# On Integrability of Hirota-Kimura-Type Discretizations: Experimental Study of the Discrete Clebsch System 

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}
R. Hirota and K. Kimura discovered integrable discretizations of the Euler and the Lagrange tops, given by birational maps. Their method is a specialization to the integrable context of a general discretization scheme introduced by W . Kahan and applicable to any vector field with a quadratic dependence on phase variables. According to a proposal by T. Ratiu, discretizations of Hirota-Kimura type can be considered for numerous integrable systems of classical mechanics. Due to a remarkable and not well understood mechanism, such discretizations seem to inherit the integrability for all algebraically completely integrable systems. We introduce an experimental method for a rigorous study of integrability of such discretizations.

Application of this method to the Hirota-Kimura-type discretization of the Clebsch system leads to the discovery of four functionally independent integrals of motion of this discrete-time system, which turn out to be much more complicated than the integrals of the continuous-time system. Further, we prove that every orbit of the discrete-time Clebsch system lies in an intersection of four quadrics in the six-dimensional phase space. Analogous results hold for the Hirota-Kimura-type discretizations for all commuting flows of the Clebsch system, as well as for the so(4) Euler top.

## 1. INTRODUCTION

The discretization method studied in this paper seems to have been introduced into the geometric integration literature by W. Kahan in the unpublished notes [Kahan 93]. It is applicable to any system of ordinary differential equations for $x: \mathbb{R} \rightarrow \mathbb{R}^{n}$ with a quadratic vector field

$$
\dot{x}=Q(x)+B x+c,
$$

where each component of $Q: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is a quadratic form, while $B \in \operatorname{Mat}_{n \times n}$ and $c \in \mathbb{R}^{n}$. Kahan's dis-
cretization reads as

$$
\begin{equation*}
\frac{\widetilde{x}-x}{2 \epsilon}=Q(x, \widetilde{x})+\frac{1}{2} B(x+\widetilde{x})+c, \tag{1-1}
\end{equation*}
$$

where

$$
Q(x, \widetilde{x})=\frac{1}{2}(Q(x+\widetilde{x})-Q(x)-Q(\widetilde{x}))
$$

is the symmetric bilinear form corresponding to the quadratic form $Q$. Here and below we use the following notational convention, which will allow us to omit a great many indices: for a sequence $x: \mathbb{Z} \rightarrow \mathbb{R}$ we write $x$ for $x_{k}$ and $\widetilde{x}$ for $x_{k+1}$. Equation (1-1) is linear with respect to $\widetilde{x}$ and therefore defines a rational map $\widetilde{x}=f(x, \epsilon)$. Clearly, this map approximates the time( $2 \epsilon$ )-shift along the solutions of the original differential system, so that $x_{k} \approx x(2 k \epsilon)$. (We have chosen a slightly unusual notation $2 \epsilon$ for the time step, in order to avoid the appearance of various powers of 2 in numerous formulas; a more standard choice would lead to changing $\epsilon$ to $\epsilon / 2$ everywhere.) Since (1-1) remains invariant under the interchange $x \leftrightarrow \widetilde{x}$ with the simultaneous sign inversion $\epsilon \mapsto-\epsilon$, one has the reversibility property

$$
f^{-1}(x, \epsilon)=f(x,-\epsilon)
$$

In particular, the map $f$ is birational.
W. Kahan applied this discretization scheme to the famous Lotka-Volterra system and showed that in this case it possesses a very remarkable nonspiraling property. We will briefly discuss this example in Section 2. Some further applications of this discretization have been explored in [Kahan and Li 97].

The next, even more intriguing, appearance of this discretization can be found in two papers by R. Hirota and K. Kimura, who (being apparently unaware of Kahan's work) applied it to two famous integrable systems of classical mechanics: the Euler top and the Lagrange top [Hirota and Kimura 00, Kimura and Hirota 00]. For the purposes of the present text, integrability of a dynamical system is synonymous with the existence of a sufficient number of functionally independent conserved quantities, or integrals of motion, that is, functions constant along the orbits. We leave aside other aspects of the multifaceted notion of integrability, such as Hamiltonian or explicit solution. Surprisingly, the Kahan-Hirota-Kimura discretization scheme produced integrable maps in both the Euler and the Lagrange cases of rigid-body motion. Even more surprisingly, the mechanism that ensures integrability in these two cases seems to be rather different from the majority of examples known in the area of integrable discretizations, and, more generally, integrable maps; cf. [Suris 03].

The case of the discrete-time Euler top is relatively simple, and the proof of its integrability given in [Hirota and Kimura 00] is rather straightforward and easy to verify by hand. As often happens, no explanation was given in [Hirota and Kimura 00] as to how this result was discovered. The "derivation" of integrals of motion for the discrete-time Lagrange top in [Kimura and Hirota $00]$ is rather cryptic and almost incomprehensible.

The present paper aims at clarifying the HirotaKimura integrability mechanism and its application to further integrable systems. We use the term Hirota-Kimura-type discretization for Kahan's discretization in the context of integrable systems. In Section 3 we propose a formalization of the Hirota-Kimura mechanism from [Kimura and Hirota 00], which will, we hope, reveal its main idea and contribute to a demystification of at least some of its aspects.

We introduce a notion of a Hirota-Kimura basis for a given map $f$. Such a basis $\Phi$ is a set of simple (often monomial) functions $\Phi=\left(\varphi_{l}, \ldots, \varphi_{l}\right)$ such that for every orbit $\left\{f^{i}(x)\right\}$ of the map $f$ there is a certain linear combination $c_{1} \varphi_{1}+\cdots+c_{l} \varphi_{l}$ of functions from $\Phi$ vanishing on this orbit. As explained in Section 3, this is a new mathematical notion, not reducible to that of integrals of motion, although closely related to it. In Section 4 we establish a theoretical foundation for the search for Hirota-Kimura bases for a given discrete-time system, and give a number of practical recipes and tricks for doing this.

We dare to claim that the results of [Hirota and Kimura 00] concerning the discrete-time Euler top were originally discovered using the mechanism of HirotaKimura bases, and we present in Section 5 an attempt to reconstruct the way in which this discovery was made. Section 6 contains the main results of this paper, namely the proof of integrability of the Hirota-Kimura-type discretization for a further famous integrable system of classical mechanics, namely for the Clebsch case of the motion of a rigid body in an ideal fluid.

Our investigations are based mainly on computer experiments, which are used for both the discovery of new results and their rigorous proof. A search for HirotaKimura bases can be accomplished with the help of numerical experiments based on the recipe $\mathcal{N}$ formulated in Section 4, which has a theoretical justification in Theorem 4.1.

If the search has been successful and a certain set of functions $\Phi$ has been identified as a Hirota-Kimura basis for a given map $f$, then numerical experiments can provide very convincing evidence in favor of such a state-


FIGURE 1. Left: a spiraling orbit of the explicit Euler method with time step $\epsilon=0.01$ applied to the Lotka-Volterra system. Right: three orbits of Kahan's discretization with $\epsilon=0.1$.
ment. A rigorous proof of such a statement turns out to be much more demanding. At present, we are not in possession of any theoretical proof strategies and are forced to verify the corresponding statements by means of symbolic computations. However, direct and simple-minded symbolic computations turn out to be infeasible, due to complexity issues.

As detailed in Section 6, the sheer size of the explicit expressions for the second iterate $f^{2}$ of the discrete-time Clebsch system precludes symbolic manipulations, such as solution of linear systems, as soon as these involve $f^{2}$. Therefore, our main effort has been devoted to finding a strategy for a complete and rigorous symbolic proof that would avoid using $f^{2}$ and would stay within the memory and performance restrictions of the available software and hardware. The resulting proofs are computer assisted and are based on symbolic computation with Maple, Singular, and Form. ${ }^{1}$

Our work was stimulated by a talk presented by T. Ratiu at the Oberwolfach workshop "Geometric Integration" [Ratiu 06], where an extension of the HirotaKimura approach to the Clebsch system and to the Kovalevski top was proposed. However, no valid derivation of integrals was presented in T. Ratiu's talk, so that the question of the integrability of these discretizations remained open. Our work answers this question in the affirmative for the Clebsch system (indeed, for a whole family of Hamiltonian flows generated by commuting integrals of the Clebsch system). In the concluding Section 7, we discuss further perspectives of this approach and

[^0]formulate a general conjecture about the integrability of the Hirota-Kimura-type discretizations.

## 2. KAHAN'S DISCRETIZATION OF THE LOTKA-VOLTERRA SYSTEM

As already mentioned in Section 1, W. Kahan applied his general discretization scheme to the famous LotkaVolterra system modeling the interaction of predator and prey populations:

$$
\begin{equation*}
\dot{x}=x(1-y), \quad \dot{y}=y(x-1) \tag{2-1}
\end{equation*}
$$

Solutions of this system lie on closed curves in (the first quadrant of) the phase plane $\mathbb{R}^{2}$, because of the presence of the integral (conserved quantity)

$$
H(x, y)=x+y-\log (x y)
$$

In fact, system $(2-1)$ is Hamiltonian with respect to the Poisson bracket

$$
\begin{equation*}
\{x, y\}=x y \tag{2-2}
\end{equation*}
$$

with the Hamilton function $H$ :

$$
\dot{x}=-x y \frac{\partial H}{\partial y}, \quad \dot{y}=x y \frac{\partial H}{\partial x}
$$

When applied to (2-1), the majority of conventional discretization schemes produce spiraling solutions. Compared with solutions of the original system, this is a qualitatively different behavior; see Figure 1 (left).

The discretization proposed by Kahan reads

$$
\begin{align*}
(\widetilde{x}-x) / \epsilon & =(\widetilde{x}+x)-(\widetilde{x} y+x \widetilde{y}),  \tag{2-3}\\
(\widetilde{y}-y) / \epsilon & =(\widetilde{x} y+x \widetilde{y})-(\widetilde{y}+y) .
\end{align*}
$$

Equations (2-3) can be written as a linear system for $(\widetilde{x}, \widetilde{y})$,

$$
\left(\begin{array}{cc}
1-\epsilon+\epsilon y & \epsilon x \\
-\epsilon y & 1+\epsilon-\epsilon x
\end{array}\right)\binom{\widetilde{x}}{\widetilde{y}}=\binom{(1+\epsilon) x}{(1-\epsilon) y}
$$

which can be immediately solved, thus yielding an explicit map $(\widetilde{x}, \widetilde{y})=f(x, y, \epsilon)$ :

$$
\left\{\begin{array}{l}
\widetilde{x}=x \frac{(1+\epsilon)^{2}-\epsilon(1+\epsilon) x-\epsilon(1-\epsilon) y}{1-\epsilon^{2}-\epsilon(1-\epsilon) x+\epsilon(1+\epsilon) y}  \tag{2-4}\\
\widetilde{y}=y \frac{(1-\epsilon)^{2}+\epsilon(1+\epsilon) x+\epsilon(1-\epsilon) y}{1-\epsilon^{2}-\epsilon(1-\epsilon) x+\epsilon(1+\epsilon) y}
\end{array}\right.
$$

A remarkable property of Kahan's discretization is that it apparently does not suffer from spiraling; solutions seem to fill out closed curves in the phase plane; see Figure 1 (right). A (partial) explanation of this behavior was given in [Sanz-Serna 94], where it was shown that the map $f$ is Poisson with respect to the invariant Poisson bracket (2-2) of the system (2-1). It is unknown whether the map (2-4) possesses an integral of motion, thus forcing all orbits to lie on smooth closed curves, as suggested by Figure 1 (right). Some numerical experiments, via a deep zoom-in into certain domains of the phase plane, indicate that the map might be nonintegrable, but a rigorous proof of a nonexistence statement seems to be rather difficult. It might be possible with the use of technology described in [Gelfreich and Lazutkin 01].

## 3. HIROTA-KIMURA BASES AND INTEGRALS

In this section a general formulation of a remarkable mechanism will be given that seems to be responsible for the integrability of the Hirota-Kimura-type (or Kahan-type) discretizations of algebraically completely integrable systems. This mechanism is thus far not well understood. In fact, at the moment we do not know what mathematical structures make it actually work.

Throughout this section, $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is a birational map, while $h_{i}, \varphi_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}$ stand for rational, usually polynomial, functions on the phase space. We start by recalling a well-known definition.

Definition 3.1. A function $h: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is called an integral, or a conserved quantity, of the map $f$ if for every $x \in \mathbb{R}^{n}$,

$$
h(f(x))=h(x)
$$

so that

$$
h \circ f^{i}(x)=h(x) \quad \forall i \in \mathbb{Z}
$$

Convention 3.2. In the last formula and everywhere in the sequel, we use the expression $h \circ f^{i}(x)$ for the evaluation of the function $h \circ f^{i}$ at the point $x$. This is equivalent to $h\left(f^{i}(x)\right)$ and is used to spare some parentheses.

Thus, each orbit of the map $f$ lies on a certain level set of its integral $h$. As a consequence, if one knows $d$ functionally independent integrals $h_{1}, \ldots, h_{d}$ of $f$, one can claim that each orbit of $f$ is confined to an $(n-d)$ dimensional invariant set, which is a common level set of the functions $h_{1}, \ldots, h_{d}$.

Definition 3.3. A set of functions $\Phi=\left(\varphi_{1}, \ldots, \varphi_{l}\right)$, linearly independent over $\mathbb{R}$, is called a Hirota-Kimura basis (or HK basis for short) if for every $x_{0} \in \mathbb{R}^{n}$ there exists a vector $c=\left(c_{1}, \ldots, c_{l}\right) \neq 0$ such that

$$
\begin{equation*}
\left(c_{1} \varphi_{1}+\cdots+c_{l} \varphi_{l}\right) \circ f^{i}(x)=0 \quad \forall i \in \mathbb{Z} \tag{3-1}
\end{equation*}
$$

For a given $x \in \mathbb{R}^{n}$, the set of all vectors $c \in \mathbb{R}^{l}$ with this property will be denoted by $K_{\Phi}(x)$ and called the null space of the basis $\Phi$ (at the point $x$ ). This set clearly is a vector space.

Thus, for an HK basis $\Phi$ and for $c \in K_{\Phi}(x)$ the function $h=c_{1} \varphi_{1}+\cdots+c_{l} \varphi_{l}$ vanishes along the $f$ orbit of $x$. Let us stress that we cannot claim that $h=c_{1} \varphi_{1}+\cdots+c_{l} \varphi_{l}$ is an integral of motion, since vectors $c \in K_{\Phi}(x)$ do not have to belong to $K_{\Phi}(y)$ for initial points $y$ not lying on the orbit of $x$. However, for any $x$ the orbit $\left\{f^{i}(x)\right\}$ is confined to the common zero level set of $d$ functions

$$
h_{j}=c_{1}^{(j)} \varphi_{1}+\cdots+c_{l}^{(j)} \varphi_{l}=0, \quad j=1, \ldots, d
$$

where the vectors $c^{(j)}=\left(c_{1}^{(j)}, \ldots, c_{l}^{(j)}\right) \in \mathbb{R}^{l}$ form a basis of $K_{\Phi}(x)$. Thus, knowledge of an HK basis with the null space of dimension $d$ leads to a similar conclusion as knowledge of $d$ independent integrals of $f$, namely to the conclusion that the orbits lie on $(n-d)$-dimensional invariant sets. Note, however, that an HK basis gives no immediate information on how these invariant sets foliate the phase space $\mathbb{R}^{n}$, since the vectors $c^{(j)}$, and therefore the functions $h_{j}$, change from one initial point $x$ to another.

Although the notions of integrals and HK bases cannot be immediately translated into one another, they turn out to be closely related.

The simplest situation for an HK basis corresponds to $l=2, \operatorname{dim} K_{\Phi}(x)=d=1$. In this case we immediately see that $h=\varphi_{1} / \varphi_{2}$ is an integral of motion of
the map $f$. Conversely, for any rational integral of motion $h=\varphi_{1} / \varphi_{2}$, its numerator and denominator $\varphi_{1}, \varphi_{2}$ satisfy

$$
\left(c_{1} \varphi_{1}+c_{2} \varphi_{2}\right) \circ f^{i}(x)=0, \quad i \in \mathbb{Z}
$$

with $c_{1}=1, c_{2}=-h(x)$, and thus build an HK basis with $l=2$. Thus, the notion of an HK basis generalizes (for $l \geq 3$ ) the notion of integrals of motion.

On the other hand, knowing an HK basis $\Phi$ with $\operatorname{dim} K_{\Phi}(x)=d \geq 1$ allows one to find integrals of motion for the map $f$. Indeed, from Definition 3.3 there follows immediately the following result:

Proposition 3.4. If $\Phi$ is an HK basis for a map $f$, then

$$
K_{\Phi}(f(x))=K_{\Phi}(x)
$$

Thus, the $d$-dimensional null space $K_{\Phi}(x) \in \operatorname{Gr}(d, l)$, regarded as a function of the initial point $x \in \mathbb{R}^{n}$, is constant along trajectories of the map $f$, i.e., it is a $\operatorname{Gr}(d, l)$ valued integral. One can extract from this fact a number of scalar integrals.

Corollary 3.5. Let $\Phi$ be an HK basis for $f$ with $\operatorname{dim} K_{\Phi}(x)=d$ for all $x \in \mathbb{R}^{n}$. Take a basis of $K_{\Phi}(x)$ consisting of $d$ vectors $c^{(i)} \in \mathbb{R}^{l}$ and put them into the columns of an $l \times d$ matrix $C(x)$. For any $d$-index $\alpha=\left(\alpha_{1}, \ldots, \alpha_{d}\right) \subset\{1,2, \ldots, n\}$ let $C_{\alpha}=C_{\alpha_{1}, \ldots, \alpha_{d}}$ denote the $d \times d$ minor of the matrix $C$ built from the rows $\alpha_{1}, \ldots, \alpha_{d}$. Then for any two $d$-indices $\alpha, \beta$, the function $C_{\alpha} / C_{\beta}$ is an integral of $f$.

Proof: The functions $C_{\alpha}$ are nothing other than the Grassmann-Plücker coordinates of the $d$-space $K_{\Phi}(x)$ in the Grassmannian $\operatorname{Gr}(d, l)$, which are defined up to a common factor. More precisely, any basis of $K_{\Phi}(f(x))$ is obtained from the given basis of $K_{\Phi}(x)$ via a right multiplication of $C$ by a nondegenerate $d \times d$ matrix $D$. This yields a simultaneous multiplication of all $C_{\alpha}$ by the common factor $\operatorname{det} D$. This operation does not change the quotients $C_{\alpha} / C_{\beta}$.

Especially simple is the situation in which the null space of an HK basis has dimension $d=1$.

Corollary 3.6. Let $\Phi$ be an HK basis for $f$ with $\operatorname{dim} K_{\Phi}(x)=1$ for all $x \in \mathbb{R}^{n}$. Let $K_{\Phi}(x)=\left[c_{1}(x)\right.$ : $\left.\cdots: c_{l}(x)\right] \in \mathbb{R}^{l-1}$. Then the functions $c_{j} / c_{k}$ are integrals of motion for $f$.

An interesting (and difficult) question is the number of functionally independent integrals obtained from a given HK basis according to Corollaries 3.5 and 3.6. We will see later that it is possible for an HK basis with a onedimensional null space to produce more than one independent integral (see Theorem 6.5).

The first examples of this mechanism (with $d=1$ ) were given in [Kimura and Hirota 00] and (somewhat implicitly) in [Hirota and Kimura 00].

## 4. FINDING HIROTA-KIMURA BASES

At present, we cannot give any theoretically sufficient conditions for the existence of a Hirota-Kimura basis $\Phi$ for a given map $f$, and the only way to find such a basis remains the experimental one. Definition 3.3 requires the verification of condition (3-1) for all $i \in \mathbb{Z}$, which is, of course, impractical. We now show that it is enough to check this condition for a finite number of iterates $f^{i}$.

For a given set of functions $\Phi=\left(\varphi_{1}, \ldots, \varphi_{l}\right)$ and for any interval $[j, k] \subset \mathbb{Z}$, we define

$$
X_{[j, k]}(x)=\left(\begin{array}{ccc}
\varphi_{1}\left(f^{j}(x)\right) & \cdots & \varphi_{l}\left(f^{j}(x)\right)  \tag{4-1}\\
\varphi_{1}\left(f^{j+1}(x)\right) & \cdots & \varphi_{l}\left(f^{j+1}(x)\right) \\
\cdots & & \cdots \\
\varphi_{1}\left(f^{k}(x)\right) & \cdots & \varphi_{l}\left(f^{k}(x)\right)
\end{array}\right)
$$

In particular, $X_{(-\infty, \infty)}(x)$ will denote the doubly infinite matrix of type (4-1). Obviously,

$$
\operatorname{ker} X_{(-\infty, \infty)}(x)=K_{\Phi}(x)
$$

Thus, Definition 3.3 requires $\operatorname{dim} \operatorname{ker} X_{(-\infty, \infty)}(x) \geq 1$. Our algorithm for detecting this situation is based on the following observation.

## Theorem 4.1. Let

$$
\operatorname{dim} \operatorname{ker} X_{[0, s-1]}(x)= \begin{cases}l-s & \text { for } 1 \leq s \leq l-d  \tag{4-2}\\ d & \text { for } s=l-d+1\end{cases}
$$

hold with some $d$ for all $x \in \mathbb{R}^{n}$. Then for any $x \in \mathbb{R}^{n}$,

$$
\operatorname{ker} X_{(-\infty, \infty)}(x)=\operatorname{ker} X_{[0, l-d-1]}(x)
$$

and in particular,

$$
\operatorname{dim} \operatorname{ker} X_{(-\infty, \infty)}(x)=d
$$

Proof: By definition, $X_{[j, k]}(x)=X_{[0, k-j]}\left(f^{j}(x)\right)$. Therefore, applying condition (4-2) to iterates $f^{j}(x)$ instead of $x$ itself, we see that the kernel of any submatrix of
$X_{(-\infty, \infty)}(x)$ with $l-d$ rows, as well as the kernel of any submatrix with $l-d+1$ rows, is $d$-dimensional:

$$
\begin{aligned}
\operatorname{dim} \operatorname{ker} X_{[j, j+l-d-1]}(x) & =\operatorname{dim} \operatorname{ker} X_{[j, j+l-d]}(x) \\
& =\operatorname{dim} \operatorname{ker} X_{[j+1, j+l-d]}(x)
\end{aligned}
$$

Since, obviously,

$$
\begin{aligned}
\operatorname{ker} X_{[j, j+l-d-1]}(x) & \supset \operatorname{ker} X_{[j, j+l-d]}(x) \\
& \subset \operatorname{ker} X_{[j+1, j+l-d]}(x)
\end{aligned}
$$

we find that all three kernels coincide; in particular,

$$
\operatorname{ker} X_{[j, j+l-d-1]}(x)=\operatorname{ker} X_{[j+1, j+l-d]}(x)
$$

By induction, all ker $X_{[j, j+l-d-1]}(x), j \in \mathbb{Z}$, coincide, and therefore they coincide with $\operatorname{ker} X_{(-\infty, \infty)}(x)$ as well.

These results lead us to formulate the following numerical algorithm for the estimation of $\operatorname{dim} K_{\Phi}(x)$ for a hypothetical HK basis $\Phi=\left(\varphi_{1}, \ldots, \varphi_{l}\right)$ :
$\mathcal{N}:$ For several randomly chosen initial points $x \in \mathbb{R}^{n}$, compute $\operatorname{dim} \operatorname{ker} X_{[0, s-1]}(x)$ for $1 \leq s \leq l$. If for every $x$ condition (4-2) is satisfied with one and the same $d \geq 1$, then $\Phi$ is likely to be an HK basis for $f$, with $\operatorname{dim} K_{\Phi}(x)=d$.

We stress once again that generally (for general maps $f$ and general monomial sets $\Phi$ ) one will find that the $l \times l$ matrix $X_{[0, l-1]}(x)$ is nondegenerate for a typical $x$, so that $\operatorname{dim} K_{\Phi}(x)=0$. Finding (a candidate for) an HK basis $\Phi$ is a highly nontrivial task.

Having found an HK basis $\Phi$ with $\operatorname{dim} K_{\Phi}(x)=d$ numerically, one faces the next problem: to prove this fact, that is, to prove that the system of equations (3-1) with $i=i_{0}, i_{0}+1, \ldots, i_{0}+l-d$ admits (for some, and then for all, $i_{0} \in \mathbb{Z}$ ) a $d$-dimensional space of solutions. For the sake of clarity, we restrict our following discussion to the most important case $d=1$. Thus, one has to prove that the homogeneous system

$$
\begin{gather*}
\left(c_{1} \varphi_{1}+\cdots+c_{l} \varphi_{l}\right) \circ f^{i}(x)=0  \tag{4-3}\\
\quad i=i_{0}, i_{0}+1, \ldots, i_{0}+l-1
\end{gather*}
$$

admits for every $x \in \mathbb{R}^{n}$ a one-dimensional vector space of nontrivial solutions.

The main obstruction to a symbolic solution of the system $(4-3)$ is the growing complexity of the iterates $f^{i}(x)$. While the expression for $f(x)$ is typically of moderate size, already the second iterate $f^{2}(x)$ becomes typically prohibitively big. In such a situation a symbolic solution
of the linear system (4-3) should be considered impossible as soon as $f^{2}(x)$ is involved, for instance, if $l \geq 3$ and one considers the linear system with $i=0,1, \ldots, l-1$.

Therefore it becomes crucial to reduce the number of iterates involved in $(4-3)$ as far as possible. A reduction of this number by 1 becomes in many cases crucial! One can imagine several ways to accomplish this.
$\mathcal{A}$ : Take into account that because of the reversibility $f^{-1}(x, \epsilon)=f(x,-\epsilon)$, the negative iterates $f^{-i}$ are of the same complexity as $f^{i}$. Therefore, one can reduce the complexity of the functions involved in (4-3) by choosing $i_{0}=-[l / 2]$ instead of the naive choice $i_{0}=0$.

For instance, in the case $l=3$ one should consider the system (4-3) with $i=-1,0,1$, and not with $i=0,1,2$.

However, already in the case $l=4$ this simple recipe does not allow us to avoid considering $f^{2}$. In this case, the following way of dealing with the system (4-3) becomes useful.
$\mathcal{B}$ : Set $c_{l}=-1$ and consider instead of the homogeneous system (4-3) of $l$ equations the inhomogeneous system

$$
\begin{array}{r}
\left(c_{1} \varphi_{1}+\cdots+c_{l-1} \varphi_{l-1}\right) \circ f^{i}(x)=\varphi_{l} \circ f^{i}(x)  \tag{4-4}\\
i=i_{0}, i_{0}+1, \ldots, i_{0}+l-2
\end{array}
$$

of $l-1$ equations. Having found the (unique) solution $\left(c_{1}(x), \ldots, c_{l-1}(x)\right)$, prove that these functions are integrals of motion, that is,

$$
\begin{equation*}
c_{1}(f(x))=c_{1}(x), \quad \ldots, \quad c_{l-1}(f(x))=c_{l-1}(x) \tag{4-5}
\end{equation*}
$$

Thus, for instance, in the case $l=4$ one has to deal with the inhomogeneous system of equations (4-4) with $i=$ $-1,0,1$. Unfortunately, even if one is able to solve this system symbolically, the task of a symbolic verification of $(4-5)$ might become very hard due to the complexity of the solutions $\left(c_{1}(x), \ldots, c_{l-1}(x)\right)$.

This is the way taken, for instance, in [Kimura and Hirota 00]. In that paper, the task of verifying equations of type $(4-5)$ for the discrete-time Lagrange top is performed with the following method.
$\mathcal{G}:$ In order to verify that a rational function $c(x)=$ $p(x) / q(x)$ is an integral of motion of the map $\widetilde{x}=$ $f(x)$ coming from a system (1-1):
(i) find a Gröbner basis $G$ of the ideal $I$ generated by the components of $(1-1)$, considered as multilinear polynomials of $2 n$ variables $x, \widetilde{x}$ of total degree 2;
(ii) check, via polynomial division through elements of $G$, whether the polynomial $\delta(x, \widetilde{x})=$ $p(\widetilde{x}) q(x)-p(x) q(\widetilde{x})$ belongs to the ideal $I$.

An advantage of this method is that neither of its two steps needs the complicated explicit expressions for the map $f$. Nevertheless, both steps might be very demanding, especially the second step in the case of a complicated integral $c(x)$.

Sometimes, the task of verifying equations (4-5) can be circumvented by means of the following tricks.
$\mathcal{C}$ : Solve system (4-4) for two different but overlapping ranges $i \in\left[i_{0}, i_{0}+l-2\right]$ and $i \in\left[i_{1}, i_{1}+l-2\right]$. If the solutions coincide, then (4-5) holds automatically.

Indeed, in this situation the functions

$$
\left(c_{1}(x), \ldots, c_{l-1}(x)\right)
$$

solve the system with

$$
i \in\left[i_{0}, i_{0}+l-2\right] \cup\left[i_{1}, i_{1}+l-2\right]
$$

consisting of more than $l-1$ equations.
A clever modification of this idea, which allows one to avoid solving the second system, is as follows:
$\mathcal{D}$ : Suppose that the index range $i \in\left[i_{0}, i_{0}+l-2\right]$ in (4-4) contains 0 but is asymmetric. If the solution of this system $\left(c_{1}(x, \epsilon), \ldots, c_{l-1}(x, \epsilon)\right)$ is even with respect to $\epsilon$, then equations (4-5) hold automatically.
Indeed, the reversibility of the map $f^{-1}(x, \epsilon)=f(x,-\epsilon)$ yields in this case that equations of the system (4-4) are satisfied for $i \in\left[-\left(i_{0}+l-2\right),-i_{0}\right]$ as well, and the intervals $\left[i_{0}, i_{0}+l-2\right]$ and $\left[-\left(i_{0}+l-2\right),-i_{0}\right]$ overlap but do not coincide.

Finally, the most powerful method of reducing the number of iterations to be considered is as follows:
$\mathcal{E}:$ Often, the solutions $\left(c_{1}(x), \ldots, c_{l-1}(x)\right)$ satisfy some linear relations with constant coefficients. Find (observe) such relations numerically. Each such (still hypothetical) relation can be used to replace one equation in the system (4-4). Solve the resulting system symbolically, and proceed as in recipe $\mathcal{C}$ or $\mathcal{D}$ in order to verify $(4-5)$.
In some (rare) cases the integrals found by this approach are nice and simple enough to enable one to verify $(4-5)$ directly. Of course, it would be highly desirable to find some structures, like a Lax representation or biHamiltonian structure, that would allow one to check the conservation of integrals in a cleverer way, but up to now no such structures have been found for any of the Hirota-Kimura-type discretizations.

## 5. HIROTA-KIMURA DISCRETIZATION OF THE EULER TOP

We now illustrate the Hirota-Kimura mechanism by its application to the Euler top. This three-dimensional system is simple enough to enable one to perform all necessary computations symbolically, even by hand. At the same time, it provides a perfect illustration for many of the issues mentioned in the previous section.

### 5.1 Euler Top

The differential equations of motion of the Euler top read

$$
\begin{equation*}
\dot{x}_{1}=\alpha_{1} x_{2} x_{3}, \quad \dot{x}_{2}=\alpha_{2} x_{3} x_{1}, \quad \dot{x}_{3}=\alpha_{3} x_{1} x_{2} \tag{5-1}
\end{equation*}
$$

with $\alpha_{i}$ being real parameters of the system. This is one of the most famous integrable systems of classical mechanics, with an extensive literature devoted to it. We mention only that this system can be explicitly integrated in terms of elliptic functions, and it admits two functionally independent integrals of motion. Indeed, a quadratic function $H(x)=\gamma_{1} x_{1}^{2}+\gamma_{2} x_{2}^{2}+\gamma_{3} x_{3}^{2}$ is an integral for (5-1) if $\langle\gamma, \alpha\rangle=\gamma_{1} \alpha_{1}+\gamma_{2} \alpha_{2}+\gamma_{2} \alpha_{2}=0$. In particular, the following three functions are integrals of motion:

$$
\begin{aligned}
& H_{1}=\alpha_{3} x_{2}^{2}-\alpha_{2} x_{3}^{2} \\
& H_{2}=\alpha_{1} x_{3}^{2}-\alpha_{3} x_{1}^{2} \\
& H_{3}=\alpha_{2} x_{1}^{2}-\alpha_{1} x_{2}^{2}
\end{aligned}
$$

Clearly, only two of these are functionally independent because of $\alpha_{1} H_{1}+\alpha_{2} H_{2}+\alpha_{3} H_{3}=0$.

### 5.2 Discrete Equations of Motion

The Hirota-Kimura discretization of the Euler top introduced in [Hirota and Kimura 00] reads as follows:

$$
\left\{\begin{array}{l}
\widetilde{x}_{1}-x_{1}=\epsilon \alpha_{1}\left(\widetilde{x}_{2} x_{3}+x_{2} \widetilde{x}_{3}\right)  \tag{5-2}\\
\widetilde{x}_{2}-x_{2}=\epsilon \alpha_{2}\left(\widetilde{x}_{3} x_{1}+x_{3} \widetilde{x}_{1}\right) \\
\widetilde{x}_{3}-x_{3}=\epsilon \alpha_{3}\left(\widetilde{x}_{1} x_{2}+x_{1} \widetilde{x}_{2}\right)
\end{array}\right.
$$

Thus, the map $f: x \mapsto \widetilde{x}$ obtained by solving (5-2) for $\widetilde{x}$, is given by

$$
\begin{align*}
\widetilde{x} & =f(x, \epsilon)=A^{-1}(x, \epsilon) x  \tag{5-3}\\
A(x, \epsilon) & =\left(\begin{array}{ccc}
1 & -\epsilon \alpha_{1} x_{3} & -\epsilon \alpha_{1} x_{2} \\
-\epsilon \alpha_{2} x_{3} & 1 & -\epsilon \alpha_{2} x_{1} \\
-\epsilon \alpha_{3} x_{2} & -\epsilon \alpha_{3} x_{1} & 1
\end{array}\right) .
\end{align*}
$$

It might be instructive to have a look at the explicit formulas for this map:

$$
\begin{align*}
\widetilde{x}_{1}=\frac{1}{\Delta(x, \epsilon)}[ & x_{1}+2 \epsilon \alpha_{1} x_{2} x_{3} \\
& \left.+\epsilon^{2} x_{1}\left(-\alpha_{2} \alpha_{3} x_{1}^{2}+\alpha_{3} \alpha_{1} x_{2}^{2}+\alpha_{1} \alpha_{2} x_{3}^{2}\right)\right] \\
\widetilde{x}_{2}=\frac{1}{\Delta(x, \epsilon)}[ & x_{2}+2 \epsilon \alpha_{2} x_{3} x_{1}  \tag{5-4}\\
& \left.+\epsilon^{2} x_{2}\left(\alpha_{2} \alpha_{3} x_{1}^{2}-\alpha_{3} \alpha_{1} x_{2}^{2}+\alpha_{1} \alpha_{2} x_{3}^{2}\right)\right] \\
\widetilde{x}_{3}=\frac{1}{\Delta(x, \epsilon)}[ & x_{3}+2 \epsilon \alpha_{3} x_{1} x_{2} \\
& \left.+\epsilon^{2} x_{3}\left(\alpha_{2} \alpha_{3} x_{1}^{2}+\alpha_{3} \alpha_{1} x_{2}^{2}-\alpha_{1} \alpha_{2} x_{3}^{2}\right)\right]
\end{align*}
$$

where

$$
\begin{aligned}
\Delta(x, \epsilon)= & \operatorname{det} A(x, \epsilon) \\
= & 1-\epsilon^{2}\left(\alpha_{2} \alpha_{3} x_{1}^{2}+\alpha_{3} \alpha_{1} x_{2}^{2}+\alpha_{1} \alpha_{2} x_{3}^{2}\right) \\
& -2 \epsilon^{3} \alpha_{1} \alpha_{2} \alpha_{3} x_{1} x_{2} x_{3} .
\end{aligned}
$$

As always with HK-type discretizations, this map is birational, and one has the reversibility property

$$
f^{-1}(x, \epsilon)=f(x,-\epsilon)
$$

Apart from the Lax representation that is still missing, the discretization ( $5-3$ ) exhibits all the usual features of an integrable map: an invariant volume form, a biHamiltonian structure (that is, two compatible invariant Poisson structures), two functionally independent conserved quantities in involution, and solutions in terms of elliptic functions. The difference between its qualitative behavior and that of nonintegrable discretizations is striking; see Figure 2. For further details about the properties of this discretization we refer to [Hirota and Kimura 00] and [Petrera and Suris 07].

The integrals were first found in [Hirota and Kimura 00], apparently with the help of the approach discussed in the present work. However, since the resulting integrals are sufficiently simple and nice, their conservation can be easily verified by hand. Therefore, the paper [Hirota and Kimura 00] presents them in an ad hoc form, without explaining how they were discovered. We now try to reconstruct the way the results of [Hirota and Kimura 00 ] were originally found. To this end, we apply to the map (5-3) the method described in Section 3.

### 5.3 Hirota-Kimura Bases

Since all integrals of the Euler top are linear combinations of the functions $x_{k}^{2}$, it is natural to try the set

$$
\begin{equation*}
\Phi=\left(x_{1}^{2}, x_{2}^{2}, x_{3}^{2}, 1\right) \tag{5-5}
\end{equation*}
$$

as an HK basis for the discrete-time Euler top. An application of the numerical algorithm $\mathcal{N}$ suggests that the following statement holds:

Theorem 5.1. The set $(5-5)$ is an HK basis for the map (5-3) with $\operatorname{dim} K_{\Phi}(x)=2$. Therefore, any orbit of this map lies on the intersection of two quadrics in $\mathbb{R}^{3}$.

We will prove this theorem by finding two smaller HK bases with $d=1$. Namely, application of the numerical algorithm $\mathcal{N}$ suggests that omitting any one of the four functions $1, x_{k}^{2}$ from the basis $\Phi$ leads to an HK basis with $d=1$. In other words, for every $x \in \mathbb{R}^{3}$ there exists a one-dimensional space of vectors $\left(c_{1}, c_{2}, c_{3}\right)$ such that

$$
\left(c_{1} x_{1}^{2}+c_{2} x_{2}^{2}+c_{3} x_{3}^{2}\right) \circ f^{i}(x)=0, \quad i \in \mathbb{Z}
$$

as well as a one-dimensional space of vectors $\left(d_{1}, d_{2}, d_{4}\right)$ such that

$$
\left(d_{1} x_{1}^{2}+d_{2} x_{2}^{2}+d_{4}\right) \circ f^{i}(x)=0, \quad i \in \mathbb{Z}
$$

These numerical results can be now proven analytically.

Proposition 5.2. The set

$$
\Phi_{0}=\left(x_{1}^{2}, x_{2}^{2}, x_{3}^{2}\right)
$$

is an HK basis for the map (5-3) with $\operatorname{dim} K_{\Phi_{0}}(x)=1$. At each point $x \in \mathbb{R}^{3}$,

$$
\begin{aligned}
K_{\Phi_{0}}(x) & =\left[c_{1}: c_{2}: c_{3}\right] \\
& =\left[\alpha_{3} x_{2}^{2}-\alpha_{2} x_{3}^{2}: \alpha_{1} x_{3}^{2}-\alpha_{3} x_{1}^{2}: \alpha_{2} x_{1}^{2}-\alpha_{1} x_{2}^{2}\right]
\end{aligned}
$$

With $c_{3}=-1$, the functions

$$
\begin{equation*}
c_{1}(x)=\frac{\alpha_{3} x_{2}^{2}-\alpha_{2} x_{3}^{2}}{\alpha_{1} x_{2}^{2}-\alpha_{2} x_{1}^{2}}, \quad c_{2}(x)=\frac{\alpha_{1} x_{3}^{2}-\alpha_{3} x_{1}^{2}}{\alpha_{1} x_{2}^{2}-\alpha_{2} x_{1}^{2}} \tag{5-6}
\end{equation*}
$$

are integrals of motion of the map (5-3).
Proof: We proceed according to recipe $\mathcal{B}$. Set $c_{3}=-1$, and solve symbolically the system

$$
\begin{equation*}
\left(c_{1} x_{1}^{2}+c_{2} x_{2}^{2}\right) \circ f^{i}(x)=x_{3}^{2} \circ f^{i}(x), \quad i=0,1 \tag{5-7}
\end{equation*}
$$

which involves two inhomogeneous equations in two unknowns. System (5-7) can be written as

$$
\left\{\begin{array}{l}
c_{1} x_{1}^{2}+c_{2} x_{2}^{2}=x_{3}^{2}  \tag{5-8}\\
c_{1} \widetilde{x}_{1}^{2}+c_{2} \widetilde{x}_{2}^{2}=\widetilde{x}_{3}^{2}
\end{array}\right.
$$

where, of course, explicit formulas (5-4) have to be used for $\widetilde{x}_{k}$. The solution of this system is given by formulas


FIGURE 2. Left: a spiraling orbit of the explicit Euler method with time step $\epsilon=0.3$ applied to the Euler top. Right: a single orbit of the Hirota-Kimura discretization with the same time step, lying on an invariant spatial elliptic curve (intersection of two quadrics).
(5-6). The components of the solution do not depend on $\epsilon$. Therefore, according to recipe $\mathcal{D}$, we conclude that functions (5-6) are integrals of motion of the map (5-3).

It should be mentioned that the independence of the solution ( $c_{1}, c_{2}$ ) from $\epsilon$, or more generally, the dependence through even powers of $\epsilon$ only, which will be mentioned on many occasions below, starting with Proposition 5.3, is not given by any well-understood mechanism. Rather, it is just an instance of some very remarkable and miraculous cancellation of noneven polynomials. We illustrate this phenomenon by providing additional details to the previous proof. The solution of $(5-8)$ by Cramer's rule is given by ratios of determinants of type

$$
\begin{align*}
\left|\begin{array}{cc}
x_{i}^{2} & x_{j}^{2} \\
\widetilde{x}_{i}^{2} & \widetilde{x}_{j}^{2}
\end{array}\right|=\frac{1}{\Delta^{2}(x, \epsilon)} & {\left[4 \epsilon\left(\alpha_{j} x_{i}^{2}-\alpha_{i} x_{j}^{2}\right)\left(x_{1}+\epsilon \alpha_{1} x_{2} x_{3}\right)\right.} \\
& \left.\times\left(x_{2}+\epsilon \alpha_{2} x_{3} x_{1}\right)\left(x_{3}+\epsilon \alpha_{3} x_{1} x_{2}\right)\right] . \tag{5-9}
\end{align*}
$$

In the ratios of such determinants everything cancels out, except for the factors $\alpha_{j} x_{i}^{2}-\alpha_{i} x_{j}^{2}$. The cancellation of the denominators $\Delta^{2}(x, \epsilon)$ is, of course, no wonder, but the cancellation of the noneven factors in the numerators is rather miraculous.

One more typical phenomenon occurs in Proposition 5.2: although we have found apparently two integrals of motion (5-6), they turn out to be functionally dependent. Indeed, we have the identity

$$
\alpha_{1} c_{1}(x)+\alpha_{2} c_{2}(x)=\alpha_{3}
$$

so that for each $x \in \mathbb{R}^{3}$ the space $K_{\Phi_{0}}(x)$ is orthogonal to the constant vector $\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right)$. If one had guessed this relation numerically, one could have simplified the computation of the integrals $c_{1}, c_{2}$ by considering the system

$$
\left\{\begin{array}{r}
c_{1} x_{1}^{2}+c_{2} x_{2}^{2}=x_{3}^{2}  \tag{5-10}\\
c_{1} \alpha_{1}+c_{2} \alpha_{2}=\alpha_{3}
\end{array}\right.
$$

instead of (5-8).
Observe that existence of a linear relation allows one to reduce the number of iterates of $f$ involved in the linear system. (In the present situation, the system (5-10) contains no iterates of $f$ at all!) The latter system would lead to the same formulas (5-6); however, in this case one could not argue as in $\mathcal{D}$ and would be forced to prove that the functions (5-6) are integrals of motion directly, by verifying for them equations (4-5).

Anyhow, the existence of the HK basis $\Phi_{0}$ yields the existence of only one independent integral of the map $f$, which is not enough to ensure the integrability of $f$.

Proposition 5.3. The set

$$
\Phi_{1}=\left(x_{1}^{2}, x_{2}^{2}, 1\right)
$$

is an HK basis for the map (5-3) with $\operatorname{dim} K_{\Phi_{1}}(x)=1$. At each point $x \in \mathbb{R}^{3}$ one has

$$
K_{\Phi_{1}}(x)=\left[d_{1}: d_{2}:-1\right]
$$

where
$d_{1}(x)=\frac{\alpha_{2}\left(1-\epsilon^{2} \alpha_{3} \alpha_{1} x_{2}^{2}\right)}{\alpha_{2} x_{1}^{2}-\alpha_{1} x_{2}^{2}}, \quad d_{2}(x)=\frac{\alpha_{1}\left(1-\epsilon^{2} \alpha_{2} \alpha_{3} x_{1}^{2}\right)}{\alpha_{1} x_{2}^{2}-\alpha_{2} x_{1}^{2}}$.
These functions are integrals of motion of the map (5-3).

Proof: Following again prescription $\mathcal{B}$, we set $d_{4}=-1$, and solve symbolically the inhomogeneous system
or

$$
\left(d_{1} x_{1}^{2}+d_{2} x_{2}^{2}\right) \circ f^{i}(x)=1, \quad i=0,1
$$

$$
\left\{\begin{array}{l}
d_{1} x_{1}^{2}+d_{2} x_{2}^{2}=1 \\
d_{1} \widetilde{x}_{1}^{2}+d_{2} \widetilde{x}_{2}^{2}=1
\end{array}\right.
$$

The solution is given by $(5-11)$, due to (5-9) and

$$
\begin{aligned}
\left|\begin{array}{cc}
1 & x_{i}^{2} \\
1 & \widetilde{x}_{i}^{2}
\end{array}\right|=\frac{1}{\Delta^{2}(x, \epsilon)} & {\left[4 \epsilon \alpha_{i}\left(1-\epsilon^{2} \alpha_{j} \alpha_{k} x_{i}^{2}\right)\left(x_{1}+\epsilon \alpha_{1} x_{2} x_{3}\right)\right.} \\
& \left.\times\left(x_{2}+\epsilon \alpha_{2} x_{3} x_{1}\right)\left(x_{3}+\epsilon \alpha_{3} x_{1} x_{2}\right)\right]
\end{aligned}
$$

This time its components do depend on $\epsilon$, but are manifestly even functions of $\epsilon$. Everything noneven luckily cancels, again. Therefore, the argument $\mathcal{D}$ is still applicable, so that the functions $(5-11)$ are integrals of motion of the map $f$.

Functions (5-11) are again functionally dependent, because of

$$
\alpha_{1} d_{1}(x)+\alpha_{2} d_{2}(x)=\epsilon^{2} \alpha_{1} \alpha_{2} \alpha_{3}
$$

However, they are, clearly, functionally independent on the previously found functions (5-6), because $c_{1}, c_{2}$ depend on $x_{3}$, while $d_{1}, d_{2}$ do not.

Of course, the permutational symmetry yields that each of the sets of monomials $\Phi_{2}=\left(x_{2}^{2}, x_{3}^{2}, 1\right)$ and $\Phi_{3}=$ $\left(x_{1}^{2}, x_{3}^{2}, 1\right)$ is an HK basis as well, with $\operatorname{dim} K_{\Phi_{2}}(x)=$ $\operatorname{dim} K_{\Phi_{3}}(x)=1$. Any two of the four one-dimensional null spaces found span the full null space $K_{\Phi}(x)$. In particular, $K_{\Phi_{0}}(x)$ lies in $K_{\Phi_{1}}(x) \oplus K_{\Phi_{2}}(x)$.

Summarizing, we have found an HK basis with a twodimensional null space, as well as two functionally independent conserved quantities for the HK-discretization of the Euler top. Both results yield integrability of this discretization, in the sense that its orbits are confined to closed curves in $\mathbb{R}^{3}$. Moreover, each such curve is an intersection of two quadrics, which in the general-position case is an elliptic curve.

## 6. HIROTA-KIMURA-TYPE DISCRETIZATION OF THE CLEBSCH SYSTEM

### 6.1 The Clebsch System

The motion of a rigid body in an ideal fluid can be described by the so-called Kirchhoff equations
[Kirchhoff 70]:

$$
\left\{\begin{array}{l}
\dot{m}=m \times \frac{\partial H}{\partial m}+p \times \frac{\partial H}{\partial p}  \tag{6-1}\\
\dot{p}=p \times \frac{\partial H}{\partial m}
\end{array}\right.
$$

with $H$ being a quadratic form in $m=\left(m_{1}, m_{2}, m_{3}\right) \in \mathbb{R}^{3}$ and $p=\left(p_{1}, p_{2}, p_{3}\right) \in \mathbb{R}^{3}$; here $\times$ denotes the vector product in $\mathbb{R}^{3}$. The physical meaning of $m$ is the total angular momentum, whereas $p$ represents the total linear momentum of the system. System (6-1) is Hamiltonian with the Hamilton function $H(m, p)$, with respect to the Poisson bracket

$$
\begin{equation*}
\left\{m_{i}, m_{j}\right\}=m_{k}, \quad\left\{m_{i}, p_{j}\right\}=p_{k} \tag{6-2}
\end{equation*}
$$

where $(i, j, k)$ is a cyclic permutation of $(1,2,3)$ (all other pairwise Poisson brackets of the coordinate functions are obtained from these by skew-symmetry, or otherwise vanish). A detailed introduction to the general context of rigid-body dynamics and its mathematical foundations can be found in [Marsden and Ratiu 99].

A famous integrable case of the Kirchhoff equations was presented in [Clebsch 70] and is characterized by the Hamilton function $H=\frac{1}{2} \sum_{i=1}^{3}\left(m_{i}^{2}+\omega_{i} p_{i}^{2}\right)$. The corresponding equations of motion read

$$
\left\{\begin{array}{l}
\dot{m}=p \times \Omega p \\
\dot{p}=p \times m
\end{array}\right.
$$

where $\Omega=\operatorname{diag}\left(\omega_{1}, \omega_{2}, \omega_{3}\right)$ is the matrix of parameters, or in components,

$$
\begin{aligned}
\dot{m}_{1} & =\left(\omega_{3}-\omega_{2}\right) p_{2} p_{3} \\
\dot{m}_{2} & =\left(\omega_{1}-\omega_{3}\right) p_{3} p_{1} \\
\dot{m}_{3} & =\left(\omega_{2}-\omega_{1}\right) p_{1} p_{2} \\
\dot{p}_{1} & =m_{3} p_{2}-m_{2} p_{3} \\
\dot{p}_{2} & =m_{1} p_{3}-m_{3} p_{1} \\
\dot{p}_{3} & =m_{2} p_{1}-m_{1} p_{2} .
\end{aligned}
$$

This is the system that will be called the Clebsch system hereinafter. For an embedding of this system into the modern theory of integrable systems see [Perelomov 90, Reyman and Semenov-Tian-Shansky 94]. The Clebsch system possesses four independent quadratic integrals:

$$
\begin{align*}
H_{1}= & m_{1}^{2}+m_{2}^{2}+m_{3}^{2}+\omega_{1} p_{1}^{2}+\omega_{2} p_{2}^{2}+\omega_{3} p_{3}^{2},  \tag{6-3}\\
H_{2}= & \omega_{1} m_{1}^{2}+\omega_{2} m_{2}^{2}+\omega_{3} m_{3}^{2}-\omega_{2} \omega_{3} p_{1}^{2}-\omega_{3} \omega_{1} p_{2}^{2} \\
& -\omega_{1} \omega_{2} p_{3}^{2}  \tag{6-4}\\
H_{3}= & p_{1}^{2}+p_{2}^{2}+p_{3}^{2},  \tag{6-5}\\
H_{4}= & m_{1} p_{1}+m_{2} p_{2}+m_{3} p_{3} . \tag{6-6}
\end{align*}
$$

|  | $\operatorname{deg}$ | $\operatorname{deg}_{\boldsymbol{p}_{1}}$ | $\operatorname{deg}_{\boldsymbol{p}_{\mathbf{2}}}$ | $\operatorname{deg}_{\boldsymbol{p}_{\mathbf{3}}}$ | $\operatorname{deg}_{\boldsymbol{m}_{\boldsymbol{1}}}$ | $\operatorname{deg}_{\boldsymbol{m}_{\boldsymbol{2}}}$ | $\operatorname{deg}_{m_{\boldsymbol{3}}}$ |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| Common denominator of $\boldsymbol{f}^{\mathbf{2}}$ | 27 | 24 | 24 | 24 | 12 | 12 | 12 |
| Numerator of $\boldsymbol{p}_{\boldsymbol{1}}$-component of $\boldsymbol{f}^{\mathbf{2}}$ | 27 | 25 | 24 | 24 | 12 | 12 | 12 |
| Numerator of $\boldsymbol{p}_{\boldsymbol{2}}$-component of $\boldsymbol{f}^{\mathbf{2}}$ | 27 | 24 | 25 | 24 | 12 | 12 | 12 |
| Numerator of $\boldsymbol{p}_{\mathbf{3}}$-component of $\boldsymbol{f}^{\mathbf{2}}$ | 27 | 24 | 24 | 25 | 12 | 12 | 12 |
| Numerator of $\boldsymbol{m}_{\boldsymbol{1}}$-component of $\boldsymbol{f}^{\mathbf{2}}$ | 33 | 28 | 28 | 28 | 15 | 14 | 14 |
| Numerator of $\boldsymbol{m}_{\boldsymbol{2}}$-component of $\boldsymbol{f}^{\mathbf{2}}$ | 33 | 28 | 28 | 28 | 14 | 15 | 14 |
| Numerator of $\boldsymbol{m}_{\boldsymbol{3}}$-component of $\boldsymbol{f}^{\mathbf{2}}$ | 33 | 28 | 28 | 28 | 14 | 14 | 15 |

TABLE 1. Degrees of the numerators and the denominator of the second iterate $f^{2}(m, p)$.

These integrals are in involution with respect to the bracket (6-2). Moreover, $H_{3}, H_{4}$ are its Casimir functions (are in involution with any function on the phase space). However, the Hamiltonian structure will not play any role in the present paper. The set of linear combinations of the quadratic Hamiltonians $H_{1}, H_{2}, H_{3}$ coincides with the set of linear combinations of the functions

$$
\begin{align*}
& I_{1}=p_{1}^{2}+\frac{m_{2}^{2}}{\omega_{1}-\omega_{3}}+\frac{m_{3}^{2}}{\omega_{1}-\omega_{2}} \\
& I_{2}=p_{2}^{2}+\frac{m_{1}^{2}}{\omega_{2}-\omega_{3}}+\frac{m_{3}^{2}}{\omega_{2}-\omega_{1}}  \tag{6-7}\\
& I_{3}=p_{3}^{2}+\frac{m_{1}^{2}}{\omega_{3}-\omega_{2}}+\frac{m_{2}^{2}}{\omega_{3}-\omega_{1}}
\end{align*}
$$

For instance,

$$
\begin{aligned}
H_{1} & =\omega_{1} I_{1}+\omega_{2} I_{2}+\omega_{3} I_{3}, \\
H_{2} & =-\omega_{2} \omega_{3} I_{1}-\omega_{3} \omega_{1} I_{2}-\omega_{1} \omega_{2} I_{3} \\
H_{3} & =I_{1}+I_{2}+I_{3}
\end{aligned}
$$

### 6.2 Discrete Equations of Motion

Applying the Hirota-Kimura (or Kahan) approach to the Clebsch system, we arrive at the following discretization, proposed in [Ratiu 06]:

$$
\begin{align*}
\widetilde{m}_{1}-m_{1} & =\epsilon\left(\omega_{3}-\omega_{2}\right)\left(\widetilde{p}_{2} p_{3}+p_{2} \widetilde{p}_{3}\right),  \tag{6-8}\\
\widetilde{m}_{2}-m_{2} & =\epsilon\left(\omega_{1}-\omega_{3}\right)\left(\widetilde{p}_{3} p_{1}+p_{3} \widetilde{p}_{1}\right), \\
\widetilde{m}_{3}-m_{3} & =\epsilon\left(\omega_{2}-\omega_{1}\right)\left(\widetilde{p}_{1} p_{2}+p_{1} \widetilde{p}_{2}\right), \\
\widetilde{p}_{1}-p_{1} & =\epsilon\left(\widetilde{m}_{3} p_{2}+m_{3} \widetilde{p}_{2}\right)-\epsilon\left(\widetilde{m}_{2} p_{3}+m_{2} \widetilde{p}_{3}\right), \\
\widetilde{p}_{2}-p_{2} & =\epsilon\left(\widetilde{m}_{1} p_{3}+m_{1} \widetilde{p}_{3}\right)-\epsilon\left(\widetilde{m}_{3} p_{1}+m_{3} \widetilde{p}_{1}\right), \\
\widetilde{p}_{3}-p_{3} & =\epsilon\left(\widetilde{m}_{2} p_{1}+m_{2} \widetilde{p}_{1}\right)-\epsilon\left(\widetilde{m}_{1} p_{2}+m_{1} \widetilde{p}_{2}\right) .
\end{align*}
$$

In matrix form this can be written as

$$
M(m, p, \epsilon)\binom{\widetilde{m}}{\widetilde{p}}=\binom{m}{p}
$$

where

$$
\begin{aligned}
& M(m, p, \epsilon) \\
& \quad=\left(\begin{array}{cccccc}
1 & 0 & 0 & 0 & \epsilon \omega_{23} p_{3} & \epsilon \omega_{23} p_{2} \\
0 & 1 & 0 & \epsilon \omega_{31} p_{3} & 0 & \epsilon \omega_{31} p_{1} \\
0 & 0 & 1 & \epsilon \omega_{12} p_{2} & \epsilon \omega_{12} p_{1} & 0 \\
0 & \epsilon p_{3} & -\epsilon p_{2} & 1 & -\epsilon m_{3} & \epsilon m_{2} \\
-\epsilon p_{3} & 0 & \epsilon p_{1} & \epsilon m_{3} & 1 & -\epsilon m_{1} \\
\epsilon p_{2} & -\epsilon p_{1} & 0 & -\epsilon m_{2} & \epsilon m_{1} & 1
\end{array}\right),
\end{aligned}
$$

and the abbreviation $\omega_{i j}=\omega_{i}-\omega_{j}$ is used. The solution of this $6 \times 6$ linear system yields the birational map $f$ : $\mathbb{R}^{6} \rightarrow \mathbb{R}^{6}$,

$$
\begin{equation*}
\binom{\widetilde{m}}{\widetilde{p}}=f(m, p, \epsilon)=M^{-1}(m, p, \epsilon)\binom{m}{p}, \tag{6-9}
\end{equation*}
$$

called hereinafter the discrete Clebsch system. As usual, the reversibility property holds:

$$
\begin{equation*}
f^{-1}(m, p, \epsilon)=f(m, p,-\epsilon) \tag{6-10}
\end{equation*}
$$

A remark on the complexity of the iterates of $f$ is in order here. Each component of $(\widetilde{m}, \widetilde{p})=f(m, p)$ is a rational function whose numerator and denominator are polynomials in $m_{k}, p_{k}$ of total degree 6 . The numerators of $\widetilde{p}_{k}$ consist of 31 monomials; the numerators of $\widetilde{m}_{k}$ consist of 41 monomials; the common denominator consists of 28 monomials. It should be taken into account that the coefficients of all these polynomials depend, in turn, polynomially on $\epsilon$ and $\omega_{k}$, which additionally increases their complexity for a symbolic manipulator.

Expressions for the second iterate swell to astronomical length, prohibiting naive attempts to compute them symbolically. However, using Maple's LargeExpressions package [Carette et al. 06] and an appropriate veiling strategy, it is possible to obtain $f^{2}(m, p)$ with a reasonable amount of memory. Some impression of the complexity can be obtained from Table 1. The resulting expressions are too big to be used in further symbolic computations. Consider, for instance, the numerator of the $p_{1}$-component of $f^{2}(m, p)$. As a polynomial
of $m_{k}, p_{k}$, it contains 64,056 monomials; their coefficients are, in turn, polynomials in $\epsilon$ and $\omega_{k}$, and considered as a polynomial of the phase variables and the parameters, this expression contains $1,647,595$ terms.

### 6.3 Phase Portrait and Integrability

We now address the problem whether the discrete Clebsch system is integrable. Figures 3 and 4 show plots of the discrete Clebsch system (6-9), produced with MatLAB, for two different sets of parameter values. These plots indicate a quite regular behavior of the orbits of the discrete Clebsch system. Each orbit seems to fill out a two-dimensional surface in the 6-dimensional phase space. Leaving aside the Hamiltonian aspects of integrability, we are interested just in this simpler issue: do orbits of the map (6-9) lie on two-dimensional surfaces in $\mathbb{R}^{6}$ ?

A usual way to establish such a property would be to establish the existence of four functionally independent conserved quantities for this map. (We note in passing that plots of orbits are not very reliable in deciding about integrability. For instance, there are indications that Kahan's discretization (2-3) of the Lotka-Volterra system is nonintegrable, even if its orbits visually lie on closed curves in the phase plane. A strong magnification unveils the existence of very small regions in the phase plane with chaotic behavior.)

We will show that the answer to the above question is in the affirmative. To this end, we apply the approach based on the notion of HK basis. As a first step, we apply the numerical algorithm $\mathcal{N}$ to the maximal set of monomials, which includes all monomials of which the integrals $(6-3)-(6-6)$ of the continuous Clebsch system are built:

$$
\begin{array}{ll}
\varphi_{1}(m, p)=p_{1}^{2}, & \varphi_{2}(m, p)=p_{2}^{2}, \\
\varphi_{3}(m, p)=p_{3}^{2}, & \varphi_{4}(m, p)=m_{1}^{2}, \\
\varphi_{5}(m, p)=m_{2}^{2}, & \varphi_{6}(m, p)=m_{3}^{2}, \\
\varphi_{7}(m, p)=m_{1} p_{1}, & \varphi_{8}(m, p)=m_{2} p_{2}, \\
\varphi_{9}(m, p)=m_{3} p_{3}, & \varphi_{10}(m, p)=1 .
\end{array}
$$

We arrive at the following result:

Theorem 6.1. The set of functions

$$
\Phi=\left(p_{1}^{2}, p_{2}^{2}, p_{3}^{2}, m_{1}^{2}, m_{2}^{2}, m_{3}^{2}, m_{1} p_{1}, m_{2} p_{2}, m_{3} p_{3}, 1\right)
$$

is an HK basis for the map (6-9), with $\operatorname{dim} K_{\Phi}(m, p)=4$. Thus, any orbit of the map (6-9) lies on an intersection of four quadrics in $\mathbb{R}^{6}$.

At this point, Theorem 6.1 remains a numerical result, based on the algorithm $\mathcal{N}$. A direct symbolic proof of this statement is impossible, since it requires dealing with $f^{i}$, $i \in[-4,4]$, and the fourth iterate $f^{4}$ is a forbiddingly large expression. In order actually to prove Theorem 6.1 and to extract from it four independent integrals of motion, it is desirable to find HK (sub)bases with a smaller number of monomials, corresponding to some (preferably one-dimensional) subspaces of $K_{\Phi}(m, p)$. Much more detailed information on HK bases is provided by the following statement.

Theorem 6.2. The following four sets of functions are HK bases for the map (6-9) with one-dimensional null spaces:

$$
\begin{align*}
& \Phi_{0}=\left(p_{1}^{2}, p_{2}^{2}, p_{3}^{2}, 1\right)  \tag{6-11}\\
& \Phi_{1}=\left(p_{1}^{2}, p_{2}^{2}, p_{3}^{2}, m_{1}^{2}, m_{2}^{2}, m_{3}^{2}, m_{1} p_{1}\right)  \tag{6-12}\\
& \Phi_{2}=\left(p_{1}^{2}, p_{2}^{2}, p_{3}^{2}, m_{1}^{2}, m_{2}^{2}, m_{3}^{2}, m_{2} p_{2}\right)  \tag{6-13}\\
& \Phi_{3}=\left(p_{1}^{2}, p_{2}^{2}, p_{3}^{2}, m_{1}^{2}, m_{2}^{2}, m_{3}^{2}, m_{3} p_{3}\right) \tag{6-14}
\end{align*}
$$

If all the null spaces are considered as subspaces of $\mathbb{R}^{10}$, so that

$$
\begin{aligned}
K_{\Phi_{0}} & =\left[c_{1}: c_{2}: c_{3}: 0: 0: 0: 0: 0: 0: c_{10}\right] \\
K_{\Phi_{1}} & =\left[\alpha_{1}: \alpha_{2}: \alpha_{3}: \alpha_{4}: \alpha_{5}: \alpha_{6}: \alpha_{7}: 0: 0: 0\right], \\
K_{\Phi_{2}} & =\left[\beta_{1}: \beta_{2}: \beta_{3}: \beta_{4}: \beta_{5}: \beta_{6}: 0: \beta_{8}: 0: 0\right], \\
K_{\Phi_{3}} & =\left[\gamma_{1}: \gamma_{2}: \gamma_{3}: \gamma_{4}: \gamma_{5}: \gamma_{6}: 0: 0: \gamma_{9}: 0\right],
\end{aligned}
$$

then

$$
K_{\Phi}=K_{\Phi_{0}} \oplus K_{\Phi_{1}} \oplus K_{\Phi_{2}} \oplus K_{\Phi_{3}}
$$

This statement, too, was first found with the help of numerical experiments based on the algorithm $\mathcal{N}$. In what follows, we will discuss how these claims can be given a rigorous (computer-assisted) proof, and how much additional information (for instance, about conserved quantities for the map (6-9)) can be extracted from such a proof.

### 6.4 First HK Basis

Theorem 6.3. The set (6-11) is an HK basis for the map (6-9) with $\operatorname{dim} K_{\Phi_{0}}(m, p)=1$. At each point


FIGURE 3. An orbit of the discrete Clebsch system with $\omega_{1}=1, \omega_{2}=0.2, \omega_{3}=30$, and $\epsilon=1$; initial point $\left(m_{0}, p_{0}\right)=$ $(1,1,1,1,1,1)$.


FIGURE 4. An orbit of the discrete Clebsch system with $\omega_{1}=0.1, \omega_{2}=0.2, \omega_{3}=0.3$, and $\epsilon=1$; initial point $\left(m_{0}, p_{0}\right)=(1,1,1,1,1,1)$.
$(m, p) \in \mathbb{R}^{6}$,

$$
\begin{align*}
K_{\Phi_{0}}( & m, p)=\left[c_{1}: c_{2}: c_{3}: c_{10}\right] \\
= & {\left[\frac{1+\epsilon^{2}\left(\omega_{1}-\omega_{2}\right) p_{2}^{2}+\epsilon^{2}\left(\omega_{1}-\omega_{3}\right) p_{3}^{2}}{p_{1}^{2}+p_{2}^{2}+p_{3}^{2}}:\right.} \\
& \frac{1+\epsilon^{2}\left(\omega_{2}-\omega_{1}\right) p_{1}^{2}+\epsilon^{2}\left(\omega_{2}-\omega_{3}\right) p_{3}^{2}}{p_{1}^{2}+p_{2}^{2}+p_{3}^{2}}: \\
& \left.\frac{1+\epsilon^{2}\left(\omega_{3}-\omega_{1}\right) p_{1}^{2}+\epsilon^{2}\left(\omega_{3}-\omega_{2}\right) p_{2}^{2}}{p_{1}^{2}+p_{2}^{2}+p_{3}^{2}}:-1\right] \\
= & {\left[\frac{1}{J}+\epsilon^{2} \omega_{1}: \frac{1}{J}+\epsilon^{2} \omega_{2}: \frac{1}{J}+\epsilon^{2} \omega_{3}:-1\right] } \tag{6-15}
\end{align*}
$$

where

$$
\begin{equation*}
J(m, p, \epsilon)=\frac{p_{1}^{2}+p_{2}^{2}+p_{3}^{2}}{1-\epsilon^{2}\left(\omega_{1} p_{1}^{2}+\omega_{2} p_{2}^{2}+\omega_{3} p_{3}^{2}\right)} . \tag{6-16}
\end{equation*}
$$

The function (6-16) is an integral of motion of the map (6-9).

Proof: The statement of the theorem means that for every $(m, p) \in \mathbb{R}^{6}$ the space of solutions of the homogeneous system

$$
\left(c_{1} p_{1}^{2}+c_{2} p_{2}^{2}+c_{3} p_{3}^{2}+c_{10}\right) \circ f^{i}(m, p)=0, \quad i=0, \ldots, 3
$$

is one-dimensional. This system involves the third iterate of $f$; therefore its symbolic treatment is impossible. According to the strategy $\mathcal{B}$, we set $c_{10}=-1$ and consider the inhomogeneous system

$$
\left(c_{1} p_{1}^{2}+c_{2} p_{2}^{2}+c_{3} p_{3}^{2}\right) \circ f^{i}(m, p)=1, \quad i=0,1,2
$$



FIGURE 5. Plot of the coefficients $c_{1}, c_{2}, c_{3}$.

This system involves the second iterate of $f$, which still precludes its symbolic treatment. There are now several ways in which to proceed.

- First, we could follow the recipe $\mathcal{E}$ and find further information about the solutions $c_{i}$. To this end, we plot the points $\left(c_{1}(m, p), c_{2}(m, p), c_{3}(m, p)\right)$ for different initial data $(m, p) \in \mathbb{R}^{6}$. Figure 5 shows such a plot, with 300 initial data points ( $m, p$ ) randomly chosen from the set $[0,1]^{6}$. The points $\left(c_{1}(m, p), c_{2}(m, p), c_{3}(m, p)\right)$ seem to lie on a line in $\mathbb{R}^{3}$, which means that there should be two linear dependencies among the functions $c_{1}, c_{2}$, and $c_{3}$.
In order to identify these linear dependencies, we run the PSLQ algorithm [Ferguson and Bailey 91, Ferguson et al. 99] with the vectors $\left(c_{1}, c_{2}, 1\right)$ as input (see Remark 6.4 after the end of the proof, concerning implementation of this step). In this way we obtain the conjecture

$$
c_{1}-c_{2}=\epsilon^{2}\left(\omega_{1}-\omega_{2}\right)
$$

Similarly, running the PSLQ algorithm with the vectors $\left(c_{2}, c_{3}, 1\right)$ as input leads to the conjecture

$$
c_{2}-c_{3}=\epsilon^{2}\left(\omega_{2}-\omega_{3}\right)
$$

Having identified (numerically!) these two linear relations, we use them instead of two equations in the system (6-17)), say the equations for $i=1,2$. The resulting system becomes extremely simple:

$$
\left\{\begin{array}{l}
c_{1} p_{1}^{2}+c_{2} p_{2}^{2}+c_{3} p_{3}^{2}=1 \\
c_{1}-c_{2}=\epsilon^{2}\left(\omega_{1}-\omega_{2}\right) \\
c_{2}-c_{3}=\epsilon^{2}\left(\omega_{2}-\omega_{3}\right)
\end{array}\right.
$$

It contains no iterates of $f$ at all and can be solved immediately by hand, with the result (6-15). It should be stressed that this result still remains conjectural, and one has to prove a posteriori that the functions $c_{1}, c_{2}, c_{3}$ are integrals of motion.

- Alternatively, we can combine the above approach based on the prescription $\mathcal{E}$ with the recipe $\mathcal{D}$. For this, we use just one of the linear dependencies found above to replace the equation in $(6-17)$ with $i=2$, and then let Maple solve the remaining system. The computation takes 22.33 seconds on a $1.83-\mathrm{GHz}$ Core Duo PC and consumes 32.43 MB RAM. The output is still as in (6-15), but arguing in this way, one does not need to verify a posteriori that $c_{1}, c_{2}, c_{3}$ are integrals of motion, because they are manifestly even functions of $\epsilon$, while the symmetry of the linear system with respect to $\epsilon$ has been broken.

To finish the proof along the lines of the first of the possible arguments above, we show how to verify the statement that the function $J$ in $(6-16)$ is an integral of motion, i.e., that

$$
\frac{p_{1}^{2}+p_{2}^{2}+p_{3}^{2}}{1-\epsilon^{2}\left(\omega_{1} p_{1}^{2}+\omega_{2} p_{2}^{2}+\omega_{3} p_{3}^{2}\right)}=\frac{\widetilde{p}_{1}^{2}+\widetilde{p}_{2}^{2}+\widetilde{p}_{3}^{2}}{1-\epsilon^{2}\left(\omega_{1} \widetilde{p}_{1}^{2}+\omega_{2} \widetilde{p}_{2}^{2}+\omega_{3} \widetilde{p}_{3}^{2}\right)}
$$

This is equivalent to

$$
\begin{aligned}
& \widetilde{p}_{1}^{2}-p_{1}^{2}+\widetilde{p}_{2}^{2}-p_{2}^{2}+\widetilde{p}_{3}^{2}-p_{3}^{2} \\
& =\epsilon^{2}\left[\left(\omega_{2}-\omega_{1}\right)\left(\widetilde{p}_{1}^{2} p_{2}^{2}-\widetilde{p}_{2}^{2} p_{1}^{2}\right)+\left(\omega_{3}-\omega_{2}\right)\left(\widetilde{p}_{2}^{2} p_{3}^{2}-\widetilde{p}_{3}^{2} p_{2}^{2}\right)\right. \\
& \quad \\
& \left.\quad+\left(\omega_{1}-\omega_{3}\right)\left(\widetilde{p}_{3}^{2} p_{1}^{2}-\widetilde{p}_{1}^{2} p_{3}^{2}\right)\right]
\end{aligned}
$$

On the left-hand side of this equation we replace $\widetilde{p}_{i}-$ $p_{i}$ with the expressions from the last three equations of motion (6-8); on the right-hand side we replace $\epsilon\left(\omega_{k}-\omega_{j}\right)\left(\widetilde{p}_{j} p_{k}+p_{j} \widetilde{p}_{k}\right)$ by $\widetilde{m}_{i}-m_{i}$, according to the first three equations of motion (6-8). This brings the equation we want to prove into the form

$$
\begin{aligned}
& \left(\widetilde{p}_{1}+p_{1}\right)\left(\widetilde{m}_{3} p_{2}+m_{3} \widetilde{p}_{2}-\widetilde{m}_{2} p_{3}-m_{2} \widetilde{p}_{3}\right) \\
& \quad+\left(\widetilde{p}_{2}+p_{2}\right)\left(\widetilde{m}_{1} p_{3}+m_{1} \widetilde{p}_{3}-\widetilde{m}_{3} p_{1}-m_{3} \widetilde{p}_{1}\right) \\
& \quad+\left(\widetilde{p}_{3}+p_{3}\right)\left(\widetilde{m}_{2} p_{1}+m_{2} \widetilde{p}_{1}-\widetilde{m}_{1} p_{2}-m_{1} \widetilde{p}_{2}\right) \\
& \quad=\left(\widetilde{p}_{1} p_{2}-p_{1} \widetilde{p}_{2}\right)\left(\widetilde{m}_{3}-m_{3}\right) \\
& \quad+\left(\widetilde{p}_{2} p_{3}-p_{2} \widetilde{p}_{3}\right)\left(\widetilde{m}_{1}-m_{1}\right) \\
& \quad+\left(\widetilde{p}_{3} p_{1}-p_{3} \widetilde{p}_{1}\right)\left(\widetilde{m}_{2}-m_{2}\right) .
\end{aligned}
$$

But the latter equation is an algebraic identity in twelve variables $m_{k}, p_{k}, \widetilde{m}_{k}, \widetilde{p}_{k}$. This finishes the proof.

Remark 6.4. In the above proof and on many occasions below we make use of the PSLQ algorithm in order to identify possible linear relations among conserved quantities. Its applications are well documented in the literature on experimental mathematics [Borwein and Bailey 03, Borwein et al. 04], so that we restrict ourselves here to a couple of minor remarks.

We apply the PSLQ algorithm to the numerical values of (the candidates for) the conserved quantities obtained from the algorithm $\mathcal{N}$. We note that it is crucial to apply the PSLQ algorithm with many different initial data; from the large amount of possible linear relations one should, of course, filter out those relations that remain unaltered for different initial data.

It proved useful to perform these computations with rational data (initial values of phase variables and parameters of the map) as well as with high-precision floatingpoint numbers. In our experiments we have been able to automate this task to a large extent. All computations of this kind were performed on an Apple MacBook with a $1.83-\mathrm{GHz}$ Intel Core Duo processor and 2 GB of RAM.

### 6.5 Remaining HK Bases

We now consider the remaining HK bases $\Phi_{1}, \Phi_{2}$, and $\Phi_{3}$. Here we are dealing with the three linear systems

$$
\begin{align*}
& \left(\alpha_{1} p_{1}^{2}+\alpha_{2} p_{2}^{2}+\alpha_{3} p_{3}^{2}+\alpha_{4} m_{1}^{2}+\alpha_{5} m_{2}^{2}+\alpha_{6} m_{3}^{2}\right) \circ f^{i}(m, p) \\
& \quad=m_{1} p_{1} \circ f^{i}(m, p),  \tag{6-18}\\
& \left(\beta_{1} p_{1}^{2}+\beta_{2} p_{2}^{2}+\beta_{3} p_{3}^{2}+\beta_{4} m_{1}^{2}+\beta_{5} m_{2}^{2}+\beta_{6} m_{3}^{2}\right) \circ f^{i}(m, p) \\
& \quad=m_{2} p_{2} \circ f^{i}(m, p),  \tag{6-19}\\
& \left(\gamma_{1} p_{1}^{2}+\gamma_{2} p_{2}^{2}+\gamma_{3} p_{3}^{2}+\gamma_{4} m_{1}^{2}+\gamma_{5} m_{2}^{2}+\gamma_{6} m_{3}^{2}\right) \circ f^{i}(m, p) \\
& \quad=m_{3} p_{3} \circ f^{i}(m, p), \tag{6-20}
\end{align*}
$$

already made inhomogeneous by normalizing the last coefficient in each system, as in recipe $\mathcal{B}$, with $l=7$. The claim about each of the systems is that it admits a unique solution for $i \in \mathbb{Z}$. It is enough to solve each system for two different but intersecting ranges of $l-1=6$ consecutive indices $i$, such as $i \in[-2,3]$ and $i \in[-3,2]$, and to show that solutions coincide for both ranges (recipe $\mathcal{C}$ ). Actually, since the index range $i \in[-2,3]$ is asymmetric, it would be enough to consider the system for this one range and to show that the solutions $\alpha_{j}, \beta_{j}, \gamma_{j}$ are even functions with respect to $\epsilon$ (recipe $\mathcal{D}$ ). However, symbolic manipulations with the iterates $f^{i}$ for $i= \pm 2, \pm 3$ are impossible. In what follows, we will gradually extend the available information about the coefficients $\alpha_{j}, \beta_{j}, \gamma_{j}$, which at the end will allow us to get the analytic expres-
sions for all of them and to prove that they are indeed integrals.

### 6.6 First Additional HK Basis

Theorem 6.2 shows that after finding the HK basis $\Phi_{0}$ with $\operatorname{dim} K_{\Phi_{0}}(x)=1$ it is enough to concentrate on (sub)bases not containing the constant function $\varphi_{10}(m, p)=1$. It turns out to be possible to find an HK basis without $\varphi_{10}$ and with a one-dimensional null space, which is more amenable to a symbolic treatment than $\Phi_{1}, \Phi_{2}, \Phi_{3}$. Numerical algorithm $\mathcal{N}$ suggests that the following set of functions is an HK basis with $d=1$ :

$$
\begin{equation*}
\Psi=\left(p_{1}^{2}, p_{2}^{2}, p_{3}^{2}, m_{1} p_{1}, m_{2} p_{2}, m_{3} p_{3}\right) \tag{6-21}
\end{equation*}
$$

Theorem 6.5. The set (6-21) is an HK basis for the map (6-9) with $\operatorname{dim} K_{\Psi}(m, p)=1$. At every point $(m, p) \in$ $\mathbb{R}^{6}$,

$$
K_{\Psi}(m, p)=\left[-1:-1:-1: d_{7}: d_{8}: d_{9}\right]
$$

with

$$
\begin{equation*}
d_{k}=\frac{\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right)\left(1+\epsilon^{2} d_{k}^{(2)}+\epsilon^{4} d_{k}^{(4)}+\epsilon^{6} d_{k}^{(6)}\right)}{\Delta} \tag{6-22}
\end{equation*}
$$

$$
\begin{align*}
& k=7,8,9, \text { and } \\
& \Delta=m_{1} p_{1}+m_{2} p_{2}+m_{3} p_{3}+\epsilon^{2} \Delta^{(4)}+\epsilon^{4} \Delta^{(6)}+\epsilon^{6} \Delta^{(8)}, \tag{6-23}
\end{align*}
$$

where $d_{k}^{(2 q)}$ and $\Delta^{(2 q)}$ are homogeneous polynomials of degree $2 q$ in phase variables. In particular,

$$
\begin{aligned}
d_{7}^{(2)}= & m_{1}^{2}+m_{2}^{2}+m_{3}^{2}+\left(\omega_{2}+\omega_{3}-2 \omega_{1}\right) p_{1}^{2}+\left(\omega_{3}-\omega_{2}\right) p_{2}^{2} \\
& +\left(\omega_{2}-\omega_{3}\right) p_{3}^{2}, \\
d_{8}^{(2)}= & m_{1}^{2}+m_{2}^{2}+m_{3}^{2}+\left(\omega_{3}-\omega_{1}\right) p_{1}^{2}+\left(\omega_{3}+\omega_{1}-2 \omega_{2}\right) p_{2}^{2} \\
& +\left(\omega_{1}-\omega_{3}\right) p_{3}^{2}, \\
d_{9}^{(2)}= & m_{1}^{2}+m_{2}^{2}+m_{3}^{2}+\left(\omega_{2}-\omega_{1}\right) p_{1}^{2}+\left(\omega_{1}-\omega_{2}\right) p_{2}^{2} \\
& +\left(\omega_{1}+\omega_{2}-2 \omega_{3}\right) p_{3}^{2},
\end{aligned}
$$

and

$$
\Delta^{(4)}=m_{1} p_{1} d_{7}^{(2)}+m_{2} p_{2} d_{8}^{(2)}+m_{3} p_{3} d_{9}^{(2)}
$$

(All other polynomials are too messy to be given here.)
The functions $d_{7}, d_{8}, d_{9}$ are integrals of the map (6-9).
They are dependent due to the linear relation

$$
\begin{equation*}
\left(\omega_{2}-\omega_{3}\right) d_{7}+\left(\omega_{3}-\omega_{1}\right) d_{8}+\left(\omega_{1}-\omega_{2}\right) d_{9}=0 \tag{6-24}
\end{equation*}
$$

Any two of them are functionally independent. Moreover, any two of them together with $J$ are still functionally independent.

Proof: As already mentioned, numerical experiments suggest that for any $(m, p) \in \mathbb{R}^{6}$ there exists a onedimensional space of vectors $\left(d_{1}, d_{2}, d_{3}, d_{7}, d_{8}, d_{9}\right)$ satisfying

$$
\begin{aligned}
& \left(d_{1} p_{1}^{2}+d_{2} p_{2}^{2}+d_{3} p_{3}^{2}+d_{7} m_{1} p_{1}+d_{8} m_{2} p_{2}+d_{9} m_{3} p_{3}\right) \\
& \quad \circ f^{i}(m, p)=0
\end{aligned}
$$

for $i=0,1, \ldots, 5$. According to recipe $\mathcal{A}$, one can equally well consider this system for $i=-2,-1, \ldots, 3$, which, however, still contains the third iterate of $f$ and is therefore not manageable. Therefore, we apply recipe $\mathcal{E}$ and look for linear relations between the (numerical) solutions. Two such relations can be observed immediately, namely

$$
\begin{equation*}
d_{1}=d_{2}=d_{3} \tag{6-25}
\end{equation*}
$$

Accepting these (still hypothetical) relations and applying recipe $\mathcal{B}$, i.e., setting the common value of (6-25) equal to -1 , we arrive at the inhomogeneous system of only three linear relations

$$
\begin{align*}
& \left(d_{7} m_{1} p_{1}+d_{8} m_{2} p_{2}+d_{9} m_{3} p_{3}\right) \circ f^{i}(m, p)  \tag{6-26}\\
& \quad=\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right) \circ f^{i}(m, p)
\end{align*}
$$

for $i=-1,0,1$. Fortunately, it is possible to find one more linear relation among $d_{7}, d_{8}, d_{9}$. This was discovered numerically: we produced a three-dimensional plot of the points $\left(d_{7}(m, p), d_{8}(m, p), d_{9}(m, p)\right)$, which can be seen in Figure 6 in two different projections. This figure suggests that all these points lie on a plane in $\mathbb{R}^{3}$, the second picture being a "side view" along a direction parallel to this plane. Thus, it is plausible that one more linear relation exists. With the help of the PSLQ algorithm this hypothetical relation can be then identified as (6-24). Now the ansatz (6-26) is reduced to the following system of three equations for $\left(d_{7}, d_{8}, d_{9}\right)$, which involves only one iterate of the map $f$ :

$$
\left\{\begin{array}{l}
\left(d_{7} m_{1} p_{1}+d_{8} m_{2} p_{2}+d_{9} m_{3} p_{3}\right) \circ f^{i}(m, p)  \tag{6-27}\\
\quad=\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right) \circ f^{i}(m, p), \quad i=0,1, \\
\left(\omega_{2}-\omega_{3}\right) d_{7}+\left(\omega_{3}-\omega_{1}\right) d_{8}+\left(\omega_{2}-\omega_{2}\right) d_{9}=0 .
\end{array}\right.
$$

This system can be solved by Maple, resulting in functions given in $(6-22),(6-23) .{ }^{2}$. They are manifestly even functions of $\epsilon$, while the system has no symmetry with respect to $\epsilon \mapsto-\epsilon$. This proves that they are integrals of motion for the map $f$. This argument slightly generalizes the recipes $\mathcal{D}$ and $\mathcal{E}$, and since it is used not only here

[^1]but also on several further occasions in this paper, we give here its formalization.

Proposition 6.6. Consider a map $f: \mathbb{R}^{6} \rightarrow \mathbb{R}^{6}$ depending on a parameter $\epsilon$, reversible in the sense of (6-10). Let $I(m, p, \epsilon)$ be an integral of $f$, even in $\epsilon$, and let $A_{1}, A_{2}, A_{3} \in \mathbb{R}$. Suppose that the set of functions $\Phi=\left(\varphi_{1}, \ldots, \varphi_{4}\right)$ is such that the system of three linear equations for $\left(a_{1}, a_{2}, a_{3}\right)$,
$\left\{\begin{array}{l}\left(a_{1} \varphi_{1}+a_{2} \varphi_{2}+a_{3} \varphi_{3}\right) \circ f^{i}(m, p, \epsilon)=\varphi_{4} \circ f^{i}(m, p, \epsilon), \\ \quad i=0,1, \\ A_{1} a_{1}+A_{2} a_{2}+A_{3} a_{3}=I(m, p, \epsilon),\end{array}\right.$
admits a unique solution that is even with respect to $\epsilon$. Then this solution $\left(a_{1}, a_{2}, a_{3}\right)$ consists of integrals of the map $f$, and $\Phi$ is an $H K$ basis with $\operatorname{dim} K_{\Phi}(m, p)=1$.

Proof: Since $\left(a_{1}, a_{2}, a_{3}\right)$ are even functions of $\epsilon$, they satisfy also the system (6-28) with $\epsilon \mapsto-\epsilon$, which, due to the reversibility $(6-10)$, can be represented as
$\left\{\begin{array}{l}\left(a_{1} \varphi_{1}+a_{2} \varphi_{2}+a_{3} \varphi_{3}\right) \circ f^{i}(m, p, \epsilon)=\varphi_{4} \circ f^{i}(m, p, \epsilon), \\ i=0,-1, \\ A_{1} a_{1}+A_{2} a_{2}+A_{3} a_{3}=I(m, p, \epsilon) .\end{array}\right.$
Since the functions $\left(a_{1}, a_{2}, a_{3}\right)$ are uniquely determined by any of the systems ( $6-28$ ), ( $6-29$ ), we conclude that they remain invariant under the change $(m, p) \mapsto$ $f(m, p, \epsilon)$, or in other words, that they are integrals of motion. Finally, we can conclude that these functions satisfy the equation $\left(a_{1} \varphi_{1}+a_{2} \varphi_{2}+a_{3} \varphi_{3}\right) \circ f^{i}=\varphi_{4} \circ f^{i}$ for all $i \in \mathbb{Z}$ (and can be uniquely determined by this property), and that the linear relation $A_{1} a_{1}+A_{2} a_{2}+A_{3} a_{3}=I$ is satisfied as well.

Application of Proposition 6.6 to system (6-27) shows that $d_{7}, d_{8}, d_{9}$ are integrals of motion, since they are even in $\epsilon$. Note that here, as always in similar context, the evenness of solutions is due to "miraculous cancellation" of the equal noneven polynomials that factor out both in the numerators and denominators of the solutions. In the present computation, these common noneven factors are of degree 2 in $\epsilon$.

It remains to prove that any two of the integrals $d_{7}, d_{8}, d_{9}$ together with the previously found integral $J$ are functionally independent. To this end, we show that from such a triple of integrals one can construct another triple of integrals that yields in the limit $\epsilon \rightarrow 0$ three independent conserved quantities $H_{3}, H_{4}, H_{1}$ of the con-


FIGURE 6. Plot of the points $\left(d_{7}, d_{8}, d_{9}\right)$ for 729 values of $(m, p)$ from a six-dimensional grid around the point $(1,1,1,1,1,1)$ with a grid size of 0.01 and the parameters $\epsilon=0.1, \omega_{1}=0.1, \omega_{2}=0.2, \omega_{3}=0.3$.
tinuous Clebsch system. Indeed,

$$
\begin{aligned}
J & =p_{1}^{2}+p_{2}^{2}+p_{3}^{2}+O\left(\epsilon^{2}\right)=H_{3}+O\left(\epsilon^{2}\right) \\
\frac{J}{d_{k+6}} & =m_{1} p_{1}+m_{2} p_{2}+m_{3} p_{3}+O\left(\epsilon^{2}\right)=H_{4}+O\left(\epsilon^{2}\right) .
\end{aligned}
$$

On the other hand, it is easy to derive

$$
\begin{aligned}
\frac{d_{7}}{d_{8}} & =1+\epsilon^{2}\left(d_{7}^{(2)}-d_{8}^{(2)}\right)+O\left(\epsilon^{4}\right) \\
& =1+\epsilon^{2}\left(\omega_{2}-\omega_{1}\right)\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right)+O\left(\epsilon^{4}\right)
\end{aligned}
$$

and taking this into account and computing the terms of order $\epsilon^{4}$, one obtains

$$
\begin{aligned}
& \frac{d_{7}}{d_{8}}-1-\epsilon^{2}\left(\omega_{2}-\omega_{1}\right) J \\
& \quad=\epsilon^{4}\left(\omega_{2}-\omega_{1}\right)\left(2 H_{4}^{2}+\omega_{2} H_{3}^{2}-2 H_{3} H_{1}\right)+O\left(\epsilon^{6}\right)
\end{aligned}
$$

from which one easily extracts $H_{1}$. This proves our claim.

Remark 6.7. With the basis $\Psi$, we encounter for the first time the following interesting phenomenon: it can happen that an HK basis with a one-dimensional null space provides several (in this case two) functionally independent integrals. With Theorem 6.5, we established the existence of three independent conserved quantities and two HK bases with linearly independent null spaces. So, every orbit of the discrete Clebsch system is shown to lie in a three-dimensional manifold that belongs to an intersection of two quadrics in $\mathbb{R}^{6}$.

The aim of the following is to find one more independent integral and two more HK bases with onedimensional null spaces linearly independent of $K_{\Phi_{0}}, K_{\Psi}$.

### 6.7 Second Additional HK Basis

From the (still hypothetical) properties (6-18))-(6-20) of the bases $\Phi_{1}, \Phi_{2}, \Phi_{3}$ it follows that for any $(m, p) \in \mathbb{R}^{6}$ the system of linear equations

$$
\begin{align*}
& \left(g_{1} p_{1}^{2}+g_{2} p_{2}^{2}+g_{3} p_{3}^{2}+g_{4} m_{1}^{2}+g_{5} m_{2}^{2}+g_{6} m_{3}^{2}\right) \circ f^{i}(m, p) \\
& \quad=\left(m_{1} p_{1}+m_{2} p_{2}+m_{3} p_{3}\right) \circ f^{i}(m, p) \tag{6-30}
\end{align*}
$$

has a unique solution $\left(g_{1}, g_{2}, g_{3}, g_{4}, g_{5}, g_{6}\right)$. Indeed, the solution should be given by

$$
\begin{equation*}
g_{j}=\alpha_{j}+\beta_{j}+\gamma_{j}, \quad j=1, \ldots, 6 \tag{6-31}
\end{equation*}
$$

As for the bases $\Phi_{1}, \Phi_{2}, \Phi_{3}$, the solution of (6-30) can be determined by solving these equations for two different but intersecting ranges of six consecutive values of $i$, say for $i \in[-3,2]$ and $i \in[-2,3]$. However, it turns out that due to the existence of several linear relations among the solutions $g_{j}$, system $(6-30)$ is much easier to deal with than systems $(6-18)-(6-20)$, so that the functions $g_{j}$ can be determined and studied independently of $\alpha_{j}, \beta_{j}, \gamma_{j}$.

Theorem 6.8. The set of functions

$$
\Theta=\left(p_{1}^{2}, p_{2}^{2}, p_{3}^{2}, m_{1}^{2}, m_{2}^{2}, m_{3}^{2}, m_{1} p_{1}+m_{2} p_{2}+m_{3} p_{3}\right)
$$

is an HK basis for the map (6-9) with $\operatorname{dim} K_{\Theta}(m, p)=1$. At every point $(m, p) \in \mathbb{R}^{6}$,

$$
K_{\Theta}(m, p)=\left[g_{1}: g_{2}: g_{3}: g_{4}: g_{5}: g_{6}:-1\right]
$$

Here $g_{1}, g_{2}, g_{3}$ are integrals of the map (6-9) given by

$$
g_{k}=\frac{g_{k}^{(4)}+\epsilon^{2} g_{k}^{(6)}+\epsilon^{4} g_{k}^{(8)}+\epsilon^{6} g_{k}^{(10)}}{2\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right) \Delta}, \quad k=1,2,3
$$

where $g_{k}^{(2 q)}$ are homogeneous polynomials of degree $2 q$ in phase variables, and $\Delta$ is given in (6-23). For instance,

$$
g_{k}^{(4)}=2 H_{4}^{2}-H_{3} H_{1}+\omega_{k} H_{3}^{2}
$$

Integrals $g_{4}, g_{5}, g_{6}$ are given by

$$
\begin{equation*}
g_{4}=\frac{g_{2}-g_{3}}{\omega_{2}-\omega_{3}}, \quad g_{5}=\frac{g_{3}-g_{1}}{\omega_{3}-\omega_{1}}, \quad g_{6}=\frac{g_{1}-g_{2}}{\omega_{1}-\omega_{2}} \tag{6-32}
\end{equation*}
$$

Proof: Since system (6-30) involves too many iterates of $f$ for a symbolic treatment, we look for linear relations among the (numerical) solutions of this system. Application of the PSLQ algorithm allows us to identify three such relations, as given in (6-32). This reduces system $(6-30)$ to the following one:

$$
\begin{align*}
& {\left[g_{1}\left(p_{1}^{2}+\frac{m_{2}^{2}}{\omega_{1}-\omega_{3}}+\frac{m_{3}^{2}}{\omega_{1}-\omega_{2}}\right)\right.} \\
& \quad+g_{2}\left(p_{2}^{2}+\frac{m_{1}^{2}}{\omega_{2}-\omega_{3}}+\frac{m_{3}^{2}}{\omega_{2}-\omega_{1}}\right) \\
& \left.\quad+g_{3}\left(p_{3}^{2}+\frac{m_{1}^{2}}{\omega_{3}-\omega_{2}}+\frac{m_{2}^{2}}{\omega_{3}-\omega_{1}}\right)\right] \circ f^{i}(m, p) \\
& \quad=\left(m_{1} p_{1}+m_{2} p_{2}+m_{3} p_{3}\right) \circ f^{i}(m, p) . \tag{6-33}
\end{align*}
$$

Thus, one can say that we are dealing with a reduced Hirota-Kimura basis consisting of $l=4$ functions

$$
\widetilde{\Theta}=\left(I_{1}, I_{2}, I_{3}, H_{4}\right)
$$

see (6-7). Interestingly, this is a basis of integrals for the continuous-time Clebsch system, but we do not know whether this is just a coincidence or it has some deeper meaning. System (6-33) has to be solved for two different but intersecting ranges of $l-1=3$ consecutive indices $i$. It would be enough to show that the solution for one asymmetric range, e.g., for $i \in[0,2]$, consists of even functions of $\epsilon$. However, this asymmetric system involves of necessity the second iterate $f^{2}$. To avoid dealing with $f^{2}$, one more linear relation for $g_{1}, g_{2}, g_{3}$ would be needed. Such a relation has been found with the help of the PSLQ algorithm; it no longer has constant coefficients but involves the previously found integrals $d_{7}, d_{8}, d_{9}$ :

$$
\begin{align*}
& \left(\omega_{2}-\omega_{3}\right) g_{1}+\left(\omega_{3}-\omega_{1}\right) g_{2}+\left(\omega_{1}-\omega_{2}\right) g_{3}  \tag{6-34}\\
& \quad=\frac{1}{2}\left(\omega_{2}-\omega_{3}\right)\left(\omega_{3}-\omega_{1}\right)\left(d_{8}-d_{7}\right) .
\end{align*}
$$

Of course, due to (6-24), the right-hand side of (6-34) can be equivalently written as
$\frac{1}{2}\left(\omega_{3}-\omega_{1}\right)\left(\omega_{1}-\omega_{2}\right)\left(d_{9}-d_{8}\right)=\frac{1}{2}\left(\omega_{1}-\omega_{2}\right)\left(\omega_{2}-\omega_{3}\right)\left(d_{7}-d_{9}\right)$.

The linear system consisting of $(6-33)$ for $i=0,1$ and (6-34) can be solved by Maple with the result given in the theorem. Since $\left(d_{7}, d_{8}, d_{9}\right)$ have already been proven to be integrals of motion, and since the solutions $\left(g_{1}, g_{2}, g_{3}\right)$ are manifestly even in $\epsilon$, Proposition 6.6 yields that $\left(g_{1}, g_{2}, g_{3}\right)$ are integrals of the map $f$.

Theorem 6.8 gives us the third HK basis with a onedimensional null space for the discrete Clebsch system. Thus, it shows that every orbit lies in the intersection of three quadrics in $\mathbb{R}^{6}$. Concerning the integrals of motion, it turns out that the basis $\Theta$ does not provide us with additional ones: a numerical check with gradients shows that integrals $g_{1}, g_{2}, g_{3}$ are functionally dependent on those previously found. At this point we are lacking one more HK basis with a one-dimensional null space, linearly independent of $K_{\Phi_{0}}, K_{\Psi}, K_{\Theta}$, and one more integral of motion, functionally independent of $J$ and $d_{7}, d_{8}$.

### 6.8 Proof for the Bases $\Phi_{1}, \Phi_{2}, \Phi_{3}$

Now we return to the bases $\Phi_{1}, \Phi_{2}, \Phi_{3}$ discussed in Section 6.5. In order to be able to solve systems (6-18)-(6-20) symbolically and to prove that the solutions $\alpha_{j}, \beta_{j}, \gamma_{j}$ are indeed integrals, we have to find additional linear relations for these quantities (recipe $\mathcal{E}$ ). Within each set of coefficients we were able to identify just one relation:

$$
\begin{align*}
\left(\omega_{1}-\omega_{3}\right) \alpha_{5} & =\left(\omega_{1}-\omega_{2}\right) \alpha_{6}  \tag{6-35}\\
\left(\omega_{2}-\omega_{3}\right) \beta_{4} & =\left(\omega_{2}-\omega_{1}\right) \beta_{6}  \tag{6-36}\\
\left(\omega_{3}-\omega_{2}\right) \gamma_{4} & =\left(\omega_{3}-\omega_{1}\right) \gamma_{5} \tag{6-37}
\end{align*}
$$

This reduces the number of equations in each system by one, which, however, does not resolve our problems. A way out consists in looking for linear relations among all the coefficients $\alpha_{j}, \beta_{j}, \gamma_{j}$. Remarkably, six more independent linear relations of this kind can be identified:

$$
\begin{equation*}
\alpha_{4}=\beta_{5}=\gamma_{6} \tag{6-38}
\end{equation*}
$$

and

$$
\begin{aligned}
\frac{\alpha_{2}-\alpha_{3}-\left(\omega_{2}-\omega_{3}\right) \alpha_{4}}{\omega_{2}-\omega_{3}} & =\frac{\beta_{2}-\beta_{3}-\left(\omega_{2}-\omega_{3}\right) \beta_{4}}{\omega_{3}-\omega_{1}} \\
& =\frac{\gamma_{2}-\gamma_{3}-\left(\omega_{2}-\omega_{3}\right) \gamma_{4}}{\omega_{1}-\omega_{2}}
\end{aligned}
$$

$$
\frac{\alpha_{3}-\alpha_{1}-\left(\omega_{3}-\omega_{1}\right) \alpha_{5}}{\omega_{2}-\omega_{3}}=\frac{\beta_{3}-\beta_{1}-\left(\omega_{3}-\omega_{1}\right) \beta_{5}}{\omega_{3}-\omega_{1}}
$$

$$
\begin{equation*}
=\frac{\gamma_{3}-\gamma_{1}-\left(\omega_{3}-\omega_{1}\right) \gamma_{5}}{\omega_{1}-\omega_{2}} \tag{6-40}
\end{equation*}
$$

There are two more similar relations:

$$
\begin{aligned}
\frac{\alpha_{1}-\alpha_{2}-\left(\omega_{1}-\omega_{2}\right) \alpha_{6}}{\omega_{2}-\omega_{3}} & =\frac{\beta_{1}-\beta_{2}-\left(\omega_{1}-\omega_{2}\right) \beta_{6}}{\omega_{3}-\omega_{1}} \\
& =\frac{\gamma_{1}-\gamma_{2}-\left(\omega_{1}-\omega_{2}\right) \gamma_{6}}{\omega_{1}-\omega_{2}}
\end{aligned}
$$

but these follow from the already listed ones (6-35)-(6-40). We stress that all these linear relations were identified numerically, with the help of the PSLQ algorithm, and remain at this stage hypothetical.

With nine linear relations (6-35)-(6-40), we have to solve systems (6-18)-(6-20) simultaneously for a range of three consecutive indices $i$. Taking this range as $i=-1,0,1$, we can avoid dealing with $f^{2}$, which, however, would leave us with the problem of a proof that the solutions are integrals. Alternatively, we can choose the range $i=0,1,2$, and then the solutions are automatically integrals, as soon as it is established that they are even functions of $\epsilon$.

A symbolic solution of the system consisting of eighteen linear equations, namely $(6-18)-(6-20)$ with $i=$ $0,1,2$, along with nine simple equations (6-35)-(6-40) would require astronomical amounts of memory, because of the complexity of $f^{2}$. However, this task becomes manageable and even simple for fixed (numerical) values of the phase variables $(m, p)$ and of the parameters $\omega_{i}$, while leaving $\epsilon$ a symbolic variable. For rational values of $m_{k}, p_{k}, \omega_{k}$, all computations can be done precisely (in rational arithmetic). This means that $\alpha_{j}, \beta_{j}$, and $\gamma_{j}$ can be evaluated, as functions of $\epsilon$, at arbitrary points in $\mathbb{Q}^{9}\left(m, p, \omega_{1}, \omega_{2}, \omega_{3}\right)$. A large number of such evaluations provides us with convincing evidence in favor of the claim that these functions are even in $\epsilon$.

In order to obtain a rigorous proof without dealing with $f^{2}$, further linear relations would be necessary. Before introducing these, we present some preliminary considerations. Assuming that $\Phi_{1}, \Phi_{2}, \Phi_{3}$ are HK bases with one-dimensional null spaces, results of Theorem 6.5 on the HK basis $\Psi$ tell us that the row vector $\left(d_{7}, d_{8}, d_{9}\right)$ is the unique left null vector for the matrix

$$
M_{2}=\left(\begin{array}{ccc}
\alpha_{4} & \alpha_{5} & \alpha_{6} \\
\beta_{4} & \beta_{5} & \beta_{6} \\
\gamma_{4} & \gamma_{5} & \gamma_{6}
\end{array}\right)
$$

normalized so that

$$
\left(d_{7}, d_{8}, d_{9}\right) M_{1}=(1,1,1)
$$

where

$$
M_{1}=\left(\begin{array}{lll}
\alpha_{1} & \alpha_{2} & \alpha_{3} \\
\beta_{1} & \beta_{2} & \beta_{3} \\
\gamma_{1} & \gamma_{2} & \gamma_{3}
\end{array}\right)
$$

Note that due to (6-35)-(6-38), the matrix $M_{2}$ has at most four (linearly) independent entries. Denoting the common values in these equations by $A, B, C, D$, respectively, we obtain

$$
\begin{align*}
M_{2} & =\left(\begin{array}{ccc}
\alpha_{4} & \alpha_{5} & \alpha_{6} \\
\beta_{4} & \beta_{5} & \beta_{6} \\
\gamma_{4} & \gamma_{5} & \gamma_{6}
\end{array}\right)  \tag{6-41}\\
& =\left(\begin{array}{ccc}
D & A /\left(\omega_{1}-\omega_{3}\right) & A /\left(\omega_{1}-\omega_{2}\right) \\
B /\left(\omega_{2}-\omega_{3}\right) & D & B /\left(\omega_{2}-\omega_{1}\right) \\
C /\left(\omega_{3}-\omega_{2}\right) & C /\left(\omega_{3}-\omega_{1}\right) & D
\end{array}\right) .
\end{align*}
$$

The existence of the left null vector $\left(d_{7}, d_{8}, d_{9}\right)$ shows that $\operatorname{det}\left(M_{2}\right)=0$, or equivalently,

$$
\begin{align*}
& D^{2}-\frac{A B}{\left(\omega_{1}-\omega_{3}\right)\left(\omega_{2}-\omega_{3}\right)}-\frac{B C}{\left(\omega_{2}-\omega_{1}\right)\left(\omega_{3}-\omega_{1}\right)} \\
&-\frac{C A}{\left(\omega_{3}-\omega_{2}\right)\left(\omega_{1}-\omega_{2}\right)} \tag{6-42}
\end{align*}
$$

From (6-41) and (6-42) one easily derives that the row vector

$$
\begin{aligned}
& \left(D-\frac{B}{\omega_{2}-\omega_{3}}-\frac{C}{\omega_{3}-\omega_{2}}, D-\frac{A}{\omega_{1}-\omega_{3}}-\frac{C}{\omega_{3}-\omega_{1}}\right. \\
& \left.D-\frac{A}{\omega_{1}-\omega_{2}}-\frac{B}{\omega_{2}-\omega_{1}}\right) \\
& \quad=\left(\alpha_{4}-\beta_{4}-\gamma_{4},-\alpha_{5}+\beta_{5}-\gamma_{5},-\alpha_{6}-\beta_{6}+\gamma_{6}\right)
\end{aligned}
$$

is a left null vector of the matrix $M_{2}$, and therefore $\left(d_{7}, d_{8}, d_{9}\right)$ is proportional to this vector. The proportionality coefficient can be now determined with the help of the PSLQ algorithm, and it turns out to be extremely simple. Namely, the following relations hold:

$$
\begin{align*}
\alpha_{4}-\beta_{4}-\gamma_{4} & =D-\frac{B-C}{\omega_{2}-\omega_{3}} \tag{6-43}
\end{align*}=\frac{1}{2} d_{7}, ~\left(\alpha_{5}+\beta_{5}-\gamma_{5}=D-\frac{C-A}{\omega_{3}-\omega_{1}}=\frac{1}{2} d_{8}, ~=D-\frac{A-B}{\omega_{1}-\omega_{2}}=\frac{1}{2} d_{9} .\right.
$$

Only two of them are independent, because of (6-24). We note also that according to (6-31), one has

$$
\begin{align*}
& \alpha_{4}+\beta_{4}+\gamma_{4}=D+\frac{B-C}{\omega_{2}-\omega_{3}}=g_{4},  \tag{6-46}\\
& \alpha_{5}+\beta_{5}+\gamma_{5}=D+\frac{C-A}{\omega_{3}-\omega_{1}}=g_{5},  \tag{6-47}\\
& \alpha_{6}+\beta_{6}+\gamma_{6}=D+\frac{A-B}{\omega_{1}-\omega_{2}}=g_{6} . \tag{6-48}
\end{align*}
$$

Equations (6-43)-(6-48) and (6-42) are already enough to determine all four integrals $A, B, C, D$, that is, all
$\alpha_{j}, \beta_{j}, \gamma_{j}$ with $j=4,5,6$, provided it is proven that they are indeed integrals. These (conditional) results read

$$
\begin{align*}
& A=\frac{1+\epsilon^{2} A^{(2)}+\epsilon^{4} A^{(4)}+\epsilon^{6} A^{(6)}+\epsilon^{8} A^{(8)}}{2 \epsilon^{2} \Delta}  \tag{6-49}\\
& B=\frac{1+\epsilon^{2} B^{(2)}+\epsilon^{4} B^{(4)}+\epsilon^{6} B^{(6)}+\epsilon^{8} B^{(8)}}{2 \epsilon^{2} \Delta}  \tag{6-50}\\
& C=\frac{1+\epsilon^{2} C^{(2)}+\epsilon^{4} C^{(4)}+\epsilon^{6} C^{(6)}+\epsilon^{8} C^{(8)}}{2 \epsilon^{2} \Delta}  \tag{6-51}\\
& D=\frac{p_{1}^{2}+p_{2}^{2}+p_{3}^{2}+\epsilon^{2} D^{(4)}+\epsilon^{4} D^{(6)}+\epsilon^{6} D^{(8)}}{2 \Delta} \tag{6-52}
\end{align*}
$$

where $A^{(2 q)}, B^{(2 q)}, C^{(2 q)}, D^{(2 q)}$ are homogeneous polynomials of degree $2 q$ in phase variables, for instance,

$$
\begin{aligned}
A^{(2)}= & B^{(2)}=C^{(2)} \\
= & m_{1}^{2}+m_{2}^{2}+m_{3}^{2}+\left(\omega_{2}+\omega_{3}-2 \omega_{1}\right) p_{1}^{2} \\
& +\left(\omega_{3}+\omega_{1}-2 \omega_{2}\right) p_{2}^{2}+\left(\omega_{1}+\omega_{2}-2 \omega_{3}\right) p_{3}^{2} \\
D^{(4)}= & \left(m_{1} p_{1}+m_{2} p_{2}+m_{3} p_{3}\right)^{2} \\
& +\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right) \\
& \quad \times\left(\left(\omega_{2}+\omega_{3}-2 \omega_{1}\right) p_{1}^{2}+\left(\omega_{3}+\omega_{1}-2 \omega_{2}\right) p_{2}^{2}\right. \\
& \left.\quad+\left(\omega_{1}+\omega_{2}-2 \omega_{3}\right) p_{3}^{2}\right)
\end{aligned}
$$

We remark that ( $6-42$ ) tells us that no more than three of the functions $A, B, C, D$ are actually functionally independent. Computation with gradients shows that $A, B, C$ are functionally indeed independent. Moreover, all other previously found integrals $J, d_{7}, d_{8}, d_{9}$, and $g_{1}, g_{2}, g_{3}$ are functionally dependent on these.

Theorem 6.9. The sets $(6-12)-(6-14)$ are HK bases for the map (6-9) with $\operatorname{dim} K_{\Phi_{1}}(m, p)=\operatorname{dim} K_{\Phi_{2}}(m, p)=$ $\operatorname{dim} K_{\Phi_{3}}(m, p)=1$. At each point $(m, p) \in \mathbb{R}^{6}$,

$$
\begin{aligned}
& K_{\Phi_{1}}(m, p)=\left[\alpha_{1}: \alpha_{2}: \alpha_{3}: \alpha_{4}: \alpha_{5}: \alpha_{6}:-1\right] \\
& K_{\Phi_{2}}(m, p)=\left[\beta_{1}: \beta_{2}: \beta_{3}: \beta_{4}: \beta_{5}: \beta_{6}:-1\right] \\
& K_{\Phi_{3}}(m, p)=\left[\gamma_{1}: \gamma_{2}: \gamma_{3}: \gamma_{4}: \gamma_{5}: \gamma_{6}:-1\right],
\end{aligned}
$$

where $\alpha_{j}, \beta_{j}$, and $\gamma_{j}$ are rational functions of $(m, p)$, even with respect to $\epsilon$. They are integrals of motion for the map (6-9) and satisfy linear relations (6-35)-(6-40). For $j=4,5,6$ they are given by $(6-41),(6-51),(6-52)$. For $j=1,2,3$ they are of the form

$$
\begin{equation*}
h=\frac{h^{(2)}+\epsilon^{2} h^{(4)}+\epsilon^{4} h^{(6)}+\epsilon^{6} h^{(8)}+\epsilon^{8} h^{(10)}+\epsilon^{10} h^{(12)}}{2 \epsilon^{2}\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right) \Delta} \tag{6-53}
\end{equation*}
$$

where $h$ stands for any of the functions $\alpha_{j}, \beta_{j}, \gamma_{j}, j=$ $1,2,3$, and the corresponding $h^{(2 q)}$ are homogeneous polynomials in phase variables of degree $2 q$. For instance,

$$
\begin{gather*}
\alpha_{1}^{(2)}=H_{3}-I_{1}, \quad \alpha_{2}^{(2)}=-I_{1}, \quad \alpha_{3}^{(2)}=-I_{1} \\
\beta_{1}^{(2)}=-I_{2}, \quad \beta_{2}^{(2)}=H_{3}-I_{2}, \quad \beta_{3}^{(2)}=-I_{2}  \tag{6-54}\\
\gamma_{1}^{(2)}=-I_{3}, \quad \gamma_{2}^{(2)}=-I_{3}, \quad \gamma_{3}^{(2)}=H_{3}-I_{3}
\end{gather*}
$$

The four functions $J, \alpha_{1}, \beta_{1}$, and $\gamma_{1}$ are functionally independent.

Proof: The proof consists of several steps.
Step 1. Consider the system for the 18 unknowns $\alpha_{j}, \beta_{j}, \gamma_{j}, j=1, \ldots, 6$, consisting of 17 linear equations: $(6-18)-(6-20)$ with $i=0,1,(6-35)-(6-40)$, and $(6-43)$, (6-44). This system is underdetermined, so that in principle it admits a one-parameter family of solutions.

Remarkably, the symbolic Maple solution shows that all variables $\alpha_{j}, \beta_{j}, \gamma_{j}$ with $j=4,5,6$ are determined by this system uniquely, the results coinciding with (6-41), (6-51)-(6-52). (Actually, the Maple answers are much more complicated, and their simplification was carried out with Singular, which was used to cancel out common factors from the huge expressions in numerators and denominators of these rational functions.) Since these uniquely determined $\alpha_{j}, \beta_{j}, \gamma_{j}$ with $j=4,5,6$ are even functions of $\epsilon$, this proves that they (i.e., $A, B, C, D$ ) are integrals of motion.

Step 2. Having determined $\alpha_{j}, \beta_{j}, \gamma_{j}$ with $j=4,5,6$, we are in a position to compute $\alpha_{j}, \beta_{j}, \gamma_{j}$ with $j=1,2,3$. For instance, to obtain the values of $\alpha_{j}$ with $j=1,2,3$, we consider the symmetric linear system (6-18) with $i=-1,0,1$ (and with already found $\alpha_{4}, \alpha_{5}, \alpha_{6}$ ). This system was solved by Maple. The solutions are huge rational functions, which, however, turn out to admit massive cancellations. These cancellations were carried out with the help of Singular. The resulting expressions for $\alpha_{1}, \alpha_{2}, \alpha_{3}$ turn out to satisfy the ansatz (6-53) with the leading terms given in the first line of $(6-54) .^{3}$

However, this computation does not prove that the functions so obtained are indeed integrals of motion. To prove this, one could, in principle, either check directly the identities $\alpha_{j} \circ f=\alpha_{j}, j=1,2,3$, or verify equation (6-18) with $i=2$. Both approaches are prohibitively expensive, so that we have to look for an alternative one.

[^2]Step 3. The results of Step 2 yield an explicit expression for the function

$$
F=\left(\omega_{2}-\omega_{3}\right) \alpha_{1}+\left(\omega_{3}-\omega_{1}\right) \alpha_{2}+\left(\omega_{1}-\omega_{2}\right) \alpha_{3}, \quad(6-55)
$$

which is of the form

$$
F=\frac{\left(\omega_{2}-\omega_{3}\right)\left(1+\epsilon^{2} F^{(2)}+\epsilon^{4} F^{(4)}+\epsilon^{6} F^{(6)}+\epsilon^{8} F^{(8)}\right)}{2 \epsilon^{2} \Delta}
$$

It is of a crucial importance for our purposes that it can be proven directly that $F$ is an integral of motion. We have proved this with the method $\mathcal{G}$ based on the Gröbner basis for the ideal generated by discrete equations of motion. The application of this method to $F$ is more feasible than to any of the $\alpha_{j}, j=1,2,3$, because of the cancelation of the huge polynomial coefficient of $\epsilon^{10}$ in the numerator of $F$.

In fact, more is true: $F$ is not only an integral, but is functionally dependent on the previously found ones, say on $J, d_{7}, d_{8}$. For a proof of this claim, it would be most favorable to find the explicit dependence $F=F\left(J, d_{7}, d_{8}\right)$, but it remains unknown to us. Instead, we have chosen the way of verification that

$$
\nabla F \in \operatorname{span}\left(\nabla J, \nabla d_{7}, \nabla d_{8}\right)
$$

This is easily checked numerically for arbitrarily many (rational) values of the data involved. For a symbolic check, one has to prove the existence of three scalar functions $\lambda_{1}, \lambda_{2}, \lambda_{3}$ such that

$$
\nabla F=\lambda_{1} \nabla J+\lambda_{2} \nabla d_{7}+\lambda_{3} \nabla d_{8}
$$

This is a system of six equations in three unknowns. Since $J$ does not depend on $m_{k}$, one can determine $\lambda_{2}$, $\lambda_{3}$ from a system of only three equations:

$$
\nabla_{m} F=\lambda_{2} \nabla_{m} d_{7}+\lambda_{3} \nabla_{m} d_{8}
$$

After that, it remains to check that $\nabla_{p} F-\lambda_{2} \nabla_{p} d_{7}-$ $\lambda_{3} \nabla_{p} d_{8}$ is proportional to $\nabla_{p} J$. Clearly, these computations can be arranged so as to verify the vanishing of certain (very big) polynomials. We have been able to perform these computations with the help of Singular for symbolic $m_{k}, p_{k}$ but with (several sets of) numeric values of coefficients $\omega_{k}$ only.
Step 4. The result of Step 3 allows us to proceed as follows. Consider the system of three linear equations for $\alpha_{1}, \alpha_{2}, \alpha_{3}$, consisting of (6-18) with $i=0,1$ and

$$
\left(\omega_{2}-\omega_{3}\right) \alpha_{1}+\left(\omega_{3}-\omega_{1}\right) \alpha_{2}+\left(\omega_{1}-\omega_{2}\right) \alpha_{3}=F
$$

where $F$ is the explicit expression obtained and proven to be an integral in Step 3. This system can now be solved
by Maple; the results, again simplified with singular, are even functions of $\epsilon$ (indeed, the same ones obtained in Step 1 from the symmetric system). Noneven polynomials in $\epsilon$ of degree 7 cancel in a miraculous way from the numerators and the denominator. Now Proposition 6.6 ensures that these solutions are integrals of motion.

Step 5. Finally, in order to find $\beta_{1}, \beta_{2}, \beta_{3}$ and $\gamma_{1}, \gamma_{2}, \gamma_{3}$, we solve the two systems consisting of $(6-19)$ and $(6-20)$ with $i=0,1$, and respectively the first and the second linear relations in (6-39). The results are even functions of $\epsilon$ satisfying the ansatz ( $6-53$ ) with the leading terms given in (6-54). Proposition 6.6 yields that these functions are also integrals of motion. ${ }^{4}$

### 6.9 Preliminary Results on the Hirota-Kimura-Type Discretization of the General Flow of the Clebsch System

The general flow of the Clebsch system, depending on three real parameters $b_{1}, b_{2}, b_{3}$ (or rather on their differences $b_{i}-b_{j}$, which gives two independent real parameters), reads as follows:

$$
\left\{\begin{array}{l}
\dot{m}=m \times C m+p \times B p  \tag{6-56}\\
\dot{p}=p \times C m
\end{array}\right.
$$

where $B=\operatorname{diag}\left(b_{1}, b_{2}, b_{3}\right)$ and $C=\operatorname{diag}\left(c_{1}, c_{2}, c_{3}\right)$ with

$$
\begin{equation*}
c_{1}=\frac{b_{2}-b_{3}}{\omega_{2}-\omega_{3}}, \quad c_{2}=\frac{b_{3}-b_{1}}{\omega_{3}-\omega_{1}}, \quad c_{3}=\frac{b_{1}-b_{2}}{\omega_{1}-\omega_{2}} \tag{6-57}
\end{equation*}
$$

This flow is Hamiltonian with the quadratic Hamilton function

$$
\begin{aligned}
H & =\frac{1}{2}\langle m, C m\rangle+\frac{1}{2}\langle p, B p\rangle=\frac{1}{2} \sum_{k=1}^{3}\left(c_{k} m_{k}^{2}+b_{k} p_{k}^{2}\right) \\
& =\frac{1}{2}\left(b_{1} I_{1}+b_{2} I_{2}+b_{3} I_{3}\right)
\end{aligned}
$$

In components, system (6-56) reads

$$
\begin{aligned}
\dot{m}_{1} & =\left(c_{3}-c_{2}\right) m_{2} m_{3}+\left(b_{3}-b_{2}\right) p_{2} p_{3}, \\
\dot{m}_{2} & =\left(c_{1}-c_{3}\right) m_{3} m_{1}+\left(b_{1}-b_{3}\right) p_{3} p_{1}, \\
\dot{m}_{3} & =\left(c_{2}-c_{1}\right) m_{1} m_{2}+\left(b_{2}-b_{1}\right) p_{1} p_{2}, \\
\dot{p}_{1} & =c_{3} m_{3} p_{2}-c_{2} m_{2} p_{3}, \\
\dot{p}_{2} & =c_{1} m_{1} p_{3}-c_{3} m_{3} p_{1}, \\
\dot{p}_{3} & =c_{2} m_{2} p_{1}-c_{1} m_{1} p_{2} .
\end{aligned}
$$

[^3]The KH discretization of the flow $(6-56)$ reads
$\left\{\begin{array}{l}\widetilde{m}-m=\epsilon(\widetilde{m} \times C m+m \times C \widetilde{m}+\widetilde{p} \times B p+p \times B \widetilde{p}), \\ \widetilde{p}-p=\epsilon(\widetilde{p} \times C m+p \times C \widetilde{m}) .\end{array}\right.$
In components,

$$
\begin{align*}
\widetilde{m}_{1}-m_{1}= & \epsilon\left(c_{3}-c_{2}\right)\left(\widetilde{m}_{2} m_{3}+m_{2} \widetilde{m}_{3}\right)  \tag{6-58}\\
& +\epsilon\left(b_{3}-b_{2}\right)\left(\widetilde{p}_{2} p_{3}+p_{2} \widetilde{p}_{3}\right), \\
\widetilde{m}_{2}-m_{2}= & \epsilon\left(c_{1}-c_{3}\right)\left(\widetilde{m}_{3} m_{1}+m_{3} \widetilde{m}_{1}\right) \\
& +\epsilon\left(b_{1}-b_{3}\right)\left(\widetilde{p}_{3} p_{1}+p_{3} \widetilde{p}_{1}\right), \\
\widetilde{m}_{3}-m_{3}= & \epsilon\left(c_{2}-c_{1}\right)\left(\widetilde{m}_{1} m_{2}+m_{1} \widetilde{m}_{2}\right) \\
& +\epsilon\left(b_{2}-b_{1}\right)\left(\widetilde{p}_{1} p_{2}+p_{1} \widetilde{p}_{2}\right), \\
\widetilde{p}_{1}-p_{1}= & \epsilon c_{3}\left(\widetilde{m}_{3} p_{2}+m_{3} \widetilde{p}_{2}\right)-\epsilon c_{2}\left(\widetilde{m}_{2} p_{3}+m_{2} \widetilde{p}_{3}\right), \\
\widetilde{p}_{2}-p_{2}= & \epsilon c_{1}\left(\widetilde{m}_{1} p_{3}+m_{1} \widetilde{p}_{3}\right)-\epsilon c_{3}\left(\widetilde{m}_{3} p_{1}+m_{3} \widetilde{p}_{1}\right), \\
\widetilde{p}_{3}-p_{3}= & \epsilon c_{2}\left(\widetilde{m}_{2} p_{1}+m_{2} \widetilde{p}_{1}\right)-\epsilon c_{1}\left(\widetilde{m}_{1} p_{2}+m_{1} \widetilde{p}_{2}\right) .
\end{align*}
$$

In what follows, we will use the abbreviations $b_{i j}=$ $b_{i}-b_{j}$ and $c_{i j}=c_{i}-c_{j}$. The linear system (6-58) defines an explicit birational map $f: \mathbb{R}^{6} \rightarrow \mathbb{R}^{6}$,

$$
\begin{equation*}
\binom{\widetilde{m}}{\widetilde{p}}=f(m, p, \epsilon)=M^{-1}(m, p, \epsilon)\binom{m}{p}, \tag{6-59}
\end{equation*}
$$

where

$$
\begin{aligned}
& M(m, p, \epsilon)= \\
& \left(\begin{array}{cccccc}
1 & \epsilon c_{23} m_{3} & \epsilon c_{23} m_{2} & 0 & \epsilon b_{23} p_{3} & \epsilon b_{23} p_{2} \\
\epsilon c_{31} m_{3} & 1 & \epsilon c_{31} m_{1} & \epsilon b_{31} p_{3} & 0 & \epsilon b_{31} p_{1} \\
\epsilon c_{12} m_{2} & \epsilon c_{12} m_{1} & 1 & \epsilon b_{12} p_{2} & \epsilon b_{12} p_{1} & 0 \\
0 & \epsilon c_{2} p_{3} & -\epsilon c_{3} p_{2} & 1 & -\epsilon c_{3} m_{3} & \epsilon c_{2} m_{2} \\
-\epsilon c_{1} p_{3} & 0 & \epsilon c_{3} p_{1} & \epsilon c_{3} m_{3} & 1 & -\epsilon c_{1} m_{1} \\
\epsilon c_{1} p_{2} & -\epsilon c_{2} p_{1} & 0 & -\epsilon c_{2} m_{2} & \epsilon c_{1} m_{1} & 1
\end{array}\right) .
\end{aligned}
$$

As usual, the map (6-59) possesses the reversibility property

$$
f^{-1}(m, p, \epsilon)=f(m, p,-\epsilon)
$$

Conjecture 6.10. All claims of Theorems 6.1, 6.2 hold also for the discretization (6-59) of the general flow of the Clebsch system, with the HK basis $\Phi_{0}$ being replaced by

$$
\begin{equation*}
\Phi_{0}=\left(p_{1}^{2}, p_{2}^{2}, p_{3}^{2}, m_{1}^{2}, m_{2}^{2}, m_{3}^{2}, 1\right) \tag{6-60}
\end{equation*}
$$

This conjecture is supported by numerical results based on the algorithm $\mathcal{N}$. The claim concerning the set $\Phi_{0}$ given in (6-60) is proven symbolically. In order to keep the notation compact, we give here this proof for the second flow of the Clebsch system only.

Recall that the first flow of the Clebsch system, considered in Section 6, corresponds to $b_{i}=\omega_{i}$ and $c_{i}=1$. The
second flow is characterized by the Hamilton function

$$
\begin{aligned}
H= & \frac{1}{2} H_{2} \\
= & \frac{1}{2}\left(\omega_{1} m_{1}^{2}+\omega_{2} m_{2}^{2}+\omega_{3} m_{3}^{2}-\omega_{2} \omega_{3} p_{1}^{2}-\omega_{1} \omega_{3} p_{2}^{2}\right. \\
& \left.\quad-\omega_{1} \omega_{2} p_{3}^{2}\right) .
\end{aligned}
$$

In other words, the choice of parameters $b_{k}$ characterizing the second flow is

$$
\begin{equation*}
b_{1}=-\omega_{2} \omega_{3}, \quad b_{2}=-\omega_{3} \omega_{1}, \quad b_{3}=-\omega_{1} \omega_{2} \tag{6-61}
\end{equation*}
$$

so that

$$
\begin{equation*}
c_{1}=\omega_{1}, \quad c_{2}=\omega_{2}, \quad c_{3}=\omega_{3} \tag{6-62}
\end{equation*}
$$

For the HK discretization of the second Clebsch flow, we give a more concrete formulation of our findings concerning the HK basis $\Phi_{0}$, including a "nice" integral.

Theorem 6.11. For the map (6-59) the set of functions $(6-60)$ is an HK basis with $\operatorname{dim} K_{\Phi_{0}}(m, p)=1$. At each point $(m, p) \in \mathbb{R}^{6}$, we have

$$
K_{\Phi_{0}}(m, p)=\left[e_{1}: e_{2}: e_{3}: e_{4}: e_{5}: e_{6}:-1\right]
$$

where all $e_{i}$ are fractional linear functions of a single integral $L=L(m, p, \epsilon)$ of the map (6-59), which is a quotient of two quadratic polynomials in $m_{k}, p_{k}$.

If the coefficients $b_{k}, c_{k}$ are as in (6-61)), (6-62), then the integral $L$ can be taken as

$$
L=\frac{E_{1} G_{1}+E_{2} G_{2}+E_{3} G_{3}}{1+\epsilon^{2} \omega_{1} G_{1}+\epsilon^{2} \omega_{2} G_{2}+\epsilon^{2} \omega_{3} G_{3}}
$$

with

$$
\begin{align*}
& E_{1}=\omega_{3} \omega_{1}+\omega_{1} \omega_{2}-\omega_{2} \omega_{3}, \\
& E_{2}=\omega_{1} \omega_{2}+\omega_{2} \omega_{3}-\omega_{3} \omega_{1},  \tag{6-63}\\
& E_{3}=\omega_{2} \omega_{3}+\omega_{3} \omega_{1}-\omega_{1} \omega_{2}
\end{align*}
$$

and

$$
\begin{aligned}
G_{1} & =\omega_{1} m_{1}^{2}+\omega_{2} \omega_{3} p_{1}^{2} \\
G_{2} & =\omega_{2} m_{2}^{2}+\omega_{3} \omega_{1} p_{2}^{2} \\
G_{3} & =\omega_{3} m_{3}^{2}+\omega_{1} \omega_{2} p_{3}^{2}
\end{aligned}
$$

In this case,

$$
\begin{aligned}
\frac{e_{1}}{\omega_{2} \omega_{3}} & =\frac{e_{4}}{\omega_{1}}=\frac{E_{1}}{L}-\epsilon^{2} \omega_{1} \\
& =\frac{E_{1}+\epsilon^{2}\left(\omega_{1}-\omega_{2}\right) E_{3} G_{2}+\epsilon^{2}\left(\omega_{1}-\omega_{3}\right) E_{2} G_{3}}{E_{1} G_{1}+E_{2} G_{2}+E_{3} G_{3}}, \\
\frac{e_{2}}{\omega_{3} \omega_{1}} & =\frac{e_{5}}{\omega_{2}}=\frac{E_{2}}{L}-\epsilon^{2} \omega_{2} \\
& =\frac{E_{2}+\epsilon^{2}\left(\omega_{2}-\omega_{3}\right) E_{1} G_{3}+\epsilon^{2}\left(\omega_{2}-\omega_{1}\right) E_{3} G_{1}}{E_{1} G_{1}+E_{2} G_{2}+E_{3} G_{3}}, \\
\frac{e_{3}}{\omega_{1} \omega_{2}} & =\frac{e_{6}}{\omega_{3}}=\frac{E_{3}}{L}-\epsilon^{2} \omega_{3} \\
& =\frac{E_{3}+\epsilon^{2}\left(\omega_{3}-\omega_{1}\right) E_{2} G_{1}+\epsilon^{2}\left(\omega_{3}-\omega_{2}\right) E_{1} G_{2}}{E_{1} G_{1}+E_{2} G_{2}+E_{3} G_{3}}
\end{aligned}
$$

The numerator of $L$,

$$
\begin{aligned}
L(m, p, 0)= & E_{1}\left(\omega_{1} m_{1}^{2}+\omega_{2} \omega_{3} p_{1}^{2}\right)+E_{2}\left(\omega_{2} m_{2}^{2}+\omega_{3} \omega_{1} p_{2}^{2}\right) \\
& +E_{3}\left(\omega_{3} m_{3}^{2}+\omega_{1} \omega_{2} p_{3}^{2}\right)
\end{aligned}
$$

is a linear combination of quadratic integrals of motion of the continuous Clebsch system.

Proof: We will present the proof for only the second Clebsch flow. The claim of the theorem refers to the linear system
$\left(e_{1} p_{1}^{2}+e_{2} p_{2}^{2}+e_{3} p_{3}^{2}+e_{4} m_{1}^{2}+e_{5} m_{2}^{2}+e_{6} m_{3}^{2}\right) \circ f^{i}(m, p)=1$
for $i$ from the ranges containing six consecutive numbers, such as $i \in[-2,3]$ or $i \in[-3,2]$. Since the solution of such a system clearly requires more iterates of the map $f$ than could be handled symbolically, we follow recipe $\mathcal{E}$ and look for linear relations between the $e_{i}$. It turns out to be possible to identify the following five relations:

$$
\begin{align*}
\omega_{1} e_{1}-\omega_{2} \omega_{3} e_{4} & =0  \tag{6-64}\\
\omega_{2} e_{2}-\omega_{3} \omega_{1} e_{5} & =0  \tag{6-65}\\
\omega_{3} e_{3}-\omega_{1} \omega_{2} e_{6} & =0  \tag{6-66}\\
e_{1}-e_{2}-\left(\omega_{1}-\omega_{2}\right) e_{6} & =\epsilon^{2} \omega_{3}^{2}\left(\omega_{1}-\omega_{2}\right),  \tag{6-67}\\
e_{3}-e_{1}-\left(\omega_{3}-\omega_{1}\right) e_{5} & =\epsilon^{2} \omega_{2}^{2}\left(\omega_{3}-\omega_{1}\right) . \tag{6-68}
\end{align*}
$$

Of course, there is also a third inhomogeneous relation:

$$
e_{2}-e_{3}-\left(\omega_{2}-\omega_{3}\right) e_{4}=\epsilon^{2} \omega_{1}^{2}\left(\omega_{2}-\omega_{3}\right)
$$

but it is actually a consequence of the previous five. As usual, these (at this point conjectural) identities can be (and have been) found using the PSLQ algorithm. Now we obtain the six functions $e_{i}$ by solving a simple system of six linear equations that involves no iterates of the map $f$ at all and consists of

$$
e_{1} p_{1}^{2}+e_{2} p_{2}^{2}+e_{3} p_{3}^{2}+e_{4} m_{1}^{2}+e_{5} m_{2}^{2}+e_{6} m_{3}^{2}=1
$$

along with the relations $(6-64)-(6-68)$. The solution is given in the formulation of the theorem. To prove that the function $L$ is an integral of motion, one can use a straightforward computation using Maple. Also, a proof based on the equations of motion alone can be given, similar to the proof for $L$ (see the proof of Theorem 6.3). The last claim of the theorem about $L(x, 0)$ follows in the limit $\epsilon \rightarrow 0$, but can be also easily checked directly, by verifying conditions (6-57) for $b_{i}=\omega_{j} \omega_{k} E_{j}$ and $c_{i}=$ $\omega_{i} E_{i}$ with $E_{i}$ from (6-63). These conditions are satisfied due to the identities

$$
\omega_{j} E_{i}-\omega_{i} E_{j}=\left(\omega_{i}-\omega_{j}\right) E_{k}
$$

where $(i, j, k)$ is any permutation of $(1,2,3)$.

## 7. CONCLUSIONS

We have established the integrability of the Hirota-Kimura-type discretization of the Clebsch system, in the following senses:

- the existence for every initial point $(m, p) \in \mathbb{R}^{6}$ of a four-dimensional pencil of quadrics containing the orbit of this point; in our terminology, this can be formulated as existence of an HK basis with a fourdimensional null space consisting of quadratic monomials;
- the existence of four functionally independent integrals of motion (conserved quantities).

Numerical experiments show that this remains true also for an arbitrary flow of the Clebsch system. It is interesting to note that the maps generated by HirotaKimura discretizations of various flows do not commute with each other. It would be important to understand whether some analogue of commutativity of the continuous flows survives in the discrete situation.

Our investigations were based mainly on computer experiments. Our proofs are computer assisted and were obtained with the help of symbolic calculations with Maple, Singular, and Form. A general structure behind these facts that would provide us with more-systematic and less-computational proofs and with more insight remains unknown. In particular, nothing like a Lax representation has been found. Nothing is known about the existence of an invariant Poisson structure for these maps. (For a simpler system, Hirota-Kimura discretization of the Euler top, an invariant volume measure as well as a bi-Hamiltonian structure have been found in [Petrera and Suris 07].)

Hirota and Kimura demonstrated that their discretization leads to an integrable map also for the Lagrange top [Kimura and Hirota 00]. Our preliminary investigations have shown remarkable features pointing toward the integrability of the Hirota-Kimura discretizations of the following systems: Zhukovsky-Volterra gyrostat; so(4) Euler top and its commuting flows; Volterra and Toda lattices; classical Gaudin magnet. Based on these observations, we formulate the following conjecture.

Conjecture 7.1. For any algebraically completely integrable system with a quadratic vector field, its HirotaKimura discretization remains algebraically completely integrable.

If true, this statement could be related to addition theorems for multidimensional theta functions. Such a relation has already been established for the Hirota-Kimura discretization of the Euler top, which can be solved explicitly in elliptic functions [Suris 08]. In our ongoing investigations, we hope to establish integrability of the above-mentioned discrete-time systems and to uncover general mechanisms behind it.

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[^0]:    ${ }^{1}$ Singular is available online at http://www.singular.uni-kl.de/; Form can be found at http://www.nikhef.nl/~form/.

[^1]:    ${ }^{2}$ These (long) expressions can be found at http://www-m8.ma. tum.de/personen/suris/Worksheets.zip.

[^2]:    ${ }^{3}$ All further terms can be found at http://www-m8.ma.tum.de/ personen/suris/Worksheets.zip.

[^3]:    ${ }^{4}$ Maple worksheets for all computations used in this section can be found at http://www-m8.ma.tum.de/personen/suris/ Worksheets.zip.

