \overline{1} \cdots \bar{m}\}$ (see Fig. 21).


Figure 2: The red dots represent the subset $\mathcal{I}=\left\{i_{1}=1, i_{2}=\overline{2}, i_{3}=\overline{3}, i_{4}=4, i_{5}=\overline{5}\right\}$. The green dots represent the complementary (and conjugate) subset $\overline{\mathcal{I}}=\left\{\overline{i_{1}}=\overline{1}, \overline{i_{2}}=2, \overline{i_{3}}=3, \overline{i_{4}}=\overline{4}, \overline{i_{5}}=5\right\}$.

It is useful to relate this definition of the components $\mathcal{C}_{\mathcal{S}(\mathcal{I})}^{\text {contracting }}$ of the tangent cone with the one based on the normal coordinates introduced in Section 3, As we have said, either $j \in \mathcal{I}$ or $\bar{j} \in \mathcal{I}$. If $j \in \mathcal{I}, z_{j} \rightarrow 0$, so according to eqs. (22), (23), $\left|C\left(E_{\bar{j}}\right)\right| \ll\left|C\left(E_{j}\right)\right|$. If $\bar{j} \in \mathcal{I}$, the converse holds and $\left|C\left(E_{j}\right)\right| \ll\left|C\left(E_{\bar{j}}\right)\right|$. To check that this is indeed true, we need to know the asymptotic expression for $C(\lambda)$ when $\left|X_{j}\right| \rightarrow \infty$ for any $j \in \mathcal{I}$. Using eq. (33) and the asymptotic expressions eqs. (95), (96), (97) given in Section 10 below, we see that on the component $\mathcal{C}_{\mathcal{I}}$ of tangent cone, $C(\lambda)$ behaves as:

$$
\begin{equation*}
C(\lambda) \simeq \frac{2}{\prod_{i}\left(\lambda-\epsilon_{i}\right)} \prod_{i \in \mathcal{I}}\left(\lambda-E_{\bar{i}}\right) \sum_{j \in \mathcal{I}}\left(\frac{E_{\bar{j}}-E_{j}}{X_{j}} \prod_{\substack{k \in \mathcal{I} \\ k \neq j}} \frac{E_{\bar{k}}-E_{j}}{E_{k}-E_{j}}\left(\lambda-E_{k}\right)\right) \tag{35}
\end{equation*}
$$

This expression shows that $C(\lambda)$ goes to zero when all the $X_{j}$ 's go to infinity, for $j \in \mathcal{I}$. This shows that we are indeed describing a $2 m$-dimensional subspace (over $\mathbb{R}$ ) of the tangent space at the equilibrium point. Next, we see that, at this leading order, $C(\lambda)$ vanishes at $\lambda=E_{\bar{j}}$. In reality, $C\left(E_{\bar{j}}\right)$ is not exactly zero, because of the subleading terms included in the determinant $D_{1}$ defined in Section 10. If the typical value of the $X_{j}$ 's scales like $X$, the ratio of these subleading terms over the dominant one scales like $X^{-2}$. From this, we conclude that $C\left(E_{\bar{j}}\right) / C\left(E_{j}\right) \rightarrow 0$ when $X \rightarrow \infty$, as expected. More precisely, if $j \in \mathcal{I}$, to leading order in the scale $X$, we have from eq. (35):

$$
\begin{equation*}
C\left(E_{j}, X_{i \in \mathcal{I}} \rightarrow \infty\right) \simeq-\frac{2}{X_{j} \prod_{i}\left(E_{j}-\epsilon_{i}\right)} \prod_{k \in \mathcal{I}}\left(E_{j}-E_{\bar{k}}\right)^{2}, \quad C\left(E_{\bar{j}}\right)=O\left(X^{-3}\right) \tag{36}
\end{equation*}
$$

In the limit $a_{j} \rightarrow-\infty$, the dominant variables become the $X_{\bar{j}}$ 's. The leading behavior for $C(\lambda)$ is now given by eq. (35) where the set $\mathcal{I}$ is replaced by the complementary set $\overline{\mathcal{I}}$ (in which all indices are replaced by $i \rightarrow \bar{i}$ ). The state of the system is now labelled by the coordinates $B\left(E_{j}\right)=\overline{C\left(E_{\bar{j}}\right)}$. Using the reality condition (32) and replacing $E_{j}$ 's by $E_{\bar{j}}$ 's in eq. (36) gives for $j \in \mathcal{I}$ :

$$
\begin{equation*}
B\left(E_{j}, X_{i \in \mathcal{I}} \rightarrow 0\right) \simeq \frac{8 X_{j}}{\prod_{i}\left(E_{j}-\epsilon_{i}\right)} \prod_{k \in \mathcal{I}}\left(E_{j}-E_{\bar{k}}\right)^{2}, \quad B\left(E_{\bar{j}}\right)=O\left(X^{3}\right) \tag{37}
\end{equation*}
$$

where the scale $X$ goes to zero.

## 5 Periodic flows, symplectic invariants and monodromy

### 5.1 Motivation

In an integrable system, the action integrals $S_{\gamma}=\int_{\gamma} \alpha$, where the symplectic form $\omega=d \alpha$ and $\gamma$ denotes any closed cycle on regular Arnold-Liouville tori, are smooth functions of the conserved quantities in any open set containing only regular values of the moment map. These functions are also the generators of periodic flows on Arnold-Liouville tori. Inspired by previous works [3, 4, 5], we wish to examine the construction of periodic flows in the vicinity of a singular torus. We know already that $m$ such flows are generated by the $L_{j}$ 's. As we have seen, on the singular torus, the flow generated by $K_{j}$ corresponds to a large motion, moving away from the
critical point $x^{(0)}$ along the corresponding expanding manifold in the remote past and returning to $x^{(0)}$ along the contracting manifold in the far future. As suggested by Fig. 3] it is possible to use these large solitonic trajectories to construct periodic orbits, at least for regular tori close to the singular one.



Figure 3: The red curves represent the large motion of the $\alpha_{1} K_{1}+\beta_{1} L_{1}$ flow projected on the $z_{1}$, $w_{1}$ manifold in the soliton case (outer curve) and on a nearby Liouville torus (inner curve). The motion starts from a position $\left(z_{1}^{i n}, w_{1}^{i n} \simeq 0\right)$ on the expanding manifold and returns to a position $\left(z_{1}^{\text {out }} \simeq 0, w_{1}^{\text {out }}\right)$ on the contracting manifold. At the same time on the $j^{\text {th }}>1$ component, there is an induced motion from $\left(z_{j}^{i n}, w_{j}^{i n}\right)$ to $\left(z_{j}^{o u t}, w_{j}^{o u t}\right)$. One can close the trajectory by adding an extra motion $\alpha_{1} K_{1}+\beta_{1} L_{1}$ from $\left(z_{1}^{\text {out }}, w_{1}^{\text {out }}\right)$ to $\left(z_{1}^{i n}, w_{1}^{\text {in }}\right)$ and a motion $\alpha_{j} K_{j}+\beta_{j} L_{j}$ from $\left(z_{j}^{o u t}, w_{j}^{\text {out }}\right)$ to $\left(z_{j}^{i n}, w_{j}^{i n}\right)$ (blue segments). These extra motions lie entirely in the domain of validity of the normal form.

### 5.2 Periodic flows.

The Hamiltonian generating the periodic flow is a function of the $K_{l}, L_{l}$ :

$$
\mathcal{H}=\mathcal{H}\left(K_{1}, L_{1}, \cdots, K_{m}, L_{m}\right)
$$

This function however is not defined globally and this gives rise to the monodromy phenomenon. The flow itself is a linear combinations of the flows generated by $K_{l}, L_{l}$

$$
\begin{equation*}
\partial_{t}=\sum_{l} \alpha_{l} \partial_{a_{l}}+\beta_{l} \partial_{\theta_{l}}, \quad \alpha_{l}=\partial \mathcal{H} / \partial K_{l}, \quad \beta_{l}=\partial \mathcal{H} / \partial L_{l} \tag{38}
\end{equation*}
$$

The $X$ 's become functions of time

$$
\begin{equation*}
X_{k}=X_{k}(0) e^{\alpha_{k} t+i \beta_{k} t}, \quad X_{\bar{k}}=X_{\bar{k}}(0) e^{-\alpha_{k} t+i \beta_{k} t} \tag{39}
\end{equation*}
$$

We want to find conditions on $\alpha_{k}, \beta_{k}$ in order to get a closed trajectory.

To simplify notations, we will consider the large motion generated by $\alpha_{1} K_{1}+\beta_{1} L_{1}$. Fig. 3 shows that we first have to study the corresponding symplectic map for an initial condition $\left(z_{j}^{i n}, w_{j}^{i n}\right)$ close to the expanding manifold associated to $K_{1}$ i.e. such that $w_{1}^{i n}$ is small. The initial conditions and the time lapse $2 T$ are chosen in such a way that the image $\left(z_{j}^{\text {out }}, w_{j}^{\text {out }}\right)$ is close to the contracting manifold, i.e. $z_{1}^{\text {out }}$ is small. In general, $\left(z_{j}^{\text {out }}, w_{j}^{\text {out }}\right)$ is different from $\left(z_{j}^{i n}, w_{j}^{i n}\right)$ for $j \neq 1$, reflecting the impossibility to express our dynamical system in the vicinity of the pinched torus as a product of $m$ independent integrable systems with 2 degrees of freedom, each of them exhibiting an isolated focus-focus singularity. Although such factorization holds at the topological level [2], it does not hold at the level of symplectic manifolds. The obstruction to such a symplectic product decomposition is encoded in the non-diagonal symplectic invariants, which we will now compute. For this purpose we will use the fact recalled in Section 7 , that the $\operatorname{map}\left(z_{j}^{\text {in }}, w_{j}^{\text {in }}\right) \rightarrow\left(z_{j}^{\text {out }}, w_{j}^{\text {out }}\right)$ is strongly constrained by the condition that $K_{1}$ and $L_{1}$ commute with all the generators of the moment map.

It is first necessary to identify the components $\mathcal{C}_{\mathcal{S}(\mathcal{I})}^{\text {contracting }}$ of the cone corresponding to the in and out states for a large motion generated by $\alpha_{1} K_{1}+\beta_{1} L_{1}$. This flow sends the solitonic coordinate $X_{1}$ into $e^{\alpha_{1} a_{1}+i \beta_{1} \theta_{1}} X_{1}$ and $X_{\overline{1}}$ into $e^{-\alpha_{1} a_{1}+i \beta_{1} \theta_{1}} X_{\overline{1}}$, and leaves unchanged all the other coordinates $X_{k}, X_{\bar{k}}$. Let us choose an initial condition on an expanding subspace for this flow. This means that $\left|X_{\overline{1}}^{\mathrm{in}}\right| \gg\left|X_{1}^{\mathrm{in}}\right|$. In principle, the other coordinates can be chosen arbitrarily, but it is convenient to assume also that $\left|X_{\bar{k}}^{\mathrm{in}}\right| \gg\left|X_{k}^{\mathrm{in}}\right|$ for $k>1$, so that the initial point is close to the critical point. Let us denote by $\mathcal{I}_{0}$ the set of indices $\{2 \cdots m\}$. The manifold $\mathcal{C}^{\text {in }}$ is obtained by sending $a_{1}$ to $-\infty$ and the manifold $\mathcal{C}^{\text {out }}$ by sending $a_{1}$ to $+\infty$, the other $j \in \overline{\mathcal{I}}_{0}$ being spectators. Then $\mathcal{C}^{\text {in }}=\mathcal{C}_{\mathcal{S}\left(\{\overline{1}\} \cup \overline{\mathcal{I}}_{0}\right)}^{\text {contracting }}$ and $\mathcal{C}^{\text {out }}=\mathcal{C}_{\mathcal{S}\left(\{1\} \cup \overline{\mathcal{I}}_{0}\right)}^{\text {contracting }}$.

The symplectic invariants are smooth functions of the conserved quadratic Hamiltonians $K_{j}, L_{j}(1 \leq j \leq m)$ defined in the vicinity of the singular torus [3, 4, 5]. The knowledge of solitonic trajectories allows us to compute them exactly on the singular torus, that is for $K_{j}=L_{j}=0$.

### 5.2.1 Diagonal case.

As apparent of Fig. [3, the trajectory is composed of two pieces. The first one is the large solitonic motion generated by $\alpha_{1} K_{1}+\beta_{1} L_{1}$ during a time lapse $2 T$ which yields a relation between $\left(z_{1}^{\text {in }}, w_{1}^{\text {in }}\right) \rightarrow\left(z_{1}^{\text {out }}, w_{1}^{\text {out }}\right)$ as explained in Section 7 This part of the trajectory is complemented by a further action of the flow $\alpha_{1} K_{1}+\beta_{1} L_{1}$ during a time lapse $\tau$ relating $\left(z_{1}^{\text {out }}, w_{1}^{\text {out }}\right) \rightarrow\left(z_{1}^{\text {in }}, w_{1}^{\text {in }}\right)$ so as to get a closed trajectory. This last part of the analysis lies entirely in the domain of validity of the normal forms. One should therefore have

$$
\begin{equation*}
z_{1}^{i n}=e^{\alpha_{1} \tau+i \beta_{1} \tau} z_{1}^{o u t}, \quad w_{1}^{i n}=e^{-\alpha_{1} \tau+i \beta_{1} \tau} w_{1}^{\text {out }} \tag{40}
\end{equation*}
$$

Using the leading order expressions for $z_{1}^{\text {out }}, w_{1}^{\text {out }}$, eqs. (63|64) these two equations are equivalent to

$$
\begin{equation*}
\bar{z}_{1}^{i n} w_{1}^{i n}=K_{1}+i L_{1}=e^{-\alpha_{1} \tau+i \beta_{1} \tau} \Phi_{1}^{0}(T) \tag{41}
\end{equation*}
$$

As shown in Section 7, one has, again to leading order

$$
w_{1}^{o u t}=\frac{\Phi_{1}^{0}(T)}{K_{1}+i L_{1}} w_{1}^{i n}=\frac{\Phi_{1}^{0}(T)}{\bar{z}_{1}^{i n}}
$$

The variables $z_{1}^{i n}, w_{1}^{i n}$ are obtained from $C\left(E_{1}\right), C\left(E_{\overline{1}}\right)$ evaluated at $t=-T$ and the variables $z_{1}^{\text {out }}, w_{1}^{\text {out }}$ are obtained from $C\left(E_{1}\right), C\left(E_{\overline{1}}\right)$ evaluated at $t=T$. Explicitely we have:

$$
\left\{\begin{array}{l}
\left.C\left(E_{1}\right)\right|_{-T}=\bar{w}_{1}^{\text {in }} \simeq 0 \\
\left.C\left(E_{\overline{1}}\right)\right|_{-T}=i a^{\prime}\left(\bar{E}_{1}\right) \bar{z}_{1}^{\text {in }}
\end{array}, \quad\left\{\begin{array}{l}
\left.C\left(E_{1}\right)\right|_{T}=\bar{w}_{1}^{\text {out }} \\
\left.C\left(E_{\overline{1}}\right)\right|_{T}=i a^{\prime}\left(\bar{E}_{1}\right) \bar{z}_{1}^{\text {out }} \simeq 0
\end{array}\right.\right.
$$

so that:

$$
\Phi_{1}^{0}=w_{1}^{o u t} \bar{z}_{1}^{\text {in }}=\left.\left.\frac{1}{i a^{\prime}\left(\bar{E}_{1}\right)} C\left(E_{\overline{1}}\right)\right|_{-T} \overline{C\left(E_{1}\right)}\right|_{T}
$$

When $t \rightarrow-T$ the set of $X$ 's tending to $\infty$ is $\mathcal{I}=\overline{\mathcal{I}}_{0} \cup\{\overline{1}\}$, then:

$$
\left.C\left(E_{\overline{1}}\right)\right|_{X_{\overline{1}} \rightarrow \infty} \simeq-2 \frac{\left(E_{\overline{1}}-E_{1}\right)^{2}}{\prod_{i=1}^{n}\left(E_{\overline{1}}-\epsilon_{i}\right)} \frac{\prod_{k \in \mathcal{I}_{0}}\left(E_{\overline{1}}-E_{k}\right)^{2}}{X_{\overline{1}}}, \quad C\left(E_{1}\right) \simeq 0
$$

Similarly, when $t \rightarrow T$ the set of $X$ 's tending to $\infty$ is $\mathcal{I}=\overline{\mathcal{I}}_{0} \cup\{1\}$, then:

$$
\left.C\left(E_{1}\right)\right|_{X_{1} \rightarrow \infty} \simeq-2 \frac{\prod_{k \in \mathcal{I}_{0}}\left(E_{1}-E_{k}\right)^{2}}{\prod_{i=1}^{n}\left(E_{1}-\epsilon_{i}\right)} \frac{\left(E_{1}-E_{\overline{1}}\right)^{2}}{X_{1}}, \quad C\left(E_{\overline{1}}\right) \simeq 0
$$

Hence, we get:

$$
\Phi_{1}^{0}=\frac{4}{i a^{\prime}\left(E_{\overline{1}}\right)} \frac{\left(E_{\overline{1}}-E_{1}\right)^{4} \prod_{k \in \mathcal{I}_{0}}\left(E_{\overline{1}}-E_{\bar{k}}\right)^{2}\left(E_{\overline{1}}-E_{k}\right)^{2}}{\prod_{i}\left(E_{\overline{1}}-\epsilon_{i}\right)^{2}} \frac{1}{X_{\overline{1}}(-T) \overline{X_{1}}(T)}
$$

Next, we have

$$
\left\{\begin{array}{l}
X_{1}(T)=X_{1}(0) e^{\alpha_{1} T+i \beta_{1} T} \\
X_{\overline{1}}(-T)=X_{\overline{1}}(0) e^{\alpha_{1} T-i \beta_{1} T}
\end{array} \quad \Longrightarrow \quad X_{\overline{1}}(-T) \overline{X_{1}(T)}=-\frac{1}{4} e^{2 \alpha_{1} T-2 i \beta_{1} T}\right.
$$

so that eq.(41) becomes

$$
c_{1} \equiv \bar{z}_{1} w_{1}=e^{-\alpha_{1} \tau+i \beta_{1} \tau} \Phi_{1}^{0}(T)=\rho_{1} e^{i \gamma_{1}} e^{-\alpha_{1}(2 T+\tau)+i \beta_{1}(2 T+\tau)}
$$

where we have set

$$
\begin{equation*}
\rho_{1} e^{i \gamma_{1}}=-\frac{16}{i a^{\prime}\left(E_{\overline{1}}\right)} \frac{\left(E_{\overline{1}}-E_{1}\right)^{4} \prod_{k>1}\left(E_{\overline{1}}-E_{\bar{k}}\right)^{2}\left(E_{\overline{1}}-E_{k}\right)^{2}}{\prod_{k}\left(E_{\overline{1}}-\epsilon_{k}\right)^{2}} \tag{42}
\end{equation*}
$$

The quantity $2 T+\tau$ is the total time it takes to run once along the closed trajectory, that is the period of the motion. Normalizing it to $2 \pi$, the above equation determines $\alpha_{1}$ and $\beta_{1}$ :

$$
\begin{equation*}
\alpha_{1}=-\frac{1}{2 \pi} \log \frac{\left|c_{1}\right|}{\rho_{1}}, \quad \beta_{1}=\frac{1}{2 \pi}\left(\arg \left(c_{1}\right)-\gamma_{1}\right) \tag{43}
\end{equation*}
$$

### 5.2.2 Non diagonal case.

Similarly, one can compensate the non-diagonal motion induced by the large $\alpha_{1} K_{1}+\beta_{1} L_{1}$ flow by adding flows $\alpha_{j} K_{j}+\beta_{j} L_{j}$. The evolution of $z_{j}$ is now the composition of this new flow and the indirect evolution due to the $\alpha_{1} K_{1}+\beta_{1} L_{1}$ flow. Hence

$$
z_{j}^{\text {out }}=e^{\left(\alpha_{j}+i \beta_{j}\right)(2 T+\tau)} \rho_{z j}^{(1)} z_{j}^{i n}, \quad j \neq 1
$$

Periodicity requires

$$
\begin{equation*}
e^{-2 \pi\left(\alpha_{j}+i \beta_{j}\right)}=\rho_{z j}^{(1)} \tag{44}
\end{equation*}
$$

where we have set the period $2 T+\tau$ to $2 \pi$. This determines the remaining constants $\alpha_{j}, \beta_{j}$ in eq.(38), once we have computed $\rho_{z j}^{(1)}$ and $\rho_{w j}^{(1)}(j \neq 1)$ as defined in Section 7

The initial conditions were chosen such that when $t=-T$ the system lies on a subspace of the tangent cone whose natural coordinates are the $B\left(E_{j}\right)^{\text {in }}$ 's given by eq. (37) with $\mathcal{I}=\mathcal{I}_{0} \cup\{1\}$. In the long time limit $t=T$, i.e. when $a_{1}=\alpha_{1} t$ is positive and large, we get $\left|X_{\overline{1}}^{\text {out }}\right| \ll\left|X_{1}^{\text {out }}\right|$, and the system lies on a different subspace of the tangent cone. The $B\left(E_{j}\right)^{\text {out }}$ are again given by eq. (37) with $\mathcal{I}=\mathcal{I}_{0} \cup\{\overline{1}\}$. From the discussion in Section [7, we expect that:

$$
B\left(E_{j}\right)^{\text {out }}=\rho_{z j}^{(1)} B\left(E_{j}\right)^{\text {in }}, \quad j \neq 1
$$

By using eq.(37), this is indeed the case and we find :

$$
\begin{equation*}
\rho_{z j}^{(1)}=\left(\frac{E_{j}-E_{1}}{E_{j}-E_{\overline{1}}}\right)^{2} \tag{45}
\end{equation*}
$$

As explained in Section 7, we expect the general relation $\rho_{z j}^{(1)} \bar{\rho}_{w j}^{(1)}=1$ between $\rho_{z j}^{(1)}$ and $\rho_{w j}^{(1)}$. As a consistency check, we may evaluate $\rho_{w j}^{(1)}$ directly. Using that $w_{j}=\overline{C\left(E_{j}\right)}$, we have

$$
C\left(E_{j}\right)^{\text {out }}=\bar{\rho}_{w j}^{(1)} C\left(E_{j}\right)^{\text {in }}, \quad j \neq 1
$$

Now $C\left(E_{j}\right)^{\text {in }}$ is obtained by eq.(36) with the index set $\mathcal{I}=\overline{\mathcal{I}_{0}} \cup\{\overline{1}\}$ and $C\left(E_{j}\right)^{\text {out }}$ is obtained from the same formula with an index set $\mathcal{I}=\overline{\mathcal{I}_{0}} \cup\{1\}$. So we have:

$$
\begin{equation*}
\bar{\rho}_{w j}^{(1)}=\frac{C\left(E_{j}\right)^{\mathrm{out}}}{C\left(E_{j}\right)^{\mathrm{in}}}=\left(\frac{E_{j}-E_{\overline{1}}}{E_{j}-E_{1}}\right)^{2} \tag{46}
\end{equation*}
$$

Clearly, the relation $\rho_{z j}^{(1)} \bar{\rho}_{w j}^{(1)}=1$ is satisfied. Finally the constants $\alpha_{j}, \beta_{j}$ are determined by:

$$
\begin{equation*}
e^{-2 \pi\left(\alpha_{j}+i \beta_{j}\right)}=\left(\frac{E_{j}-E_{1}}{E_{j}-E_{\overline{1}}}\right)^{2} \tag{47}
\end{equation*}
$$

### 5.3 Action variables and symplectic invariants.

Let us define the 1 -form on the base (i.e. the image of the moment map)

$$
\Omega^{(1)}=\sum_{j} \Omega_{j}^{(1)}, \quad \Omega_{j}^{(1)}=\alpha_{j} d K_{j}+\beta_{j} d L_{j}
$$

Remembering that

$$
c_{1}=\bar{z}_{1}^{\text {in }} w_{1}^{\text {in }}=\bar{z}_{1}^{\text {out }} w_{1}^{\text {out }}=K_{1}+i L_{1}
$$

we see that $\Omega^{(1)}$ is singular when $c_{1} \rightarrow 0$. Isolating the singular part, we may write

$$
\Omega^{(1)}=-\frac{1}{4 \pi}\left(\log c_{1} d c_{1}+\log \bar{c}_{1} d \bar{c}_{1}\right)+\Omega^{(1) r e g}
$$

where

$$
\Omega^{(1) r e g}=\frac{1}{2 \pi} \log \rho_{1} d K_{1}-\frac{1}{2 \pi} \gamma_{1} d L_{1}+\sum_{j>1} \alpha_{j} d K_{j}+\beta_{j} d L_{j}
$$

The form $\Omega^{\text {reg }}$ contains all the regularized symplectic invariants of San Vũ Ngoc [4], computed here on the singular fiber. In the case of one spin, we find

$$
\rho_{1} e^{i \gamma_{1}}=-\frac{8}{s}\left(2 s-\epsilon_{1}\right)^{3 / 2}\left(\epsilon_{1}+i \sqrt{2 s-\epsilon_{1}^{2}}\right)
$$

Setting $s=1$ and $\epsilon_{1}=0$ yields

$$
\Omega^{(1) r e g}=\frac{1}{2 \pi}\left(5 \log 2 d K_{1}+\frac{\pi}{2} d L_{1}\right)
$$

and we recover the invariants computed in [6] (up to the normalization of the period to $2 \pi$ ).
To define the form $\Omega^{(1)}$, we need to introduce a cut in the $K_{1}, L_{1}$ plane. In this cut plane, the form $\Omega^{(1)}$ is closed and remembering eq.(38) we have

$$
\Omega^{(1)}=d \mathcal{H}^{(1)}
$$

The Hamiltonian $\mathcal{H}^{(1)}$ is the action associated to the closed trajectory we have constructed. It is defined on the cut plane $K_{1}, L_{1}$.

We have computed the action $\mathcal{H}^{(1)}$ to leading order. However following [5] we can give a more global interpretation using the spectral curve. First the normal actions $K_{j}, L_{j}$ are just the integrals over the vanishing cycles $A_{j}$ (encircling pairs of roots of the spectral polynomial which coalesce into double roots on the critical fiber) as shown in Fig. 5 .

$$
\begin{equation*}
K_{j}+i L_{j}=\frac{1}{2 i \pi} \oint_{A_{j}} \sqrt{\Lambda(\lambda)} d \lambda \tag{48}
\end{equation*}
$$

This can already be seen at the quadratic level

$$
\begin{aligned}
\frac{1}{2 i \pi} \oint_{A_{j}} \sqrt{\Lambda} d \lambda & =\frac{1}{2 i \pi} \oint_{A_{j}} \sqrt{a^{2}(\lambda)+\sum_{k} \frac{a(\lambda)}{a^{\prime}\left(E_{k}\right)\left(\lambda-E_{k}\right)} B_{k} C_{k}} d \lambda \\
& =\frac{1}{2 i \pi} \oint_{A_{j}}\left(a(\lambda)+\frac{1}{2} \sum_{j} \frac{1}{a^{\prime}\left(E_{k}\right)\left(\lambda-E_{k}\right)} B_{k} C_{k}\right) d \lambda=\frac{1}{2 a^{\prime}\left(E_{j}\right)} B_{j} C_{j}=K_{j}+i L_{j}
\end{aligned}
$$

Eq.(48) has a meaning beyond the quadratic approximation and can be viewed as the inverse of the function $\psi$ in Eliasson normal form [10, 5]. It is a non linear extension of eqs.(27, 28).


Figure 4: The motion of the three separated variables for the soliton solution in the 3 -spins model. The flow is given by eq.(39) with $\alpha_{k}, \beta_{k}$ given by eqs.(43|47). The limiting flow is obtained by setting $X_{j} \rightarrow X_{j}\left(\frac{t}{2 T}\right)$ where $T \simeq-\frac{1}{4 \pi} \log \left|c_{1}\right|$ is large, and $-2 \pi T \leq t \leq 2 \pi T$.

Concerning the action $\mathcal{H}^{(1)}$ we see that it has the unique property of being singular when $c_{1}=K_{1}+i L_{1} \rightarrow 0$ but it remains regular when $c_{j}=K_{j}+i L_{j} \rightarrow 0, j \neq 1$. To get an indication of its global definition, we may consider the example of the 3 -spins model and the motion of the separated variables of the soliton solution (see Appendix (9) corresponding to the limiting motion generated by $\mathcal{H}^{(1)}$ on the singular fiber. On Fig. 4 we show the motion of the separated coordinates $\lambda_{k}$ defined in eq.(69). Here $\lambda_{1}$ is associated to the sign (+) in eq.(78), while $\lambda_{2}, \lambda_{3}$ are associated to the sign ( - ). Hence the corresponding points ( $\mu_{k}, \lambda_{k}$ ) belong to different sheets in the representation of the spectral curve as a covering of the complex plane. On a fiber close to the singular one, the action integral is

$$
\mathcal{H}^{(1)}=\frac{1}{2 i \pi} \sum_{k} \oint \mu_{k} d \lambda_{k}=\frac{1}{2 i \pi} \sum_{k} \oint \sqrt{\Lambda\left(\lambda_{k}\right)} d \lambda_{k}
$$

where the $\lambda_{k}$ describe trajectories close to the critical ones. When $k=3$ we integrate over a trivial cycle and the contribution vanishes. The $k=1,2$ contributions combine to make a non trivial cycle $B_{1}$ as indicated on Fig. 5 .

$$
\begin{equation*}
\mathcal{H}^{(1)}=\frac{1}{2 i \pi} \oint_{B_{1}} \sqrt{\Lambda(\lambda)} d \lambda \tag{49}
\end{equation*}
$$

We can eliminate the $H_{i}$ between eq.(48) and eq.(49) thereby obtaining the full Taylor expansion


Figure 5: The spectral curve with the system of cuts and cycles corresponding to the limiting motion of the separated variables $\lambda_{k}$.
of $\mathcal{H}^{(1)}$ in terms of the $K_{j}, L_{j}$.
We believe that this can be generalized to the $n$ spin model. We have $2 m$ vanishing cycles $A_{j}$ and $2 m$ cycles $B_{j}$. The action integrals on the cycles $A_{j}$ give directly the $c_{j}=K_{j}+i L_{j}$ and the action integral on the $B_{j}$ give the periodic Hamiltonian $\mathcal{H}^{(j)}$ (see section 11) which become singular when $c_{j}$ tends to zero but remain regular when $c_{i}, i \neq j$ tends to zero. The important ingredient here is the identification of the $B_{j}$ cycles through the motion of the divisor of the proper soliton solution. A complete analysis of this motion has not yet been achieved, however using the asymptotic expressions for determinants given in section 10 it is possible to show that as the time $t_{1}$ goes to $-\infty$, we have one separated variable $\lambda_{1}^{+}$located at $E_{\overline{1}}$, and another one, $\lambda_{1}^{-}$, at $E_{1}$. When $t_{1}$ goes to $+\infty$, there is one separated variable $\lambda_{i}^{+}$located at $E_{1}$ and one $\lambda_{j}^{+}$at $E_{\overline{1}}$, The other remaining separated variables, $m-2$ of type $\lambda^{+}$and $m-1$ of type $\lambda^{-}$are found to occupy the same sets of points in the complex plane in both limits $t_{1} \rightarrow-\infty$ and $t_{1} \rightarrow \infty$ due to the periodicity conditions eq.(47). The difficult task remains to show that this possibly complicated motion of the separated variables can be continuously deformed into one of the $B_{j}$ cycles. If one succeeds in doing this, the periodic Hamiltonians can then be obtained more globally by the construction of section [11. An alternative approach would be to use Picard-Fuchs equations for the action variables as in 7 .

### 5.4 Monodromy

The orbits of the periodic motions defined in the previous section provide a basis of cycles on the Liouville tori close to the singular fiber. It is well known that when we follow these cycles by continuity when we perform a non trivial closed contour in the image of regular values of the moment map, the initial and final basis of cycles may be different and are related by an element of $G L(n, \mathbb{Z})$.

In order to compute the monodromy, we have first to determine the fundamental group $\pi_{1}$ of the manifold of regular values of the moment map. The manifold of singular values of the moment map is a stratified manifold by the rank of the moment map. At the equilibrium point
the rank is zero, however the equilibrium point belongs to bigger strata where the rank is $r \leq 2 m$. Eliasson's theorem allows an easy description of this stratification in a local neigborhood $U$ of the equilibrium point. The moment map is given by

$$
(x, p) \in U \rightarrow\left(K_{1}, L_{1}, K_{2}, L_{2}, \cdots, K_{m}, L_{m}\right)
$$

The derivatives of each pair $K_{j}, L_{j}$ generically span a two dimensional subspace of the tangent space of the image of the moment map near the equilibrium point and these spaces are linearly independent for two different pairs. For each pair the only way to span a space of dimension smaller than 2 is when the two tangent vectors are identically zero. In this case the rank of the momentum map drops to $r=2 m-2$ and the equation of this stratum is $K_{j}=L_{j}=0$. The singular values of the momentum map in a neighborhood of the equilibrium point consists therefore of $m$ strata $\mathcal{O}_{j}$ of dimension $2 m-2$ given by the equations $K_{j}=L_{j}=0$, for $i \in$ $\{1, \cdots m\}$. Of course each one of these strata is itself stratified by the intersection of $p$ such manifolds $\mathcal{O}_{j_{1}} \cap \mathcal{O}_{j_{2}} \cap \cdots \mathcal{O}_{j_{p}}$ where the rank is $r=2 m-2 p$. The equilibrium point itself is obtained when $p=m$.

Because the co-dimension of $\mathcal{O}_{j}$ is 2 , we can draw non trivial circles $S_{1}$ around it. Indeed $\mathcal{O}_{j}$ intersects the plane $K_{j}, L_{j}$ precisely at the origin and such an $S_{1}$ can be taken to be $K_{j}+i L_{j} \rightarrow$ $e^{i \alpha_{j}}\left(K_{j}+i L_{j}\right)$. The fundamental group of the image of the momentum map minus the $\mathcal{O}_{j}$ 's is a product of $m$ factors $S_{1}: \pi_{1}=S_{1} \times S_{1} \times \cdots S_{1}$.

We can now compute the monodromy. Locally we have $\Omega^{(1)}=d \mathcal{H}^{(1)}$. Recall that

$$
\bar{z}_{1} w_{1}=K_{1}+i L_{1}=\rho_{1} e^{i \gamma_{1}}
$$

The regular values of the moment map are such that $\rho_{1} \neq 0$. Monodromy is defined in the open set of the regular values of the moment map. Let us take as a non trivial closed path a circle around the origin in the $K_{1}, L_{1}$ plane: $c_{1}\left(\alpha_{1}\right)=e^{i \alpha_{1}} \rho_{1}, \gamma_{1} \leq \alpha_{1} \leq \gamma_{1}+2 \pi$. Integrating the form $\Omega^{(1)}$ over this circle one gets a non trivial contribution from $\Omega_{1}^{(1)}$

$$
\begin{aligned}
\oint_{C} \Omega^{(1)}=\oint_{C} \Omega_{1}^{(1)} & =-\frac{i}{4 \pi} \oint_{\gamma_{1}}^{2 \pi+\gamma_{1}}\left[c_{1}\left(\alpha_{1}\right) \log c_{1}\left(\alpha_{1}\right)-\bar{c}_{1}\left(\alpha_{1}\right) \log \bar{c}_{1}\left(\alpha_{1}\right)\right] d \alpha_{1} \\
& =\frac{1}{2 \pi} \rho_{1} \oint_{\gamma_{1}}^{2 \pi+\gamma_{1}} \alpha_{1} \cos \alpha_{1} d \alpha_{1}=\rho_{1} \sin \gamma_{1}=L_{1}
\end{aligned}
$$

This means that when we go around that circle, the Hamiltonian $\mathcal{H}^{(1)}$ becomes $\mathcal{H}^{(1)}+L_{1}$. The new periodic trajectory contains an extra $S^{1}$ factor generated by $L_{1}$. This is exactly the monodromy phenomenon. Note that when the image of the moment map winds once around the stratum $\mathcal{O}_{j}, \mathcal{H}^{(1)}$ remains single-valued when $j \neq 1$, so for $\mathcal{H}^{(1)}$ only paths which wind around the stratum $\mathcal{O}_{1}$ give rise to a non-trivial monodromy.

Following up the transformation of the action integrals $\mathcal{H}^{(j)}$ defined in the previous section along such a non trivial path would also lead to an interpretation of the monodromy in the framework of the Picard-Lefschetz theory as in [12].

## 6 Appendix : Normal Forms.

Normal coordinates simultaneously diagonalize all the flows $H_{j}$ in the quadratic approximation.

$$
\begin{equation*}
\left\{H_{j}, a_{\alpha}\right\}=E_{j, \alpha} a_{\alpha} \tag{50}
\end{equation*}
$$

To construct them, we return to eqs.(14- 16). They imply

$$
\left\{\frac{1}{2} \operatorname{Tr} L^{2}(\lambda), C(\mu)\right\}=\frac{2 i}{\lambda-\mu}(A(\lambda) C(\mu)-A(\mu) C(\lambda))
$$

When we expand around a configuration given by eq.(21), all the quantities $\left(b, \bar{b}, s_{j}^{+}, s_{j}^{-}\right)$are first order, but $s_{j}^{z}$ is second order because $s_{j}^{z}=s e_{j}-\frac{e_{j}}{2 s} s_{j}^{+} s_{j}^{-}+\cdots, \quad e_{j}= \pm 1$. In this expansion, the Hamiltonians are quadratic while $C(\mu)$ is first order. Therefore the Poisson bracket in the left-hand side is linear. Now $A(\lambda)$ is constant plus second order, so that in the right-hand side we can replace $A(\lambda)$ and $A(\mu)$ by their zeroth order expression :

$$
A(\lambda) \simeq a(\lambda)=2 \lambda+\sum_{j=1}^{n} \frac{s e_{j}}{\lambda-\epsilon_{j}}
$$

and we arrive at

$$
\begin{equation*}
\left\{\frac{1}{2} \operatorname{Tr} L^{2}(\lambda), C(\mu)\right\}=\frac{2 i}{\lambda-\mu}(a(\lambda) C(\mu)-a(\mu) C(\lambda)) \tag{51}
\end{equation*}
$$

Eq.(51) will be precisely of the form of eq.(50) if we can kill the unwanted term $C(\lambda)$. This is achieved by imposing the condition

$$
\begin{equation*}
a(\mu)=0, \quad \text { "Classical Bethe Equation" } \tag{52}
\end{equation*}
$$

This is an equation of degree $n+1$ for $\mu$. For non degenerate singularities it has only simple roots. Calling $E_{i}$ its solutions, we construct in this way $n+1$ variables $C\left(E_{i}\right)$. Clearly, they all Poisson commute

$$
\begin{equation*}
\left\{C\left(E_{i}\right), C\left(E_{j}\right)\right\}=0 \tag{53}
\end{equation*}
$$

Since phase space has dimension $2(n+1)$ this is half what we need. To construct the conjugate variables, we consider eq.(16). In our linear approximation it reads

$$
\left\{B\left(E_{i}\right), C\left(E_{j}\right)\right\}=\frac{2 i}{E_{i}-E_{j}}\left(a\left(E_{i}\right)-a\left(E_{j}\right)\right)
$$

If $E_{i}$ and $E_{j}$ are different solutions of eq.(52), then obviously

$$
\begin{equation*}
\left\{B\left(E_{i}\right), C\left(E_{j}\right)\right\}=0, \quad E_{i} \neq E_{j} \tag{54}
\end{equation*}
$$

If however $E_{j}=E_{i}$ then

$$
\begin{equation*}
\left\{B\left(E_{i}\right), C\left(E_{i}\right)\right\}=2 i a^{\prime}\left(E_{i}\right) \tag{55}
\end{equation*}
$$

Finally, we also have

$$
\begin{equation*}
\left\{B\left(E_{i}\right), B\left(E_{j}\right)\right\}=0 \tag{56}
\end{equation*}
$$

Up to normalisation, we have indeed constructed canonical coordinates.
It is simple to express the quadratic Hamiltonians in theses coordinates:

$$
\begin{equation*}
\frac{1}{2} \operatorname{Tr} L^{2}(\lambda)=a^{2}(\lambda)+\sum_{j} \frac{a(\lambda)}{a^{\prime}\left(E_{j}\right)\left(\lambda-E_{j}\right)} B\left(E_{j}\right) C\left(E_{j}\right) \tag{57}
\end{equation*}
$$

This has the correct analytical properties in $\lambda$ and together with the Poisson brackets eqs. 5354 55. (56) we reproduce eq.(51). Note that there is no pole at $\lambda=E_{j}$ because $a\left(E_{j}\right)=0$. Expanding around $\lambda=\infty$ we get

$$
H_{n+1}=s \sum_{k} e_{k}+\sum_{i} \frac{1}{2 a^{\prime}\left(E_{i}\right)} B\left(E_{i}\right) C\left(E_{i}\right)
$$

and computing the residue at $\lambda=\epsilon_{j}$, we find

$$
H_{j}=s e_{j}\left[2 \epsilon_{j}+\sum_{k} \frac{s e_{k}}{\epsilon_{j}-\epsilon_{k}}\right]+\sum_{i} \frac{1}{2 a^{\prime}\left(E_{i}\right)} \frac{s e_{j}}{\epsilon_{j}-E_{i}} B\left(E_{i}\right) C\left(E_{i}\right)
$$

We can invert these formulae: devide eq.(57) by $\lambda-E_{j}$ and take the residue at $\lambda=E_{j}$. Since $a\left(E_{j}\right)=0$ we get

$$
\frac{1}{2} \operatorname{Tr} L^{2}\left(E_{j}\right)=B\left(E_{j}\right) C\left(E_{j}\right)
$$

or explicitly

$$
B\left(E_{j}\right) C\left(E_{j}\right)=4 E_{j}^{2}+4 H_{n+1}+\sum_{i=1}^{n} \frac{2 H_{i}}{E_{j}-\epsilon_{i}}+\sum_{k=1}^{n} \frac{s^{2}}{\left(E_{j}-\epsilon_{i}\right)^{2}}
$$

We can now make contact with the Williamson classification theorem.
If $E_{j}$ is real we have $B\left(E_{j}\right)=\overline{C\left(E_{j}\right)}$ and we set

$$
C\left(E_{j}\right)=\sqrt{\left|a^{\prime}\left(E_{j}\right)\right|}\left(p_{j}+i \eta_{j} q_{j}\right), \quad B\left(E_{j}\right)=\sqrt{\left|a^{\prime}\left(E_{j}\right)\right|}\left(p_{j}-i \eta_{j} q_{j}\right)
$$

where $\eta_{j}=-\operatorname{sign}\left(a^{\prime}\left(E_{j}\right)\right)$ and $p_{j}, q_{j}$ are real canonical coordinates. Then eq.(55) is satisfied. Moreover

$$
B\left(E_{j}\right) C\left(E_{j}\right)=\left|a^{\prime}\left(E_{j}\right)\right|\left(p_{j}^{2}+q_{j}^{2}\right)
$$

i.e. we have an elliptic singularity.

If $E_{j}$ is complex, there is another root $E_{j+1}=\bar{E}_{j}$ which is its complex conjugate. Then $\overline{B\left(E_{j}\right)}=C\left(E_{j+1}\right)$. We introduce real canonical coordinates $p_{j}, q_{j}, p_{j+1}, q_{j+1}$ and set

$$
\begin{aligned}
& B\left(E_{j}\right)=-i a^{\prime}\left(E_{j}\right)\left(q_{j}+i q_{j+1}\right), \quad B\left(E_{j+1}\right)=p_{j}+i p_{j+1} \\
& C\left(E_{j}\right)=p_{j}-i p_{j+1}, \quad C\left(E_{j+1}\right)=i a^{\prime}\left(E_{j+1}\right)\left(q_{j}-i q_{j+1}\right)
\end{aligned}
$$

then eqs. $5315455(56)$ are satisfied and

$$
\begin{aligned}
B\left(E_{j}\right) C\left(E_{j}\right) & =-i a^{\prime}\left(E_{j}\right)\left(K_{j}+i L_{j}\right) \\
B\left(E_{j+1}\right) C\left(E_{j+1}\right) & =i a^{\prime}\left(E_{j+1}\right)\left(K_{j}-i L_{j}\right)
\end{aligned}
$$

i.e. we have a focus-focus singularity.

## 7 Appendix: In-Out transformation.

We consider a flow generated by the Hamiltonian $\alpha_{1} K_{1}+\beta_{1} L_{1}$. Such linear combinations were needed in Section [5] although here only the $K_{1}$ part really matters. We assume $\alpha_{1}>0$. Let us fix some $\delta>0$ with $\delta$ small enough that Eliasson map $\Phi$ is well defined in the ball of radius $2 \delta$ around the origin. Consider the point $U:\left(z_{1}=\delta, w_{1}=0\right)$ on the expanding manifold and the trajectory of the flow starting at that point. Since the trajectory stays on the global unstable manifold, which is equal to the global stable manifold, it will reappear after a time lapse $2 T$ on the contracting manifold at a point $S:\left(z_{1}=0,\left|w_{1}\right|=\delta\right)$. The flow evaluated at time $2 T$ is well defined in a neighborhood $\Sigma_{U}$ of the point $U$ and defines a symplectic mapping $F_{T}^{(1)}$ of $\Sigma_{U}$, $\left|w_{1}^{\text {in }}\right| \ll\left|z_{1}^{\text {in }}\right|$, to a neighborhood $\Sigma_{S}$ of the point $S$ in the contracting subspace $\left|z_{1}^{\text {out }}\right| \ll\left|w_{1}^{\text {out }}\right|$ :

$$
\Sigma_{U}:\left(z_{1}^{\text {in }}, w_{1}^{\text {in }}, z_{2}^{i n}, w_{2}^{\text {in }}, \cdots, z_{m}^{\text {in }}, w_{m}^{\text {in }}\right) \xrightarrow{F_{T}^{(1)}} \Sigma_{S}:\left(z_{1}^{\text {out }}, w_{1}^{\text {out }}, z_{2}^{\text {out }}, w_{2}^{\text {out }}, \cdots, z_{m}^{\text {out }}, w_{m}^{\text {out }}\right)
$$

We also assume that the other components $z_{i}^{\text {in }}, w_{i}^{\text {in }}, z_{i}^{\text {out }}, w_{i}^{\text {out }}, i=2, \cdots, m$ are small (this will be satisfied in our case) so that the initial and final points lie in a small neighborhood of the critical point. One can find the precise form of this mapping by writing that it commutes with the flows generated by $K_{j}, L_{j}, j=1 \cdots m$. Let us fix the level set by imposing:

$$
\begin{equation*}
\bar{z}_{j}^{i n} w_{j}^{i n}=\bar{z}_{j}^{o u t} w_{j}^{o u t}=K_{j}+i L_{j} \tag{58}
\end{equation*}
$$

Let us also assume that $K_{j}+i L_{j} \neq 0$. We may then choose $\left(w_{1}^{i n}, \cdots, w_{m}^{i n}\right)$ as coordinates on the intersection of $\Sigma_{U}$ with the level set and write:

$$
w_{j}^{o u t}=F_{j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}, w_{1}^{i n}, \bar{w}_{1}^{i n}, \cdots, w_{m}^{i n}, \bar{w}_{m}^{i n}\right)
$$

The condition that the transformation $F^{(1)}$ commutes with the flows $K_{1}, L_{1}, \cdots, K_{m}, L_{m}$ means:

$$
\begin{aligned}
& e^{-a_{j}+i \theta_{j}} F_{j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}, w_{1}, \bar{w}_{1}, \cdots, w_{m}, \bar{w}_{m}\right)= \\
& F_{j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}, e^{-a_{1}+i \theta_{1}} w_{1}, e^{-a_{1}-i \theta_{1}} \bar{w}_{1}, \cdots, e^{-a_{m}+i \theta_{m}} w_{m}, e^{-a_{m}-i \theta_{m}} \bar{w}_{m}\right)
\end{aligned}
$$

To solve this equation we choose the parameters $a_{k}, \theta_{k}$ such that

$$
e^{-a_{k}+i \theta_{k}} w_{k}=\eta_{k}, \quad w_{k}, \eta_{k} \neq 0, \quad 1 \leq k \leq m
$$

where $\eta_{k}$ are constants. Then the equation becomes:

$$
F_{j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}, w_{1}, \bar{w}_{1}, \cdots, w_{m}, \bar{w}_{m}\right)=\rho_{w j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}\right) w_{j}
$$

where $\rho_{w j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}\right)=\eta_{j}^{-1} F_{j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}, \eta_{1}, \bar{\eta}_{1}, \cdots, \eta_{m}, \bar{\eta}_{m}\right)$. Hence, we have shown that

$$
\begin{equation*}
w_{j}^{o u t}=\rho_{w j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}\right) w_{j}^{i n} \tag{59}
\end{equation*}
$$

where the super script (1) is to remind that this is a transformation under the $\alpha_{1} K_{1}+\beta_{1} L_{1}$ flow. In exactly the same way, we show

$$
\begin{equation*}
z_{j}^{o u t}=\rho_{z j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}\right) z_{j}^{i n} \tag{60}
\end{equation*}
$$

The conditions eq.(58) then imply

$$
\begin{equation*}
\rho_{w j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}\right) \overline{\rho_{z j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}\right)}=1 \tag{61}
\end{equation*}
$$

When $j \neq 1$, the functions $\rho_{w j}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}\right)$ do not depend on the choice of the time lapse $2 T$ for the flow, nor on the choice of the parameters $\alpha_{1}, \beta_{1}$. They are smooth in a neighborhood of the critical torus. As shown before [4, 5] they are symplectic invariants attached to the Liouville foliation in the vicinity of the singular torus. Let us now show that for $j=1$, $\rho_{w 1}^{(1)}\left(K_{1}, L_{1}, \cdots K_{m}, L_{m}\right)$ is singular as $K_{1}+i L_{1} \rightarrow 0$.

To see this, we first consider the critical torus, on which the dilating stratum is given by $w_{1}=0$ and the contracting stratum by $z_{1}=0$. We start with an initial condition on the dilating stratum $\left(w_{1, \text { sol }}^{i n}=0, z_{1, \text { sol }}^{i n}\right)$. We let the system evolve with respect to the flow . After a large motion during the time lapse $2 T$, the point reappears on the contracting stratum $\left(w_{1, \text { sol }}^{\text {out }}, z_{1, \text { sol }}^{\text {out }}=\right.$ $0)$. The same reasoning as before on the invariance under the flows $K_{j}, L_{j}$ shows that:

$$
\begin{equation*}
w_{1, s o l}^{o u t}=\frac{\Phi_{1}^{0}}{\bar{z}_{1, \text { sol }}^{i n}} \tag{62}
\end{equation*}
$$

where $\Phi_{1}^{0}$ is a constant, which depends only on the combination $\alpha_{1} T$. Let us now consider a torus close to the critical torus. The point $\left(z_{1, \text { sol }}^{i n}, w_{1, \text { sol }}^{i n}=0\right)$ is displaced to $\left(z_{1}^{i n}, w_{1}^{i n}\right)$ and $\left(z_{1, \text { sol }}^{\text {out }}=0, w_{1, \text { sol }}^{\text {out }}\right)$ is displaced to $\left(z_{1}^{\text {out }}, w_{1}^{\text {out }}\right)$. The small parameters can be chosen to be $w_{j}^{\text {in }}$ ( $1 \leq j \leq m$ ) and we can write:

$$
\begin{align*}
& z_{1}^{\text {out }}=\left(\frac{z_{1}^{i n} \bar{w}_{1}^{\text {in }}}{\overline{\Phi_{1}^{0}}}\right) z_{1}^{i n}+O\left(w^{2}\right)  \tag{63}\\
& w_{1}^{\text {out }}=\frac{\Phi_{1}^{0}}{\bar{z}_{1}^{\text {in }}+O(1) w_{1}^{\text {in }}+O\left(w^{2}\right)=\left(\frac{\Phi_{1}^{0}}{\bar{z}_{1}^{i n} w_{1}^{\text {in }}}+O(1)\right) w_{1}^{\text {in }}+O\left(w^{2}\right)} \tag{64}
\end{align*}
$$

where we have imposed the condition $\bar{z}_{1}^{\text {in }} w_{1}^{\text {in }}=\bar{z}_{1}^{\text {out }} w_{1}^{\text {out }}=K_{1}+i L_{1}$. Connecting these equations with the definition (59) we see that:

$$
\rho_{w 1}^{(1)} \simeq \frac{\Phi_{1}^{0}}{K_{1}+i L_{1}}+O(1)
$$

in the vicinity of the singular torus. As explained in 4], a regularization procedure is needed in order to define the diagonal symplectic invariant.

## 8 Appendix: Symplectic invariance.

In this Appendix we explain why the quantity $\Omega^{\text {reg }}$ computed in Section 5.3 has an intrinsic meaning and is invariant under symplectic transformations preserving the singular Liouville foliation. We extend to higher dimensions the arguments developed in [4.

Recall that two singular foliations $\mathcal{F}$ and $\mathcal{F}^{\prime}$ in the symplectic manifolds $M, \omega$ and $M^{\prime}, \omega^{\prime}$ are equivalent if there exists a symplectic diffeomorphism $\varphi: \mathcal{F} \rightarrow \mathcal{F}^{\prime}$ sending leaves to leaves.

Having two such foliations, by Eliasson theorem, we can define for each of them normal coordinates $\left(z_{j}, w_{j}\right)$ and $\left(z_{j}^{\prime}, w_{j}^{\prime}\right), j=1 \cdots m$, valid in neighborhoods $U$ and $U^{\prime}$ of the focus-focus point. We can also define the moment maps $\mu^{(0)}:\left(z_{j}, w_{j}\right) \rightarrow\left(K_{j}, L_{j}\right)$ and $\mu^{(0)^{\prime}}:\left(z_{j}^{\prime}, w_{j}^{\prime}\right) \rightarrow$ $\left(K_{j}^{\prime}, L_{j}^{\prime}\right), j=1 \cdots m$ where the Hamiltonians $K_{i}, L_{i}$ and $K_{i}^{\prime}, L_{i}^{\prime}$ are as in Section 3. In the neighborhoods $U$ and $U^{\prime}$, the study of the symplectic diffeomorphism $\varphi$ therefore reduces to the study of the commutative diagram

where $\Phi$ is a symplectomorphism and $\Psi$ is a diffeomorphism. We will show that the diffeomorphism $\Psi$ is extremely constrained and such that $L_{i}=\zeta_{i} L_{i}^{\prime}$ with $\zeta_{i}= \pm 1$ and $K_{i}=K_{i}^{\prime}$, up to a flat function (vanishing at the origin together with all its derivatives) so that the Taylor expansion of $\Omega^{r e g}$ has an invariant meaning.

The first observation is that the rank of the moment map is preserved under $\Phi$. Manifolds where the rank of $\mu^{(0)}$ is $r$ are mapped to manifolds where the rank of $\mu^{(0)^{\prime}}$ is $r$. Manifolds of rank 2 are the (complex) coordinate planes $z_{i}, w_{i}$, characterized by the equations $c_{i} \equiv K_{i}+i L_{i} \neq 0$ and $c_{j}=0$ for $j \neq i$. Therefore the map $\Psi$ sends the coordinate axis (in complex notation) $c_{i} \neq 0$ to a coordinate axis $c_{i^{\prime}}^{\prime} \neq 0$. This defines a map $\sigma: i^{\prime}=\sigma(i)$ of the set $\{1,2 \cdots m\}$ to itself which is a permutation. We can assume that this permutation is the identity.

One can generalize this argument. Manifolds of rank 4 are the characterized by $c_{i} \neq 0, c_{j} \neq 0$ and all other $c_{k}=0$. Therefore the map $\Psi$ sends the plane $c_{i} \neq 0, c_{j} \neq 0$ and all other $c_{k}=0$ to a similar plane $c_{r}^{\prime} \neq 0, c_{s}^{\prime} \neq 0$ and all other $c_{l}^{\prime}=0$. But that plane must contain the axes $c_{i}^{\prime} \neq 0$ and $c_{j}^{\prime} \neq 0$, hence it is the whole plane $c_{i}^{\prime} \neq 0, c_{j}^{\prime} \neq 0$ and all other $c_{k}^{\prime}=0$. Repeating the argument, manifolds of rank $2 m-2$ are hyperplanes $z_{i}=0, w_{i}=0$ characterized by the equations $c_{i}=K_{i}+i L_{i}=0$ and $c_{j} \neq 0$ for $j \neq i$ and the map $\Psi$ sends the hyperplane $c_{i}=0$ to the hyperplane $c_{i}^{\prime}=0$.

The symplectomorphism $\Phi$ sends a periodic flow to a periodic flow of the same period. In particular, the flow generated by $L_{i}$ is periodic. Let us denote by $\left(K_{i}^{\prime}(K, L), L_{i}^{\prime}(K, L)\right)$ the coordinates of the image by $\Psi$ of the point of coordinates $\left(K_{i}, L_{i}\right)$. The Hamiltonian vector field of $L_{i}^{\prime}$ is

$$
X_{L_{i}^{\prime}}=\sum_{k} \frac{\partial L_{i}^{\prime}}{\partial L_{k}} X_{L_{k}}+\frac{\partial L_{i}^{\prime}}{\partial K_{k}} X_{K_{k}}
$$

However, on the open sets $U$ and $U^{\prime}$, the only periodic flows are those generated by the $L_{k}$ and $L_{k}^{\prime}$. Hence $\frac{\partial L_{i}^{\prime}}{\partial K_{k}}=0$ and $\frac{\partial L_{i}^{\prime}}{\partial L_{k}}=A_{i k}$ where $A_{i k}$ is an invertible matrix of $G L_{m}(\mathbb{Z})$. Therefore

$$
L_{i}^{\prime}=\sum_{k} A_{i k} L_{k}
$$

On the other hand, we have just seen that if only $c_{i}$ is different from zero, then only $c_{i}^{\prime}$ is different from zero. This implies that the matrix $A$ is diagonal, and since it is an element of $G L_{m}(\mathbb{Z})$, it is of the form $A_{i k}=\zeta_{i} \delta_{i j}, \zeta_{i}= \pm 1$. We have shown that

$$
\begin{equation*}
L_{i}^{\prime}=\zeta_{i} L_{i}, \quad \zeta_{i}= \pm 1 \tag{66}
\end{equation*}
$$

Next we want to analyse the functions

$$
K_{i}^{\prime}=K_{i}^{\prime}\left(K_{1}, L_{1} \cdots K_{m}, L_{m}\right)
$$

The remarks on the preservation of the rank of the moment map put constraints on the form of the Taylor expansions of these functions. First, because the manifold $c_{i}=0, c_{j} \neq 0, j \neq i$ is mapped to the manifold $c_{i}^{\prime}=0, c_{j}^{\prime} \neq 0, j \neq i$, we can write

$$
\begin{equation*}
K_{i}^{\prime}=A_{i}\left(K_{1}, L_{1} \cdots K_{m}, L_{m}\right) K_{i}+B_{i}\left(K_{1}, L_{1} \cdots K_{m}, L_{m}\right) L_{i} \tag{67}
\end{equation*}
$$

where the functions $A_{i}, B_{i}$ are regular at 0 . Moreover, since the axis $c_{i} \neq 0, c_{j}=0, j \neq i$ is mapped on the axis $c_{i}^{\prime} \neq 0, c_{j}^{\prime}=0, j \neq i$ we have

$$
A_{i}\left(K_{1}, L_{1} \cdots K_{m}, L_{m}\right)=\alpha_{i}+O(K, L), \quad B_{i}\left(K_{1}, L_{1} \cdots K_{m}, L_{m}\right)=\beta_{i}+O(K, L)
$$

where $\alpha_{i}, \beta_{i}$ are constants such that $\alpha_{i} \beta_{i} \neq 0$. We show below the much stronger results $\alpha_{i}+O(K, L)=1+$ flat and $\beta_{i}+O(K, L)=0+$ flat.

Let us assume for simplicity that $i=1$. Let us consider, in the plane $\left(z_{1}^{\prime}, w_{1}^{\prime}\right)$, the points on the critical torus $m_{0}^{\prime \text { in }}:\left(w_{1}^{\prime}=\delta, z_{1}^{\prime}=0\right)$ and $m_{0}^{\prime o u t}:\left(w_{1}^{\prime}=0, z_{1}^{\prime}=\delta\right)$ (all other coordinates vanish). The points $m_{0}^{\text {in }}=\Phi^{-1}\left({m^{\prime \prime}}_{0}^{\prime \text { in }}\right)$ and $m_{0}^{\text {out }}=\Phi^{-1}\left(m_{0}^{\prime o u t}\right)$ are in the plane $\left(z_{1}, w_{1}\right)$ on the critical torus and we can assume $m_{0}^{\text {in }}:\left(w_{1}=a, z_{1}=0\right), m_{0}^{\text {out }}:\left(w_{1}=0, z_{1}=b\right)$ (all other coordinate vanish).

Let us now consider a small neighborhood of $m_{0}^{\prime i n}$ and let us apply at each of its points the flow generated by $\left(K_{1}^{\prime}, L_{1}^{\prime}\right)$ during times $\left(a_{1}^{\prime}, \theta_{1}^{\prime}\right)$ chosen in such a way that the point $m^{\prime \text { in }}:\left(z_{1}^{\prime \text { in }}=\bar{c}_{1}^{\prime} / \delta, w_{1}^{\prime \text { in }}=\delta, z_{j}^{\prime \text { in }}, w_{j}^{\prime \text { in }}\right)$ is sent to $m^{\prime \text { out }}:\left(z_{1}^{\text {out }}=\delta, w_{1}^{\prime \text { out }}=c_{1}^{\prime} / \delta, z_{j}^{\prime \text { out }}=\right.$ \left.${z_{j}^{\prime}}^{\text {in }}, w_{j}^{\prime \text { out }}=w_{j}^{\prime \text { in }}\right)$ i.e.

$$
e^{-a_{1}^{\prime}+i \theta_{1}^{\prime}}=\frac{c_{1}^{\prime}}{\delta^{2}}
$$

In particular, when $\left|c_{1}^{\prime}\right| \rightarrow 0$, we find

$$
a_{1}^{\prime} \rightarrow-\log \left|c_{1}^{\prime}\right| \simeq-\log \left|c_{1}\right|
$$

where the last relation comes from eqs. 66[67). This map sends $m_{0}^{\prime \text { in }}$ to $m_{0}^{\prime \text { out }}$ and extends to a smooth map $F^{\prime}$ from a neighborhood of $m_{0}^{\prime \text { in }}$ to a neighborhood of $m_{0}^{\prime o u t}$. The map $F=\Phi^{-1} \circ F^{\prime} \circ \Phi$ sends a neighborhood of $m_{0}^{i n}$ to a neighborhood of $m_{0}^{\text {out }}$ and this map, being a composition of $C^{\infty}$ maps, should be $C^{\infty}$. On the other hand, the map $F$ is obtained by applying the flow

$$
\sum_{j=1}^{m} a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial K_{j}} X_{K_{j}}+\theta_{1}^{\prime} X_{L_{1}}+\sum_{j=1}^{m} a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial L_{j}} X_{L_{j}}
$$

On $\left(z_{1}, w_{1}\right)$, this gives

$$
\begin{aligned}
z_{1}^{\text {out }} & =\exp \left[a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial K_{1}}+i\left(\theta_{1}^{\prime}+a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial L_{1}}\right)\right] z_{1}^{i n} \\
w_{1}^{\text {out }} & =\exp \left[-a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial K_{1}}+i\left(\theta_{1}^{\prime}+a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial L_{1}}\right)\right] w_{1}^{i n}
\end{aligned}
$$

When $\left|c_{1}\right| \rightarrow 0$, we should have $z^{i n} \rightarrow \frac{\bar{c}_{1}}{\bar{a}}, z^{\text {out }} \rightarrow b$. Inserting into the first equation, we find

$$
\bar{a} b=\exp \left[\log \bar{c}_{1}+a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial K_{1}}+i\left(\theta_{1}^{\prime}+a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial L_{1}}\right)\right]
$$

(for the other equation $w^{i n} \rightarrow a w^{o u t} \rightarrow \frac{c_{1}}{b}$, we obtain the same condition). Using the values of $a_{1}^{\prime}, \theta_{1}^{\prime}$, we get in the limit $\left|c_{1}\right| \rightarrow 0$

$$
\left.\bar{a} b=\exp \left[\log \mid c_{1}\right]\left(1-\frac{\partial K_{1}^{\prime}}{\partial K_{1}}-i \frac{\partial K_{1}^{\prime}}{\partial L_{1}}\right)\right]
$$

Finiteness when $\left.\mid c_{1}\right] \rightarrow 0$ requires $\alpha_{1}=1, \beta_{1}=0$. One can push the argument and show that smoothness of the map in the limit $c_{1} \rightarrow 0$ gives

$$
\frac{\partial K_{1}^{\prime}}{\partial K_{1}}=1+\text { flat }, \quad \frac{\partial K_{1}^{\prime}}{\partial L_{1}}=\text { flat }
$$

the argument is the same as in [4].
On $\left(z_{j}, w_{j}\right), j \neq 1$, we get

$$
\begin{aligned}
z_{j}^{\text {out }} & =\exp \left[a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial K_{j}}+i\left(a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial L_{j}}\right)\right] z_{j}^{\text {in }} \\
w_{j}^{\text {out }} & =\exp \left[-a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial K_{j}}+i\left(a_{1}^{\prime} \frac{\partial K_{1}^{\prime}}{\partial L_{j}}\right)\right] w_{j}^{i n}
\end{aligned}
$$

We require that $\left(z_{j}^{\text {out }}, w_{j}^{\text {out }}\right)$ be finite when $c_{1}$ tends to zero with $\left(z_{j}^{i n}, w_{j}^{i n}\right)$ finite. This gives

$$
\frac{\partial K_{1}^{\prime}}{\partial K_{j}}=\text { flat }, \quad \frac{\partial K_{1}^{\prime}}{\partial L_{j}}=\text { flat }
$$

## 9 Appendix: Solitons.

In this appendix we recall the construction of solitons in the Jaynes-Cummings-Gaudin model 13 . We present it in some details because we need to generalize it to include all times $t_{i}$ and we wish to present in a self contained way the formulae we use in our calculation of the symplectic invariants.

The Lax form eq. (10) of the equation of motion implies that the so-called spectral curve $\Gamma$, defined by $\operatorname{det}(L(\lambda)-\mu)=0$ is a constant of motion. Specifically:

$$
\begin{equation*}
\Gamma: \mu^{2}-A^{2}(\lambda)-B(\lambda) C(\lambda)=0, \text { i.e. } \mu^{2}=\frac{Q_{2 n+2}(\lambda)}{\prod_{j}\left(\lambda-\epsilon_{j}\right)^{2}} \tag{68}
\end{equation*}
$$

Defining $y=\mu \prod_{j}\left(\lambda-\epsilon_{j}\right)$, the equation of the curve becomes $y^{2}=Q_{2 n+2}(\lambda)$ which is an hyperelliptic curve. Since the polynomial $Q_{2 n+2}(\lambda)$ has degree $2 n+2$, the genus of the curve in $n$.

The separated variables are $g=n$ points on the curve whose coordinates $\left(\lambda_{k}, \mu_{k}\right)$ can be taken as coordinates on phase space. They are defined as follows. Let us write

$$
\begin{equation*}
C(\lambda)=2 \bar{b}+\sum_{j=1}^{n} \frac{s_{j}^{+}}{\lambda-\epsilon_{j}} \equiv 2 \bar{b} \frac{\prod_{k=1}^{n}\left(\lambda-\lambda_{k}\right)}{\prod_{j=1}^{n}\left(\lambda-\epsilon_{j}\right)} \tag{69}
\end{equation*}
$$

the separated coordinates are the collection $\left(\lambda_{k}, \mu_{k}=A\left(\lambda_{k}\right)\right)$. They have canonical Poisson brackets

$$
\begin{equation*}
\left\{\lambda_{k}, \mu_{k^{\prime}}\right\}=-i \delta_{k, k^{\prime}} \tag{70}
\end{equation*}
$$

Notice however that if $\bar{b}=0$ theses variables are not well defined.
There are only $2 n$ such coordinates which turn out to be invariant under the global $U(1)$ rotation generated by $H_{n+1}$ :

$$
b \rightarrow e^{i \theta} b, \quad \bar{b} \rightarrow e^{-i \theta} \bar{b}, \quad s_{j}^{-} \rightarrow e^{i \theta} s_{j}^{-}, \quad s_{j}^{+} \rightarrow e^{-i \theta} s_{j}^{+}
$$

So they describe the reduced model obtained by fixing the value of $H_{n+1}$ and taking into consideration only the dynamical variables invariant under this $U(1)$ action. The initial dynamical model can be recovered by adding the phase of the oscillator coordinates $\bar{b}, b$ to the separated variables.

We show now how to reconstruct spin coordinates from the separated variables $\left(\lambda_{k}, \mu_{k}\right)$ and the oscillator coordinates $(\bar{b}, b)$. For $C(\lambda)$ we have eq.(69). It is a rational fraction of $\lambda$ which has simple poles at $\lambda=\epsilon_{j}$ whose residue is $s_{j}^{+}$. For $A(\lambda)$, we write:

$$
A(\lambda)=\frac{P_{n+1}(\lambda)}{\prod_{j=1}^{n}\left(\lambda-\epsilon_{j}\right)}
$$

The polynomial $P_{n+1}(\lambda)-2 \lambda \prod_{j=1}^{n}\left(\lambda-\epsilon_{j}\right)$ is of degree $n-1$, and we know its value at the $n$ points $\lambda_{j}$ because:

$$
\begin{equation*}
P_{n+1}\left(\lambda_{j}\right)=\mu_{j} \prod_{k=1}^{n}\left(\lambda_{j}-\epsilon_{k}\right) \tag{71}
\end{equation*}
$$

Therefore, we can write:

$$
\begin{equation*}
P_{n+1}(\lambda)=2 \lambda \prod_{j=1}^{n}\left(\lambda-\epsilon_{j}\right)+\sum_{i}\left(\mu_{i}-2 \lambda_{i}\right) \prod_{k=1}^{n}\left(\lambda_{i}-\epsilon_{k}\right) \frac{\prod_{l \neq i}\left(\lambda-\lambda_{l}\right)}{\prod_{l \neq i}\left(\lambda_{i}-\lambda_{l}\right)} \tag{72}
\end{equation*}
$$

Once $A(\lambda)$ and $C(\lambda)$ are known, $B(\lambda)$ is determined by:

$$
\begin{equation*}
B(\lambda)=\frac{Q_{2 n+2}(\lambda)-P_{n+1}^{2}(\lambda)}{C(\lambda) \prod_{k}\left(\lambda-\epsilon_{k}\right)^{2}}=\frac{1}{2 \bar{b}} \frac{Q_{2 n+2}(\lambda)-P_{n+1}^{2}(\lambda)}{\prod_{i}\left(\lambda-\lambda_{i}\right) \prod_{k}\left(\lambda-\epsilon_{k}\right)} \tag{73}
\end{equation*}
$$

The polynomial in the numerator is of degree $2 n$, and moreover it is divisible by $\prod_{i}\left(\lambda-\lambda_{i}\right)$ because $P_{n+1}^{2}\left(\lambda_{i}\right)=Q_{2 n+2}\left(\lambda_{i}\right)$, so we can write

$$
\begin{equation*}
B(\lambda)=2 b \frac{\prod_{i}\left(\lambda-\bar{\lambda}_{i}\right)}{\prod_{k}\left(\lambda-\epsilon_{k}\right)} \tag{74}
\end{equation*}
$$

In this construction, the variables $\bar{\lambda}_{i}$ and the corresponding $\bar{\mu}_{i}$ are complicated functions of the $\left(\lambda_{i}, \mu_{i}\right)$. The set of physical configurations (often refered to here as the real slice) is obtained by writing that the set $\left(\overline{\lambda_{i}}, \overline{\mu_{i}}\right)$ is the complex conjugate of the set $\left(\lambda_{i}, \mu_{i}\right)$. This leads to a set of complicated relations whose solution is not known in general, but that we will solve in the soliton case.

The Hamiltonians $H_{j}$ of the reduced model are obtained by writing that the points $\left(\lambda_{k}, \mu_{k}\right)$ belong to the spectral curve. We get a linear system of equations $\sum_{j} \mathcal{B}_{k j} H_{j}=V_{k}$ where

$$
\mathcal{B}_{k j}=\frac{1}{\lambda_{k}-\epsilon_{j}}, \quad V_{k}=\frac{1}{2}\left(\mu_{k}^{2}-\lambda_{k}^{2}-4 H_{n+1}-\sum_{j} \frac{s^{2}}{\left(\lambda_{k}-\epsilon_{j}\right)^{2}}\right)
$$

Its solution is

$$
H_{i}=\sum_{k}\left(\mathcal{B}^{-1}\right)_{i k} V_{k}, \quad\left(\mathcal{B}^{-1}\right)_{j p}=\frac{\prod_{l \neq p}\left(\epsilon_{j}-\lambda_{l}\right) \prod_{i}\left(\lambda_{p}-\epsilon_{i}\right)}{\prod_{i \neq j}\left(\epsilon_{j}-\epsilon_{i}\right) \prod_{l \neq p}\left(\lambda_{p}-\lambda_{l}\right)}
$$

Using the Poisson bracket eq.(70), the equation of motion of the variable $\lambda_{k}$ with respect to $H_{i}$
is then (no summation over $k$ ):

$$
\begin{equation*}
\partial_{t_{i}} \lambda_{k}=i \mu_{k}\left(\mathcal{B}^{-1}\right)_{i k} \tag{75}
\end{equation*}
$$

We are interested in the system with prescribed real values of the conserved quantities, $H_{1}, \ldots, H_{n+1}$, i.e. we take as coordinates the $\left(\lambda_{j}, H_{j}\right)$ instead of the $\left(\lambda_{j}, \mu_{j}\right)$. This amounts to fixing the Liouville torus we work with, or equivalently the spectral polynomial $Q_{2 n+2}(\lambda)$. In this setting, the $\mu_{j}$ are determined by the equations

$$
\prod_{k=1}^{n}\left(\lambda_{j}-\epsilon_{k}\right) \mu_{j}=( \pm)_{j} \sqrt{Q_{2 n+2}\left(\lambda_{j}\right)}
$$

The equations of motion eq.(75) must be complemented by the equation of motion for the phase of $b$.

We consider now the level set containing the critical point eq.(21). It corresponds to a maximally degenerate spectral curve:

$$
\begin{equation*}
Q_{2 n+2}(\lambda)=4 \prod_{l=1}^{n+1}\left(\lambda-E_{l}\right)^{2} \tag{76}
\end{equation*}
$$

where the zeroes of $Q_{2 n+2}(\lambda)$ are the roots of the classical Bethe equation:

$$
\begin{equation*}
2 E+\sum_{j=1}^{n} \frac{s e_{j}}{E-\epsilon_{j}}=0 \tag{77}
\end{equation*}
$$

The separated variables $\lambda_{i}, \mu_{i}$ now satisfy

$$
\begin{equation*}
\mu_{i}=2( \pm)_{i} \frac{\prod_{l}\left(\lambda_{i}-E_{l}\right)}{\prod_{j}\left(\lambda_{i}-\epsilon_{j}\right)} \tag{78}
\end{equation*}
$$

The choice of sign here plays a crucial role in the description of the various strata of the level set.

We first examine the case where a $\lambda_{i}$ is frozen at a root $E_{l}$. Let us assume that $Q_{2 n+2}(\lambda)$ has a real root at $\lambda=E$. This means that $A^{2}(\lambda)+B(\lambda) C(\lambda)$ vanishes when $\lambda=E$. But for real $\lambda, A(\lambda)$ is real and one has $C(\lambda)=\overline{B(\lambda)}$ so that $A(\lambda), B(\lambda), C(\lambda)$ must all vanish at $\lambda=E$. In particular, recalling eq.(69), this means that one of the separated variables, say $\lambda_{1}$, is frozen at the value $E$, provided $\bar{b} b \neq 0$. This implies also that $\lambda-E$ divides simultaneously $A(\lambda), B(\lambda)$ and $C(\lambda)$, and therefore $(\lambda-E)^{2}$ divides $Q_{2 n+2}(\lambda)$. A real root of $Q_{2 n+2}(\lambda)$ is necessarily a double root and one $\lambda_{i}$ must be frozen at $E$.

Nothing simple can be said in the case $Q_{2 n+2}(\lambda)$ has a simple complex root. So, let us assume that it has a double complex root $E$. Of course the complex conjugate $\bar{E}$ is also a double root. Since $\Lambda(\lambda)=A^{2}(\lambda)+B(\lambda) C(\lambda)$ we have

$$
\Lambda^{\prime}(\lambda)=2 A(\lambda) A^{\prime}(\lambda)+B^{\prime}(\lambda) C(\lambda)+B(\lambda) C^{\prime}(\lambda)
$$

So, if $\Lambda(\lambda)$ has a double zero $E$ we have

$$
\begin{aligned}
& 0=A^{2}(E)+B(E) C(E) \\
& 0=2 A(E) A^{\prime}(E)+B^{\prime}(E) C(E)+B(E) C^{\prime}(E)
\end{aligned}
$$

In contrast to the real case, we cannot infer from these equations that $C(E)=0$. But if $C(E)=0$ that is if one $\lambda_{i}$ freezes at $E$, the first equation implies $A(E)=0$ and the second equation implies $B(E) C^{\prime}(E)=0$. Therefore if the zero of $C(\lambda)$ at $\lambda=E$ is simple, then necessarily $B(E)=0$. But $B(E)$ is the complex conjugate of $C(\bar{E})$ which must therefore vanish provided $\bar{b} b \neq 0$. From this we conclude that another separated variable $\lambda_{k}$ freezes at $\bar{E}$. In contrast to the real case, freezing is not compulsory, but it is the possibility to freeze the $\lambda_{i}$ by complex conjugated pairs that leads to the description of the real slice and the stratification of the level set.

Coming back to eq.(78), we see from eqs.(7176) that:

$$
P_{n+1}\left(\lambda_{i}\right)=( \pm)_{i} \sqrt{Q_{2 n+2}\left(\lambda_{i}\right)}=2( \pm)_{i} \prod_{l}\left(\lambda_{i}-E_{l}\right), \quad P_{n+1}(\lambda)=2\left(\lambda^{n+1}-\sigma_{1}(\epsilon) \lambda^{n}+\cdots\right)
$$

If we take the $+\operatorname{sign}$ for all $i$, then obviously $P_{n+1}(\lambda)=2 \prod_{l}\left(\lambda-E_{l}\right)$, and taking the residue at $\lambda=\epsilon_{j}$ in $P_{n+1}(\lambda) / \prod_{i}\left(\lambda-\epsilon_{i}\right)$, we find $s_{j}^{z}=s e_{j}$. This is the static solution corresponding to the critical point.

To go beyond this trivial solution, we divide the $\lambda_{i}$ into three sets $\lambda_{i}^{+} \in \mathcal{E}^{+}$and $\lambda_{i}^{-} \in \mathcal{E}^{-}$ depending on the sign in this formula, and $\lambda_{i}^{0} \in \mathcal{E}^{0}$ is the set of $\lambda_{i}$ frozen at some $E_{l}$.

We define the polynomials:

$$
\mathcal{P}(\lambda, t)=\prod_{j}\left(\lambda-\lambda_{j}(t)\right)=\mathcal{P}^{(+)}(\lambda) \mathcal{P}^{(0)}(\lambda) \mathcal{P}^{(-)}(\lambda)
$$

The equations of motion eq.(75) read in the soliton case

$$
\partial_{t_{i}} \lambda_{k}=\sqrt{-1} \mu_{k}\left(\mathcal{B}^{-1}\right)_{i k}=2 \sqrt{-1}( \pm)_{k} \frac{\prod_{l^{\prime}}\left(\lambda_{k}-E_{l^{\prime}}\right) \prod_{l \neq k}\left(\epsilon_{i}-\lambda_{l}\right)}{\prod_{j \neq i}\left(\epsilon_{i}-\epsilon_{j}\right) \prod_{l \neq k}\left(\lambda_{k}-\lambda_{l}\right)}
$$

hence, for $E_{l} \notin \mathcal{E}^{0}$,

$$
\sum_{k}^{\prime}( \pm)_{k} \frac{1}{\lambda_{k}-E_{l}} \partial_{t_{i}} \lambda_{k}=-2 \sqrt{-1} \frac{\prod_{l}\left(\epsilon_{i}-\lambda_{l}\right)}{\prod_{j \neq i}\left(\epsilon_{i}-\epsilon_{j}\right)} \sum_{k}^{\prime} \frac{\prod_{l^{\prime} \neq l}^{\prime}\left(\lambda_{k}-E_{l^{\prime}}\right)}{\left(\lambda_{k}-\epsilon_{i}\right) \prod_{l \neq k}^{\prime}\left(\lambda_{k}-\lambda_{l}\right)}
$$

where $\sum_{k}{ }^{\prime}$ means that the frozen $\lambda_{k} \in \mathcal{E}^{0}$ are excluded. We can rewrite this equation as

$$
\begin{aligned}
\partial_{t_{i}} \log \frac{\mathcal{P}^{(+)}\left(E_{l}\right)}{\mathcal{P}^{(-)}\left(E_{l}\right)} & =-2 \sqrt{-1} \frac{\prod_{l}\left(\epsilon_{i}-\lambda_{l}\right)}{\prod_{j \neq i}\left(\epsilon_{i}-\epsilon_{j}\right)} \sum_{k}^{\prime} \operatorname{Res}_{\lambda_{k}} \frac{\prod_{l^{\prime} \neq l}\left(z-E_{l^{\prime}}\right)}{\left(z-\epsilon_{i}\right) \prod_{l}\left(z-\lambda_{l}\right)} \\
& =2 \sqrt{-1} \frac{\prod_{l}\left(\epsilon_{i}-\lambda_{l}\right)}{\prod_{j \neq i}\left(\epsilon_{i}-\epsilon_{j}\right)}\left(\operatorname{Res}_{\infty}+\operatorname{Res}_{\epsilon_{i}}\right) \frac{\prod_{l^{\prime} \neq l}^{\prime}\left(z-E_{l^{\prime}}\right)}{\left(z-\epsilon_{i}\right) \prod_{l}^{\prime}\left(z-\lambda_{l}\right)} \\
& =2 \sqrt{-1} \frac{\prod_{l}\left(\epsilon_{i}-\lambda_{l}\right)}{\prod_{j \neq i}\left(\epsilon_{i}-\epsilon_{j}\right)}\left(-1+\frac{\prod_{l^{\prime} \neq l}^{\prime}\left(\epsilon_{i}-E_{l^{\prime}}\right)}{\prod_{l}^{\prime}\left(\epsilon_{i}-\lambda_{l}\right)}\right)
\end{aligned}
$$

so that

$$
\partial_{t_{i}} \log \frac{\mathcal{P}^{(+)}\left(E_{l}\right)}{\mathcal{P}^{(-)}\left(E_{l}\right)}=-2 \sqrt{-1} \frac{\prod_{l}\left(\epsilon_{i}-\lambda_{l}\right)}{\prod_{j \neq i}\left(\epsilon_{i}-\epsilon_{j}\right)}+2 \frac{\sqrt{-1}}{\epsilon_{i}-E_{l}} \frac{\prod_{l^{\prime}}\left(\epsilon_{i}-E_{l^{\prime}}\right)}{\prod_{j \neq i}\left(\epsilon_{i}-\epsilon_{j}\right)}
$$

Remembering the identities:

$$
\begin{equation*}
2 \frac{\prod_{k}\left(\epsilon_{j}-E_{k}\right)}{\prod_{k \neq j}\left(\epsilon_{j}-\epsilon_{k}\right)}=s e_{j}, \quad \sum_{l} E_{l}=\sigma_{1}(\epsilon) \tag{79}
\end{equation*}
$$

we can rewrite

$$
\partial_{t_{i}} \log \frac{\mathcal{P}^{(+)}\left(E_{l}\right)}{\mathcal{P}^{(-)}\left(E_{l}\right)}=-\sqrt{-1} \frac{s_{i}^{+}}{\bar{b}}+\sqrt{-1} \frac{s e_{i}}{\epsilon_{i}-E_{l}}
$$

But the equation of motion for $\bar{b}$ is

$$
\partial_{t_{i}} \bar{b}=\left\{H_{i}, \bar{b}\right\}=\sqrt{-1} s_{i}^{+}
$$

so that finally

$$
\partial_{t_{i}} \log \frac{\mathcal{P}^{(+)}\left(E_{l}\right)}{\mathcal{P}^{(-)}\left(E_{l}\right)}=-\partial_{t_{i}} \log \bar{b}+\sqrt{-1} \frac{s e_{i}}{\epsilon_{i}-E_{l}}, \quad E_{l} \notin \mathcal{E}^{0}
$$

Integrating these equations we get

$$
\begin{equation*}
\mathcal{P}^{(-)}\left(E_{l}\right)=\bar{b}(\{t\}) X_{l}(\{t\}) \mathcal{P}^{(+)}\left(E_{l}\right), \quad X_{l}(\{t\})=X_{l}(0) e^{\sqrt{-1}\left(\sum_{i} \frac{s e_{i}}{E_{l}-\epsilon_{i}} t_{i}-t_{n+1}\right)}, \quad E_{l} \notin \mathcal{E}^{0} \tag{80}
\end{equation*}
$$

The $t_{n+1}$ dependence of $X_{l}$ is fixed by noticing that $\mathcal{P}^{( \pm)}\left(E_{l}\right)$ are independent of $t_{n+1}$, while $\partial_{t_{n+1}} \bar{b}=\sqrt{-1} \bar{b}$.

Denoting $n_{ \pm}, n_{0}$ the number of elements in $\mathcal{E}^{( \pm)}, \mathcal{E}^{0}$ respectively, we get a set of $n-n_{0}+1$ linear equations for the $n_{+}+n_{-}+1=n-n_{0}+1$ unknown coefficients of the polynomials $\mathcal{P}_{-}(\lambda)$ and $\bar{b}(t) \mathcal{P}_{+}(\lambda)$.

These real soliton solutions live on the pre-images of focus-focus singularities These are composed of various strata, depending on the number of $n_{0}$ of separated variables frozen by pairs at the double roots ( $E_{l}, E_{\bar{l}}=E_{l}^{*}$ ) of the spectral polynomial. Separated variables are necessarily frozen at the real double roots of the spectral polynomial. These strata are themselves products of $\left(n+1-n_{0}\right) / 2$ two dimensional pinched tori.

We then consider the remaining complex roots of the spectral polynomial where no separated variables are frozen. This is a self conjugate set. We write eq. 80 ) as a linear system for the symmetric functions $\sigma_{i}^{(-)}\left(\left\{\lambda_{j}^{(-)}\right\}\right)$and $\sigma_{i}^{(+)}\left(\left\{\lambda_{j}^{(+)}\right\}\right)$and $\bar{b}$. We introduce the vector

$$
V=\left(\begin{array}{c}
(-1)^{n_{-}} \sigma_{n_{-}}^{(-)} / \bar{b} \\
\vdots \\
-\sigma_{1}^{(-)} / \bar{b} \\
1 / \bar{b} \\
(-1)^{1+n_{+}} \sigma_{n_{+}}^{(+)} \\
\vdots \\
\sigma_{1}^{(+)}
\end{array}\right)
$$

then eq.(80) reads

$$
\left(1, E, \cdots, E^{n_{-}}, X, X E, \cdots, X E^{n_{+}-1}\right) V=X E^{n_{+}}
$$

where we have defined the column vectors $\left(E^{j}\right)_{l}=E_{l}^{j}$ and $\left(X E^{j}\right)_{l}=X_{l} E_{l}^{j}$ where $X_{l}$ is given by eq. (80). The formulae used in the text are straightforward consequences Cramer's solution of this linear system.

In the vector $X$ we should incorporate the reality conditions

$$
\overline{X_{l}} X_{\bar{l}}=-\frac{1}{4}, \quad E_{\bar{l}} \equiv \overline{E_{l}}
$$

which we derive below.

### 9.1 Reality conditions.

It remains to impose the condition that the set of zeros $\left\{\bar{\lambda}_{k}\right\}$ of $B(\lambda)$ in eq.(74) is the complex conjugate of the set of zeros $\left\{\lambda_{k}\right\}$ of $C(\lambda)$. Going back to eq. (73), the first step is to build the polynomial $P_{n+1}(\lambda)$. Using the fact that $P_{n+1}\left(\lambda_{i}\right)=2( \pm)_{i} \prod_{l}\left(\lambda_{i}-E_{l}\right)$ we may write :

$$
P_{n+1}(\lambda)=2 \prod_{l}\left(\lambda-E_{l}\right)-4 \sum_{\lambda_{i}^{-}} \prod_{E_{l} \notin \mathcal{E}^{0}}\left(\lambda_{i}^{-}-E_{l}\right) \frac{\prod_{\lambda_{l}^{-} \neq \lambda_{i}^{-}}\left(\lambda-\lambda_{l}^{-}\right)}{\prod_{\lambda_{l}^{-} \neq \lambda_{i}^{-}}\left(\lambda_{i}^{-}-\lambda_{l}^{-}\right)} \frac{\mathcal{P}_{+}(\lambda) \mathcal{P}_{0}(\lambda)}{\mathcal{P}_{+}\left(\lambda_{i}^{-}\right)}
$$

To derive this formula, we used the fact that the polynomials $P_{n+1}(\lambda)$ and $2 \prod_{l}\left(\lambda-E_{l}\right)$ have the same terms of degrees $n+1$ and $n$; the last statement comes from the relation $\sum_{j} \epsilon_{j}=\sum_{l} E_{l}$
which is a direct consequence of the classical Bethe equation. Alternatively, we can also write:

$$
\begin{aligned}
P_{n+1}(\lambda)= & -2 \prod_{l}\left(\lambda-E_{l}\right)+4 \sum_{\lambda_{i}^{+}} \prod_{E_{l} \notin \mathcal{E}^{0}}\left(\lambda_{i}^{+}-E_{l}\right) \frac{\prod_{\lambda_{l}^{+} \neq \lambda_{i}^{+}}\left(\lambda-\lambda_{l}^{+}\right)}{\prod_{\lambda_{l}^{+} \neq \lambda_{i}^{+}}\left(\lambda_{i}^{+}-\lambda_{l}^{+}\right)} \frac{\mathcal{P}_{-}(\lambda) \mathcal{P}_{0}(\lambda)}{\mathcal{P}_{-}\left(\lambda_{i}^{+}\right)} \\
& +4\left(\lambda+\Sigma_{1}^{\prime}-\sigma_{1}^{\prime}(E)\right) \mathcal{P}_{+}(\lambda) \mathcal{P}_{-}(\lambda) \mathcal{P}_{0}(\lambda)
\end{aligned}
$$

Note that the last term in the right hand side is necessary to adjust the terms of degree $n+1$ and $n$ in $\lambda$ between the two sides of the equation. Again, we used the relation $\sum_{j} \epsilon_{j}=\sum_{l} E_{l}$. These two expressions for $P_{n+1}(\lambda)$ motivate the following definitions of polynomials $\mathcal{S}_{ \pm}(\lambda)$ :

$$
\begin{aligned}
& \mathcal{S}_{+}(\lambda)=\sum_{\lambda_{i}^{-}} \prod_{E_{l} \notin \mathcal{E}^{0}}\left(\lambda_{i}^{-}-E_{l}\right) \frac{\prod_{\lambda_{l}^{-} \neq \lambda_{i}^{-}}\left(\lambda-\lambda_{l}^{-}\right)}{\mathcal{P}_{-}^{\prime}\left(\lambda_{i}^{-}\right) \mathcal{P}_{+}\left(\lambda_{i}^{-}\right)} \\
& \mathcal{S}_{-}(\lambda)=\sum_{\lambda_{i}^{+}} \prod_{E_{l} \notin \mathcal{E}^{0}}\left(\lambda_{i}^{+}-E_{l}\right) \frac{\prod_{\lambda_{l}^{+} \neq \lambda_{i}^{+}}\left(\lambda-\lambda_{l}^{+}\right)}{\mathcal{P}_{+}^{\prime}\left(\lambda_{i}^{+}\right) \mathcal{P}_{-}\left(\lambda_{i}^{+}\right)}+\left(\lambda+\Sigma_{1}^{\prime}-\sigma_{1}^{\prime}(E)\right) \mathcal{P}_{+}(\lambda)
\end{aligned}
$$

so that we can write:

$$
\begin{align*}
& P_{n+1}(\lambda)=2 \prod_{l}\left(\lambda-E_{l}\right)-4 \mathcal{S}_{+}(\lambda) \mathcal{P}_{+}(\lambda) \mathcal{P}_{0}(\lambda)  \tag{81}\\
& P_{n+1}(\lambda)=-2 \prod_{l}\left(\lambda-E_{l}\right)+4 \mathcal{S}_{-}(\lambda) \mathcal{P}_{-}(\lambda) \mathcal{P}_{0}(\lambda) \tag{82}
\end{align*}
$$

Note that $\mathcal{S}_{+}(\lambda)$ has degree $n_{-}-1$ and $\mathcal{S}_{-}(\lambda)$ has degree $n_{+}+1$. Now, we have:

$$
4 \prod_{l}\left(\lambda-E_{l}\right)^{2}-P_{n+1}^{2}(\lambda)=16 \mathcal{S}_{-}(\lambda) \mathcal{S}_{+}(\lambda) \mathcal{P}_{0}(\lambda) \mathcal{P}(\lambda)=4 \bar{b} b \mathcal{P}(\lambda) \overline{\mathcal{P}}(\lambda)
$$

where $\mathcal{P}(\lambda)=\prod_{j}\left(\lambda-\lambda_{j}\right)=\mathcal{P}_{-}(\lambda) \mathcal{P}_{0}(\lambda) \mathcal{P}_{+}(\lambda)$ and $\overline{\mathcal{P}}(\lambda)=\prod_{j}\left(\lambda-\bar{\lambda}_{j}\right)$ is the complex conjugate of $\mathcal{P}(\lambda)$, that is $\overline{\mathcal{P}}(\lambda)=\overline{\mathcal{P}(\bar{\lambda})}$. Therefore:

$$
\begin{equation*}
\bar{b} b \overline{\mathcal{P}}(\lambda)=4 \mathcal{S}_{-}(\lambda) \mathcal{S}_{+}(\lambda) \mathcal{P}_{0}(\lambda) \tag{83}
\end{equation*}
$$

So the zeroes $\bar{\lambda}_{i}$ of $\overline{\mathcal{P}}(\lambda)$ split into the zeroes $\bar{\lambda}_{i}^{+}, \bar{\lambda}_{i}^{-}$and $\bar{\lambda}_{i}^{0}$ of $\mathcal{S}_{+}(\lambda), \mathcal{S}_{-}(\lambda)$ and $\mathcal{P}_{0}(\lambda)$ respectively. A direct consequence of these definitions and of eqs. (81) and (82) is that:

$$
\begin{align*}
& P_{n+1}\left(\bar{\lambda}_{i}^{+}\right)=+2 \prod_{l}\left(\bar{\lambda}_{i}^{+}-E_{l}\right)  \tag{84}\\
& P_{n+1}\left(\bar{\lambda}_{i}^{-}\right)=-2 \prod_{l}\left(\bar{\lambda}_{i}^{-}-E_{l}\right) \tag{85}
\end{align*}
$$

By the definition of the $\bar{\lambda}_{i}^{0}$ 's, we see that the set $\mathcal{E}^{0}$ is self conjugate and that $\mathcal{P}_{0}(\lambda)=\overline{\mathcal{P}}_{0}(\lambda)$. The above definition of $\mathcal{S}_{-}(\lambda)$ shows that the coefficient of its term of highest degree is equal to one (this is not the case for $\mathcal{S}_{+}(\lambda)$ ). Because of this:

$$
\begin{equation*}
\mathcal{S}_{-}(\lambda)=\overline{\mathcal{P}}_{-}(\lambda) \tag{86}
\end{equation*}
$$

Combining this with eq. (83) we get also:

$$
\begin{equation*}
\mathcal{S}_{+}(\lambda)=\frac{1}{4} \bar{b} b \overline{\mathcal{P}}_{+}(\lambda) \tag{87}
\end{equation*}
$$

Comparing the terms of highest degrees in $\mathcal{S}_{+}(\lambda)$ and $\overline{\mathcal{P}}_{+}(\lambda)$ gives:

$$
\begin{equation*}
\bar{b} b=4 \sum_{\lambda_{i}^{-}} \frac{\prod_{E_{l} \notin \mathcal{E}^{0}}\left(\lambda_{i}^{-}-E_{l}\right)}{\mathcal{P}_{-}^{\prime}\left(\lambda_{i}^{-}\right) \mathcal{P}_{+}\left(\lambda_{i}^{-}\right)} \tag{88}
\end{equation*}
$$

It is expressed only in terms of $\lambda_{i}^{-}$'s. So, if $n_{-}=0$, we recover the fact already mentioned that $\bar{b} b=0$ and the system remains at the critical point.

At this stage, we are ready to enforce the reality condition. As discussed above, the real slice is obtained by imposing that the set $\left\{\bar{\lambda}_{i}\right\}$ be the same as the set $\left\{\lambda_{i}^{*}\right\}$, where in the rest of this section we denote by $z^{*}$ the complex conjugate of $z$. Equivalently:

$$
\overline{\mathcal{P}}\left(\lambda^{*}\right)=\mathcal{P}(\lambda)^{*}
$$

for any $\lambda$. From the discussion at the beginning of this section, we know that the frozen variables $\lambda_{i}^{0}$ appear in complex conjugate pairs so that $\mathcal{P}_{0}\left(\lambda^{*}\right)=\mathcal{P}_{0}(\lambda)^{*}$. We also know that $\mathcal{P}_{0}(\lambda)=\overline{\mathcal{P}}_{0}(\lambda)$ so that $\overline{\mathcal{P}}_{0}\left(\lambda^{*}\right)=\mathcal{P}_{0}(\lambda)^{*}$. The above reality condition becomes then:

$$
\overline{\mathcal{P}}_{-}\left(\lambda^{*}\right) \overline{\mathcal{P}}_{+}\left(\lambda^{*}\right)=\mathcal{P}_{-}(\lambda)^{*} \mathcal{P}_{+}(\lambda)^{*}
$$

It is clearly sufficient to impose simultaneously:

$$
\begin{equation*}
\overline{\mathcal{P}}_{-}\left(\lambda^{*}\right)=\mathcal{P}_{-}(\lambda)^{*}, \quad \overline{\mathcal{P}}_{+}\left(\lambda^{*}\right)=\mathcal{P}_{+}(\lambda)^{*} \tag{89}
\end{equation*}
$$

but we claim that this condition is also necessary. This comes from the fact already noted that the sign of $P_{n+1}(\lambda) / \prod_{l}\left(\lambda-E_{l}\right)$ is positive for $\lambda=\lambda_{i}^{+}$or $\lambda=\bar{\lambda}_{i}^{+}$and negative for $\lambda=\lambda_{i}^{-}$or $\lambda=\bar{\lambda}_{i}^{-}$. So the roots $\bar{\lambda}_{i}^{+}$have to be complex conjugates of $\lambda_{i}^{+}$and likewise, the roots $\bar{\lambda}_{i}^{-}$have to be complex conjugates of $\lambda_{i}^{-}$. An interesting and useful consequence of this is that we must have deg $\mathcal{P}_{ \pm}=\operatorname{deg} \mathcal{S}_{ \pm}=\operatorname{deg} \overline{\mathcal{P}}_{ \pm}$, which requires $n_{+}=n_{-}-1$. Since $n_{+}+n_{0}+n_{-}=n$ we find:

$$
n_{+}=\frac{1}{2}\left(n-1-n_{0}\right), \quad n_{-}=\frac{1}{2}\left(n+1-n_{0}\right)
$$

Because the coefficients of highest degrees of $\mathcal{P}_{-}(\lambda)$ and $\mathcal{P}_{+}(\lambda)$ are set equal to one, the above constraints (89) are equivalent to (assuming of course that $\overline{\mathcal{P}}_{+}(\lambda)$ and $\overline{\mathcal{P}}_{-}(\lambda)$ are mutually prime and similarly for $\mathcal{P}_{+}(\lambda)$ and $\left.\mathcal{P}_{-}(\lambda)\right)$.

$$
\frac{\overline{\mathcal{P}}_{-}\left(\lambda^{*}\right)}{\overline{\mathcal{P}}_{+}\left(\lambda^{*}\right)}=\left(\frac{\mathcal{P}_{-}(\lambda)}{\mathcal{P}_{+}(\lambda)}\right)^{*}
$$

We also note that we should add the constraint $\bar{b}=b^{*}$. This plus the fact that these two polynomials involve a total of $n_{+}+n_{-}$unknown coefficients, shows that it is necessary and sufficient to enforce the following conditions for the $n+1-n_{0}=n_{+}+n_{-}+1$ roots $E_{l}$ which don't belong to $\mathcal{E}^{0}$ :

$$
\begin{equation*}
\frac{\overline{\mathcal{P}}_{-}\left(E_{l}^{*}\right)}{b \overline{\mathcal{P}}_{+}\left(E_{l}^{*}\right)}=\left(\frac{\mathcal{P}_{-}\left(E_{l}\right)}{\bar{b} \mathcal{P}_{+}\left(E_{l}\right)}\right)^{*} \tag{90}
\end{equation*}
$$

As we have seen, the general solution of the Hamiltonian evolution on the critical torus, eq. (80) implies:

$$
\begin{equation*}
\frac{\mathcal{P}_{-}\left(E_{l}\right)}{\bar{b} \mathcal{P}_{+}\left(E_{l}\right)}=X_{l} \tag{91}
\end{equation*}
$$

To evaluate the left-hand side of the conditions (90), we set $\lambda=E_{l}$ in eqs. (81, (82), which gives:

$$
P_{n+1}\left(E_{l}\right)=-4 \mathcal{S}_{+}\left(E_{l}\right) \mathcal{P}_{+}\left(E_{l}\right) \mathcal{P}_{0}\left(E_{l}\right), \quad P_{n+1}\left(E_{l}\right)=4 \mathcal{S}_{-}\left(E_{l}\right) \mathcal{P}_{-}\left(E_{l}\right) \mathcal{P}_{0}\left(E_{l}\right), \quad E_{l} \notin \mathcal{E}^{0}
$$

and therefore:

$$
\begin{equation*}
\mathcal{S}_{+}\left(E_{l}\right) \mathcal{P}_{+}\left(E_{l}\right)=-\mathcal{S}_{-}\left(E_{l}\right) \mathcal{P}_{-}\left(E_{l}\right) \tag{92}
\end{equation*}
$$

Since $\mathcal{P}_{+}\left(E_{l}\right) \neq 0$, and remembering eq. (80) this implies:

$$
\frac{\mathcal{S}_{-}\left(E_{l}\right)}{\mathcal{S}_{+}\left(E_{l}\right)}=-\frac{1}{\bar{b} X_{l}}
$$

But using eqs. (86) and (87), we get:

$$
\begin{equation*}
\frac{\overline{\mathcal{P}}_{-}\left(E_{l}^{*}\right)}{b \overline{\mathcal{P}}_{+}\left(E_{l}^{*}\right)}=-\frac{1}{4 X_{\bar{l}}} \tag{93}
\end{equation*}
$$

where we define the index $\bar{l}$ to be such that $E_{\bar{l}}=E_{l}^{*}$. Given eqs. (91) and (93), the reality conditions (90) and the condition $\bar{b}=b^{*}$ are satisfied if and only if:

$$
\begin{equation*}
X_{l}^{*} X_{\bar{l}}=-\frac{1}{4} \tag{94}
\end{equation*}
$$

As expected, times disappears from these conditions so that they reduce to constraints on the integration constants:

$$
X_{l}(0)^{*} X_{\bar{l}}(0)=-\frac{1}{4}
$$

This characterizes a stratum of dimension $n-n_{0}+1$ on the real slice.

## 10 Appendix : Asymptotic behavior of solitonic formulae

Let us assume that the $m$ solitonic amplitudes $X_{j}$ become large while keeping fixed ratios. Because of the reality condition (94), the $m$ remaining amplitudes $X_{\bar{j}}$ go to zero with fixed ratios. With no loss of generality, let us assume that $j$ runs from 1 to $m$ and $\bar{j}$ from $m+1$ to $2 m$. The complete expression (33) for $C(\lambda)$ involves three determinants $D_{1}, D_{2}$ and $D$ defined by:

$$
\begin{aligned}
D_{1} & =\operatorname{det}\left(\begin{array}{cccccccc}
1 & \lambda & \cdots & \lambda^{m} & 0 & 0 & \cdots & 0 \\
1 & E & \cdots & E^{n} & X & X E & \cdots & X E^{m-1}
\end{array}\right) \\
D_{2} & =\operatorname{det}\left(\begin{array}{cccccccc}
0 & 0 & \cdots & 0 & 1 & \lambda & \cdots & \lambda^{m-1} \\
1 & E & \cdots & E^{m} & X & X E & \cdots & X E^{m-1}
\end{array}\right) \\
D & =\operatorname{det}\left(\begin{array}{cccccccc}
1 & E & \cdots & E^{m-1} & X & X E & \cdots & X E^{m-1}
\end{array}\right)
\end{aligned}
$$

In the limit where $\left|X_{j}\right| \rightarrow \infty$ and $X_{\bar{j}} \rightarrow 0$ for $1 \leq j \leq m$, the dominant term in $D_{1}$ is given by:

$$
D_{1} \simeq \operatorname{det}\left(\begin{array}{cccccccc}
1 & \lambda & \cdots & \lambda^{m} & 0 & 0 & \cdots & 0 \\
0 & 0 & \cdots & 0 & X_{1} & X_{1} E_{1} & \cdots & X_{1} E_{1}^{m-1} \\
. & . & & . & . & . & & . \\
. & . & & . & . & . & & . \\
0 & 0 & \cdots & 0 & X_{m} & X_{m} E_{m} & \cdots & X_{m} E_{m}^{m-1} \\
1 & E_{\overline{1}} & \cdots & E_{\overline{1}}^{m} & 0 & 0 & \cdots & 0 \\
. & . & & . & \cdot & . & & . \\
\cdot & . & & . & . & . & & . \\
1 & E_{\bar{m}} & \cdots & E_{\bar{m}}^{m} & 0 & 0 & \cdots & 0
\end{array}\right)
$$

So, we get:

$$
\begin{equation*}
D_{1} \simeq \prod_{i<j}\left|E_{j}-E_{i}\right|^{2} \prod_{j=1}^{m}\left(\lambda-E_{\bar{j}}\right) X_{1} \cdots X_{m} \tag{95}
\end{equation*}
$$

It is also important to note that the subleading terms of $D_{1}$ vanish. We obtain such terms either by removing one of the $X_{j}$ 's or by adding one of the $X_{\bar{j}}$ 's. In the first case, the $m$ last columns to the right have only $m-1$ non-vanishing lines, so they are linearly dependent. Likewise, in the second case, the first $m+1$ columns on the left have only $m$ non-vanishing lines. So the first corrections to $D_{1}$ originate from the terms where we simultaneously remove one $X_{j}$ and add one $X_{\bar{k}}$.

The leading term in $D_{2}$ has the form:

$$
\operatorname{det}\left(\begin{array}{cccccccc}
0 & 0 & \cdots & 0 & 1 & \lambda & \cdots & \lambda^{m-1} \\
0 & 0 & \cdots & 0 & X_{1} & X_{1} E_{1} & \cdots & X_{1} E_{1}^{m-1} \\
. & \cdot & & . & \cdot & \cdot & & . \\
. & . & & . & . & \cdot & & . \\
0 & 0 & \cdots & 0 & X_{m} & X_{m} E_{m} & \cdots & X_{m} E_{m}^{m-1} \\
1 & E_{\overline{1}} & \cdots & E_{\overline{1}}^{m} & 0 & 0 & \cdots & 0 \\
. & \cdot & & \cdot & \cdot & . & & . \\
\cdot & . & & . & . & . & & . \\
1 & E_{\bar{m}} & \cdots & E_{\bar{m}}^{m} & 0 & 0 & \cdots & 0
\end{array}\right)
$$

which is equal to zero. One has therefore to consider subleading terms. They can be obtained by removing one of the $X_{j}$ 's or by adding one of the $X_{\bar{j}}$ 's. In fact, only the first option leads to
non-zero terms. Therefore, we may write:

$$
D_{2} \simeq \sum_{j=1}^{m} \operatorname{det}\left(\begin{array}{cccccccc}
0 & 0 & \cdots & 0 & 1 & \lambda & \cdots & \lambda^{m-1} \\
0 & 0 & \cdots & 0 & X_{1} & X_{1} E_{1} & \cdots & X_{1} E_{1}^{m-1} \\
\cdot & \cdot & & \cdot & \cdot & \cdot & & \cdot \\
\cdot & \cdot & & \cdot & \cdot & \cdot & & \cdot \\
0 & 0 & \cdots & 0 & X_{j-1} & X_{j-1} E_{j-1} & \cdots & X_{j-1} E_{j-1}^{m-1} \\
1 & E_{j} & \cdots & E_{j}^{m} & 0 & 0 & \cdots & 0 \\
0 & 0 & \cdots & 0 & X_{j+1} & X_{j+1} E_{j+1} & \cdots & X_{j+1} E_{j+1}^{m-1} \\
\cdot & \cdot & & \cdot & \cdot & \cdot & & \cdot \\
\cdot & \cdot & & \cdot & \cdot & \cdot & & \cdot \\
0 & 0 & \cdots & 0 & X_{m} & X_{m} E_{m} & \cdots & X_{m} E_{m}^{m-1} \\
1 & E_{\overline{1}} & \cdots & E_{\overline{1}}^{m} & 0 & 0 & \cdots & 0 \\
\cdot & \cdot & & \cdot & \cdot & \cdot & & \cdot \\
\cdot & \cdot & & \cdot & \cdot & \cdot & & . \\
1 & E_{\bar{m}} & \cdots & E_{\bar{m}}^{m} & 0 & 0 & \cdots & 0
\end{array}\right)
$$

Explicitely:

$$
\begin{equation*}
D_{2} \simeq \prod_{i<j}\left|E_{j}-E_{i}\right|^{2} \prod_{j=1}^{m}\left(\lambda-E_{j}\right) \sum_{j=1}^{m}\left(\frac{1}{X_{j}} \frac{E_{\bar{j}}-E_{j}}{\lambda-E_{j}} \prod_{k \neq j} \frac{E_{\bar{k}}-E_{j}}{E_{k}-E_{j}}\right) \quad X_{1} \cdots X_{m} \tag{96}
\end{equation*}
$$

The same reasoning as for $D_{1}$ shows that the next corrections to this estimate of $D_{2}$ are obtained by removing one more of the $X_{j}$ 's and by adding simultaneously one of the $X_{\bar{k}}$ 's.

The determinant $D$ has a non-vanishing leading order term:

$$
D \simeq \operatorname{det}\left(\begin{array}{cccccccc}
0 & 0 & \cdots & 0 & X_{1} & X_{1} E_{1} & \cdots & X_{1} E_{1}^{m-1} \\
\cdot & \cdot & & \cdot & \cdot & \cdot & & \cdot \\
. & \cdot & & . & . & . & & . \\
0 & 0 & \cdots & 0 & X_{m} & X_{m} E_{m} & \cdots & X_{m} E_{m}^{m-1} \\
1 & E_{\overline{1}} & \cdots & E_{\overline{1}}^{m} & 0 & 0 & \cdots & 0 \\
\cdot & \cdot & & \cdot & \cdot & \cdot & & \cdot \\
\cdot & . & & . & . & . & & . \\
1 & E_{\bar{m}} & \cdots & E_{\bar{m}}^{m} & 0 & 0 & \cdots & 0
\end{array}\right)
$$

so that:

$$
\begin{equation*}
D \simeq(-1)^{m} \prod_{i<j}\left|E_{j}-E_{i}\right|^{2} X_{1} \cdots X_{m} \tag{97}
\end{equation*}
$$

## 11 Appendix: Periodic flows.

There is a general procedure in the algebro-geometric setting to construct Hamiltonians generating periodic flows. From the equation of the spectral curve we have

$$
\delta \mu d \lambda=\sum_{i} \frac{d \lambda}{\mu\left(\lambda-\epsilon_{i}\right)} \delta H_{i}=\sum_{i} \sigma_{i}(\lambda) d \lambda \delta H_{i}
$$

where

$$
\sigma_{i}(\lambda) d \lambda=\frac{d \lambda}{\mu\left(\lambda-\epsilon_{i}\right)}=\frac{\prod_{k \neq i}\left(\lambda-\epsilon_{k}\right)}{\sqrt{Q_{2 n+2}(\lambda)}} d \lambda
$$

form a basis of holomorphic differentials. Recalling the flows generated by $H_{i}$, eq. (75), we obtain the important relation

$$
\sigma_{j}\left(\lambda_{k}\right) \sum_{k} \partial_{t_{i}} \lambda_{k}=i \sum_{k} \mu_{k}\left(\mathcal{B}^{-1}\right)_{i k} \frac{\mathcal{B}_{k j}}{\mu_{k}}=i \delta_{i j}
$$

To find the Hamiltonian generating the periodic flows, we should need the description of the periodic trajectory of the $n$ coordinates $\lambda_{k}$ of the divisor as $n$ ovals on the spectral curve. However only the homotopy class of the trajectory will be important in the following discussion.

Let us assume for the time being that it is homotopic to $n$ independent non intersecting ovals $B_{i}$. Let us define

$$
\mathcal{H}^{(i)}=\oint_{B_{i}} \sqrt{\Lambda(\lambda)} d \lambda
$$

To the $B_{i}$-cycles we associate a basis of normalized Abelian differentials $\omega_{i}(\lambda) d \lambda$. Of course we can expand it on the $\sigma_{j}(\lambda) d \lambda$ basis

$$
\omega_{i}(\lambda) d \lambda=\sum_{j} \mathcal{N}_{i j} \sigma_{j}(\lambda) d \lambda
$$

Equipped with these Abelian differentials, we define the angles through the Abel transformation

$$
\theta_{j}=\sum_{k} \int^{\lambda_{k}} \omega_{j}(\lambda) d \lambda
$$

Let us compute the equation of motion of these angles under the Hamiltonians $\mathcal{H}^{(i)}$.

$$
\begin{aligned}
\partial_{\tau_{i}} \theta_{j} & =\sum_{k} \omega_{j}\left(\lambda_{k}\right) \partial_{\tau_{i}} \lambda_{k}=\sum_{k, l} \omega_{j}\left(\lambda_{k}\right) \frac{\partial \mathcal{H}^{(i)}}{\partial H_{l}} \partial_{t_{l}} \lambda_{k} \\
& =\sum_{l, n} \frac{\partial \mathcal{H}^{(i)}}{\partial H_{l}} \mathcal{N}_{j n} \sum_{k} \sigma_{n}\left(\lambda_{k}\right) \partial_{t_{l}} \lambda_{k} \\
& =i \sum_{l, n} \frac{\partial \mathcal{H}^{(i)}}{\partial H_{l}} \mathcal{N}_{j n} \delta_{l n}=\oint_{B_{i}} \mathcal{N}_{j l} \frac{\partial}{\partial H_{l}} \sqrt{\Lambda} d \lambda \\
& =i \oint_{B_{i}} \mathcal{N}_{j l} \sigma_{l}(\lambda) d \lambda=i \oint_{B_{i}} \omega_{j}(\lambda) d \lambda=2 i \pi \delta_{i j}
\end{aligned}
$$

Hence

$$
\partial_{\tau_{i}} \theta_{j}=2 i \pi \delta_{i j}
$$

and the flows of the Hamiltonians $\mathcal{H}^{(i)}$ are $2 \pi$-periodic. So the problem of finding the Hamiltonians generating the periodic flows reduces to the description of the homotopy class of the $n$ non intersecting ovals of the periodic trajectories.

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