EXPONENTIALLY SMALL ASYMPTOTIC ESTIMATES FOR
THE SPLITTING OF SEPARATRICES TO WHISKERED TORI
WITH QUADRATIC AND CUBIC FREQUENCIES

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ABSTRACT. We study the splitting of invariant manifolds of whiskered tori
with two or three frequencies in nearly-integrable Hamiltonian systems, such
that the hyperbolic part is given by a pendulum. We consider a 2-dimensional
torus with a frequency vector \( \omega = (1, \Omega) \), where \( \Omega \) is a quadratic irrational
number, or a 3-dimensional torus with a frequency vector \( \omega = (1, \Omega, \Omega^2) \), where
\( \Omega \) is a cubic irrational number. Applying the Poincaré–Melnikov method,
we find exponentially small asymptotic estimates for the maximal splitting
distance between the stable and unstable manifolds associated to the invariant
torus, and we show that such estimates depend strongly on the arithmetic
properties of the frequencies. In the quadratic case, we use the continued
fractions theory to establish a certain arithmetic property, fulfilled in 24 cases,
which allows us to provide asymptotic estimates in a simple way. In the cubic
case, we focus our attention to the case in which \( \Omega \) is the so-called cubic golden
number (the real root of \( x^3 + x - 1 = 0 \)), obtaining also asymptotic estimates.
We point out the similitudes and differences between the results obtained for
both the quadratic and cubic cases.

1. INTRODUCTION

1.1. Background and objectives. The aim of this paper is to introduce a methodology for measuring the exponentially small splitting of separatrices in a perturbed Hamiltonian system, associated to an \( \ell \)-dimensional whiskered torus (invariant hyperbolic torus) with an algebraic frequency vector, quadratic in the case \( \ell = 2 \), and cubic in the case \( \ell = 3 \).

As the unperturbed system, we consider an integrable Hamiltonian \( H_0 \) with \( \ell + 1 \) degrees of freedom having \( \ell \)-dimensional whiskered tori with coincident stable and unstable whiskers. In general, for a perturbed Hamiltonian \( H = H_0 + \mu H_1 \), where \( \mu \neq 0 \) is small, the whiskers do not coincide anymore, giving rise to the phenomenon called splitting of separatrices, discovered by Poincaré [46]. In order
to give a measure for the splitting, one often describes it by a periodic vector function $M(\theta)$, $\theta \in \mathbb{T}^\ell$, usually called splitting function, giving the distance between the invariant manifolds in the complementary directions, on a transverse section. The most popular tool to measure the splitting is the Poincaré–Melnikov method, introduced in [46] and rediscovered later by Melnikov and Arnold [40, 1]. This method provides a first order approximation

$$M(\theta) = \mu M(\theta) + O(\mu^2), \quad (1)$$

where $M(\theta)$ is called the Melnikov function and is defined by an integral. In fact, it was established [17, 8] that both the splitting and the Melnikov functions are the gradients of scalar functions: the splitting potential and the Melnikov potential, denoted $L(\theta)$ and $L(\theta)$, respectively. This result implies the existence of homoclinic orbits (i.e., intersections between the stable and unstable whiskers) in the perturbed system.

We focus our attention on a concrete torus with an $\ell$-dimensional frequency vector of fast frequencies:

$$\omega_\varepsilon = \frac{\omega}{\sqrt{\varepsilon}}, \quad \varepsilon > 0, \quad (2)$$

with a relation between the parameters, of the form $\mu = \varepsilon^p$ for some $p > 0$. Thus, we have a singular perturbation problem, and the interest for this situation lies in its relation to the normal form of a nearly-integrable Hamiltonian, with $\varepsilon$ as the perturbation parameter, in the vicinity of a simple resonance [44, 9]. In such a singular problem, one can give upper bounds for the splitting, showing that it is exponentially small with respect to $\varepsilon$. The first of such upper bounds was obtained by Neishtadt [43] in one and a half degrees of freedom i.e., for 1 frequency (see also [29, 18, 19, 20]), and later this was extended to the case of 2 or more frequencies (see for instance [50, 21, 4, 3, 12, 13]).

The problem of establishing asymptotic estimates for the exponentially small splitting is much more difficult, but some results have also been obtained by several methods, for explicit perturbations. The difficulty lies in the fact that the Melnikov function is exponentially small in $\varepsilon$ and the error of the method could overcome the main term in (1). Then, an additional study is required in order to validate the Poincaré–Melnikov method. In the case of 1 frequency, the first result providing an asymptotic estimate for the exponentially small splitting was obtained by Lazutin [34] in 1984, for the Chirikov standard map, using complex parameterizations of the invariant manifolds. The same technique was used to justify the Poincaré–Melnikov method in a Hamiltonian with one and a half degrees of freedom [15, 16, 22] or an area-preserving map [14].

For 2 or more frequencies, it turns out that small divisors appear in the splitting function and, as first noticed by Lochak [37], the arithmetic properties of the frequency vector $\omega$ play an important role. This was established by Simó [50], and rigorously proved in [12] for the quasi-periodically forced pendulum, assuming a trigonometric polynomial perturbation in the coordinates associated to the pendulum. Recently, a more general (meromorphic) perturbation was considered in [26]. It is worth remarking that, in some cases, the Poincaré–Melnikov method does not predict correctly the size of the splitting, as shown for instance in [5]. Indeed, when the Poincaré–Melnikov approach cannot be validated, other techniques can be applied to get exponentially small estimates, such as complex matching [2, 45, 41, 42], “beyond all orders” asymptotic methods [38], or continuous averaging [52, 47].
A different approach, followed by Lochak, Marco and Sauzin [49, 36], and Rudnev and Wiggins [48], is the parametrization of the whiskers as solutions of Hamilton–Jacobi equation, to obtain exponentially small estimates of the splitting, and the existence of transverse homoclinic orbits for some intervals of the perturbation parameter \( \varepsilon \). Following a similar approach, it was shown in [11] the continuation of the exponentially small estimates and the transversality of the splitting for all sufficiently small values of \( \varepsilon \) hold under a quite general condition on the phases of the perturbation. Otherwise, homoclinic bifurcations can occur, studied by Simó and Valls [51] in Arnold’s example. The quoted papers considered the case of 2 frequencies, and assumed in most cases that the frequency ratio is the famous golden mean \( \Omega_1 = (\sqrt{5} - 1)/2 \). A generalization to some other quadratic frequency ratios was studied in [10].

We also stress that, for some purposes, it is not necessary to establish the transversality of the splitting, and can be enough to provide asymptotic estimates of the maximal splitting distance. Indeed, such estimates imply the existence of splitting between the invariant manifolds, which provides a strong indication of the non-integrability of the system near the given torus, and opens the door to the application of topological methods [24, 25] for the study of Arnold diffusion in such systems. For a more complete background and references concerning exponentially small splitting, the reader is referred to [36, 13].

The main objective of this paper is to develop a unified methodology in order to generalize the results on exponentially small splitting to frequency vectors \( \omega \) in \( \mathbb{R}^2 \) or \( \mathbb{R}^3 \), in order to obtain asymptotic estimates for the maximal splitting distance, and to emphasize their dependence on the arithmetic properties of the frequencies. Namely, we consider two possibilities:

- **quadratic frequencies**: \( \ell = 2 \) and \( \omega = (1, \Omega) \), where \( \Omega \) is a quadratic irrational number;
- **cubic frequencies**: \( \ell = 3 \) and \( \omega = (1, \Omega, \Omega^2) \), where \( \Omega \) is a cubic irrational number whose two conjugates are not real.

Such frequency vectors satisfy a Diophantine condition,

\[
|\langle k, \omega \rangle| \geq \frac{\gamma}{|k|^\ell - 1}, \quad \forall k \in \mathbb{Z}^\ell \setminus \{0\}
\]  

with some \( \gamma > 0 \), in both the quadratic and the cubic cases. We point out that \( \ell - 1 \) is the minimal possible exponent for Diophantine inequalities in \( \mathbb{R}^\ell \) (see for instance [35, ap. 4]).

One of the goals of this paper is to show, for the above frequencies, that we can detect the integer vectors \( k \in \mathbb{Z}^\ell \setminus \{0\} \) providing an approximate equality in (3), i.e., giving the “least” small divisors (relatively to the size of \( |k| \)). We call such vectors \( k \) the primary resonances of \( \omega \), and other vectors the secondary resonances. We show that, if a certain arithmetic condition is fulfilled (see the separation condition (37)), then the harmonics associated to such primary vectors \( k \) are the dominant ones in the splitting function \( M(\theta) \), for each small enough value of the perturbation parameter \( \varepsilon \). This possibility of detecting the dominant harmonics in the splitting function in terms of primary resonances is the main reason for our choice of the frequency vectors.

In the quadratic case, the required arithmetic condition (37) can be formulated in terms of the continued fraction of \( \Omega \), which is (eventually) periodic, and in fact, we can restrict ourselves to the case of purely periodic continued fractions. There
are 24 numbers satisfying (37), all of them having 1-periodic or 2-periodic continued fractions,
\[ \Omega_a = [a], \quad a = 1, \ldots, 13, \quad \text{and} \quad \Omega_{1,a} = [1, a], \quad a = 2, \ldots, 12 \quad (4) \]
(this includes the golden number \( \Omega_1 = [1, 1, 1, \ldots] = (\sqrt{5} - 1)/2 \)).

In the cubic case, there is no standard continued fraction theory, but a particular study can be carried out for each cubic irrational \( \Omega \). We consider in this paper the cubic golden number (see for instance [27]):
\[ \Omega \approx 0.6823, \quad \text{the real root of} \quad x^3 + x - 1 = 0, \quad (5) \]
although we stress that a similar approach could be carried out for other cubic irrationals (see some examples in [7]).

In the main result of this paper (see Theorem 1), we establish exponentially small asymptotic estimates for the maximal splitting distance, valid in all the cases (4–5).

In this way, the results provided in [10] for some quadratic frequencies are extended to other cases, including a particular case of cubic frequencies. As far as we know, this is the first result providing asymptotic estimates for the exponentially small splitting of separatrices with 3 frequencies. To avoid technicalities, we put emphasis on the constructive part of the proofs, using the arithmetic properties of the frequencies in order to provide a unified methodology which can be applied to both the quadratic and the cubic cases, stressing the similarities and differences between them. We determine, for every \( \epsilon \) small enough, the dominant harmonic of the Melnikov function \( M(\theta) \), associated to a primary resonance and, consequently, we obtain an estimate for the maximal value of this function.

In a further step, the first order approximation has to be validated showing that the dominant harmonics of the splitting function \( M(\theta) \) correspond to the dominant harmonics of the Melnikov function, as done in [11]. Moreover, one can show in the cases (4–5) that the invariant manifolds intersect along transverse homoclinic orbits, with an exponentially small angle. To obtain this, one needs to consider the “next” dominant harmonics (at least 2 ones in the quadratic case and at least 3 ones in the cubic case, provided their associated vectors \( k \) are linearly independent), which can be carried out for the frequency vectors considered. Nevertheless, in some cases the secondary resonances have to be taken into account giving rise to more involved estimates. We only provide here the main ideas, and rigorous proofs will be published elsewhere.

\section*{1.2. Setup and main result} In order to formulate our main result, let us describe the Hamiltonian considered, which is analogous to the one considered in [13] and other related works. In symplectic coordinates \((x, y, \varphi, I) \in T \times \mathbb{R} \times T^\ell \times \mathbb{R}^\ell\),
\[ H(x, y, \varphi, I) = H_0(x, y, I) + \mu H_1(x, \varphi), \]
\[ H_0(x, y, I) = \langle \omega_x, I \rangle + \frac{1}{2} \langle \Lambda I, I \rangle + \frac{y^2}{2} + \cos x - 1, \quad H_1(x, \varphi) = h(x)f(\varphi), \quad (6) \]
with
\[ h(x) = \cos x, \quad f(\varphi) = \sum_{\substack{k \in \mathbb{Z}^\ell \\text{even} \quad \sigma_k \geq 0}} e^{-\rho|k|} \cos(\langle k, \varphi \rangle - \sigma_k), \quad (7) \]
where the restriction in the sum is introduced in order to avoid repetitions. The Hamiltonian (6–7) provides a particular model for the behavior of a nearly-integrable
Hamiltonian system $\mathcal{K} = \mathcal{K}_0 + \varepsilon \mathcal{K}_1$ in the vicinity of a simple resonance, after carrying out one step of resonant normal form (see for instance [44, 9]). In this way, our unperturbed Hamiltonian $H_0$ and the perturbation $\mu H_1$ play the role of the truncated normal form and the remainder, respectively. In its turn, the truncated normal form is an $O(\varepsilon)$-perturbation of the initial Hamiltonian $\mathcal{K}_0$, making the hyperbolicity appear (a rescaling leads to the fast frequencies (2)). The parameters $\varepsilon$ and $\mu$ should not be considered as independent but rather linked by a relation of the type $\mu = \varepsilon^p$.

To justify the form of $H_1$ chosen in (6–7), let us first stress that, without some restrictive conditions on the perturbation, it is not possible (at the present moment) to establish the existence of splitting, since it is necessary to compute explicitly the Melnikov approximation in order to show that it dominates the error term in (1). Although our methodology can be directly applied to rather general perturbations $H_1$, the reason of our choice of $H_1$ as in (6–7) is that it makes easier the explicit computation of the Melnikov potential. Moreover, the fact that all harmonics in the Fourier expansion with respect to $\varphi$ are non-zero, having an exponential decay, ensures that the study of the dominant harmonics of the Melnikov potential can be carried out directly from the arithmetic properties of the frequency vector $\omega$ (see Section 3.1). In fact, the particular form chosen in (6–7) is significative, since it is a generalization of the Arnold example (introduced in [1] to illustrate the transition chain mechanism in Arnold diffusion) and has often been considered in the literature (see for instance [23, 36]).

Notice that the unperturbed system $H_0$ consists of a classical pendulum and $\ell$ rotors with fast frequencies:

$$P(x, y) = \frac{y^2}{2} + \cos x - 1, \quad \dot{\varphi} = \omega_x + \Lambda I, \quad \dot{I} = 0. \quad (8)$$

The pendulum has a hyperbolic equilibrium at the origin, and the (upper) separatrix can be parameterized by $(x_0(s), y_0(s)) = (4 \arctan e^s, 2/\cosh s)$, $s \in \mathbb{R}$. The rotors system $(\varphi, I)$ has the solutions $\varphi = \varphi_0 + (\omega_x + \Lambda_0) t, I = I_0$. Consequently, $H_0$ has an $\ell$-parameter family of $\ell$-dimensional whiskered invariant tori, with coincident stable and unstable whiskers. Among the family of whiskered tori, we will focus our attention on the torus located at $I = 0$, whose frequency vector is $\omega_x$ as in (2), in our case a quadratic (for $\ell = 2$) or cubic frequency vector (for $\ell = 3$). We also assume the condition of isoenergetic nondegeneracy

$$\det \left( \begin{array}{cc} \Lambda & \omega \\ \omega^\top & 0 \end{array} \right) \neq 0. \quad (9)$$

When adding the perturbation $\mu H_1$, the hyperbolic KAM theorem can be applied (see for instance [44]) thanks to the Diophantine condition (3) and the isoenergetic nondegeneracy (9). For $\mu$ small enough, the whiskered torus persists with some shift and deformation, as well as its local whiskers.

In general, for $\mu \neq 0$ the (global) whiskers do not coincide anymore, and one can introduce a splitting function $\mathcal{M}(\theta), \theta \in \mathbb{T}^\ell$, giving the distance between the whiskers in the complementary directions, on a transverse section. Namely, taking the action coordinates $I \in \mathbb{R}^\ell$ as the complementary directions, and introducing parameterizations $\mathcal{J}^+(\theta)$ and $\mathcal{J}^-(\theta), \theta \in \mathbb{T}^\ell$, for the intersection of the stable and unstable manifolds with the transverse section, one can define the splitting function:

$$\mathcal{M}(\theta) = \mathcal{J}^-(\theta) - \mathcal{J}^+(\theta), \quad \theta \in \mathbb{T}^\ell$$
where

\[ C \]

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approximation (1) is given by the (vector) Melnikov function details, and also [17]). Applying the Poincaré–Melnikov method, the first order (in fact, a normal form procedure is carried out previously, see [8, §5.2] for more details, and also [17]).

\[ \text{Melnikov function} \]

\[ L \]

\[ \sigma \]

\[ \text{function. The phases obtain asymptotic estimates for the splitting. Notice that the constant } \rho > 0 \text{ already considered in [11], and allows us to deal with the Melnikov function and we have chosen for the perturbation the special form given in (7). This form was the Fourier expansion of } f \text{.}

\[ \text{the Fourier expansion of } L \text{ to use the method of residues in order to compute the coefficients of the Fourier expansion of } L(\theta) \text{ (see their expression in Section 3.1). However, our approach can also be directly applied to other classical 1-degree-of-freedom Hamiltonians } P(x, y) = y^2/2 + V(x), \text{ with a potential } V(x) \text{ having a unique nondegenerate maximum, although the use of residues becomes more cumbersome when the separatrix has poles of higher orders (see some examples in [16]).}

\[ \text{In order to emphasize the role played by the arithmetic properties of the splitting, we have chosen for the perturbation the special form given in (7). This form was already considered in [11], and allows us to deal with the Melnikov function and obtain asymptotic estimates for the splitting. Notice that the constant } \rho > 0 \text{ in the Fourier expansion of } f(\varphi) \text{ in (7) gives the complex width of analyticity of this function. The phases } \sigma_k \text{ can be chosen arbitrarily for the purpose of this paper.}

\[ \text{Now we can formulate our main result, providing asymptotic estimates for the maximal splitting distance in both the quadratic and cubic cases.}

\textbf{Theorem 1} (main result). \textit{For the Hamiltonian system (6–7) with } \ell + 1 \text{ degrees of freedom, satisfying the isoenergetic condition (9), assume that } \varepsilon \text{ is small enough and } \mu = \varepsilon^p \text{ with } p > 3. \text{ For } \ell = 2, \text{ if } \Omega \text{ is one of the } 2^d \text{ quadratic numbers (4), and for } \ell = 3, \text{ if } \Omega \text{ is the cubic golden number (5), the following asymptotic estimate holds:}

\[ \max_{\theta \in \mathbb{T}} |M(\theta)| \sim \frac{\mu}{\varepsilon^{1/\ell}} \exp \left\{ \frac{-C_0 h_1(\varepsilon)}{\varepsilon^{1/2\ell}} \right\}, \]

\text{where } C_0 = C_0(\Omega, \rho) \text{ is a positive constant, defined in (28). Concerning the function } h_1(\varepsilon),

\[ \begin{align*}
\text{(a) in the quadratic case } \ell = 2, & \text{ it is periodic in } \ln \varepsilon, \text{ with } \min h_1(\varepsilon) = 1 \text{ and } \\
\text{max } h_1(\varepsilon) = A_1, & \text{ where the constant } A_1 = A_1(\Omega) \text{ is defined in (34);} \\
\text{(b) in the cubic case } \ell = 3, & \text{ it satisfies the bound } 0 < A_0^- \leq h_1(\varepsilon) \leq A_1^+, \text{ where the constants } A_0^- = A_0^- (\Omega) \text{ and } A_1^+ = A_1^+ (\Omega) \text{ are defined in (35).}
\end{align*} \]

\textbf{Remarks.}

\[ \begin{align*}
\text{(1) The periodicity in } \ln \varepsilon \text{ of the function } h_1(\varepsilon) \text{ in the quadratic case (a) was first established in [12] for the quasi-periodically forced pendulum, assuming that the frequency ratio is the golden number } \Omega_1. \text{ Previously, the existence of an oscillatory behavior with lower and upper bounds had been shown in [50].}
\end{align*} \]

\[ \begin{align*}
\text{(2) In contrast to the quadratic case, it turns out in the cubic case (b) that the function } h_1(\varepsilon) \text{ is not periodic in } \ln \varepsilon \text{ and has a more complicated form (see Figure 3, where one can conjecture that } h_1(\varepsilon) \text{ is a quasiperiodic function).}
\end{align*} \]
(3) The exponent $p > 3$ in the relation $\mu = \varepsilon^p$ can be improved in some special cases. For instance, if in (7) one considers $h(x) = \cos x - 1$, then the asymptotic estimates are valid for $p > 2$. This is related to the fact that, in this case, the invariant torus remains fixed under the perturbation and only the whiskers deform [13].

This paper is organized as follows. In Section 2 we study the arithmetic properties of quadratic and cubic frequencies, and in Section 3 we find, for the frequencies considered in (4–5), an asymptotic estimate of the dominant harmonic of the splitting potential, together with a bound of the remaining harmonics which allows us to provide an asymptotic estimate for the maximal splitting distance, as established in Theorem 1.

2. Arithmetic properties of quadratic and cubic frequencies

2.1. Iteration matrices and resonant sequences. We review in this section the technic developed in [10] for studying the resonances of quadratic frequencies ($\ell = 2$), showing that it admits a direct generalization to the case of cubic frequencies ($\ell = 3$).

In the 2-dimensional case, we consider a quadratic frequency vector $\omega \in \mathbb{R}^2$, i.e., its frequency ratio is a quadratic irrational number. Of course, we can assume without loss of generality that the vector has the form $\omega = (1, \Omega)$.

On the other hand, in the 3-dimensional case we consider a cubic frequency vector $\omega \in \mathbb{R}^3$, i.e., the frequency ratios generate a cubic field (an algebraic number field of degree 3). In order to simplify our exposition, we assume that the vector has the form $\omega = (1, \Omega, \Omega^2)$, where $\Omega$ is a cubic irrational number. Hence, the cubic field is $\mathbb{Q}(\Omega)$.

Any quadratic or cubic frequency vector $\omega \in \mathbb{R}^\ell$ satisfies the Diophantine condition (3), with the minimal exponent $\ell - 1$, see for instance [6]. With this in mind, we define the “numerators”

$$\gamma_k := \langle (k, \omega) \rangle \cdot |k|^\ell - 1, \quad k \in \mathbb{Z}^\ell \setminus \{0\},$$

provided a norm $|\cdot|$ for integer vectors has been chosen (for quadratic vectors, the norm $|\cdot| = |\cdot|_1$, i.e., the sum of absolute values of the components of the vector, was used in [10]; for cubic vectors it will be more convenient to use the Euclidean norm $|\cdot| = |\cdot|_2$). Our goal is to provide a classification of the integer vectors $k$, according to the size of $\gamma_k$, in order to find the primary resonances (i.e., the integer vectors $k$ for which $\gamma_k$ is smallest and, hence, fitting best the Diophantine condition (3)) and study their separation with respect to the secondary resonances.

The key point is to use a result by Koch [32]: for a vector $\omega \in \mathbb{R}^\ell$ whose frequency ratios generate an algebraic field of degree $\ell$, there exists a unimodular matrix $T$ (a square matrix with integer entries and determinant $\pm 1$) having the eigenvector $\omega$ with associated eigenvalue $\lambda$ of modulus $> 1$, and such the other $\ell - 1$ eigenvalues are simple and of modulus $< 1$. This result is valid for any dimension $\ell$, and is usually applied in the context of renormalization theory (see for instance [32, 39]), since the iteration of the matrix $T$ provides successive rational approximations to the direction of the vector $\omega$. Notice that the matrix $T$ satisfying the conditions above is not unique (for instance, any power $T^j$, with $j$ positive, also satisfies them).

We will assume without loss of genericity that $\lambda$ is positive ($\lambda > 1$).
In this paper, we are not interested in finding approximations to \( \omega \), but rather to the \textit{quasi-resonances} of \( \omega \), which lie close to the orthogonal hyperplane \( \langle \omega \rangle^\perp \).

With this aim, we consider the matrix \( U = (T^{-1})^\top \), which satisfies the following fundamental equality:

\[
\langle Uk, \omega \rangle = \langle k, U^\top \omega \rangle = \frac{1}{\lambda} \langle k, \omega \rangle.
\] (12)

We say that an integer vector \( k \) is \textit{admissible} if \( |\langle k, \omega \rangle| < \frac{1}{2} \). We restrict ourselves to the set \( A \) of admissible vectors, since for any \( k \not\in A \) we have \( |\langle k, \omega \rangle| > \frac{1}{2} \) and \( \gamma_k > |k|^{\ell - 1}/2 \). We see from (12) that if \( k \in A \), then also \( Uk \in A \). We say that \( k \) is \textit{primitive} if \( k \in A \) but \( U^{-1}k \not\in A \). We also deduce from (12) that \( k \) is primitive if and only if

\[
\frac{1}{2\lambda} < |\langle k, \omega \rangle| < \frac{1}{2}.
\] (13)

Since the first component of \( \omega \) is equal to 1, it is clear that any admissible vectors can be presented in the form

\[
k^0(j) = (-\text{rint}(j\Omega), j), \quad j = \mathbb{Z} \setminus \{0\} \quad (\ell = 2),
\]

\[
k^0(j) = (-\text{rint}(j_1\Omega + j_2\Omega^2), j_1, j_2), \quad j = (j_1, j_2) \in \mathbb{Z}^2 \setminus \{0\} \quad (\ell = 3)
\]

(we denote \( \text{rint}(x) \) the closest integer to \( x \)). If \( k^0(j) \in \mathbb{Z}^\ell \) is primitive, we also say that \( j \in \mathbb{Z}^{\ell - 1} \) is primitive, and denote \( P \) be the set of such primitives. Now we define, for each \( j \in P \), the following \textit{resonant sequence} of integer vectors:

\[
s(j, n) := U^nk^0(j), \quad n = 0, 1, 2, \ldots
\] (14)

It turns out that such resonant sequences cover the whole set \( A \) of admissible vectors, providing a classification of them. The properties of such a classification follow from Proposition 2 for the case of quadratic frequencies and from Proposition 3 for the case of cubic frequencies.

### 2.2. Properties of quadratic frequencies.

It is well-known that all \textit{quadratic irrational numbers} \( \Omega \in (0, 1) \), i.e., the real roots of quadratic polynomials with rational coefficients, have the continued fraction

\[
\Omega = [a_1, a_2, a_3, \ldots], \quad a_i \in \mathbb{Z}^+,
\]

that is \textit{eventually periodic}, i.e., periodic starting with some element \( a_i \). In fact, as we see below we can restrict ourselves to the numbers with \textit{purely periodic} continued fractions and denote them according to their periodic part; for an \( m \)-periodic continued fraction, we write \( \Omega_{a_1, \ldots, a_m} = [\overline{a_1, \ldots, a_m}] \). For example, the famous \textit{golden number} is \( \Omega_1 = [\overline{1}] = (\sqrt{5} - 1)/2 \), and the \textit{silver number} is \( \Omega_2 = [\overline{2}] = \sqrt{2} - 1 \).

For a quadratic frequency \( \omega = (1, \Omega) \), the matrix \( T \) provided by Koch’s result [32] can be constructed directly from the continued fraction of \( \Omega \). The quadratic numbers (4), considered in this paper, have 1-periodic or 2-periodic continued fractions.

Let us write their matrix \( T = T(\Omega) \) with \( \omega \) as an eigenvector, and the associated eigenvalue \( \lambda = \lambda(\Omega) > 1 \):

\[
\begin{align*}
\text{for } \Omega &= \Omega_a, & T &= \begin{pmatrix} a & 1 \\ 1 & 0 \end{pmatrix}, & \lambda &= \frac{1}{\Omega}; \\
\text{for } \Omega &= \Omega_{1,a}, & T &= \begin{pmatrix} a + 1 & 1 \\ a & 1 \end{pmatrix}, & \lambda &= \frac{1}{1 - \Omega}.
\end{align*}
\]
Remark. In what concerns the contents of this paper, it is enough to consider quadratic numbers with purely periodic continued fractions, due to the equivalence of any quadratic number \( \Omega \), with an eventually periodic continued fraction, to some \( \Omega \) with a purely periodic one: \( \hat{\Omega} = \frac{c + d\Omega}{a + b\Omega} \), with integers \( a, b, c, d \) such that \( ad - bc = \pm 1 \). Then, it can be shown that the same results apply to both numbers \( \Omega \) and \( \hat{\Omega} \) for \( \varepsilon \) small enough. For instance, the results for the golden number \( \Omega_1 \) also apply to the noble numbers \( \Omega = [b_1, \ldots, b_n, \overline{1}] \). We point out that this happens for all the cases we have considered in this paper, would depend on the non-periodic part of the continued fraction.

Now we consider the resonant sequences defined in (14). For the matrix \( T \), let \( v_2 \) be a second eigenvector (with eigenvalue \( \sigma/\lambda \) of modulus < 1, where \( \sigma = \det T = \pm 1 \)). Then, the vectors \( \omega, v_2 \) are a basis of eigenvectors. For the matrix \( U = (T^{-1})^\top \), let \( u_1, u_2 \) be a basis of eigenvectors with eigenvalues \( 1/\lambda \) and \( \sigma \lambda \), respectively. It is well-known that \( \langle u_2, \omega \rangle = \langle u_1, v_2 \rangle = 0 \). For any primitive integer \( j \in \mathcal{P} \), we define the quantities

\[
\gamma_j^r := \lim_{n \to \infty} \gamma_j s_{j(n)} = |r_j| K_j, \quad K_j = \left| \frac{p_j}{\langle u_2, v_2 \rangle} u_2 \right| = \left| k^0(j) - \frac{r_j}{\langle u_1, \omega \rangle} u_1 \right|,
\]

and one has:

1. \( \gamma_j s_{j(n)} = \gamma_j^r + \mathcal{O}(\lambda^{-2n}), \quad n \geq 0 \);
2. \( |s(j, n)| = K_j \lambda^n + \mathcal{O}(\lambda^{-n}), \quad n \geq 0 \);
3. \( \gamma_j^r > \frac{(1 + \Omega) |j| - a}{2\lambda}, \quad a = \frac{1}{2} \left( 1 + \frac{|u_1|}{|\langle u_1, \omega \rangle|} \right) \).

Since the lower bounds given in (c) for the “limit numerators” \( \gamma_j^r \) are increasing with respect to the primitive \( j \), we can select the minimal of them, corresponding to some primitive \( j_0 \). We denote

\[
\gamma^* := \liminf_{|k| \to \infty} \gamma_k = \min_{j \in \mathcal{P}} \gamma_j^* = \gamma_{j_0}^* > 0.
\]

The corresponding sequence \( s_0(n) := s(j_0, n) \) gives us the primary resonances, and we call secondary resonances the integer vectors belonging to any of the remaining sequences \( s(j, n), j \neq j_0 \).

We introduce normalized numerators \( \tilde{\gamma}_k \) and their limits \( \tilde{\gamma}_j^r, j \in \mathcal{P} \), after dividing by \( \gamma^* \), and in this way \( \tilde{\gamma}_{j_0}^r = 1 \). We also define a parameter \( B_0 \) measuring the separation between primary and secondary resonances:

\[
\tilde{\gamma}_k := \frac{\gamma_k}{\gamma^*}, \quad \tilde{\gamma}_j^r := \frac{\gamma_j^r}{\gamma^*}, \quad B_0 := \min_{j \in \mathcal{P}\setminus\{j_0\}} \left( \tilde{\gamma}_j^r \right)^{1/2},
\]

where we included the square root for convenience, see (37). We are implicitly assuming the hypothesis that the primitive \( j_0 \) is unique. In other words, we are assuming that \( B_0 > 1 \). We point out that this happens for all the cases we have explored.
2.3. Properties of cubic frequencies. Now, we consider a frequency vector of the form \( \omega = (1, \Omega, \Omega^2) \), where \( \Omega \) is a cubic irrational number. If we consider the matrix \( T \) given by Koch’s result [32], mentioned in Section 2.1, we can distinguish two possible cases for its three eigenvalues \( \lambda, \lambda_2, \lambda_3 \) (recall that \( \lambda > 1 \) is the eigenvalue with eigenvector \( \omega \)):

- **the real case**: the three eigenvalues \( \lambda, \lambda_2, \lambda_3 \) are real;
- **the complex case**: only the eigenvalue \( \lambda \) is real, and the other two ones \( \lambda_2, \lambda_3 \) are a pair of complex conjugate numbers.

These two cases are often called **totally real** and **non-totally real**, respectively. In this paper we only consider cubic frequency vectors in the complex or non-totally real case.

**Remark.** The reason to restrict ourselves to the complex case is that the remaining two (complex) eigenvalues have the same modulus. As we see below, it is natural to extend the results for quadratic frequencies to cubic frequencies of complex type. Instead, the study of the real case would require a different approach, since the behavior of the associated small divisors turns out to be different from the complex case considered here.

Unlike the 2-dimensional quadratic frequencies, in the case of 3-dimensional cubic frequencies there is no standard theory of continued fractions providing a direct construction of the matrix \( T \). However, there are some multidimensional continued fractions algorithms, which applied to the pair \( (\Omega, \Omega^2) \) could be helpful to provide \( T \). Such algorithms have been studied in the context of renormalization and its connections to KAM theory, see for instance [33, 27, 28, 30, 31]. Fortunately, for a given concrete cubic frequency vector it is not hard to find the matrix \( T \) by inspection, as we do in Section 2.4 for the cubic golden number. Other examples of cubic frequencies and their associated matrices are given in [7] (see also [37] for an account of examples and results concerning cubic frequencies).

As in Section 2.2, we are going to establish the properties of the resonant sequences (14). Let us consider a basis of eigenvectors of \( T \), writing the two complex ones in terms of real and imaginary parts: \( \omega, v_2 + iv_3, v_2 - iv_3 \), with eigenvalues \( \lambda, \lambda_2, \lambda_3 = \bar{\lambda}_2 \), respectively. Notice that \( |\lambda_2| = \lambda^{-1/2} \); we denote \( \phi := \arg(\lambda_2) \).

In a similar way, we consider for the matrix \( U = (T^{-1})^\top \) a basis \( u_1, u_2 + iu_3, u_2 - iu_3 \) with eigenvalues \( \lambda^{-1}, \lambda_2^{-1} \), and \( \lambda_3^{-1} = \bar{\lambda}_2^{-1} \), respectively. In this way, we avoid working with complex vectors. One readily sees that \( \langle u_2, \omega \rangle = \langle u_3, \omega \rangle = 0 \), i.e., \( u_2 \) and \( u_3 \) span the resonant plane \( \langle \omega \rangle \). Other useful equalities are: \( \langle u_1, v_2 \rangle = \langle u_1, v_3 \rangle = 0, \langle u_2, v_2 \rangle = -\langle u_3, v_3 \rangle, \langle u_2, v_3 \rangle = \langle u_3, v_2 \rangle \). We define \( Z_1, Z_2 \) and \( \theta \) through the formulas

\[
\frac{1}{2}(|u_2|^2 + |u_3|^2) = Z_1, \quad \frac{1}{2}(|u_2|^2 - |u_3|^2) = Z_2 \cos \theta, \quad \langle u_2, u_3 \rangle = Z_2 \sin \theta. \tag{17}
\]

For any primitive \( j \), we define the quantities

\[
r_j := \langle k^0(j), \omega \rangle, \quad p_j := \langle k^0(j), v_2 \rangle, \quad q_j := \langle k^0(j), v_3 \rangle, \tag{18}
\]

and \( E_j, \psi_j \) through the formulas

\[
\frac{(v_2, u_2)p_j + (v_2, u_3)q_j}{(v_2, u_2)^2 + (v_2, u_3)^2} = E_j \cos \psi_j, \quad \frac{(v_2, u_3)p_j - (v_2, u_2)q_j}{(v_2, u_2)^2 + (v_2, u_3)^2} = E_j \sin \psi_j. \tag{19}
\]

The following proposition extends the results, given in Proposition 2 for the quadratic case, to the complex cubic case.
Proposition 3. For any primitive \(j = (j_1, j_2)\) \(\in \mathcal{P}\), the sequence of numerators \(\gamma_{s(j,\cdot)}\) oscillates as \(n \to \infty\) between two values,

\[
\gamma_j^- = \gamma_j^+ (1 - \delta), \quad \gamma_j^+ = \gamma_j^+ (1 + \delta),
\]

where we define

\[
\gamma_j^+ = |r_j| K_j, \quad K_j = E_j^2 Z_1, \quad \delta = \frac{Z_2}{Z_1} < 1.
\]

We also have

\[
\begin{align*}
(a) \quad & \gamma_{s(j,n)} = \gamma_j^+ (1 + \delta \cos[2n\phi + 2\psi_j - \theta]) + \mathcal{O}(\lambda^{-3n/2}); \\
(b) \quad & |s(j, n)|^2 = K_j (1 + \delta \cos[2n\phi + 2\psi_j - \theta]) \cdot \lambda^n + \mathcal{O}(\lambda^{-n/2}); \\
(c) \quad & \gamma_j^- \geq \frac{1 - \delta}{2\lambda(1 + \delta)} \left[ |j| - \frac{|u_1|}{2(|u_1, \omega)|} \right]^2.
\end{align*}
\]

Proof. We present the primitive vector associated to \(j\) in the basis \(u_1, u_2, u_3\):

\[
k^0(j) = c_1 u_1 + c_2 u_2 + c_3 u_3,
\]

and taking scalar products with \(\omega, v_2\) and \(v_3\) and solving a linear system, one can obtain the values of the coefficients:

\[
c_1 = \frac{r_j}{(u_1, \omega)}, \quad c_2 = E_j \cos \psi_j, \quad c_3 = E_j \sin \psi_j,
\]

where the definitions (18–19) have been taken into account. Now, we apply the iteration matrix \(U\). Using the identities

\[
U^n u_2 = \lambda^{n/2} [\cos(n\phi) u_2 + \sin(n\phi) u_3],
\]

\[
U^n u_3 = \lambda^{n/2} [-\sin(n\phi) u_2 + \cos(n\phi) u_3],
\]

we find

\[
s(j, n) = U^n k^0(j) = \lambda^{n/2} E_j [\cos(n\phi + \psi_j) u_2 + \sin(n\phi + \psi_j) u_3] + \mathcal{O}(\lambda^{-n}),
\]

and we deduce, according to the definitions (17),

\[
|s(j, n)|^2 = \lambda^n E_j^2 (Z_1 + Z_2 \cos[2n\phi + 2\psi_j - \theta]) + \mathcal{O}(\lambda^{-n/2}),
\]

which gives (b). Multiplying by \(|s(j, n), \omega| = |r_j| \lambda^{-n}\), we obtain \(\gamma_{s(j,n)}\) as given in (a). This implies the asymptotic bounds introduced in (20).

Finally, one easily sees that

\[
|k^0(j) - c_1 u_1|^2 = |c_2 u_2 + c_3 u_3|^2 = E_j^2 (Z_1 + Z_2 \cos[2\psi_j - \theta]) \leq E_j^2 (Z_1 + Z_2)
\]

and, hence, using (21) and (13) and also that \(|k^0(j)| \geq |j|\), we get

\[
K_j = E_j^2 Z_1 \geq \frac{|k^0(j) - c_1 u_1|^2}{1 + \delta} \geq \frac{1}{1 + \delta} \left[ |j| - \frac{|u_1|}{2(|u_1, \omega)|} \right]^2,
\]

which implies the lower bound (c). □

As we can see in (a), the existence of limit of the sequences \(\gamma_{s(j,n)}\) stated in Proposition 2 for the quadratic case is replaced here by an oscillatory limit behavior, with a lower limit \(\liminf_{n \to \infty} \gamma_{s(j,n)} \geq \gamma_j^-\) and an upper limit \(\limsup_{n \to \infty} \gamma_{s(j,n)} \leq \gamma_j^+\).

Notice that such values of the limits are exact if the phase \(\phi/2\pi\), that appears in (a), is irrational.

Selecting the primitive \(j_0 \in \mathcal{P}\) which gives the minimal limits, we have the primary resonances, and we denote them by \(s_0(n) := s(j_0, n)\), and we call secondary
resonances the integer vectors belonging to any of the remaining sequences \( s(j, n) \), \( j \neq j_0 \). Such primary resonances can easily be detected thanks to Proposition 3(c): although \( \gamma_j^\pm \) are not increasing in general with respect to \( |j| \), we have an increasing lower bound, which implies that \( \lim_{|j| \to \infty} \gamma_j^\pm = \infty \), and then one has to check only a finite number of primitive vectors \( j \) in order to find the minimal \( \gamma_j^- \) and \( \gamma_j^+ \), corresponding to the primary resonances.

As in Section 2.2, we define normalized values \( \tilde{\gamma}_k \), \( \tilde{\gamma}_j^* \), \( \tilde{\gamma}_j^\pm \), after dividing by the minimal among the values \( \gamma_j^* \),

\[
\tilde{\gamma}_k := \frac{\gamma_k}{\gamma^*}, \quad \tilde{\gamma}_j^* := \frac{\gamma_j^*}{\gamma^*}, \quad \tilde{\gamma}_j^\pm := \frac{\gamma_j^\pm}{\gamma^*}, \quad \text{where } \gamma^* := \min_{j \in P} \gamma_j^* = \gamma_{j_0}^*.
\]  

We also introduce a parameter \( B_0^- \), as a measure for the separation between primary and secondary resonances:

\[
B_0^- := \min_{j \in P \setminus \{j_0\}} \left( \frac{\tilde{\gamma}_j^-}{\tilde{\gamma}_j^+} \right)^{1/3}
\]  

(compare with (16), and see also (37)). Notice that the distinction between primary and secondary resonances makes sense if \( B_0^- > 1 \), i.e., the interval \([\tilde{\gamma}_j^-, \tilde{\gamma}_j^+]\) has no intersection with any other interval \([\tilde{\gamma}_{j'}, \tilde{\gamma}_{j'}^+], j \neq j_0 \) (as happens in the cubic golden case, see the next section).

2.4. The cubic golden frequency vector. Now, we assume that \( \Omega \) is the cubic golden number: the real root of \( x^3 + x - 1 = 0 \). We have \( \Omega \approx 0.6823 \). In this case, the matrix \( T \) can easily be found by inspection. We have

\[
T = \begin{pmatrix}
1 & 0 & 1 \\
1 & 0 & 0 \\
0 & 1 & 0
\end{pmatrix}, \quad U = (T^{-1})^\top = \begin{pmatrix}
0 & 0 & 1 \\
1 & 0 & -1 \\
0 & 1 & 0
\end{pmatrix},
\]

with the eigenvalue \( \lambda = 1/\Omega \approx 1.4656 \).

It is not hard to compute the data provided by Proposition 3 in this concrete case. In particular, we have

\[
\phi = \arg(\lambda_2) = -\arctan \frac{4\sqrt{3\Omega}}{\Omega(6\Omega^2 + 9\Omega + 4)} + \pi \approx \frac{13\pi}{22}
\]  

and, from Proposition 3(b), we have the following approximately periodic behaviors:

\[
\gamma_{s(j, n + 22)} \approx \gamma_{s(j, n)}, \quad \text{and } |s(j, n + 22)| \approx \lambda^{11} |s(j, n)|.
\]

Other relevant parameters are

\[
\gamma^* = \frac{2}{\Omega} (5 + \Omega + 4\Omega^2) \approx 0.4867 \quad \text{and } \quad \delta = (3 - 2\Omega)\sqrt{2 - \Omega + \Omega^2}/(5 + \Omega + 4\Omega^2) \approx 0.2895.
\]

In Table 1, we write down the values \( \gamma_j^* \), and the bounds \( \gamma_j^- \) and \( \gamma_j^+ \), for the resonant

\[
\begin{array}{|c|c|c|c|}
\hline
k_0(j) & \gamma_j^- & \gamma_j^* & \gamma_j^+ \\
\hline
[0, 0, 1] & 0.3459 & 0.4867 & 0.6276 \\
[-1, 2, 0] & 1.0376 & 1.4602 & 1.8829 \\
[-2, 1, 2] & 3.1127 & 4.3807 & 5.6488 \\
|j| \geq 3 & \geq 1.2742 \\
\hline
\end{array}
\]

Table 1. Numerical data for the cubic golden frequency vector.
sequences induced by a few primitives $k_0(j)$, as well as a lower bound for all other primitives. The smallest ones correspond to the primitive vector $k_0([0, 1]) = [0, 0, 1]$ (primary resonances). The parameter introduced in (23), indicating the separation between the primary and the secondary resonances, is $B_0^- \approx 1.1824$.

Additionally, it is interesting to visualize such a separation in the following way. Taking logarithm in both hands of the Diophantine condition (3), we can write it as

$$-\ln |\langle k, \omega \rangle| \leq 2 \ln |k| - \ln \gamma.$$ 

If we draw all the points with coordinates $(\ln |k|, -\ln |\langle k, \omega \rangle|)$ (see Figure 1), we can see a sequence of points lying between the two straight lines $-\ln |\langle k, \omega \rangle| = 2 \ln |k| - \ln \gamma(\beta)$. Such points correspond to integer vectors belonging to the sequence of primary resonances: $k = s_0(n)$, $n \geq 0$.

![Figure 1. Points $(\ln |k|, -\ln |\langle k, \omega \rangle|)$ for the cubic golden frequency vector](image)

### 3. Asymptotic estimates for the maximal splitting distance

In order to provide asymptotic estimates for the splitting, we start with the first order approximation, given by the Poincaré–Melnikov method. It is convenient for us to work with the (scalar) Melnikov potential $L$ and the splitting potential $L$, but we state our main result in terms of the splitting function $M = \nabla L$, which gives a measure of the splitting distance between the invariant manifolds of the whiskered torus. Notice also that the nondegenerate critical points of $L$ correspond to simple zeros of $M$, and give rise to transverse homoclinic orbits to the whiskered torus.

In the present paper, we restrict ourselves to present the constructive part of the proofs, which corresponds to find, for every sufficiently small $\varepsilon$, the dominant harmonics of the Fourier expansion of the Melnikov potential $L(\theta)$, as well as to provide bounds for the sum of the remaining terms of that expansion. The final step, to ensure that the Poincaré–Melnikov method (1) predicts correctly the size of splitting in our singular case $\mu = \varepsilon^p$, can be worked out simply by showing that the asymptotic estimates of the dominant harmonics are large enough to overcome the harmonics of the error term. This final step is analogous to the one done in
for the case of the golden number $\Omega_1$ (using the upper bounds for the error term provided in [13]), and will be published elsewhere for all cases considered in Theorem 1.

First, we are going to find in Section 3.1 an exponentially small asymptotic estimate for the dominant harmonic among the ones associated to primary resonances, given by a function $h_1(\varepsilon)$ in the exponent. We also provide an estimate for the sum of all other (primary or secondary) harmonics. This can be done jointly for both the quadratic ($\ell = 2$) and cubic ($\ell = 3$) cases. In Section 3.2, we establish a condition ensuring that the dominant harmonic among all harmonics is given by a primary resonance. This condition is fulfilled for the frequencies (4–5). To complete the proof of Theorem 1, we show that the different arithmetic properties of quadratic and cubic frequencies lead to different properties of the function $h_1(\varepsilon)$: periodic (with respect to $\ln \varepsilon$) in the quadratic case, and a more complicated bounded function in the cubic case.

3.1. **Dominant harmonics of the splitting potential.** We put our functions $f$ and $h$, defined in (7), into the integral (10) and get the Fourier expansion of the Melnikov potential, where the coefficients can be obtained using residues:

$$L(\theta) = \sum_{k \notin \mathbb{Z}, k \geq 0} L_k \cos(\langle k, \theta \rangle - \sigma_k), \quad L_k = \frac{2\pi \langle [k, \omega_\varepsilon] \rangle e^{-\rho|k|}}{\sinh \left[ \frac{\pi}{2} \langle [k, \omega_\varepsilon] \rangle \right]}.$$  

Using (2) and (11), we present the coefficients in the form

$$L_k = \alpha_k e^{-\beta_k}, \quad \alpha_k \approx \frac{4\pi \gamma_k}{|k|^{\ell-1}} \frac{\varepsilon^{|k|}}{\sqrt{\varepsilon}}, \quad \beta_k = \rho |k| + \frac{\pi \gamma_k}{2|k|^{\ell-1}} \sqrt{\varepsilon},$$  

where an exponentially small term has been neglected in the denominator of $\alpha_k$.

For any given $\varepsilon$, the harmonics with largest coefficients $L_k(\varepsilon)$ correspond essentially to the smallest exponents $\beta_k(\varepsilon)$. Thus, we have to study the dependence on $\varepsilon$ of such exponents.

With this aim, we introduce for any $X, Y$ (and a fixed $\ell = 2, 3$) the function

$$G(\varepsilon; X, Y) := \frac{Y^{1/\ell}}{\ell} \left[ (\ell - 1) \left( \frac{\varepsilon}{X} \right)^{1/2\ell} + \left( \frac{X}{\varepsilon} \right)^{(\ell-1)/2\ell} \right].$$  

One easily checks that this function has its minimum at $\varepsilon = X$, and the corresponding minimum value is $G(X; X, Y) = Y^{1/\ell}$. Then, the exponents $\beta_k(\varepsilon)$ in (25) can be presented in the form

$$\beta_k(\varepsilon) = \frac{C_0}{\varepsilon^{1/2\ell}} g_k(\varepsilon),$$  

where we define

$$g_k(\varepsilon) := G(\varepsilon; \varepsilon_k, \tilde{\gamma}_k), \quad \varepsilon_k := D_0 \frac{\tilde{\gamma}_k^2}{|k|^{2\ell}},$$  

$$C_0 = \ell \left( \frac{\rho}{\ell - 1} \right)^{(\ell-1)/\ell} \left( \frac{\pi \gamma^*}{2} \right)^{1/\ell}, \quad D_0 = \left( \frac{\ell - 1}{2\rho} \right)^2,$$  

with $\gamma^* = \gamma^*_b$ and $\tilde{\gamma}_k$ given in (15–16) for $\ell = 2$, and in (22) for $\ell = 3$. Consequently, for all $k$ we have $\beta_k \geq \frac{C_0 \tilde{\gamma}_k^{1/\ell}}{\varepsilon^{1/2\ell}}$, which provides the maximum value of the coefficient.
As we show in the next result, the function \( h \) where we define, in order to unify the notation, \( L_k(\varepsilon) \) of the harmonic given by the integer vector \( k \). Recall that for \( k = s(j,n) \), belonging to the resonant sequence generated by a given primitive \( j \in \mathcal{P} \) (see definition (14)), the (normalized) numerators \( \hat{\gamma}_k \) tend to a limit \( \hat{\gamma}_j^* \) (see Proposition 2 for the quadratic case), or oscillate between two limit values \( \hat{\gamma}_j^- \), with \( \hat{\gamma}_j^- < \hat{\gamma}_j^* < \hat{\gamma}_j^+ \) (see Proposition 3 for the cubic case).

The primary integer vectors \( k \), belonging to the sequence \( s_0(n) = s(j_0,n) \), play an important role here, since they give the smallest limit \( \hat{\gamma}_j^* = 1 \). Consequently, they give the dominant harmonics of the Melnikov potential, at least for \( \varepsilon \) close to their minimum points \( \varepsilon_k^* \). Our aim is to show that this happens also for any \( \varepsilon \) (small enough) not necessarily close to \( \varepsilon_k^* \), under a condition ensuring that the separation between the primary and secondary resonances is large enough (see the separation condition (37)). For the sequence of primary resonances, the asymptotic behavior of the functions \( g_k(\varepsilon) \) in (27), \( k = s_0(n) \), is obtained from the main terms, as \( n \to \infty \), given by Propositions 2 and 3. Thus, we can write

\[
g_{s_0(n)}(\varepsilon) \approx g_n^*(\varepsilon) := G(\varepsilon; \varepsilon_n^*, b_n), \quad \varepsilon_n^* := \frac{D_0}{K_{j_0}^{2\ell/(\ell-1)} \left( \frac{1}{b_n \lambda^{\ell n}} \right)^{2/(\ell-1)}},
\]

where we define, in order to unify the notation,

\[
\begin{align*}
b_n &= 1 \quad \text{ in the quadratic case;} \\
b_n &= 1 + \delta \cos[2n\phi + 2\psi_{j_0} - \theta] \quad \text{ in the cubic case.}
\end{align*}
\]

Notice that (29) relies on the approximations \( \hat{\gamma}_{s_0(n)} \approx b_n \) and \( |s_0(n)|^{\ell-1} \approx K_{j_0} b_n \lambda^n \), as well as the fact that \( \hat{\gamma}_{j_0}^* = 1 \).

Now we define, for any given \( \varepsilon \), the function \( h_1(\varepsilon) \) as the minimum of the values \( g_n^*(\varepsilon) \), which takes place for some index \( N = N(\varepsilon) \),

\[
h_1(\varepsilon) := \min_{n \geq 0} g_n^*(\varepsilon) = g_N^*(\varepsilon).
\]

As we show in the next result, the function \( h_1(\varepsilon) \) indicates for any \( \varepsilon \) the size of the dominant harmonic among the primary resonances, given by the integer vector \( k = s_0(N) \). Furthermore, we are going to establish an asymptotic estimate for the sum of all remaining coefficients in the Fourier expansion of the splitting function. This second estimate is written in terms of the function

\[
h_2(\varepsilon) := \min_{k \neq s_0(N)} g_k(\varepsilon).
\]

See Figure 2 as an illustration for the functions \( h_1(\varepsilon) \) and \( h_2(\varepsilon) \), corresponding to the case of the quadratic vector given by \( \Omega_{1,2} \), and Figure 3 for \( h_1(\varepsilon) \) in the case of the cubic golden vector.

Notice that this result is stated in terms of the Fourier coefficients of the splitting function \( \mathcal{M} = \nabla \mathcal{L} \). We write, for the splitting potential,

\[
\mathcal{L}(\theta) = \sum_{k \in \mathbb{Z} \setminus \{0\} \atop k \geq 0} \mathcal{L}_k \cos(\langle k, \theta \rangle - \tau_k),
\]

with upper bounds for \( |\mathcal{L}_k - \mu L_k| \) and \( |\tau_k - \sigma_k| \). Then, for the splitting function we have \( |\mathcal{M}_k| = |k| \mathcal{L}_k \).
Figure 2. Graphs of the functions $g_k(\varepsilon)$, $h_1(\varepsilon)$, $h_2(\varepsilon)$ for the quadratic number $\Omega_{1,2}$, using a logarithmic scale for $\varepsilon$.

Figure 3. Graphs of the functions $g_*^n(\varepsilon)$ and $h_1(\varepsilon)$ for the cubic golden number $\Omega \approx 0.6823$, using a logarithmic scale for $\varepsilon$.

Proposition 4. For $\varepsilon$ small enough and $\mu = \varepsilon^p$ with $p > 3$, one has:

(a) $|M_{s_0(n)}| \sim |s_0(N)| L_{s_0(n)} \sim \frac{\mu}{\varepsilon^{1/\ell}} \exp \left\{ -\frac{C_0 h_1(\varepsilon)}{\varepsilon^{1/2\ell}} \right\}$;

(b) $\sum_{k \neq s_0(n)} |M_k| \sim \frac{\mu}{\varepsilon^{1/\ell}} \exp \left\{ -\frac{C_0 h_2(\varepsilon)}{\varepsilon^{1/2\ell}} \right\}$.

Sketch of the proof. At first order in $\mu$, for the coefficients of the splitting function we can write

$$|M_k| \sim \mu |k| L_k = \mu |k| \alpha_k e^{-\beta_k},$$  \hspace{1cm} (33)
where we have neglected the error term in the Melnikov approximation (1) and used the expression (25) for the coefficients of the Melnikov potential. As mentioned throughout this section, the main behavior of the coefficients $L_k$ is given by the exponents $\beta_k$, which have been written in (27) in terms of the functions $g_k(\varepsilon)$. In particular, the coefficient associated to the dominant harmonic, among the primary resonances, $L_{s_0(N)}$ with $N = N(\varepsilon)$, can be expressed in terms of the function $h_1(\varepsilon)$ introduced in (31).

Now, we consider the remaining factors in (33). We see from (25) that such factors can be written as $|k| \alpha_k \sim |k|^{-(\ell-2)} \sqrt{\varepsilon}$. For $k = s_0(N)$, let us show that they turn out to be polynomial with respect to $\varepsilon$, with a concrete exponent to be determined. First, we use that in (29) we have $\varepsilon_n^* \sim \lambda^{-2n/(\ell-1)}$. Then, for a given $\varepsilon$, the coefficient $N = N(\varepsilon)$ giving the dominant harmonic is such that $\varepsilon_N^*$ (the minimum of the function $g_N^*$) is close to $\varepsilon$. Then, we have $\lambda^N \sim \varepsilon^{-(\ell-1)/2\ell}$. On the other hand, we deduce from from Propositions 2(b) and 3(b) that $|s_0(N)|^{\ell-1} \sim \lambda^N$. Putting the obtained estimates together, we get $|s_0(N)| \alpha_{s_0(N)} = \varepsilon^{-1/\ell}$, which provides the polynomial factor in part (a). The estimate obtained is valid for the dominant coefficient of the Melnikov function. To get the analogous estimate for the splitting function, one has to bound the corresponding coefficient of the error term in (1), showing that it is also exponentially small and dominated by the main term in the approximation. This works as in [11], where the case of the golden number was considered, and we omit the details here.

The proof of part (b) can be carried out in similar terms. For the second dominant harmonic, we get an exponentially small estimate with the function $h_2(\varepsilon)$, defined in (32). This estimate is also valid if one considers the whole sum in (b), since the terms of this sum can be bounded by a geometric series and, hence, it can be estimated by its dominant term (see [11] for more details).

3.2. Study of the functions $h_1(\varepsilon)$ and $h_2(\varepsilon)$. To conclude the proof of Theorem 1, we show in this section the different properties of the function $h_1(\varepsilon)$ for the quadratic and cubic cases, and establish a condition, fulfilled in all cases (4–5), ensuring that $h_2(\varepsilon) \leq h_1(\varepsilon)$ for any $\varepsilon$.

Lemma 5.

(a) In the quadratic case $\ell = 2$, the function $h_1(\varepsilon)$ is $4 \ln \lambda$-periodic in $\ln \varepsilon$, with $\min h_1(\varepsilon) = A_0$ and $\max h_1(\varepsilon) = A_1$, with constants

$$A_0 = 1, \quad A_1 = \frac{1}{2} \left( \frac{1}{\sqrt{\lambda}} + \sqrt{\lambda} \right).$$

(b) In the cubic case $\ell = 3$, the function $h_1(\varepsilon)$ satisfies the bound $0 < A_0^- \leq h_1(\varepsilon) \leq A_1^+$, with constants

$$A_0^- = (1 - \delta)^{1/3}, \quad A_1^+ = \left( \frac{1 + \delta}{3} \right)^{1/3} \left( \frac{\sqrt{\lambda} + 1}{2\lambda} \right)^{1/6} + \left( \frac{2\lambda}{\sqrt{\lambda} + 1} \right)^{1/3}.$$  

Proof. We use that the functions $g_n^*(\varepsilon)$ and the values $\varepsilon_n^*$ satisfy the following scaling properties:

$$g_{n+1}^*(\varepsilon) = \left( \frac{b_{n+1}}{b_n} \right)^{1/\ell} g_n^* \left( \frac{\varepsilon_n^*}{\varepsilon_n^* + \varepsilon} \right), \quad \varepsilon_n^* = \left( \frac{b_n}{b_{n+1} \lambda^2} \right)^{2/(\ell - 1)} \varepsilon_n^*.$$  

In the quadratic case ($\ell = 2$), we have $b_n = b_{n+1} = 1$ and, hence, the relations (36) become
\[ g^*_n(\varepsilon) = g^*_n(\lambda^4 \varepsilon), \quad \varepsilon^*_n = \frac{\varepsilon^*_n}{\lambda^4}, \]
where we have $(\varepsilon^*_n)$ as a geometric sequence, and the functions $g^*_n(\varepsilon)$ are just translations of the initial one, if we use a logarithmic scale for $\varepsilon$ (see Figure 2). It is easy to check that the intersection between the graphs of the functions $g^*_n$ and $g^*_n$ takes place at $\varepsilon^*_n := \varepsilon^*_n / \lambda^2$. Thus, for $\varepsilon \in [\varepsilon^*_n, \varepsilon^*_n - 1]$, we have $N(\varepsilon) = n$, and we deduce that $h_1(\varepsilon) = g^*_n(\varepsilon)$. We can obtain $h_1(\varepsilon)$ from any given interval $[\varepsilon^*_n, \varepsilon^*_n - 1]$ by extending it as a $4 \ln \lambda$-periodic function of $\ln \varepsilon$, and it is clear that its minimum and maximum values are $h_1(\varepsilon^*_n) = 1$ and $h_1(\varepsilon^*_n) = h_1(\varepsilon^*_n - 1) = A_1$, respectively.

The cubic case $\ell = 3$ becomes more cumbersome, because the function $h_1(\varepsilon)$ is not periodic in $\ln \varepsilon$, due to the oscillating quantities $b_n, b_{n+1}$ in (36); notice that $1 - \delta < \frac{b_n + 1}{b_n} < 1 + \delta$. Nevertheless, we are going to obtain a periodic upper bound for the function $h_1(\varepsilon)$. Let us introduce the functions
\[ g^+_{n}(\varepsilon) := G(\varepsilon; \varepsilon^+_n, 1 + \delta), \quad \varepsilon^+_n := \frac{D_0}{K_0} \cdot \frac{1}{(1 + \delta) \lambda^{3n}}, \]
obtained from (29), by replacing the oscillatory factors $b_n$ by the constant $1 + \delta$ (and taking $\ell = 3$). One can check that the graphs of the functions $g^+_n(\varepsilon)$ and $g^+_n(\varepsilon)$ have no intersection if $b_n < 1 + \delta$, and coincide if $b_n = 1 + \delta$, which implies that $g^*_n(\varepsilon) \leq g^+_n(\varepsilon)$ for any $\varepsilon$. As in (31), we can define
\[ h^+_1(\varepsilon) := \min_{n \geq 0} g^+_n(\varepsilon), \]
and it is clear that we have the upper bound $h_1(\varepsilon) \leq h^+_1(\varepsilon)$ for any $\varepsilon$. By a similar argument to the one used in the quadratic case, we can establish the periodicity in $\ln \varepsilon$ of the function $h^+_1(\varepsilon)$, and we find its maximum value. Indeed, one can check from the expression (26), with $\ell = 3$, that the graphs of the functions $g^+_n$ and $g^+_n$ intersect at
\[ \varepsilon'_n := \left( \frac{\sqrt{\lambda} + 1}{2 \lambda} \right)^2 \cdot \varepsilon^+_n, \quad \text{with} \quad g^+_n(\varepsilon'_n) = g^+_n(\varepsilon'_n) = A^+_1. \]
Thus, we have $h^+_1(\varepsilon) = g^+_n(\varepsilon)$ for $\varepsilon \in [\varepsilon'_n, \varepsilon'_n - 1]$, and we can extend it as a $3 \ln \lambda$-periodic function in $\ln \varepsilon$, whose maximum value is $A^+_1$, which provides an upper bound for $h_1(\varepsilon)$. On the other hand, since $g^+_n(\varepsilon) \geq b_n^{1/3}$ it is clear that $h_1(\varepsilon) \geq A_0^+$. □

**Remark.** In general, for the cubic case, the function $h_1(\varepsilon)$ is not periodic in $\ln \varepsilon$, but we can conjecture that it is quasiperiodic due to the oscillating quantities $b_n$ introduced in (30). In fact, using the approximation (24) for the angle $\phi$, we have $b_{n+22} \approx b_n$, which gives $66 \ln \lambda$ as an approximate period for $h_1(\varepsilon)$ (see Figure 3).

As said in Section 3.1, the function $h_1(\varepsilon)$ is related to the dominant harmonic among the primary resonances, corresponding to the integer vector $s_0(N)$, with $N = N(\varepsilon)$. In order to ensure that this harmonic provides the maximal splitting distance, we need that $h_1(\varepsilon) \leq g_\ell(\varepsilon)$ also for secondary harmonics $k$. To have this inequality, the separation between the primary and secondary resonances has to be large enough. Recalling that the separations $B_0, B_\ell^-$ for quadratic and cubic
frequencies were defined in (16) and (23), respectively, we impose the "separation condition":
\[ B_0 \geq A_1 \quad \text{in the quadratic case;} \]
\[ B_0^- \geq A_1^+ \quad \text{in the cubic case.} \]  
A numerical exploration of this condition, among quadratic frequency vectors given by purely periodic continued fractions, indicates that the 24 cases considered in (4) are all the ones satisfying it. On the other hand, the condition is fulfilled in the case of the cubic golden vector considered in (5), since \[ B_0^- \approx 1.1824 \text{ and } A_1^+ \approx 1.0909, \]
as can be checked from the numerical data given in Section 2.4.

Notice that, under the separation condition (37), we have \( h_2(\varepsilon) < h_1(\varepsilon) \), unless \( \varepsilon \) is very close to some concrete values, at which there is a change in the dominant harmonic given by \( N = N(\varepsilon) \) (in the quadratic case, this happens for \( \varepsilon \) close to the geometric sequence \( \varepsilon_n' \)); see Figure 2.

Now we can complete the proof of Theorem 1. Indeed, according to Proposition 4, the separation condition implies that for any \( \varepsilon \) sufficiently small, we have the estimate
\[ \max_{\theta \in \mathbb{T}} |M(\theta)| \sim |M_{s_0(N)}|, \]
since the coefficient of the dominant harmonic \( k = s_0(N) \), \( N = N(\varepsilon) \), is greater or equal than the sum of all other harmonics.

References


