# A FUNCTIONAL ANALYSIS APPROACH TO ARNOLD DIFFUSION ${ }^{\text {N }}$ 

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Abstract. - We discuss in the context of nearly integrable Hamiltonian systems a functional analysis approach to the "splitting of separatrices" and to the "shadowing problem". As an application we apply our method to the problem of Arnold Diffusion for nearly integrable partially isochronous systems improving known results. © 2002 Éditions scientifiques et médicales Elsevier SAS
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RÉSumé. - Nous discutons dans le contexte des systèmes hamiltoniens presque intégrables une approche basée sur l'analyse fonctionnelle des problèmes de "splitting" des séparatrices et de "shadowing". Nous appliquons notre méthode au problème de la diffusion d'Arnold pour des systèmes presque intégrables partiellement isochrones, améliorant des résultats connus. © 2002 Éditions scientifiques et médicales Elsevier SAS

## 1. Introduction

Topological instability of action variables in multidimensional nearly integrable Hamiltonian systems is known as Arnold Diffusion. This terminology was introduced by Chirikov in [17], years after Arnold discovered this phenomenon in his famous paper [3]. For autonomous Hamiltonian systems with two degrees of freedom KAM theory generically implies topological stability of the action variables (i.e. the time-evolution of the action variables for the perturbed system stay close to their initial values for all times). On the contrary, for systems with more than two degrees of freedom, outside a wide range of initial conditions (the so-called "Kolmogorov set" provided by KAM theory), the action variables may undergo a drift of order one in a very long, but finite time called the "diffusion time".

[^0]After thirty years from Arnold's seminal work [3], attention to Arnold diffusion has been renewed by [15], followed by several papers (see, e.g., [8,14,24,28,45] and references therein).

The Hamiltonian models which are usually studied (as suggested by normal form theory near simple resonances) have the form

$$
\begin{equation*}
H(I, \varphi, p, q)=\frac{1}{2} I_{1}^{2}+\omega \cdot I_{2}+\frac{1}{2} p^{2}+\varepsilon(\cos q-1)+\varepsilon \mu f(I, \varphi, p, q) \tag{1.1}
\end{equation*}
$$

where $\varepsilon$ and $\mu$ are small parameters (the "natural" order for $\mu$ being $\varepsilon^{d}$ for some positive $d) ;\left(I_{1}, I_{2}, p\right)$ and $(\varphi, q)$ are standard symplectic action-angle variables $\left(I_{i} \in \mathbf{R}^{n_{i}}\right.$, $n_{1}+n_{2}=n, \varphi=\left(\varphi_{1}, \varphi_{2}\right) \in \mathbf{T}^{n},(p, q) \in \mathbf{R} \times \mathbf{T}, \mathbf{T}$ being the standard torus $\left.\mathbf{R} / 2 \pi \mathbf{Z}\right)$. In Arnold's model $I_{1}, I_{2} \in \mathbf{R}, \omega=1, f(I, \varphi, p, q)=(\cos q-1)\left(\sin \varphi_{1}+\cos \varphi_{2}\right)$ and in [3] diffusion is proved for $\mu$ exponentially small w.r.t. $\sqrt{\varepsilon}$. Physically (1.1) describes a system of $n_{1}$ "rotators" and $n_{2}$ harmonic oscillators weakly coupled with a pendulum through a perturbation term.

The existence of Arnold diffusion is usually proved following the mechanism proposed in [3]. For $\mu=0$, Hamiltonian $H$ admits a continuous family of $n$-dimensional partially hyperbolic invariant tori $\mathcal{T}_{I}$ possessing stable and unstable manifolds $W_{0}^{s}\left(\mathcal{T}_{I}\right)=$ $W_{0}^{u}\left(\mathcal{T}_{I}\right)$, called "whiskers" by Arnold. Arnold's mechanism is then based on the following three main steps.
(i) For $\mu \neq 0$ small enough, the perturbed stable and unstable whiskers $W_{\mu}^{s}\left(\mathcal{T}_{I}^{\mu}\right)$ and $W_{\mu}^{u}\left(\mathcal{T}_{I}^{\mu}\right)$ split and intersect transversally ("splitting of the whiskers");
(ii) Prove the existence of a chain of "transition" tori connected by heteroclinic orbits ("transition chain");
(iii) Prove the existence of an orbit, "shadowing" the transition chain, for which the action variables $I$ undergo a variation of $\mathrm{O}(1)$ in a certain time $T_{d}$ called the diffusion time.
The shadowing problem (iii) has been extensively studied in the last years by geometrical (see, e.g., [15,18-20,24,37]) and by variational methods (see, e.g., [8,9,11, 12]). More rich and older literature is available on the splitting problem see, e.g., [34,33, $23,15,29,21,25,26,30,43,41,42,31,36,38,39]$ and references therein.

The aim of this paper is to provide a functional analysis approach apt to deal with Arnold diffusion, especially with "splitting" (i) and "shadowing" (iii) problems.

Rather than formulating our results in an abstract setting we shall illustrate the method (in a complete and self-contained way) on a relatively simple class of models, namely harmonic oscillators weakly coupled with a pendulum through purely spatial perturbations.

In this context defining a transition chain (ii) is a straighforward consequence of the "splitting of the whiskers" (i) since all the invariant tori are preserved by the perturbation, being just slightly deformed. This also happens, for the peculiar choice of the perturbation, in the non-isochronous system considered in [3]. Obviously this is not the case for general non-isochronous systems where the surviving perturbed tori are separated by the gaps appearing in KAM constructions, making the existence of chains of tori a more difficult matter, see [15]. We also refer to [45] for a somewhat different
mechanism of diffusion where step (ii) is bypassed for systems with three degrees of freedom using Mather theory.

We now give a rough description of the main results of this paper. We consider nearly integrable partially isochronous Hamiltonian systems given by

$$
\begin{equation*}
\mathcal{H}_{\mu}=\omega \cdot I+\frac{p^{2}}{2}+(\cos q-1)+\mu f(\varphi, q) \tag{1.2}
\end{equation*}
$$

where $(\varphi, q) \in \mathbf{T}^{n} \times \mathbf{T}^{1}:=\left(\mathbf{R}^{n} / 2 \pi \mathbf{Z}^{n}\right) \times(\mathbf{R} / 2 \pi \mathbf{Z})$ are the angle variables, $(I, p) \in$ $\mathbf{R}^{n} \times \mathbf{R}^{1}$ are the action variables and $\mu \geqslant 0$ is a small real parameter. When $\mu=0$ the energy $\omega_{i} I_{i}$ of each oscillator is a constant of the motion. The unperturbed Hamitonian possesses $n$-dimensional invariant tori $\mathcal{I}_{I_{0}}=\left\{(\varphi, I, q, p) \in \mathbf{T}^{n} \times \mathbf{R}^{n} \times \mathbf{T}^{1} \times \mathbf{R}^{1} \mid I=\right.$ $\left.I_{0}, q=p=0\right\}$ with stable and unstable manifolds $W^{s}\left(\mathcal{T}_{I_{0}}\right)=W^{u}\left(\mathcal{T}_{I_{0}}\right)=\{(\varphi, I, q, p) \in$ $\left.\mathbf{T}^{n} \times \mathbf{R}^{n} \times \mathbf{T}^{1} \times \mathbf{R}^{1} \mid I=I_{0}, p^{2} / 2+(\cos q-1)=0\right\}$. The problem of Arnold diffusion in this context is whether, for $\mu \neq 0$, there exist motions whose net effect is to transfer $\mathrm{O}(1)$-energy from one oscillator to the others. In order to exclude trivial drifts of the actions due to resonance phenomena, it is standard to assume a diophantine condition for the frequency vector $\omega$. Precisely we will always suppose that $\omega$ is $(\gamma, \tau)$-diophantine, i.e.

$$
\begin{equation*}
\exists \gamma>0, \tau \geqslant n-1 \text { such that }|\omega \cdot k| \geqslant \gamma /|k|^{\tau}, \forall k \in \mathbf{Z}^{n}, k \neq 0 \tag{H1}
\end{equation*}
$$

In the present paper we first prove general shadowing theorems (Theorems 2.3 and 3.2) which improve -for isochronous systems- known estimates on the diffusion time. Secondly we discuss a new method for the splitting of the whiskers, providing general estimates on the Fourier coefficients of some "splitting function" (see Theorem 4.2). As applications we consider the following two cases
(a) the frequencies of the harmonic oscillators form a diophantine vector $\omega$ of order 1 ("a priori-unstable case");
(b) the frequencies of the harmonic oscillators form a diophantine vector $\omega_{\varepsilon}=$ $\left(1 / \sqrt{\varepsilon}, \beta \varepsilon^{a}\right)$ with $a \geqslant 0, \mu \varepsilon^{-3 / 2}$ small and the perturbation $f(\varphi, q)=(1-$ $\cos q) f(\varphi)$ ("three-time-scales problem" with perturbations preserving all the unperturbed invariant tori). This corresponds, after a time rescaling, to $\omega=$ $\left(1, \beta \varepsilon^{a} \sqrt{\varepsilon}\right)$ in (1.1).
Case (a) highlights the improvement of our estimates on diffusion times. In this case it is easy to show, using the classical Poincaré-Melnikov function

$$
M(A)=\int_{\mathbf{R}}\left[f(\omega t+A, 0)-f\left(\omega t+A, q_{0}(t)\right)\right] \mathrm{d} t
$$

$q_{0}(t)=4 \arctan (\exp t)$ being the unperturbed separatrix of the pendulum, that the splitting of the whiskers is $\mathrm{O}(\mu)$. Then our shadowing method yields (see for a more precise statement Theorems 2.4 and 3.3)

THEOREM 1.1. - Assume ( $H 1$ ) and let M possess a proper minimum (or maximum) $A_{0}$. Then, for $\mu$ small enough, there exist orbits whose action variables undergo a drift of order one, with diffusion time $T_{d}=\mathrm{O}((1 / \mu) \log (1 / \mu))$.

Theorem 1.1 answers a question raised in [35] (Section 7) proving that, at least for isochronous systems, it is possible to reach the maximal speed of diffusion $\mu /|\log \mu|$; for related results compare also with [20] and [44]. We recall that the estimate on the diffusion time obtained in [15] is $T_{d} \gg \mathrm{O}(\exp (1 / \mu))$ and that in [24] it is improved to be $T_{d}=\mathrm{O}(\exp (1 / \mu))$; recently in [9], by means of Mather theory, the estimate on the diffusion time has been improved to be $T_{d}=\mathrm{O}\left(1 / \mu^{2 \tau+1}\right)$; in [19] it is obtained via geometric methods that $T_{d}=\mathrm{O}\left(1 / \mu^{\tau+1}\right)$. It is worth pointing out that the estimates given in [9] and [19], while providing a diffusion time polynomial in the splitting, depend on the diophantine exponent $\tau$ and hence on the number of rotators $n$. Instead our estimate (as well as that discussed in [20]) does not depend upon the number of degrees of freedom.

The three-time-scales system illustrates an application of our estimate of the splitting (Theorem 5.1), which, together with our general shadowing Theorem 2.3, yields Arnold diffusion in case (b). Roughly (see Theorem 5.2 for a precise statement) we get

THEOREM 1.2. - Under a suitable non degeneracy condition on $f$, for all $\varepsilon$ small such that $\omega_{\varepsilon}$ is $\left(\gamma_{\varepsilon}, \tau\right)$ diophantine, for $\mu \varepsilon^{-3 / 2}$ small enough, stable and unstable manifolds split and there exist orbits whose action variables undergo a drift of order one with diffusion time $T_{d}=\mathrm{O}\left((\sqrt{\varepsilon} / \mu) \mathrm{e}^{\pi /(2 \sqrt{\varepsilon})}\left[\left(\gamma_{\varepsilon}\right)^{-1}\left(\sqrt{\varepsilon} \mathrm{e}^{\pi /(2 \sqrt{\varepsilon})}\right)^{\tau}+\right.\right.$ $|\log \mu|]$ ).

The choice of three time scales frequencies enables to derive lower estimates for the splitting relatively easily. Incidentally we mention that Theorem 5.1 improves the main Theorem I in [38] which holds for $\mu=\varepsilon^{p}, p>2+a$; w.r.t. [25] (which deals with more general systems) we remark that our results hold in any dimension, while the results of [25], based on tree techniques and cancellations, are proved for $n=2$.

We also mention that, for $n \geqslant 3$, with a more careful shadowing analysis (exploiting the anisotropy of this splitting) we can prove that, along special directions, Arnold diffusion takes place in polynomial time w.r.t $1 / \varepsilon$ (see [7] and Remark 5.2).

Furthermore we also remark that our proof of Theorem 3.2 is completely selfcontained in the sense that, unlike the known approaches (excepted [45]), we do not make use of any KAM-type result for proving, under assumption (H1), the persistence of invariant tori, see Theorem 3.1.

The results of this paper have been announced in [6].
We now describe the functional analysis approach developed in this paper to prove both the results on the shadowing theorem and on the "splitting of the whiskers". It is based on a finite dimensional reduction of Lyapunov-Schmidt type, variational in nature, introduced in [2] and in [1], and later extended in [4,5] in order to construct shadowing orbits. For simplicity we describe our approach when the perturbation term $f(\varphi, q)=(1-\cos q) f(\varphi)$ so that the tori $\mathcal{T}_{I_{0}}$ are still invariant for $\mu \neq 0$. The equations of motion derived by Hamiltonian $\mathcal{H}_{\mu}$ are

$$
\begin{equation*}
\dot{\varphi}=\omega, \quad \dot{I}=-\mu(1-\cos q) \nabla f(\varphi), \quad \dot{q}=p, \quad \dot{p}=\sin q-\mu \sin q f(\varphi) \tag{1.4}
\end{equation*}
$$

The dynamics on the angles $\varphi$ is given by $\varphi(t)=\omega t+A$ so that (1.4) are reduced to the quasi-periodically forced pendulum equation

$$
\begin{equation*}
-\ddot{q}+\sin q=\mu \sin q f(\omega t+A) \tag{1.5}
\end{equation*}
$$

corresponding to the Lagrangian

$$
\begin{equation*}
\mathcal{L}_{\mu, A}(q, \dot{q}, t)=\frac{\dot{q}^{2}}{2}+(1-\cos q)+\mu(\cos q-1) f(\omega t+A) \tag{1.6}
\end{equation*}
$$

For each solution $q(t)$ of (1.5) one recovers the dynamics of the actions $I(t)$ by quadratures in (1.4).

For $\mu=0$ Eq. (1.5) is autonomous and possesses the one parameter family of homoclinic solutions $(\bmod .2 \pi) q_{\theta}(t)=4 \arctan (\exp (t-\theta)), \theta \in \mathbf{R}$. Consider the Lagrangian action functional $\Phi_{\mu, A}: q_{0}+H^{1}(\mathbf{R}) \rightarrow \mathbf{R}$ associated to the quasiperiodically forced pendulum (1.5)

$$
\begin{equation*}
\Phi_{\mu, A}(q):=\int_{\mathbf{R}} \mathcal{L}_{\mu, A}(q(t), \dot{q}(t), t) \mathrm{d} t \tag{1.7}
\end{equation*}
$$

$\Phi_{\mu, A}$ is smooth on $q_{0}+H^{1}(\mathbf{R})$ and critical points $q$ of $\Phi_{\mu, A}$ are homoclinic solutions to $0, \bmod .2 \pi$, of (1.5). These critical points $q$ are in fact smooth functions of the time $t$ and are exponentially decaying to 0 , mod. $2 \pi$, as $|t| \rightarrow+\infty$.

The unperturbed functional $\Phi_{0}:=\Phi_{0, A}$ does not depend on $A$ and possesses the 1-dimensional manifold of critical points $Z:=\left\{q_{\theta} \mid \theta \in \mathbf{R}\right\}$ with tangent space at $q_{\theta}$ spanned by $\dot{q}_{\theta}$. All the unperturbed critical points $q_{\theta}$ are degenerate since $d^{2} \Phi_{0}\left(q_{\theta}\right)\left[\dot{q}_{\theta}\right]=0$. However $q_{\theta}$ are non-degenerate critical points of the restriction $\Phi_{0 \mid E_{\theta}}$ for any subspace $E_{\theta}$ supplementary to $\left\langle\dot{q}_{\theta}\right\rangle$. It is then possible to apply a LyapunovSchmidt type reduction, based on the Implicit Function Theorem, to find near $q_{\theta}$, for $\mu$ small, critical points $q_{A, \theta}^{\mu}$ of $\Phi_{\mu, A}$ restricted to $E_{\theta}$; more precisely $q_{A, \theta}^{\mu}=q_{\theta}+w_{A, \theta}^{\mu}$ with $w_{A, \theta}^{\mu} \in E_{\theta},\left\|w_{A, \theta}^{\mu}\right\|=\mathrm{O}(\mu)$ and d $\Phi_{\mu, A}\left(q_{A, \theta}^{\mu}\right)_{\mid E_{\theta}}=0$. We will call the functions $q_{A, \theta}^{\mu, \theta}$ "1-bump pseudo-homoclinic solutions" of the quasi-periodically forced pendulum (1.5).

It turns out that the 1-dimensional manifold $Z_{\mu}=\left\{q_{A, \theta}^{\mu} \mid \theta \in \mathbf{R}\right\}$ is a "natural constraint" for the action functional $\Phi_{\mu, A}$, namely any critical point of $\Phi_{\mu, A \mid Z_{\mu}}$ is a critical point of $\Phi_{\mu, A}$, and hence a true solution of (1.5) homoclinic to $0(\bmod .2 \pi)$.

Through this paper we will perform the above finite dimensional reduction using two different supplementary spaces to $\left\langle\dot{q}_{\theta}\right\rangle$ : one is better suited for the shadowing arguments, the other is better suited for studying the splitting problem in presence of "high frequencies".

Shadowing. In Sections 2 and 3, dealing with the shadowing theorem, we choose

$$
\begin{equation*}
E_{\theta}=\{w: \mathbf{R} \rightarrow \mathbf{R} \mid w(\theta)=0\} \tag{1.8}
\end{equation*}
$$

$E_{\theta}$ and $\left\langle\dot{q}_{\theta}\right\rangle$ are supplementary spaces since $\dot{q}_{0}(0) \neq 0$. The corresponding "reduced action functional" $F_{\mu}(A, \theta):=\Phi_{\mu, A}\left(q_{A, \theta}^{\mu}\right)$, see (2.1), turns out to have a neat geometric
meaning, see Remark 2.1, related to the generating functions of the exact Lagrangian stable and unstable manifolds $W^{s, u}\left(\mathcal{T}_{I_{0}}\right)$.

From a geometrical point of view, the choice of the supplementary space defined by (1.8) means considering $W^{s}\left(\mathcal{T}_{I_{0}}\right)$ and $W^{u}\left(\mathcal{T}_{I_{0}}\right)$ at the fixed Poincaré section $\{q=\pi\}$, since $q_{A, \theta}^{\mu}(\theta)=q_{\theta}(\theta)+w_{A, \theta}^{\mu}(\theta)=q_{0}(0)=\pi$.

The supplementary space $E_{\theta}$ is very well suited to perform the shadowing theorem because the " 1 -bump pseudo solutions" $q_{A, \theta}^{\mu}(t)$ are true solutions of (1.5) except at the instant $t=\theta$ where $\dot{q}_{A, \theta}^{\mu}(t)$ may have a jump, see Lemma 2.1. Following [4,5] we will generalize in Section 2.2 the above finite dimensional reduction, finding a natural constraint for " $k$-bump pseudo homoclinic solutions" turning $k$ times near the unperturbed separatrices of the pendulum, see Lemma 2.4.

In this way under a suitable "splitting condition", satisfied for instance if $G_{\mu}(A):=$ $F_{\mu}(A, 0)$ possesses a proper minimum (see Condition 2.1 ), we can prove a general shadowing theorem with explicit estimates on the diffusion time $T_{d}$, see Theorem 2.3. This estimate (see expression (2.21)), is roughly the following: $T_{d}$ is estimated by the product of the number of heteroclinic transitions $k$ ( $=$ number of tori forming the transition chain $=$ heteroclinic jump/splitting) and of the time $T_{s}$ required for a single transition, namely $T_{d}=k T_{s}$. The time for a single transition $T_{s}$ is bounded by the maximum time between the "ergodization time" of the torus $\mathbf{T}^{n}$ run by the linear flow $\omega t$, and the time needed to "shadow" homoclinic orbits for the quasi-periodically forced pendulum.

Theorem 1.1 (for a priori unstable systems) is a straightforward consequence of the general shadowing theorem, noting that, by the classical Poincaré-Melnikov theory, the splitting of the whiskers is $\mathrm{O}(\mu)$. The main reason for which Theorem 1.1 improves the polynomial estimates $T_{d}=\mathrm{O}\left(1 / \mu^{2 \tau+1}\right)$ and $T_{d}=\mathrm{O}\left(1 / \mu^{\tau+1}\right)$, obtained respectively in [9] and [19], is that our shadowing orbit can be chosen, at each transition, to approach the homoclinic point only up to a distance $\mathrm{O}(1)$ and not $\mathrm{O}(\mu)$ like in [9] and [19]. This implies that the time spent by our diffusion orbit at each transition is $T_{s}=\mathrm{O}(\log (1 / \mu))$. Since the number of tori forming the transition chain is equal to $\mathrm{O}(1 /$ splitting $)=\mathrm{O}(1 / \mu)$ the diffusion time is finally estimated by $T_{d}=\mathrm{O}((1 / \mu) \log (1 / \mu))$.

As well, Theorem 1.2 (and Theorem 5.2) for three-time-scales systems, is a consequence of the general shadowing Theorem 2.3, but in this situation the splitting is by far less easy to measure, because it is exponentially small.

As mentioned above variational methods in the context of Arnold diffusion have been already used in [8], see Remark 2.2 for comments on that. One possible advantage of our approach is that it may be used to consider more general critical points of the reduced functional, not only minima. Another advantage is that the same shadowing arguments can be used also when the hyperbolic part is a general Hamiltonian in $\mathbf{R}^{2 m}, m \geqslant 1$, possessing one hyperbolic equilibrium and a transversal homoclinic orbit. Nevertheless we have developed all the details when the hyperbolic part is the standard one-dimensional pendulum because it is the model equation to study Arnold diffusion near a simple-resonance.

Splitting. Detecting and measuring the splitting of the whiskers is a difficult problem when the frequency vector $\omega=\omega_{\varepsilon}$ depends on some small parameter $\varepsilon$ and
contains some "fast frequencies" $\omega_{i}=\mathrm{O}\left(1 / \varepsilon^{b}\right), b>0$. Indeed, in this case, the variations of the Melnikov function along some directions turn out to be exponentially small with respect to $\varepsilon$ and then the naive Poincaré-Melnikov expansion provides a valid measure of the splitting only for $\mu$ exponentially small with respect to some power of $\varepsilon$.

The typical argument to estimate exponentially small splittings, used virtually in all papers dealing with this problem, goes back to Arnold himself and is based on Fourier analysis on complex domains.

For this reason we would like to extend analytically the "reduced action functional" $F_{\mu}(A, \theta)=\Phi_{\mu, A}\left(q_{A, \theta}^{\mu}\right)$ in a complex strip sufficiently wide in the $\theta$ variable. However $F_{\mu}(A, \theta)$ can not be easily analytically extended. Indeed, for $\theta$ complex, the supplementary space $E_{\theta}=\{w: \mathbf{R} \rightarrow \mathbf{C} \mid w(\operatorname{Re} \theta)=0\}$, appearing naturally when we try to extend the definition of $q_{A, \theta}^{\mu}$ to $\theta \in \mathbf{C}$, does not depend analytically on $\theta$. This breakdown of analyticity, arising when measuring the "splitting of the whiskers" at the fixed Poincare section $\{q=\pi\}$, is a well known difficulty and has been compensated in $[15,25,26]$ via the introduction of tree techniques which enable to prove cancellations in the power series expansions.

Our method to overcome this "loss of analyticity" is different and relies on the introduction of another supplementary space $\widetilde{E}_{\theta}$, which depends analytically on $\theta$. It is defined by

$$
\widetilde{E}_{\theta}=\left\{w: \mathbf{R} \rightarrow \mathbf{R} \mid \int_{\mathbf{R}} \psi_{\theta}(t) w(t) \mathrm{d} t=0\right\}
$$

where $\psi_{0}(t)=\cosh ^{2}(t) /(1+\cosh t)^{3} . \widetilde{E}_{\theta}$ and $\left\langle\dot{q}_{\theta}\right\rangle$ are supplementary spaces since $\int_{\mathbf{R}} \psi_{0}(t) \dot{q}_{0} \mathrm{~d} t \neq 0$ (see Remark 4.1 for the choice of $\psi_{0}$ ). The corresponding reduced functional $\widetilde{F}_{\mu}(A, \theta):=\Phi_{\mu, A}\left(Q_{A, \theta}^{\mu}\right)$, defined in (4.4), where $Q_{A, \theta}^{\mu}$ are the corresponding "1-bump pseudo-homoclinics solutions", can be analytically extended in a sufficiently large complex strip. This enables to find easily the exponentially small bounds on the Fourier coefficients for the splitting, thanks to the important invariance property $\widetilde{F}_{\mu}(A, \theta+\eta)=\widetilde{F}_{\mu}(A+\omega \eta, \theta)$, which is a consequence of the autonomy of Hamiltonian $\mathcal{H}_{\mu}$. This type of Lyapunov-Schmidt reduction was used in [2] to study the exponentially small splitting in rapidly periodically forced systems.

The crucial point is now to observe that "reduced action functionals" corresponding to different choices of the supplementary space are equivalent: it turns out (Theorem 4.1) that the reduced functionals $F_{\mu}$ and $\widetilde{F}_{\mu}$ are simply the same up to a change of variables close to the identity,

$$
F_{\mu}(A, \theta)=\widetilde{F}_{\mu}\left(A, \theta+h_{\mu}(A, \theta)\right), \quad h_{\mu}=\mathrm{O}(\mu)
$$

This fact enables to transpose the informations on $\widetilde{F}_{\mu}$ to $F_{\mu}$ and viceversa. We point out that the proof of Theorem 4.1 relies only on the standard Implicit Function Theorem. The introduction of $\widetilde{F}_{\mu}$ may be interpreted simply as measuring the splitting with a non constant Poincaré section.

A different approach for establishing exponentially small splittings for a rapidly forced pendulum has been used by [23] and recently extended to the setting considered here in [21] and $[38,39]$ by the mean of normal form theory. For the more difficult case
of the standard map, where the Melnikov method fails to give the correct prediction for the splitting, we refer to the pioneering paper [34] and to [30].

We also mention that the recent papers [43] and [36], even though quite different in spirit from ours, have also several technical similarities with our method (the use of generating functions, nonfixed Poincaré section, ...). We add that these papers deal also with non-isochronous systems.

The paper is organized as follows: in Section 2 we prove the shadowing theorem when the perturbation term is $f(\varphi, q)=(1-\cos q) f(\varphi)$. In Section 3 we show how to prove the shadowing theorem for general perturbation terms $f(\varphi, q)$. In Section 4 we provide the theorem on the Fourier coefficients of the splitting and in Section 5 we consider three-time-scales systems.

## 2. The shadowing theorem

### 2.1. 1-bump homoclinic and heteroclinic solutions

Using the Implicit Function Theorem we now prove that, near the unperturbed homoclinic solutions $q_{\theta}(t)$, there exist, for $\mu$ small enough, "1-bump pseudo-homoclinic solutions" $q_{A, \theta}^{\mu}(t)$ of Eq. (1.5). $q_{A, \theta}^{\mu}(t)$ are true solutions of (1.5) in $(-\infty, \theta)$ and $(\theta,+\infty)$; at time $t=\theta$ such pseudo-solutions are glued with continuity at value $q_{A, \theta}^{\mu}(\theta)=\pi$ and for $t \rightarrow \pm \infty$ are asymptotic to the equilibrium $0 \bmod 2 \pi$. The $q_{A, \theta}^{\mu}$ are critical points for $\Phi_{\mu, A}$ restricted to $q_{\theta}+E_{\theta}$ ( $E_{\theta}$ is defined in (1.8)).

LEMMA 2.1. - There exist $\mu_{0}, C_{0}>0$ such that $\forall 0<\mu \leqslant \mu_{0}, \forall \omega \in \mathbf{R}^{n}, \forall \theta \in \mathbf{R}$, there exists a unique function $q_{A, \theta}^{\mu}(t): \mathbf{R} \rightarrow \mathbf{R}$, smooth in $(A, \mu)$, such that
(i) $\quad q_{A, \theta}^{\mu}(t)$ is a solution of $(1.5)$ in each interval $(-\infty, \theta)$ and $(\theta,+\infty)$ and

$$
q_{A, \theta}^{\mu}(\theta)=\pi
$$

$$
\begin{equation*}
\max \left(\left|q_{A, \theta}^{\mu}(t)-q_{\theta}(t)\right|,\left|\dot{q}_{A, \theta}^{\mu}(t)-\dot{q}_{\theta}(t)\right|\right) \leqslant C_{0} \mu \exp \left(-\frac{|t-\theta|}{2}\right), \quad \forall t \in \mathbf{R} \tag{ii}
\end{equation*}
$$

$$
\begin{equation*}
q_{A, \theta}^{\mu}(t)=q_{A+k 2 \pi, \theta}^{\mu}(t), \quad \forall k \in \mathbf{Z}^{n} \tag{iii}
\end{equation*}
$$

$$
\begin{equation*}
q_{A, \theta+\eta}^{\mu}(t)=q_{A+\omega \eta, \theta}^{\mu}(t-\eta), \quad \forall \theta, \eta \in \mathbf{R} \tag{iv}
\end{equation*}
$$

$$
\begin{align*}
& \max \left(\left|\partial_{A} q_{A, \theta}^{\mu}(t)\right|,\left|\partial_{A} \dot{q}_{A, \theta}^{\mu}(t)\right|,\left|\omega \cdot \partial_{A} q_{A, \theta}^{\mu}(t)\right|,\left|\omega \cdot \partial_{A} \dot{q}_{A, \theta}^{\mu}(t)\right|\right)  \tag{v}\\
& \quad \leqslant C_{0} \mu \exp \left(-\frac{|t-\theta|}{2}\right)
\end{align*}
$$

## Proof. - In Appendix A.

We can then define the function $F_{\mu}: \mathbf{T}^{n} \times \mathbf{R} \rightarrow \mathbf{R}$ as the action functional of Lagrangian (1.6) evaluated on the 1-dimensional manifold $Z_{\mu}:=\left\{q_{A, \theta}^{\mu} \mid \theta \in \mathbf{R}\right\}$ of "1-bump pseudo-homoclinic solutions", namely

$$
F_{\mu}(A, \theta)=\Phi_{\mu, A}\left(q_{A, \theta}^{\mu}\right)
$$

$$
\begin{equation*}
=\int_{-\infty}^{\theta} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{\mu}(t), \dot{q}_{A, \theta}^{\mu}(t), t\right) \mathrm{d} t+\int_{\theta}^{+\infty} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{\mu}(t), \dot{q}_{A, \theta}^{\mu}(t), t\right) \mathrm{d} t \tag{2.1}
\end{equation*}
$$

and the "homoclinic function" $G_{\mu}: \mathbf{T}^{n} \rightarrow \mathbf{R}$ as

$$
\begin{equation*}
G_{\mu}(A)=F_{\mu}(A, 0) \tag{2.2}
\end{equation*}
$$

Since $q_{A, \theta}^{\mu}(t)$ converges exponentially fast to $0, \bmod 2 \pi$, the integrals in (2.1) are convergent.

By property (iv) of Lemma 2.1 the following invariance property holds

$$
\begin{equation*}
F_{\mu}(A, \theta+\eta)=F_{\mu}(A+\omega \eta, \theta), \quad \forall \theta, \eta \in \mathbf{R} \tag{2.3}
\end{equation*}
$$

and in particular

$$
\begin{equation*}
F_{\mu}(A, \theta)=G_{\mu}(A+\omega \theta), \quad \forall \theta \in \mathbf{R} \tag{2.4}
\end{equation*}
$$

Remark 2.1. - The homoclinic function $G_{\mu}$ is the difference between the generating functions $\mathcal{S}_{\mu, I_{0}}^{ \pm}\left(A, q_{0}\right)$ of the stable and the unstable manifolds $W_{\mu}^{s, u}\left(\mathcal{T}_{I_{0}}\right)$ (which in this case are exact Lagrangian manifolds) at section $q_{0}=\pi$, namely $G_{\mu}(A)=\mathcal{S}_{\mu, I_{0}}^{-}(A, \pi)-$ $\mathcal{S}_{\mu, I_{0}}^{+}(A, \pi)$. Indeed it can be easily verified that

$$
\begin{aligned}
& \mathcal{S}_{\mu, I_{0}}^{+}\left(A, q_{0}\right) \\
& \quad:=I_{0} \cdot A-\int_{0}^{+\infty} \frac{\left(\dot{q}_{A, q_{0}}^{\mu}(t)\right)^{2}}{2}+\left(1-\cos q_{A, q_{0}}^{\mu}(t)\right)+\mu\left(\cos q_{A, q_{0}}^{\mu}(t)-1\right) f(\omega t+A) \mathrm{d} t
\end{aligned}
$$

where $q_{A, q_{0}}^{\mu}(t)$ is the unique solution of (1.5) near $q_{0}(t)$ with $q_{A, q_{0}}^{\mu}(0)=q_{0}$ and $\lim _{t \rightarrow+\infty} q_{A, q_{0}}^{\mu}(t)=2 \pi$. Analogously

$$
\begin{aligned}
& \mathcal{S}_{\mu, I_{0}}^{-}\left(A, q_{0}\right) \\
& \quad:=I_{0} \cdot A+\int_{-\infty}^{0} \frac{\left(\dot{q}_{A, q_{0}}^{\mu}(t)\right)^{2}}{2}+\left(1-\cos q_{A, q_{0}}^{\mu}(t)\right)+\mu\left(\cos q_{A, q_{0}}^{\mu}(t)-1\right) f(\omega t+A) \mathrm{d} t
\end{aligned}
$$

where $q_{A, q_{0}}^{\mu}(t)$ is the unique solution of (1.5) near $q_{0}(t)$ with $q_{A, q_{0}}^{\mu}(0)=q_{0}$ and $\lim _{t \rightarrow-\infty} q_{A, q_{0}}^{\mu}(t)=0$.

There holds
LEMMA 2.2. $-F_{\mu}: \mathbf{T}^{n} \times \mathbf{R} \rightarrow \mathbf{R}$ is smooth and

$$
\begin{gather*}
\partial_{A} F_{\mu}(A, \theta)=\int_{\mathbf{R}} \mu\left(\cos q_{A, \theta}^{\mu}(t)-1\right) \partial_{\varphi} f(\omega t+A) \mathrm{d} t,  \tag{2.5}\\
\partial_{\theta} F_{\mu}(A, \theta)=\mu \int_{\mathbf{R}} f(\omega t+A) \sin q_{A, \theta}^{\mu}(t) \dot{q}_{A, \theta}^{\mu}(t) \mathrm{d} t=\frac{\left(\dot{q}_{A, \theta}^{\mu}\right)^{2}\left(\theta^{+}\right)}{2}-\frac{\left(\dot{q}_{A, \theta}^{\mu}\right)^{2}\left(\theta^{-}\right)}{2} . \tag{2.6}
\end{gather*}
$$

Proof. - Since $q_{A, \theta}^{\mu}$ depends smoothly on $A$ so does $F_{\mu}(A, \theta)$ and there holds

$$
\begin{aligned}
\partial_{A} F_{\mu}(A, \theta)= & \int_{-\infty}^{+\infty} \dot{q}_{A, \theta}^{\mu}(t) \partial_{A} \dot{q}_{A, \theta}^{\mu}(t)+\sin q_{A, \theta}^{\mu}(t) \partial_{A} q_{A, \theta}^{\mu}(t) \\
& -\mu \sin q_{A, \theta}^{\mu}(t) f(\omega t+A) \partial_{A} q_{A, \theta}^{\mu}(t)+\mu\left(\cos q_{A, \theta}^{\mu}(t)-1\right) \partial_{\varphi} f(\omega t+A) \mathrm{d} t
\end{aligned}
$$

Integrating by parts and using that $\lim _{t \rightarrow \pm \infty} \dot{q}_{A, \theta}^{\mu}(t)=0$ and $\partial_{A} q_{A, \theta}^{\mu}(\theta)=0$, we get

$$
\begin{aligned}
\int_{-\infty}^{+\infty} \dot{q}_{A, \theta}^{\mu}(t) \partial_{A} \dot{q}_{A, \theta}^{\mu}(t) \mathrm{d} t & =\int_{-\infty}^{\theta} \ldots+\int_{\theta}^{+\infty} \dot{q}_{A, \theta}^{\mu}(t) \partial_{A} \dot{q}_{A, \theta}^{\mu}(t) \mathrm{d} t \\
& =-\int_{-\infty}^{\theta} \ddot{q}_{A, \theta}^{\mu}(t) \partial_{A} q_{A, \theta}^{\mu}(t) \mathrm{d} t-\int_{\theta}^{+\infty} \ddot{q}_{A, \theta}^{\mu}(t) \partial_{A} q_{A, \theta}^{\mu}(t) \mathrm{d} t
\end{aligned}
$$

Hence, since $q_{A, \theta}^{\mu}(t)$ solves (1.5) on each interval $(-\infty, \theta)$ and $(\theta,+\infty)$,

$$
\begin{equation*}
\partial_{A} F_{\mu}(A, \theta)=\int_{-\infty}^{+\infty} \mu\left(\cos q_{A, \theta}^{\mu}(t)-1\right) \partial_{\varphi} f(\omega t+A) \mathrm{d} t \tag{2.7}
\end{equation*}
$$

and (2.5) is proved.
$\partial_{\theta} q_{A, \theta}^{\mu}$ is not defined in $W^{1, \infty}(\mathbf{R})$ because, $\dot{q}_{A, \theta}^{\mu}$ may have a jump at $t=\theta$. However, by the invariance property $(2.3), F_{\mu}(A, \theta)$ is smooth also in $\theta$ and we have that

$$
\begin{aligned}
\partial_{\theta} F_{\mu}(A, \theta) & =\omega \cdot \partial_{A} F_{\mu}(A, \theta) \\
& =\mu \int_{\mathbf{R}}\left(\cos q_{A, \theta}^{\mu}(t)-1\right) \frac{\mathrm{d}}{\mathrm{~d} t} f(\omega t+A) \mathrm{d} t \\
& =\mu \int_{\mathbf{R}} \sin q_{A, \theta}^{\mu}(t) \dot{q}_{A, \theta}^{\mu}(t) f(\omega t+A) \mathrm{d} t,
\end{aligned}
$$

by an integration by parts. Now since $q_{A, \theta}^{\mu}(t)$ solves (1.5) in $(\theta,+\infty)$

$$
\begin{aligned}
\int_{\theta}^{+\infty} \mu \sin q_{A, \theta}^{\mu}(t) \dot{q}_{A, \theta}^{\mu}(t) f(\omega t+A) \mathrm{d} t & =\int_{\theta}^{+\infty}-\ddot{q}_{A, \theta}^{\mu} \dot{q}_{A, \theta}^{\mu}+\sin q_{A, \theta}^{\mu}(t) \dot{q}_{A, \theta}^{\mu}(t) \mathrm{d} t \\
& =\frac{\left(\dot{q}_{A, \theta}^{\mu}\right)^{2}\left(\theta^{+}\right)}{2}+\cos q_{A, \theta}^{\mu}(\theta)
\end{aligned}
$$

In the same way

$$
\int_{-\infty}^{\theta} \mu \sin q_{A, \theta}^{\mu}(t) \dot{q}_{A, \theta}^{\mu}(t) f(\omega t+A) \mathrm{d} t=-\frac{\left(\dot{q}_{A, \theta}^{\mu}\right)^{2}\left(\theta^{-}\right)}{2}-\cos q_{A, \theta}^{\mu}(\theta) .
$$

This proves the lemma.

It follows that $Z_{\mu}:=\left\{q_{A, \theta}^{\mu} \mid \theta \in \mathbf{R}\right\}$ is a "natural constraint" for $\Phi_{\mu, A}$. Indeed, by (2.6) and Lemma 2.1(ii), if $\partial_{\theta} F_{\mu}(A, \theta)=0$ then $q_{A, \theta}^{\mu}(t)$ is a true homoclinic $(\bmod 2 \pi)$ solution of (1.5). Then, for each $I_{0} \in \mathbf{R}^{n}$,

$$
\begin{equation*}
\left(\omega t+A, I_{\mu}(t), q_{A, \theta}^{\mu}(t), \dot{q}_{A, \theta}^{\mu}(t)\right) \tag{2.8}
\end{equation*}
$$

where

$$
\begin{equation*}
I_{\mu}(t)=I_{0}-\mu \int_{-\infty}^{t}\left(1-\cos q_{A, \theta}^{\mu}(s)\right) \partial_{\varphi} f(\omega s+A) \mathrm{d} s \tag{2.9}
\end{equation*}
$$

is a solution of $\mathcal{H}_{\mu}$ emanating at $t=-\infty$ from torus $\mathcal{T}_{I_{0}}$. Since $q_{A, \theta}^{\mu}$ converges exponentially fast to the equilibrium, the "jump" in the action variables $I_{\mu}(+\infty)-I_{0}$ is finite. We shall speak of homoclinic orbit to the torus $\mathcal{T}_{I_{0}}$ when the jump is zero, and of heteroclinic orbit from $\mathcal{T}_{I_{0}}$ to $\mathcal{T}_{I_{\mu}(+\infty)}$ when the jump is not zero. As a direct consequence of (2.5) we have that such a jump is given by $\partial_{A} F_{\mu}(A, \theta)$ :

LEMMA 2.3. - Let $\partial_{\theta} F_{\mu}(A, \theta)=0$ then $I_{\mu}(t)$ given in (2.9) satisfies

$$
\begin{equation*}
\partial_{A} F_{\mu}(A, \theta)=\int_{-\infty}^{+\infty} \dot{I}_{\mu}(t) \mathrm{d} t=I_{\mu}(+\infty)-I_{0} \in \mathbf{R}^{n} \tag{2.10}
\end{equation*}
$$

In particular if $(A, \theta)$ is a critical point of $F_{\mu}(A, \theta)$ then (2.8) defines a homoclinic orbit to torus $\mathcal{T}_{I_{0}}$.

By the invariance property (2.4) if $B$ is a critical point of the homoclinic function $G_{\mu}$, then, for all $(A, \theta)$ such that $A+\omega \theta=B$, (2.8) provides, for each $I_{0}$, a homoclinic solution to the torus $\mathcal{T}_{I_{0}}$. These homoclinic orbits are not geometrically distinct since, by the autonomy of $\mathcal{H}_{\mu}$, they are all obtained by time translation of the same homoclinic orbit. By the Ljusternik-Schnirelman category theory, since cat $\mathbf{T}^{n}=n+1$, the function $G_{\mu}: \mathbf{T}^{n} \rightarrow \mathbf{R}$ has at least $n+1$ distinct critical points. This proves (see also [36])

THEOREM 2.1. - Let $0<\mu \leqslant \mu_{0} . \forall I_{0} \in \mathbf{R}^{n}$ there exist at least $n+1$ geometrically distinct homoclinic orbits to $\mathcal{T}_{I_{0}}$.

From the conservation of energy a heteroclinic orbit between $\mathcal{T}_{I_{0}}$ and $\mathcal{T}_{I_{0}^{\prime}}$, if any, must satisfy the energy relation

$$
\begin{equation*}
\omega \cdot I_{0}=\omega \cdot I_{0}^{\prime} \tag{2.11}
\end{equation*}
$$

By Lemma 2.3 a critical point of $F_{\mu, I_{0}, I_{0}^{\prime}}(A, \theta)$, defined by $F_{\mu, I_{0}, I_{0}^{\prime}}(A, \theta)=F_{\mu}(A, \theta)$ -$\left(I_{0}^{\prime}-I_{0}\right) \cdot A=G_{\mu}(A+\omega \theta)-\left(I_{0}^{\prime}-I_{0}\right) \cdot A$, gives rise to a heteroclinic solution joining the tori $\mathcal{I}_{I_{0}}$ and $\mathcal{T}_{I_{0}^{\prime}}$. If the energy condition (2.11) holds then the function $F_{\mu, I_{0}, I_{0}^{\prime}}(A, \theta)$ satisfies the invariance property

$$
\begin{equation*}
F_{\mu, I_{0}, I_{0}^{\prime}}(A, \theta)=G_{\mu}(A+\omega \theta)-\left(I_{0}^{\prime}-I_{0}\right) \cdot(A+\omega \theta)=G_{\mu, I_{0}, I_{0}^{\prime}}(A+\omega \theta) \tag{2.12}
\end{equation*}
$$

where

$$
\begin{equation*}
G_{\mu, I_{0}, I_{0}^{\prime}}(B):=G_{\mu}(B)-\left(I_{0}^{\prime}-I_{0}\right) \cdot B . \tag{2.13}
\end{equation*}
$$

Note that $G_{\mu, I_{0}, I_{0}^{\prime}}$ is not $2 \pi \mathbf{Z}^{n}$-periodic, and it might possess no critical point even for $\left|I_{0}^{\prime}-I_{0}\right|$ small. However near a homoclinic orbit to $\mathcal{T}_{I_{0}}$ satisfying some "transversality condition" there exist heteroclinic solutions connecting nearby tori $\mathcal{T}_{I_{0}^{\prime}}$. As an example, the following theorem holds, where $B_{\rho}\left(A_{0}\right)$ denotes the open ball of radius $\rho$ centered at $A_{0}$ in $\mathbf{R}^{n}$ (covering space of $\mathbf{T}^{n}$ ).

THEOREM 2.2. - Assume that there exist $A_{0} \in \mathbf{T}^{n}, \delta>0$ and $\rho>0$ such that $\inf _{\partial B_{\rho}\left(A_{0}\right)} G_{\mu} \geqslant \inf _{B_{\rho}\left(A_{0}\right)} G_{\mu}+\delta$. Then for all $I_{0}, I_{0}^{\prime} \in \mathbf{R}^{n}$ satisfying $\left(I_{0}-I_{0}^{\prime}\right) \cdot \omega=0$ and $\left|I_{0}-I_{0}^{\prime}\right|<\delta /(2 \rho)$ there exists a heteroclinic solution of $\mathcal{H}_{\mu}$ connecting $\mathcal{T}_{I_{0}}$ to $\mathcal{T}_{I_{0}^{\prime}}$.

Proof. - Let us consider $\mathcal{G}_{\mu}(A):=G_{\mu}(A)-\left(I_{0}^{\prime}-I_{0}\right) \cdot\left(A-A_{0}\right)$ which differs from $G_{\mu, I_{0}, I_{0}^{\prime}}$, defined in (2.13), only by a constant. $\mathcal{G}$ attains an absolute minimum in $\overline{B_{\rho}\left(A_{0}\right)}$ at some point $\bar{A}$. The theorem is proved if we show that $\bar{A} \in B_{\rho}\left(A_{0}\right)$. Arguing by contradiction we assume that $\bar{A} \in \partial B_{\rho}\left(A_{0}\right)$. Denote by $\widetilde{A} \in B_{\rho}\left(A_{0}\right)$ a minimum point for $G_{\mu}$. Then $\mathcal{G}_{\mu}(\bar{A}) \geqslant G_{\mu}(\bar{A})-\left|I_{0}^{\prime}-I_{0}\right| \rho>\min _{B_{\rho}\left(A_{0}\right)} G_{\mu}+\delta / 2=G_{\mu}(\widetilde{A})+\delta / 2$. However $\mathcal{G}_{\mu}(\widetilde{A}) \leqslant G_{\mu}(\widetilde{A})+\delta / 2$ and we get a contradiction.

### 2.2. The $\boldsymbol{k}$-bump pseudo-homoclinic solutions

We prove in the next lemma the existence of " $k$-bump pseudo-homoclinic solutions" $q_{A, \theta}^{L}(t)$ of the quasi-periodically forced pendulum (1.5) which turn $k$ times along the separatrices and are asymptotic to the equilibrium for $t \rightarrow \pm \infty$. Such pseudo-homoclinics $q_{A, \theta}^{L}(t)$ are found, via the Contraction Mapping Theorem, as small perturbations of a chain of "1-bump pseudo-homoclinic solutions" obtained in Lemma 2.1.

Lemma 2.4. - There exist $C_{1}, L_{1}>0$ such that $\forall \omega \in \mathbf{R}^{n}, \forall 0<\mu \leqslant \mu_{0}, \forall k \in \mathbf{N}$, $\forall L>L_{1}, \forall \theta=\left(\theta_{1}, \ldots, \theta_{k}\right) \in \mathbf{R}^{k}$ such that $\min _{i}\left(\theta_{i+1}-\theta_{i}\right)>L$, there exists a unique pseudo-homoclinic solution $q_{A, \theta}^{L}(t): \mathbf{R} \rightarrow \mathbf{R}$, smooth in $(A, \theta, \mu)$ which is a true solution of $(1.5)$ in each interval $\left(-\infty, \theta_{1}\right),\left(\theta_{i}, \theta_{i+1}\right)(i=1, \ldots, k-1),\left(\theta_{k},+\infty\right)$ and

$$
\begin{gather*}
q_{A, \theta}^{L}\left(\theta_{i}\right)=\pi(2 i-1), q_{A, \theta}^{L}(t)=q_{A, \theta_{1}}^{\mu}(t) \quad \text { in }\left(-\infty, \theta_{1}\right)  \tag{i}\\
\quad \text { and } q_{A, \theta}^{L}(t)=2 \pi(k-1)+q_{A, \theta_{k}}^{\mu}(t) \quad \text { in }\left(\theta_{k},+\infty\right)
\end{gather*}
$$

$$
\begin{align*}
& \left\|q_{A, \theta}^{L}-q_{A, \theta_{i}}^{\mu}\right\|_{W^{1, \infty}\left(J_{i}\right)} \leqslant C_{1} \exp \left(-C_{1} L\right)  \tag{ii}\\
& \quad \text { where } J_{i}=\left(\theta_{i},\left(\theta_{i}+\theta_{i+1}\right) / 2\right), \forall i=1, \ldots, k-1
\end{align*}
$$

(iii) $\left\|q_{A, \theta}^{L}-q_{A, \theta_{i+1}}^{\mu}\right\|_{W^{1, \infty}\left(J_{i}^{\prime}\right)} \leqslant C_{1} \exp \left(-C_{1} L\right)$,
where $J_{i}^{\prime}=\left(\left(\theta_{i}+\theta_{i+1}\right) / 2, \theta_{i+1}\right), \forall i=1, \ldots, k-1$;
(iv) $\quad q_{A, \theta}^{L}(t)=q_{A+k 2 \pi, \theta}^{L}(t), \quad \forall k \in \mathbf{Z}^{n} ;$
(v) $\quad q_{A, \theta+\eta e_{k}}^{L}(t)=q_{A+\omega \eta, \theta}^{L}(t-\eta), \quad \forall \eta \in \mathbf{R}, \theta \in \mathbf{R}^{k}, \quad$ where $e_{k}=(1, \ldots, 1) \in \mathbf{R}^{k}$.

Proof. - In Appendix A.
Remark 2.2. - The " $k$-bump pseudo-homoclinic solutions" $q_{A, \theta}^{L}$ of Lemma 2.4 could also be obtained minimizing the Lagrangian action functional with fixed end-points
as in [8] using the directs method of the calculus of variations. For general systems, where the homoclinic is not a minimum for the action functional, as for the Duffing equation, this approach cannot be applied. Moreover we prefer to use the Contraction Mapping Theorem since it also provides immediately the approximation estimates of Lemma 2.4(ii)-(iii).

We consider the Lagrangian action functional evaluated on the pseudo-homoclinic solutions $q_{A, \theta}^{L}$ given by lemma 2.4 depending on $n+k$ variables

$$
\begin{aligned}
& F_{\mu}^{k}\left(A_{1}, \ldots, A_{n}, \theta_{1}, \ldots, \theta_{k}\right)=\int_{-\infty}^{\theta_{1}} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right) \mathrm{d} t \\
& \quad+\sum_{i=1}^{k-1} \int_{\theta_{i}}^{\theta_{i+1}} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right) \mathrm{d} t+\int_{\theta_{k}}^{+\infty} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right) \mathrm{d} t .
\end{aligned}
$$

By Lemma 2.4(v) the following invariance property holds

$$
\begin{equation*}
F_{\mu}^{k}\left(A, \theta+\eta e_{k}\right)=F_{\mu}^{k}(A+\eta \omega, \theta), \quad \forall \eta \in \mathbf{R}, \theta \in \mathbf{R}^{k} \tag{2.14}
\end{equation*}
$$

Let $\mathcal{F}_{\mu}^{k}: \mathbf{T}^{n} \times \mathbf{R}^{k} \rightarrow \mathbf{R}$ be the " $k$-bump heteroclinic function" defined by

$$
\begin{equation*}
\mathcal{F}_{\mu}^{k}(A, \theta):=F_{\mu}^{k}(A, \theta)-\left(I_{0}^{\prime}-I_{0}\right) \cdot A \tag{2.15}
\end{equation*}
$$

It turns out that the $k$-dimensional manifold $Z_{k}:=\left\{q_{A, \theta}^{L} \mid \theta \in \mathbf{R}^{k}\right\}$ of " $k$-bump pseudohomoclinic solutions" of (1.5) is a natural constraint for $\Phi_{\mu, A}$, see [4]. Arguing as in Lemma 2.2 we have

LEMMA 2.5. - $\forall I_{0}, I_{0}^{\prime} \in \mathbf{R}^{n}$, if $(A, \theta)$ is a critical point of $\mathcal{F}_{\mu}^{k}(A, \theta)$, then $(\omega t+A$, $\left.I_{\mu}(t), q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t)\right)$ where $I_{\mu}(t)=I_{0}-\mu \int_{-\infty}^{t}\left(1-\cos q_{A, \theta}^{L}(s)\right) \partial_{\varphi} f(\omega s+A) \mathrm{d} s$ is a heteroclinic solution connecting $\mathcal{T}_{I_{0}}$ to $\mathcal{T}_{I_{0}^{\prime}}$.

By Lemma 2.5 we need to find critical points of $\mathcal{F}_{\mu}^{k}(A, \theta)$. When $\min _{i}\left(\theta_{i+1}-\theta_{i}\right) \rightarrow$ $+\infty$ the " $k$-bump homoclinic function" $F_{\mu}^{k}(A, \theta)$ is well approximated simply by the sum of $F_{\mu}\left(A, \theta_{i}\right)$ according to the following lemma. We set $\theta_{0}=-\infty$ and $\theta_{k+1}=+\infty$.

Lemma 2.6. - There exist $C_{2}, L_{2}>0$ such that $\forall \omega \in \mathbf{R}^{n}, \forall 0<\mu \leqslant \mu_{0}, \forall L>L_{2}$, $\forall \theta_{1}<\cdots<\theta_{k}$ with $\min _{i}\left(\theta_{i+1}-\theta_{i}\right)>L$

$$
\begin{equation*}
F_{\mu}^{k}\left(A, \theta_{1}, \ldots, \theta_{k}\right)=\sum_{i=1}^{k} F_{\mu}\left(A, \theta_{i}\right)+\sum_{i=1}^{k} R_{i}\left(\mu, A, \theta_{i-1}, \theta_{i}, \theta_{i+1}\right) \tag{2.16}
\end{equation*}
$$

with

$$
\left|R_{i}\left(\mu, A, \theta_{i-1}, \theta_{i}, \theta_{i+1}\right)\right| \leqslant C_{2} \exp \left(-C_{2} L\right)
$$

Proof. - We can write

$$
\begin{aligned}
F_{\mu}^{k} & \left(A, \theta_{1}, \ldots, \theta_{k}\right) \\
= & \left(\int_{-\infty}^{\theta_{1}} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right)+\int_{\theta_{1}}^{\left(\theta_{1}+\theta_{2}\right) / 2} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right)\right) \\
& +\sum_{i=2}^{k-1}\left(\int_{\left(\theta_{i-1}+\theta_{i}\right) / 2}^{\theta_{i}} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right)+\int_{\theta_{i}}^{\left(\theta_{i}+\theta_{i+1}\right) / 2} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right)\right) \\
& +\left(\int_{\left(\theta_{k-1}+\theta_{k}\right) / 2}^{\theta_{k}} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right)+\int_{\theta_{k}}^{+\infty} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right)\right) .
\end{aligned}
$$

We define

$$
\begin{aligned}
& R_{i}^{-}\left(\mu, A, \theta_{i-1}, \theta_{i}\right) \\
& \quad=\int_{\left(\theta_{i-1}+\theta_{i}\right) / 2}^{\theta_{i}} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right) \mathrm{d} t-\int_{-\infty}^{\theta_{i}} \mathcal{L}_{\mu, A}\left(q_{A, \theta_{i}}^{\mu}(t), \dot{q}_{A, \theta_{i}}^{\mu}(t), t\right) \mathrm{d} t, \\
& R_{i}^{+} \\
& \quad\left(\mu, A, \theta_{i}, \theta_{i+1}\right) \\
& \quad=\int_{\theta_{i}}^{\left(\theta_{i}+\theta_{i+1}\right) / 2} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right) \mathrm{d} t-\int_{\theta_{i}}^{+\infty} \mathcal{L}_{\mu, A}\left(q_{A, \theta_{i}}^{\mu}(t), \dot{q}_{A, \theta_{i}}^{\mu}(t), t\right) \mathrm{d} t,
\end{aligned}
$$

where $q_{A, \theta_{i}}^{\mu}$ is the 1-bump pseudo-homoclinic solution obtained in Lemma 2.1. Recalling the Definition 2.1 of $F_{\mu}(A, \theta)$ we have

$$
\begin{aligned}
F_{\mu}^{k}\left(A, \theta_{1}, \ldots, \theta_{k}\right)= & F_{\mu}\left(A, \theta_{1}\right)+R_{1}^{+}\left(\mu, A, \theta_{1}, \theta_{2}\right) \\
& +\sum_{i=2}^{k-1} F_{\mu}\left(A, \theta_{i}\right)+\left(R_{i}^{-}\left(\mu, A, \theta_{i-1}, \theta_{i}\right)+R_{i}^{+}\left(\mu, A, \theta_{i}, \theta_{i+1}\right)\right. \\
& +F_{\mu}\left(A, \theta_{k}\right)+R_{k}^{-}\left(\mu, A, \theta_{k-1}, \theta_{k}\right)
\end{aligned}
$$

Setting $R_{1}=R_{1}^{+}, R_{i}=R_{i}^{-}+R_{i}^{+}(i=2, \ldots, k-1)$ and $R_{k}=R_{k}^{-}$we derive the expression (2.16). In order to complete the proof, we have to show the existence of $C_{2}, L_{2}>0$ such that $\forall \omega \in \mathbf{R}^{n}$, for all $0<\mu \leqslant \mu_{0}, \forall L>L_{2}, \forall \theta_{1}<\cdots<\theta_{k}$ with $\min _{i}\left(\theta_{i+1}-\theta_{i}\right)>L$, for all $i=1, \ldots, k$

$$
\begin{equation*}
\left|R_{i}^{ \pm}\left(\mu, A, \theta_{i}, \theta_{i+1}\right)\right| \leqslant C_{2} \exp \left(-C_{2} L\right) \tag{2.17}
\end{equation*}
$$

We write the proof for $R_{i}^{+}$. We have

$$
\begin{aligned}
R_{i}^{+} & \left(\mu, A, \theta_{i}, \theta_{i+1}\right) \\
& =\int_{\theta_{i}}^{\left(\theta_{i}+\theta_{i+1}\right) / 2}\left(\mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right)-\mathcal{L}_{\mu, A}\left(q_{A, \theta_{i}}^{\mu}(t), \dot{q}_{A, \theta_{i}}^{\mu}(t), t\right)\right) \mathrm{d} t
\end{aligned}
$$

$$
\begin{equation*}
-\int_{\left(\theta_{i}+\theta_{i+1}\right) / 2}^{+\infty} \mathcal{L}_{\mu, A}\left(q_{A, \theta_{i}}^{\mu}(t), \dot{q}_{A, \theta_{i}}^{\mu}(t), t\right) \mathrm{d} t \tag{2.18}
\end{equation*}
$$

By Lemma 2.1(ii) the homoclinic orbit satisfies $\max \left(\left|q_{A, \theta_{i}}^{\mu}(t)\right|,\left|\dot{q}_{A, \theta_{i}}^{\mu}(t)\right|\right) \leqslant$ $C \exp \left(-\left|t-\theta_{i}\right| / 2\right)$. Hence $\mathcal{L}_{\mu, A}\left(q_{A, \theta_{i}}^{\mu}(t), \dot{q}_{A, \theta_{i}}^{\mu}(t), t\right)=\mathrm{O}\left(\exp \left(-\left|t-\theta_{i}\right|\right)\right)$ and for all $\theta_{1}<\cdots<\theta_{k}$ with $\min _{i}\left(\theta_{i+1}-\theta_{i}\right)>L$,

$$
\begin{equation*}
\left|\int_{\left(\theta_{i}+\theta_{i+1}\right) / 2}^{+\infty} \mathcal{L}_{\mu, A}\left(q_{A, \theta_{i}}^{\mu}(t), \dot{q}_{A, \theta_{i}}^{\mu}(t), t\right) \mathrm{d} t\right|=\mathrm{O}\left(\mathrm{e}^{-L / 2}\right) \tag{2.19}
\end{equation*}
$$

From Lemma 2.4(ii) we also deduce that

$$
\begin{equation*}
\left(\int_{\theta_{i}}^{\left(\theta_{i}+\theta_{i+1}\right) / 2} \mathcal{L}_{\mu, A}\left(q_{A, \theta}^{L}(t), \dot{q}_{A, \theta}^{L}(t), t\right)-\mathcal{L}_{\mu, A}\left(q_{A, \theta_{i}}^{\mu}(t), \dot{q}_{A, \theta_{i}}^{\mu}(t), t\right) \mathrm{d} t\right)=\mathrm{O}\left(\mathrm{e}^{-C L}\right) \tag{2.20}
\end{equation*}
$$

for some constant $C>0$. From (2.18), (2.19) and (2.20) we deduce (2.17).

### 2.3. The diffusion orbit

We now give an example of condition on $G_{\mu}$ which implies the existence of diffusion orbits. $B_{\alpha}\left(A_{0}\right) \subset \mathbf{R}^{n}$ denotes the open ball centered at $A_{0} \in \mathbf{R}^{n}$ and of radius $\alpha$.

Condition 2.1 ("Splitting condition"). - There exist $A_{0} \in \mathbf{T}^{n}, \alpha>0$, a bounded open set $U \subset \mathbf{R}^{n}$ (the covering space of $\mathbf{T}^{n}$ ) such that $B_{\alpha}\left(A_{0}\right) \subset U$ and a positive constant $\delta>0$ such that
(i) $\inf _{\partial U} G_{\mu} \geqslant \inf _{U} G_{\mu}+\delta$;
(ii) $\sup _{B_{\alpha}\left(A_{0}\right)} G_{\mu} \leqslant \frac{\delta}{4}+\inf _{U} G_{\mu}$;
(iii) $\quad d\left(\left\{A \in U \mid G_{\mu}(A) \leqslant \delta / 2+\inf _{U} G_{\mu}\right\},\left\{A \in U \mid G_{\mu}(A) \geqslant 3 \delta / 4+\inf _{U} G_{\mu}\right\}\right) \geqslant 2 \alpha$.

Remark 2.3. - If $G_{\mu}$ possesses a non-degenerate minimum in $A_{0}$ the "Splitting condition" 2.1 is satisfied, for example, choosing $U=B_{\rho}\left(A_{0}\right)$ with $\rho=\min \lambda_{i} /$ $\left(3 \sup _{A \in \mathbf{T}^{n}}\left|D^{3} G_{\mu}\right|\right), \quad \delta=4\left(\min \lambda_{i}\right) \rho^{2} / 9 \quad$ and $\quad \alpha=\rho((\sqrt{6}-\sqrt{5}) / 2 \sqrt{10})$ $\sqrt{\left(\min _{i} \lambda_{i}\right) /\left(\max _{i} \lambda_{i}\right)}$ where $\lambda_{i}$ are the (positive) eigenvalues of $D^{2} G^{\mu}\left(A_{0}\right)$. Details are given in Appendix A .

The following shadowing type theorem holds, where

$$
\rho_{U}:=\operatorname{diam}\left(\Pi_{\omega}(U)\right)
$$

and $\Pi_{\omega}: \mathbf{R}^{n} \rightarrow \mathbf{R}^{n}$ denotes the orthogonal projection onto $\omega^{\perp}$.
THEOREM 2.3. - Assume (H1) and the "Splitting condition" 2.1. Then $\forall I_{0}, I_{0}^{\prime}$ with $\omega \cdot I_{0}=\omega \cdot I_{0}^{\prime}$, there is a heteroclinic orbit connecting the invariant tori $\mathcal{I}_{I_{0}}$ and $\mathcal{T}_{I_{0}^{\prime}}$.

Moreover there exists $C_{3}>0$ such that $\forall \eta>0$ small enough the "diffusion time" $T_{d}$ needed to go from a $\eta$-neighbourhood of $\mathcal{I}_{I_{0}}$ to a $\eta$-neighbourhood of $\mathcal{T}_{I_{0}^{\prime}}$ is bounded by

$$
\begin{equation*}
T_{d} \leqslant C_{3} \frac{\left|I_{0}-I_{0}^{\prime}\right|}{\delta} \cdot \rho_{U} \cdot \max \left(|\log \delta|, \frac{1}{\gamma \alpha^{\tau}}\right)+C_{3}|\log (\eta)| \tag{2.21}
\end{equation*}
$$

Remark 2.4. - The meaning of (2.21) is the following: the diffusion time $T_{d}$ is estimated by the product of the number of heteroclinic transitions $k=$ ( heteroclinic jump / splitting $)=\left|I_{0}^{\prime}-I_{0}\right| / \delta$, and of the time $T_{s}$ required for a single transition, that is $T_{d}=k \cdot T_{s}$. The time for a single transition $T_{s}$ is bounded by the maximum time between the "ergodization time" $\left(1 / \gamma \alpha^{\tau}\right)$, i.e. the time needed for the flow $\omega t$ to make an $\alpha$-net of the torus, and the time $|\log \delta|$ needed to "shadow" homoclinic orbits for the forced pendulum equation. We use here that these homoclinic orbits are exponentially asymptotic to the equilibrium.

Proof. - Assume with no loss of generality that $A_{0}=0$ and $\inf _{U} G_{\mu}(A)=0$. Let us choose the number of bumps $k$ as

$$
\begin{equation*}
k=\left[\frac{24 \cdot \rho_{U} \cdot\left|I_{0}^{\prime}-I_{0}\right|}{\delta}\right]+1 \tag{2.22}
\end{equation*}
$$

By Lemma 2.4(i) and Lemma 2.1(ii), the trajectory $q_{A, \theta}^{L}$ converges exponentially fast to $\mathcal{T}_{I_{0}}$ (resp. $\mathcal{T}_{I_{0}^{\prime}}$ ) as $t \rightarrow-\infty($ resp. $+\infty)$ from $\theta_{1}$ (resp. $\theta_{k}$ ). Therefore it is enough to prove the existence of a critical point $(\bar{A}, \bar{\theta}) \in \mathbf{T}^{n} \times \mathbf{R}^{k}$ of the $k$-bump heteroclinic function $\mathcal{F}_{\mu}^{k}$, defined in (2.15), such that for some positive constant $K_{1}$

$$
\begin{equation*}
\left|\bar{\theta}_{k}-\bar{\theta}_{1}\right| \leqslant K_{1} \frac{\left|I_{0}-I_{0}^{\prime}\right|}{\delta} \cdot \rho_{U} \cdot \max \left(|\log \delta|, \frac{1}{\gamma \alpha^{\tau}}\right) \tag{2.23}
\end{equation*}
$$

More precisely we shall enforce

$$
\begin{equation*}
K_{2}|\log \delta|<\left|\bar{\theta}_{i+1}-\bar{\theta}_{i}\right|<K_{3} \max \left(|\log \delta|, \frac{1}{\gamma \alpha^{\tau}}\right) \quad \forall i=1, \ldots, k-1 \tag{2.24}
\end{equation*}
$$

for some positive constants $K_{2}, K_{3}$. Let $\left(\Omega_{1}, \ldots, \Omega_{n}\right)$ be an orthonormal basis of $\mathbf{R}^{n}$ where

$$
\Omega_{1}=\frac{\omega}{|\omega|} \quad \text { and } \quad \Omega_{2}=\frac{I_{0}^{\prime}-I_{0}}{\left|I_{0}^{\prime}-I_{0}\right|}
$$

We recall that $\omega \cdot\left(I_{0}^{\prime}-I_{0}\right)=0$. In order to find a critical point of $\mathcal{F}_{\mu}^{k}$ we introduce suitable coordinates $\left(a_{1}, \ldots, a_{n}, s_{1}, \ldots, s_{k}\right) \in \mathbf{R}^{n} \times(-\operatorname{diam} U \text {, } \operatorname{diam} U)^{k}$ defined by

$$
A=\sum_{j=1}^{n} a_{j} \Omega_{j}, \quad \theta_{i}=\frac{\eta_{i}+s_{i}-a_{1}}{|\omega|} \quad \forall i=1, \ldots, k
$$

where $\eta_{i}$ are constants to be chosen later. In these new coordinates the heteroclinic function defined in (2.15) is given by

$$
\begin{align*}
& \overline{\mathcal{F}}_{\mu}^{k}\left(a_{1}, a_{2}, \ldots, a_{n}, s_{1}, \ldots, s_{k}\right) \\
& \quad=F_{\mu}^{k}\left(\sum_{j=1}^{n} a_{j} \Omega_{j}, \frac{\eta_{1}+s_{1}-a_{1}}{|\omega|}, \ldots, \frac{\eta_{k}+s_{k}-a_{1}}{|\omega|}\right)-\left|I_{0}^{\prime}-I_{0}\right| a_{2} . \tag{2.25}
\end{align*}
$$

Using the invariance property (2.14) we see that $\overline{\mathcal{F}}_{\mu}^{k}$ does not depend on the new variable $a_{1}$ :

$$
\begin{aligned}
\overline{\mathcal{F}}_{\mu}^{k}\left(a_{1}, a_{2}, \ldots, a_{n}, s_{1}, \ldots, s_{k}\right) & =F_{\mu}^{k}\left(\sum_{j=2}^{n} a_{j} \Omega_{j}, \frac{\eta_{1}+s_{1}}{|\omega|}, \ldots, \frac{\eta_{k}+s_{k}}{|\omega|}\right)-\left|I_{0}^{\prime}-I_{0}\right| a_{2} \\
& =\overline{\mathcal{F}}_{\mu}^{k}\left(0, a_{2}, \ldots, a_{n}, s_{1}, \ldots, s_{k}\right) .
\end{aligned}
$$

For simplicity we will still use the notation $\overline{\mathcal{F}}_{\mu}^{k}\left(a_{2}, \ldots, a_{n}, s_{1}, \ldots, s_{k}\right)$ for $\overline{\mathcal{F}}_{\mu}^{k}\left(0, a_{2}, \ldots, a_{n}, s_{1}, \ldots, s_{k}\right)$. We now choose the constants $\eta_{i}$. Let

$$
\begin{equation*}
D=\frac{|\omega|}{C_{2}}\left|\log \left(\frac{24 C_{2}}{\delta}\right)\right|+2 \operatorname{diam} U \tag{2.26}
\end{equation*}
$$

where $C_{2}$ is the constant appearing in Lemma 2.6. We shall use the following fact (see [13]): there is $\bar{C}>0$ such that, for all intervals $J \subset \mathbf{R}$ of length greater or equal to $\bar{C} /\left(\gamma \alpha^{\tau}\right)$, there is $\theta \in J$ such that

$$
\begin{equation*}
d\left(\theta \omega, 2 \pi \mathbf{Z}^{n}\right)<\alpha \tag{2.27}
\end{equation*}
$$

By (2.27) there is $\left(\eta_{1}, \ldots, \eta_{k}\right) \in \mathbf{R}^{k}$ such that

$$
\begin{gather*}
\eta_{i} \Omega_{1} \equiv \chi_{i}, \bmod 2 \pi \mathbf{Z}^{n}, \quad\left|\chi_{i}\right|<\alpha \quad \text { and } \quad \chi_{i} \cdot \Omega_{1}=0, \quad \text { i.e. } \chi_{i}=\sum_{j=2}^{n} \chi_{i, j} \Omega_{j}  \tag{2.28}\\
\eta_{1}=0, \quad D<\eta_{i+1}-\eta_{i}<\left(D+\frac{\bar{C}|\omega|}{\gamma \alpha^{\tau}}\right) \tag{2.29}
\end{gather*}
$$

By (2.26), (2.29), since $s_{i} \in(-\operatorname{diam} U, \operatorname{diam} U)$ we have that $\theta_{i+1}-\theta_{i} \geqslant \frac{1}{C_{2}}\left|\log \left(\frac{24 C_{2}}{\delta}\right)\right|$; hence, by Lemma 2.6, setting

$$
\begin{aligned}
\bar{R}_{i} & =\bar{R}_{i}\left(a_{2}, \ldots, a_{n}, s_{i-1}, s_{i}, s_{i+1}\right) \\
& =R_{i}\left(\mu, \sum_{j=2}^{n} a_{j} \Omega_{j}, \frac{s_{i-1}+\eta_{i-1}}{|\omega|}, \frac{s_{i}+\eta_{i}}{|\omega|}, \frac{s_{i+1}+\eta_{i+1}}{|\omega|}\right)
\end{aligned}
$$

we get

$$
\begin{equation*}
\left|\bar{R}_{i}\left(a_{2}, \ldots, a_{n}, s_{i-1}, s_{i}, s_{i+1}\right)\right|<\frac{\delta}{24} . \tag{2.30}
\end{equation*}
$$

By Lemma 2.6, the invariance property (2.14), (2.28) and since $G_{\mu}$ is $2 \pi \mathbf{Z}^{n}$-periodic, we have

$$
\begin{align*}
& \overline{\mathcal{F}}_{\mu}^{k}\left(a_{2}, \ldots, a_{n}, s_{1}, \ldots, s_{k}\right) \\
& \quad=\sum_{i=1}^{k} F_{\mu}\left(\sum_{j=2}^{n} a_{j} \Omega_{j}, \frac{\eta_{i}+s_{i}}{|\omega|}\right)+\bar{R}_{i}-\left|I_{0}^{\prime}-I_{0}\right| a_{2} \\
& \quad=\sum_{i=1}^{k} F_{\mu}\left(\sum_{j=2}^{n} a_{j} \Omega_{j}+\chi_{i}+s_{i} \Omega_{1}, 0\right)+\bar{R}_{i}-\left|I_{0}^{\prime}-I_{0}\right| a_{2} \\
& \quad=\sum_{i=1}^{k} G_{\mu}\left(\sum_{j=2}^{n}\left(a_{j}+\chi_{i, j}\right) \Omega_{j}+s_{i} \Omega_{1}\right)+\bar{R}_{i}-\left|I_{0}^{\prime}-I_{0}\right| a_{2} \\
& \quad=\sum_{i=1}^{k} \bar{G}_{\mu}\left(a_{2}+\chi_{i, 2}, \ldots, a_{n}+\chi_{i, n}, s_{i}\right)+\bar{R}_{i}-\left|I_{0}^{\prime}-I_{0}\right| a_{2} \tag{2.31}
\end{align*}
$$

where $\bar{G}_{\mu}\left(a_{2}, \ldots, a_{n}, s\right)=G_{\mu}\left(\sum_{j=2}^{n} a_{j} \Omega_{j}+s \Omega_{1}\right)$. By (2.31), setting $\bar{\chi}_{i}:=\left(\chi_{i, 2}\right.$, $\left.\ldots, \chi_{i, n}\right) \in \mathbf{R}^{n-1}$ and $a=\left(a_{2}, \ldots, a_{n}\right)$, we get

$$
\begin{equation*}
\overline{\mathcal{F}}_{\mu}^{k}\left(a_{2}, \ldots, a_{n}, s_{1}, \ldots, s_{k}\right)=\sum_{i=1}^{k} \bar{G}_{\mu}\left(a+\bar{\chi}_{i}, s_{i}\right)+\bar{R}_{i}-\left|I_{0}^{\prime}-I_{0}\right| a_{2} \tag{2.32}
\end{equation*}
$$

Since the basis $\left(\Omega_{1}, \ldots, \Omega_{n}\right)$ is orthonormal the function $\bar{G}_{\mu}$ satisfies the same properties as $G_{\mu}$, i.e.

$$
\sup _{B_{\alpha}(0)} \bar{G}_{\mu} \leqslant \frac{\delta}{4}, \quad \inf _{\partial \widetilde{U}} \bar{G}_{\mu} \geqslant \delta
$$

and $d\left(\left\{x \in \widetilde{U} \mid \bar{G}_{\mu}(x) \leqslant \delta / 2\right\},\left\{x \in \widetilde{U} \mid \bar{G}_{\mu}(x) \geqslant 3 \delta / 4\right\}\right) \geqslant 2 \alpha$, where $\tilde{U}=\left\{\left(a_{2}, \ldots\right.\right.$, $\left.\left.a_{n}, s\right) \in \mathbf{R}^{n} \mid \sum_{j=2}^{n} a_{j} \Omega_{j}+s \Omega_{1} \in U\right\}$. Note that $\operatorname{diam} U=\operatorname{diam} \widetilde{U}$. In particular, since $(0,0) \in \widetilde{U}$, for all $(a, s) \in \widetilde{U}, s \in(-\operatorname{diam} U$, $\operatorname{diam} U)$. Moreover, for all $(a, s) \in \widetilde{U}$, $|a|=\left|\Pi_{\omega}\left(\sum_{j=2}^{n} a_{j} \Omega_{j}+s \Omega_{1}\right)\right| \leqslant \rho_{U}$.

We shall find a critical point of $\overline{\mathcal{F}}_{\mu}^{k}$ in

$$
W=\left\{(a, s) \in \mathbf{R}^{n-1} \times \mathbf{R}^{k} \mid\left(a+\bar{\chi}_{i}, s_{i}\right) \in \widetilde{U}, \forall i=1, \ldots, k\right\} .
$$

As $W$ is bounded, $\overline{\mathcal{F}}_{\mu}^{k}$ attains its minimum over $\bar{W}$ at some point $(\bar{a}, \bar{s})$. Notice that by (2.32) and (2.30)

$$
\inf _{\bar{W}} \overline{\mathcal{F}}_{\mu}^{k} \leqslant \overline{\mathcal{F}}_{\mu}^{k}(0,0)=\sum_{i=1}^{k} \bar{G}_{\mu}\left(\bar{\chi}_{i}, 0\right)+k \frac{\delta}{24} .
$$

Hence, since $\left|\bar{\chi}_{i}\right|<\alpha$ for all $i=1, \ldots, k$ and $\sup _{B_{\alpha}(0)} \bar{G}_{\mu} \leqslant \delta / 4$, we have

$$
\begin{equation*}
\inf _{\bar{W}} \overline{\mathcal{F}}_{\mu}^{k} \leqslant k \frac{\delta}{4}+k \frac{\delta}{24}=k \frac{7 \delta}{24} \tag{2.33}
\end{equation*}
$$

The theorem is proved if we show that $(\bar{a}, \bar{s}) \in W$. Arguing by contradiction we assume that $(\bar{a}, \bar{s}) \in \partial W$. Then there is some $l \in\{1, \ldots, k\}$ such that $\left(\bar{a}+\bar{\chi}_{l}, \bar{s}_{l}\right) \in \partial \widetilde{U}$, so that $\bar{G}_{\mu}\left(\bar{a}+\bar{\chi}_{l}, \bar{s}_{l}\right) \geqslant \delta$. We now prove that $\left\{\left(\bar{a}+\bar{\chi}_{l}, t\right) ; t \in(-\operatorname{diam} U, \operatorname{diam} U)\right\} \cap \widetilde{U} \subset Z:=$ $\left\{x \in \widetilde{U} \mid \bar{G}_{\mu}(x) \geqslant 3 \delta / 4\right\}$. Indeed, if not, by (2.30), for some $t \in(-\operatorname{diam} U$, $\operatorname{diam} U)$ such that $\left(\bar{a}+\bar{\chi}_{l}, t\right) \in \widetilde{U}$,

$$
\begin{aligned}
\overline{\mathcal{F}}_{\mu}^{k}\left(\bar{a}, \bar{s}_{1}, \ldots, \bar{s}_{l-1}, t, \bar{s}_{l+1}, \ldots, \bar{s}_{k}\right) \leqslant & \overline{\mathcal{F}}_{\mu}^{k}(\bar{a}, \bar{s})+\left(\bar{G}_{\mu}\left(\bar{a}+\bar{\chi}_{l}, t\right)-\bar{G}_{\mu}\left(\bar{a}+\bar{\chi}_{l}, \bar{s}_{l}\right)\right) \\
& +\left|\bar{R}_{l-1}\left(\bar{a}, \bar{s}_{l-2}, \bar{s}_{l-1}, s_{l}\right)-\bar{R}_{l-1}\left(\bar{a}, \bar{s}_{l-2}, \bar{s}_{l-1}, t\right)\right| \\
& +\left|\bar{R}_{l}\left(\bar{a}, s_{l-1}, t, \bar{s}_{l}\right)-\bar{R}_{l}\left(\bar{a}, \bar{s}_{l-1}, \bar{s}_{l}, \bar{s}_{l+1}\right)\right| \\
& +\left|\bar{R}_{l+1}\left(\bar{a}, \bar{s}_{l}, \bar{s}_{l+1}, t\right)-\bar{R}_{l+1}\left(\bar{a}, \bar{s}_{l}, \bar{s}_{l+1}, \bar{s}_{l+2}\right)\right| \\
< & \overline{\mathcal{F}}_{\mu}^{k}(\bar{a}, \bar{s})-\frac{\delta}{4}+\frac{6 \delta}{24}=\overline{\mathcal{F}}_{\mu}^{k}(\bar{a}, \bar{s}),
\end{aligned}
$$

which is wrong since $(\bar{a}, \bar{s})$ is the minimum of $\overline{\mathcal{F}}_{\mu}^{k}$ over $\bar{W}$. We deduce in particular that, for all $i$,

$$
\begin{equation*}
\left(\bar{a}+\bar{\chi}_{l}, \overline{s_{i}}\right) \in Z \cup \widetilde{U}^{c} \tag{2.34}
\end{equation*}
$$

We now prove that for all $i$,

$$
\begin{equation*}
\bar{G}_{\mu}\left(\bar{a}+\bar{\chi}_{i}, \overline{s_{i}}\right)>\frac{\delta}{2} . \tag{2.35}
\end{equation*}
$$

Indeed, assume that $\bar{G}_{\mu}\left(\bar{a}+\bar{\chi}_{i}, \overline{s_{i}}\right) \leqslant \delta / 2$. Since $\bar{G}_{\mu} \geqslant 3 \delta / 4$ in a neighbourhood of $\partial \widetilde{U}$, $d\left(x, \widetilde{U}^{c}\right) \geqslant 2 \alpha$ for all $x \in \widetilde{U}$ such that $\bar{G}_{\mu}(x) \leqslant \delta / 2$, so our splitting condition implies that

$$
d\left(Z \cup \tilde{U}^{c},\left\{x \in \widetilde{U} \mid \bar{G}_{\mu}(x) \leqslant \delta / 2\right\}\right) \geqslant 2 \alpha
$$

Therefore we get, recalling (2.34), that $\left|\left(\bar{a}+\bar{\chi}_{l}, \bar{s}_{i}\right)-\left(\bar{a}+\bar{\chi}_{i}, \bar{s}_{i}\right)\right|=\left|\bar{\chi}_{l}-\bar{\chi}_{i}\right| \geqslant 2 \alpha$. This contradicts $\left|\bar{\chi}_{i}\right|,\left|\bar{\chi}_{l}\right|<\alpha$ in (2.28).

By (2.32), (2.35), (2.30), and noting that, by (2.22), $\left|I_{0}^{\prime}-I_{0}\right|\left|a_{2}\right| \leqslant\left|I_{0}^{\prime}-I_{0}\right| \rho_{U} \leqslant$ $(k \delta) / 24$, we deduce that

$$
\overline{\mathcal{F}}_{\mu}^{k}(\bar{a}, \bar{s}) \geqslant k \frac{\delta}{2}-k \frac{\delta}{24}-k \frac{\delta}{24}=k \frac{10 \delta}{24}>k \frac{7 \delta}{24},
$$

contradicting (2.33). The proof of the theorem is complete.
Remark 2.5. - The latter proof works if $G_{\mu}$ possesses a local maximum which satisfies a non-degeneracy type condition like the "splitting condition" 2.1, while in the approaches developed in [9] and [45], based on Mather's theory, diffusion orbits are always built from local minima of $G_{\mu}$. The proof of the shadowing theorem when the homoclimic point $A_{0}$ is a saddle point requires slightly different arguments involving the topological degree. For example it holds assuming as in [24] the condition $D^{2} G_{\mu}\left(A_{0}\right) \omega \cdot \omega \neq 0$. This condition enables to prove that $\operatorname{deg}\left(\nabla \overline{\mathcal{F}}_{\mu}^{k}(a, s), W, 0\right) \neq 0$ for a suitable $W$.

When the frequency vector $\omega$ is considered as a constant, independent of any parameter ("a priori-unstable case") it is easy to justify the splitting condition 2.1 using
the first-order approximation given by the Poincaré-Melnikov primitive. With a Taylor expansion in $\mu$ we can easily prove that for $\mu$ small enough

$$
G_{\mu}(B)=b+\mu \Gamma(B)+\mathrm{O}\left(\mu^{2}\right), \quad \forall B \in \mathbf{T}^{n}
$$

where $b=\Phi_{0}\left(q_{0}\right)$ and $\Gamma: \mathbf{T}^{n} \rightarrow \mathbf{R}$ is nothing but the Poincaré-Melnikov primitive

$$
\begin{equation*}
\Gamma(B)=\int_{\mathbf{R}}\left(\cos q_{0}(t)-1\right) f(\omega t+B) \mathrm{d} t \tag{2.36}
\end{equation*}
$$

Hence, if $\Gamma$ possesses a proper minimum (resp. maximum) in $A_{0} \in \mathbf{R}^{n}$, i.e. if $\exists r>0$ such that $\inf _{\partial B_{r}\left(A_{0}\right)} \Gamma>\Gamma\left(A_{0}\right)$ (resp. $\left.\sup _{\partial B_{r}\left(A_{0}\right)} \Gamma<\Gamma\left(A_{0}\right)\right)$ then, for $\mu$ small enough, the "splitting" condition 2.1 holds with $U=B_{\rho}\left(A_{0}\right), \rho_{U}=O(1), \delta=c \mu$ and $\alpha=c^{\prime}$ for some positive constants $c$ and $c^{\prime}$ depending on $\Gamma$. We remark that $B_{r}\left(A_{0}\right)$ could be replaced by a bounded open subset $U$ of $\mathbf{R}^{n}$. Applying Theorem 2.3 we deduce

THEOREM 2.4. - Assume (H1) and let $\Gamma$ possess a proper minimum (or maximum) $A_{0}$, i.e. suppose that $\exists r>0$ such that $\inf _{\partial B_{r}\left(A_{0}\right)} \Gamma>\Gamma\left(A_{0}\right)$. Then, for $\mu$ small enough, the same statement as in Theorem 2.3 holds with a diffusion time $T_{d}=$ $\mathrm{O}((1 / \mu) \log (1 / \mu))$.

## 3. More general perturbation terms

In this section we show how to adapt the arguments of the previous section in order to deal with a more general perturbation term $f(\varphi, q)$. Regarding regularity it is sufficient to have finite large enough smoothness for $f$. The equation of motion derived by Hamiltonian $\mathcal{H}_{\mu}$ are

$$
\begin{equation*}
\dot{\varphi}=\omega, \quad \dot{I}=-\mu \partial_{\varphi} f(\varphi, q), \quad \dot{q}=p, \quad \dot{p}=\sin q-\mu \partial_{q} f(\varphi, q) \tag{3.1}
\end{equation*}
$$

corresponding to the quasi-periodically forced pendulum

$$
\begin{equation*}
-\ddot{q}+\sin q=\mu \partial_{q} f(\omega t+A, q) \tag{3.2}
\end{equation*}
$$

### 3.1. Invariant tori in the perturbed system

The first step is to prove the persistence of invariant tori for $\mu \neq 0$ small enough. It appears that no more than the standard Implicit Function Theorem is required to prove the following well known result (see for example [32] for a different proof)

THEOREM 3.1. - Let $\omega$ satisfy (H1). For $\mu$ small enough and $\forall I_{0} \in \mathbf{R}^{n}$ system $\mathcal{H}_{\mu}$ possesses $n$-dimensional invariant tori $\mathcal{T}_{I_{0}}^{\mu} \approx \mathcal{I}_{I_{0}}$ of the form

$$
\begin{equation*}
\mathcal{T}_{I_{0}}^{\mu}=\left\{I=I_{0}+a^{\mu}(\psi), p=P^{\mu}(\psi), \varphi=\psi, q=Q^{\mu}(\psi), \quad \psi \in \mathbf{T}^{n}\right\} \tag{3.3}
\end{equation*}
$$

with $Q^{\mu}(\psi), P^{\mu}(\psi)=\mathrm{O}(\mu), a^{\mu}(\psi)=\mathrm{O}(\mu)$. Moreover the dynamics on $\mathcal{T}_{I_{0}}^{\mu}$ is conjugated to the rotation of speed $\omega$ for $\psi$.

We first determine the functions $Q^{\mu}, P^{\mu}$ in (3.3). Using the standard Implicit Function Theorem we prove that there exists a unique quasi-periodic solution $q_{A}^{\mu}(t)$ for the quasi-periodically forced pendulum equation (3.2) which bifurcates from the hyperbolic equilibrium 0 .

LEMMA 3.1. - Let $f \in C^{l}\left(\mathbf{T}^{n} \times \mathbf{T}\right), l \geqslant 2$. For $\mu$ small enough there exists a unique quasi-periodic solution $q_{A}^{\mu}(t)$ of (3.2) with $\left(q_{A}^{\mu}(t), \dot{q}_{A}^{\mu}(t)\right)=\left(q_{A}^{\mu}(t), p_{A}^{\mu}(t)\right)=\mathrm{O}(\mu)$, $C^{l-1}$-smooth in A. More precisely there exist functions $Q^{\mu}, P^{\mu}: \mathbf{T}^{n} \rightarrow \mathbf{R}$ of class $C^{l-1}$, such that $\left(q_{A}^{\mu}(t), p_{A}^{\mu}(t)\right)=\left(Q^{\mu}(\omega t+A), P^{\mu}(\omega t+A)\right)$.

Proof. - Let $L$ be the Green operator of the differential operator $h \rightarrow-D^{2} h+$ $h$ with Dirichlet boundary conditions at $\pm \infty$. L is explicitely given by $L(f)=$ $\int_{\mathbf{R}} \mathrm{e}^{-|t-s|} f(s) \mathrm{d} s / 2$. It results that $L$ is a continuous linear operator in the Banach space of the continuous bounded functions from $\mathbf{R}$ to $\mathbf{R}$ endowed with the sup-norm, which we shall denote by $E$. We consider the non-linear operator $S: \mathbf{R} \times \mathbf{T}^{n} \times E \rightarrow E$

$$
S(\mu, A, q):=q-L(q-\sin q)-\mu L\left(\partial_{q} f(\omega t+A, q)\right)
$$

$S$ is of class $C^{l-1}$. We are looking for a solution $q_{A}^{\mu}$ of $S(\mu, A, q)=0$. Since $S(0, A, 0)=$ 0 and $\partial_{q} S(0, A, 0)=I d$, by the Implicit Function Theorem there exists, for $\mu$ small enough, a unique solution $q_{A}^{\mu}=\mathrm{O}(\mu)$. By (3.2) $q_{A}^{\mu} \in C^{l+1}(\mathbf{R})$; moreover it is $C^{l-1}$ smooth in $A$. We define the $C^{l-1}$-maps $Q^{\mu}(\cdot), P^{\mu}(\cdot): \mathbf{T}^{n} \rightarrow \mathbf{R}$ by

$$
Q^{\mu}(A):=q_{A}^{\mu}(0), \quad P^{\mu}(A):=\dot{q}_{A}^{\mu}(0)
$$

By uniqueness we deduce that $q_{A}^{\mu}(s+t)=q_{A+\omega s}^{\mu}(t), \forall s, t \in \mathbf{R}$. For $t=0$ this yields
$q_{A}^{\mu}(s)=q_{A+\omega s}^{\mu}(0):=Q^{\mu}(A+\omega s) \quad$ and $\quad p_{A}^{\mu}(s)=p_{A+\omega s}^{\mu}(0):=P^{\mu}(A+\omega s), \quad \forall s \in \mathbf{R}$, proving the lemma.

We now define the functions $a^{\mu}(\psi)$ of (3.3). We impose that $\left(\omega t+A, I_{0}+a^{\mu}(\omega t+\right.$ $\left.A), Q^{\mu}(\omega t+A), P^{\mu}(\omega t+A)\right)$ satisfy the equations of motions (3.1); hence the functions $a^{\mu}(\psi)$ must satisfy the following system of equations

$$
\begin{equation*}
(\omega \cdot \nabla) a^{\mu}(\psi)=\mu g^{\mu}(\psi), \quad \text { where } g^{\mu}(\psi):=-\left(\nabla_{\psi} f\right)\left(\psi, Q^{\mu}(\psi)\right) \tag{3.4}
\end{equation*}
$$

In order to solve (3.4) we expand in Fourier series the functions $a^{\mu}(\psi)=\sum_{k \in \mathbf{Z}^{n}} a_{k} \mathrm{e}^{\mathrm{i} k \cdot \psi}$, $g^{\mu}(\psi)=\sum_{k \in \mathbf{Z}^{n}} g_{k} \mathrm{e}^{\mathrm{i} k \cdot \psi}$. Each Fourier coefficient $a_{k}$ must then satisfy

$$
\begin{equation*}
i(k \cdot \omega) a_{k}=\mu g_{k}, \quad \forall k \in \mathbf{Z}^{n} \tag{3.5}
\end{equation*}
$$

It is necessary for the existence of a solution that $g_{0}=\int_{\mathbf{T}^{n}} g^{\mu}(\psi) \mathrm{d} \psi=0$. This property can be checked directly, that is

Lemma 3.2. - We have

$$
\begin{equation*}
\int_{\mathbf{T}^{n}}\left(\nabla_{\psi} f\right)\left(\psi, Q^{\mu}(\psi)\right) \mathrm{d} \psi=0 \tag{3.6}
\end{equation*}
$$

Proof. - For a proof see [16] or Appendix A12 of [15].
Since $\omega$ satisfies (H1), if $f$ is smooth enough, then the function $a^{\mu}$ defined by

$$
\begin{equation*}
a^{\mu}(\psi)=\sum_{k \in \mathbf{Z}^{n}, k \neq 0} \frac{g_{k}}{i(k \cdot \omega)} \mathrm{e}^{\mathrm{i} k \cdot \psi}, \tag{3.7}
\end{equation*}
$$

which formally solves Eq. (3.5), is well defined and smooth. Indeed since $f \in C^{l}$ the function $g^{\mu}$ defined in (3.4) is $C^{l-1}$ and there exists $M>0$ such that $\left|g_{k}\right| \leqslant M /|k|^{l-1}$, $\forall k \in \mathbf{Z}^{n}, k \neq 0$. By (H1) it follows that $\left|a_{k}\right| \leqslant M /|k|^{l-1}|\omega \cdot k| \leqslant M|k|^{\tau} /\left(\gamma|k|^{l-1}\right)$. The proof of Theorem 3.1 is complete.

### 3.2. The new symplectic coordinates

In order to reduce to the previous case we want to put the tori $\mathcal{T}_{I_{0}}^{\mu}$ at the origin by a symplectic change of variables. Recalling that the tori $\mathcal{T}_{I_{0}}^{\mu}$ are isotropic submanifolds, i.e. that the symplectic 2 form $\mathrm{d} I \wedge \mathrm{~d} \psi+\mathrm{d} p \wedge \mathrm{~d} q$ vanishes on each tangent space to $\mathcal{T}_{I_{0}}^{\mu}$ (or equivalently $\mathrm{d} a^{\mu}(\psi) \wedge \mathrm{d} \psi+\mathrm{d} P^{\mu}(\psi) \wedge \mathrm{d} Q^{\mu}(\psi)=0$ ), we can prove the following lemma

LEMMA 3.3. - The transformation of coordinates $(J, v, \psi, u) \rightarrow(I, p, \varphi, q)$ defined on the space $\mathbf{R}^{n} \times \mathbf{R} \times \mathbf{T}^{n} \times \mathbf{R}$ by

$$
\begin{gather*}
I=a^{\mu}(\psi)+u \partial_{\psi} P^{\mu}(\psi)-v \partial_{\psi} Q^{\mu}(\psi)+J, \quad p=P^{\mu}(\psi)+v, \quad \varphi=\psi \\
q=Q^{\mu}(\psi)+u \tag{3.8}
\end{gather*}
$$

is symplectic.
Proof. - Since the tori $\mathcal{T}_{I_{0}}^{\mu}$ are isotropic the 1-form on $\mathbf{T}^{n} \lambda=\sum_{i=1}^{n} a_{i}^{\mu} \mathrm{d} \varphi_{i}+P^{\mu} \mathrm{d} Q^{\mu}$ is closed. Therefore there exists a unique $c \in \mathbf{R}^{n}$ and there exists $b^{\mu}: \mathbf{T}^{n} \rightarrow \mathbf{R}$ such that $\lambda=\sum_{i=1}^{n} c_{i} \mathrm{~d} \varphi_{i}+\mathrm{d} b^{\mu}$. Then transformation (3.8) has the generating function

$$
\Phi(J, v, \varphi, q)=(J+c) \cdot \varphi+v q+b^{\mu}(\varphi)-P^{\mu}(\varphi) Q^{\mu}(\varphi)+q P^{\mu}(\varphi)-v Q^{\mu}(\varphi)
$$

and hence is symplectic.
In the new coordinates each invariant torus $\mathcal{T}_{I_{0}}^{\mu}$ is simply described by $\left\{J=I_{0}, \psi \in\right.$ $\left.\mathbf{T}^{n}, u=v=0\right\}$ and the new Hamiltonian writes

$$
\mathcal{K}_{\mu}=E_{\mu}+\omega \cdot J+\frac{v^{2}}{2}+(\cos u-1)+P_{0}(\mu, u, \psi)
$$

where

$$
\begin{aligned}
P_{0}(\mu, u, \psi)= & \left(\cos \left(Q^{\mu}+u\right)-\cos Q^{\mu}+\left(\sin Q^{\mu}\right) u+1-\cos u\right) \\
& +\mu\left(f\left(\psi, Q^{\mu}+u\right)-f\left(\psi, Q^{\mu}\right)-\partial_{q} f\left(\psi, Q^{\mu}\right) u\right)
\end{aligned}
$$

and $E_{\mu}$ is the energy of the perturbed invariant torus $\mathcal{T}_{0}^{\mu}=\left\{\left(a^{\mu}(\psi), \psi, Q^{\mu}(\psi)\right.\right.$, $\left.\left.P^{\mu}(\psi)\right) ; \psi \in \mathbf{T}^{n}\right\}$. To Hamiltonian $\left(\mathcal{K}_{\mu}\right)$ is associated the quasi-periodically forced pendulum equation

$$
\begin{equation*}
-\ddot{u}+\sin u=\partial_{u} P_{0}(\mu, u, \omega t+A) \tag{3.9}
\end{equation*}
$$

of Lagrangian

$$
\begin{equation*}
L_{\mu}=\frac{\dot{u}^{2}}{2}+(1-\cos u)-P_{0}(\mu, u, \omega t+A) \tag{3.10}
\end{equation*}
$$

Since the Hamiltonian $\mathcal{K}_{\mu}$ is no more periodic in the variable $u$ we can not directly apply Theorem 2.3 and the arguments of the previous sections require some modifications. Arguing as in Lemma 2.1 we deduce that, there exists, for $\mu$ small enough, a unique 1bump pseudo-homoclinic solution $u_{A, \theta}^{\mu}(t)$, true solution of (3.9) in $(-\infty, \theta),(\theta,+\infty)$, satisfying all the properties of Lemma 2.1. Then we define the function $\mathcal{F}_{\mu}: \mathbf{T}^{n} \times \mathbf{R} \rightarrow \mathbf{R}$ as

$$
\begin{aligned}
\mathcal{F}_{\mu}(A, \theta)= & \int_{-\infty}^{\theta} \frac{\left(\dot{u}_{A, \theta}^{\mu}\right)^{2}}{2}+\left(1-\cos u_{A, \theta}^{\mu}\right)-P_{0}\left(\mu, u_{A, \theta}^{\mu}, \omega t+A\right) \mathrm{d} t \\
& +\int_{\theta}^{+\infty} \frac{\left(\dot{u}_{A, \theta}^{\mu}\right)^{2}}{2}+\left(1-\cos u_{A, \theta}^{\mu}\right)-P_{1}\left(\mu, u_{A, \theta}^{\mu}, \omega t+A\right) \mathrm{d} t+2 \pi \dot{q}_{A}^{\mu}(\theta)
\end{aligned}
$$

where, $q_{A}^{\mu}(t)=Q^{\mu}(\omega t+A)$ and $\forall i \in \mathbf{Z}$, we have set

$$
\begin{aligned}
P_{i}(\mu, u, \omega t+A)= & \left(\cos \left(q_{A}^{\mu}(t)+u\right)-\cos q_{A}^{\mu}(t)+\sin q_{A}^{\mu}(t)(u-2 \pi i)+1-\cos u\right) \\
& +\mu\left(f\left(\omega t+A, q_{A}^{\mu}(t)+u\right)-f\left(\omega t+A, q_{A}^{\mu}(t)\right)\right. \\
& \left.-\left(\partial_{q} f\right)\left(\omega t+A, q_{A}^{\mu}(t)\right)(u-2 \pi i)\right) .
\end{aligned}
$$

Since $u_{A, \theta}^{\mu}$ converges exponentially fast to 0 for $t \rightarrow-\infty$ and to $2 \pi$ for $t \rightarrow+\infty$ the above integrals are convergent. We define the "homoclinic function" $\mathcal{G}_{\mu}: \mathbf{T}^{n} \rightarrow \mathbf{R}$ as

$$
\begin{equation*}
\mathcal{G}_{\mu}(A)=\mathcal{F}_{\mu}(A, 0) \tag{3.11}
\end{equation*}
$$

There holds also $\mathcal{F}_{\mu}(A, \theta)=\mathcal{G}_{\mu}(A+\omega \theta), \forall \theta \in \mathbf{R}$. Arguing as in Lemma 2.4 we can prove the existence of $k$-bump pseudo-homoclinic solutions $u_{A, \theta}^{L}$, which are true solutions of (3.9) in each interval $\left(-\infty, \theta_{1}\right),\left(\theta_{i}, \theta_{i+1}\right)(i=1, \ldots, k-1),\left(\theta_{k},+\infty\right)$, and satisfy all the properties of Lemma 2.4. Then we define the " $k$-bump heteroclinic function"

$$
\begin{aligned}
& \mathcal{F}_{\mu}^{k}\left(A, \theta_{1}, \ldots, \theta_{k}\right) \\
& \quad=\int_{-\infty}^{\theta_{1}} \frac{\left(\dot{u}_{A, \theta}^{L}\right)^{2}}{2}+\left(1-\cos u_{A, \theta}^{L}\right)-P_{0}\left(\mu, u_{A, \theta}^{L}, \omega t+A\right) \mathrm{d} t+2 \pi \dot{q}_{A}^{\mu}\left(\theta_{1}\right) \\
& \quad+\sum_{i=1}^{k-1} \int_{\theta_{i}}^{\theta_{i+1}} \frac{\left(\dot{u}_{A, \theta}^{L}\right)^{2}}{2}+\left(1-\cos u_{A, \theta}^{L}\right)-P_{i}\left(\mu, u_{A, \theta}^{L}, \omega t+A\right) \mathrm{d} t+2 \pi \dot{q}_{A}^{\mu}\left(\theta_{i+1}\right) \\
& \quad+\int_{\theta_{k}}^{+\infty} \frac{\left(\dot{u}_{A, \theta}^{L}\right)^{2}}{2}+\left(1-\cos u_{A, \theta}^{L}\right)-P_{k}\left(\mu, u_{A, \theta}^{\mu}, \omega t+A\right) \mathrm{d} t-\left(I_{0}^{\prime}-I_{0}\right) \cdot A
\end{aligned}
$$

If $\partial_{\theta_{i}} \mathcal{F}_{\mu}^{k}\left(A, \theta_{1}, \ldots, \theta_{k}\right)=\left(\dot{u}_{A, \theta}^{L}\right)^{2}\left(\theta_{i}^{-}\right) / 2-\left(\dot{u}_{A, \theta}^{L}\right)^{2}\left(\theta_{i}^{+}\right) / 2=0$ then $u_{A, \theta}^{L}$ is a true solution of the quasi-periodically forced pendulum equation (3.9). As in the previous
section the variation in the action variables is given by the partial derivative with respect to $A$, that is

$$
\begin{align*}
\partial_{A} \mathcal{F}_{\mu}^{k}(A, \theta)= & \int_{-\infty}^{+\infty}-\mu\left(\partial_{\varphi} f\left(\omega t+A, q_{A}^{\mu}(t)+u_{A, \theta}^{L}(t)\right)\right. \\
& \left.-\partial_{\varphi} f\left(\omega t+A, q_{A}^{\mu}(t)\right)\right) \mathrm{d} t-\left(I_{0}^{\prime}-I_{0}\right) \tag{3.12}
\end{align*}
$$

Lemma 3.4. - Let $(A, \theta)$ be a critical point of $\mathcal{F}_{\mu}^{k}$. Then there exists a heteroclinic orbit connecting the tori $\mathcal{T}_{I_{0}}^{\mu}$ and $\mathcal{T}_{I_{0}^{\prime}}^{\mu}$.

Proof. - By (3.12) it is easy to verify that the solution of (3.1) $\left(I_{\mu}(t), \omega t+A, q_{A}^{\mu}+\right.$ $u_{A, \theta}^{L}, \dot{q}_{A}^{\mu}+\dot{u}_{A, \theta}^{L}$, with $I_{\mu}(t)=C-\mu \int_{0}^{t} \partial_{\varphi} f\left(\omega s+A, q_{A}^{\mu}(s)+u_{A, \theta}^{L}(s)\right) \mathrm{d} s$ and $C=I_{0}^{\prime}+$ $a_{\mu}(A)+\mu \int_{0}^{+\infty} \partial_{\varphi} f\left(\omega t+A, q_{A}^{\mu}(t)+u_{A, \theta}^{L}(t)\right)-\partial_{\varphi} f\left(\omega t+A, q_{A}^{\mu}(t)\right) \mathrm{d} t$, is a heteroclinic solution connecting $\mathcal{T}_{I_{0}}^{\mu}$ and $\mathcal{T}_{I_{0}^{\prime}}^{\mu}$.

Finally, arguing as in the proof of Theorem 2.3, we obtain
THEOREM 3.2. - Assume (H1) and let $\mathcal{G}_{\mu}$ satisfy the "splitting condition" 2.1. Then $\forall I_{0}, I_{0}^{\prime}$ with $\omega \cdot I_{0}=\omega \cdot I_{0}^{\prime}$, there is a heteroclinic orbit connecting the invariant tori $\mathcal{T}_{I_{0}}^{\mu}$ and $\mathcal{T}_{I_{0}^{\prime}}^{\mu}$. The same estimate on the diffusion time as in Theorem 2.3 holds.

A Taylor expansion in $\mu$ gives
Lemma 3.5. - For $\mu$ small enough

$$
\begin{equation*}
\mathcal{G}_{\mu}(A)=b+\mu M(A)+\mathrm{O}\left(\mu^{2}\right), \quad \forall A \in \mathbf{T}^{n} \tag{3.13}
\end{equation*}
$$

where $b=\Phi_{0}\left(q_{0}\right)$ and $M(A)$ is the Poincaré-Melnikov primitive $M(A)=\int_{\mathbf{R}}[f(\omega t+$ $\left.A, 0)-f\left(\omega t+A, q_{0}(t)\right)\right] \mathrm{d} t$.

Proof. - We develop in $\mu$ the Lagrangian $L_{\mu, i}=\dot{u}^{2} / 2+(1-\cos u)-P_{i}(\mu, u, \omega t+A)$ for $i=0,1$.

$$
\begin{aligned}
L_{\mu, i}= & \frac{\dot{u}^{2}}{2}+(1-\cos u)-\mu((u-2 \pi i-\sin u) \gamma+f(\omega t+A, u)-f(\omega t+A, 0) \\
& \left.-\partial_{q} f(\omega t+A, 0)(u-2 \pi i)\right)+\mathcal{R}_{i}(\mu, u, t)
\end{aligned}
$$

where $\gamma(t):=\partial_{\mu \mid \mu=0} q_{A}^{\mu}(t)$ and $|\mathcal{R}(\mu, u, t)|=\mathrm{O}\left(\mu^{2}(u-2 \pi i)^{2}\right)$. Hence $\mathcal{G}_{\mu}(A)=b+$ $\mu M^{*}(A)+\mathrm{O}\left(\mu^{2}\right)$ where

$$
\begin{aligned}
M^{*}(A)= & -\int_{-\infty}^{0}\left(q_{0}(t)-\sin q_{0}(t)\right) \gamma(t)+f\left(\omega t+A, q_{0}(t)\right)-f(\omega t+A, 0) \\
& -\partial_{q} f(\omega t+A, 0) q_{0}(t) \mathrm{d} t-\int_{0}^{+\infty}\left(q_{0}(t)-2 \pi-\sin q_{0}(t)\right) \gamma(t) \\
& +f\left(\omega t+A, q_{0}(t)\right)-f(\omega t+A, 0) \mathrm{d} t
\end{aligned}
$$

$$
\begin{equation*}
+\int_{0}^{+\infty} \partial_{q} f(\omega t+A, 0)\left(q_{0}(t)-2 \pi\right) \mathrm{d} t+2 \pi \dot{\gamma}(0) \tag{3.14}
\end{equation*}
$$

The lemma is proved if we show that $M^{*}(A)=M(A)$. Integrating by parts, since $\sin q_{0}(t)=\ddot{q}_{0}(t)$ and $-\ddot{\gamma}+\gamma=\partial_{q} f(\omega t+A, 0)$, we have

$$
\begin{equation*}
\int_{-\infty}^{0}\left(q_{0}(t)-\sin q_{0}(t)\right) \gamma(t)=\int_{-\infty}^{0} q_{0}(t) \partial_{q} f(\omega t+A, 0) \mathrm{d} t-\dot{q}_{0}(0) \gamma(0)+\pi \dot{\gamma}(0) \tag{3.15}
\end{equation*}
$$

and

$$
\begin{align*}
\int_{0}^{+\infty} & \left(q_{0}(t)-2 \pi-\sin q_{0}(t)\right) \gamma(t)  \tag{3.16}\\
& =\int_{0}^{+\infty}\left(q_{0}(t)-2 \pi\right) \partial_{q} f(\omega t+A, 0) \mathrm{d} t+\dot{q}_{0}(0) \gamma(0)+\pi \dot{\gamma}(0) .
\end{align*}
$$

Finally we deduce from (3.14), (3.15) and (3.16) that $M^{*}(A)=M(A)=\int_{\mathbf{R}}[f(\omega t+$ $\left.A, 0)-f\left(\omega t+A, q_{0}(t)\right)\right] \mathrm{d} t$.

THEOREM 3.3. - Assume (H1) and let M possess a proper minimum (or maximum) $A_{0}$, i.e. $\exists r>0$ such that $\inf _{\partial B_{r}\left(A_{0}\right)} M>M\left(A_{0}\right)$. Then, for $\mu$ small enough, the same statement as in Theorem 3.2 holds where the diffusion time is $T_{d}=\mathrm{O}((1 / \mu) \log (1 / \mu))$.

Remark 3.1. - By Theorems 3.1-3.2 we obtain that, for a priori-stable, isochronous, degenerate systems considered in [9] (see also [10])

$$
\mathcal{H}_{\varepsilon}=\varepsilon \omega \cdot I+\frac{p^{2}}{2}+\varepsilon^{d}(\cos q-1)+\mu f(\varphi, q) \quad \text { with } 1<d<2
$$

for $\mu=\delta \varepsilon^{d}$, $\delta$ being a small constant, the diffusion time is bounded by $T_{d}=\mathrm{O}\left(C(\delta) / \varepsilon^{d}\right)$. This improves the result of [9], which holds for $\mu=\mathrm{O}\left(\varepsilon^{d^{\prime}}\right), d^{\prime}>d / 2+3$, and provides the upper bound on the diffusion time $T_{d}=\mathrm{O}\left(1 / \varepsilon^{C+2(\tau+1)\left(2 d^{\prime}-1-d / 2\right)}\right), C$ being a suitable positive constant.

## 4. Splitting of whiskers

If the frequency vector $\omega=\omega_{\varepsilon}$ contains some "fast frequencies" $\omega_{i}=\mathrm{O}\left(1 / \varepsilon^{b}\right)$, $b>0, \varepsilon$ being a small parameter, and if the perturbation is analytic, the oscillations of the Melnikov function along some directions turn out to be exponentially small with respect to $\varepsilon$. Hence the development (3.13) will provide a valid measure of the splitting only for $\mu$ exponentially small with respect to $\varepsilon$. In order to justify the dominance of the Poincaré-Melnikov function when $\mu=\mathrm{O}\left(\varepsilon^{p}\right)$ we need more refined estimates for the error. However it turns out that the function $F_{\mu}(A, \theta)$ can not be easily analytically extended in a sufficiently wide complex strip (roughly speaking,
the condition $q_{A, \theta}^{\mu}(\operatorname{Re} \theta)=\pi$ appearing naturally when we try to extend the definition of $q_{A, \theta}^{\mu}$ to $\theta \in \mathbf{C}$ breaks analyticity). We bypass this problem considering the action functional evaluated on different "1-bump pseudo-homoclinic solutions" $Q_{A, \theta}^{\mu}$. This new "reduced action functional" $\widetilde{F}_{\mu}(A, \theta)=\widetilde{G}_{\mu}(A+\omega \theta)$ has the advantage to have an analytic extension in $(A, \theta)$ in a wide complex strip. Moreover we will show that the homoclinic functions $G_{\mu}, \widetilde{G}_{\mu}$ corresponding to both reductions are the same up to a change of variables of the torus close to the identity. This enables to recover enough information on the homoclinic function $G_{\mu}$ to construct diffusion orbits.

We assume that $f(\varphi, q)=(1-\cos q) f(\varphi), f(\varphi)=\sum_{k \in \mathbf{Z}^{n}} f_{k} \exp (\mathrm{i} k \cdot \varphi)$ and that there are $r_{i} \geqslant 0$ such that

$$
\begin{equation*}
\forall s \in \mathbf{N}, \exists D_{s}>0 \quad \text { such that } \quad\left|f_{k}\right| \leqslant \frac{D_{s}}{|k|^{s}} \exp \left(-\sum_{i=1}^{n} r_{i}\left|k_{i}\right|\right), \forall k \in \mathbf{Z}^{n} \tag{4.1}
\end{equation*}
$$

Condition (4.1) means that $f$ has a $C^{\infty}$ extension defined in

$$
D:=\left(\mathbf{R}+\mathrm{i}\left[-r_{1}, r_{1}\right]\right) \times \cdots \times\left(\mathbf{R}+\mathrm{i}\left[-r_{n}, r_{n}\right]\right)
$$

which is holomorphic w.r.t. the variables for which $r_{i}>0$ in $\left(\mathbf{R}+\mathrm{i} I_{1}\right) \times \cdots \times\left(\mathbf{R}+\mathrm{i} I_{n}\right)$, where $I_{i}=\{0\}$ if $r_{i}=0, I_{i}=\left(-r_{i}, r_{i}\right)$ if $r_{i}>0$. We denote the supremum of $|f|$ over $D$ as

$$
\begin{equation*}
\|f\|:=\sup _{A \in D}|f(A)| \tag{4.2}
\end{equation*}
$$

Note that we can take in (4.1) $D_{0}=\sup _{\varphi \in \mathbf{T}^{n}}|f(\varphi)| \leqslant\|f\|$. It will be used starting from Section 4.2.

### 4.1. The change of coordinates

In order to obtain an analytic "reduced functional" we perform a Lyapunov-Schmidt reduction with a supplementary space $\widetilde{E}_{\theta}$ to $\left\langle\dot{q}_{\theta}\right\rangle$ which depends analytically on $\theta$. We define $\psi_{0}: \mathbf{R} \rightarrow \mathbf{R}$ by $\psi_{0}(t)=\cosh ^{2} t /(1+\cosh t)^{3}$ and set $\psi_{\theta}(t)=\psi(t-\theta)$. Note that, since,

$$
\begin{gather*}
\int_{\mathbf{R}} \psi_{0}(t) \dot{q}_{0}(t) \mathrm{d} t=\bar{K} \neq 0,  \tag{4.3}\\
\widetilde{E}_{\theta}:=\left\{w \in H^{1}(\mathbf{R}) \mid \int_{\mathbf{R}} \psi_{\theta}(t) w(t) \mathrm{d} t=0\right\}
\end{gather*}
$$

is a supplementary space to $\left\langle\dot{q}_{\theta}\right\rangle$.
Remark 4.1. - We could use any function $\psi$ decaying at infinity, analytic in the complex strip $S:=\{\theta \in \mathbf{C}| | \operatorname{Im} \theta \mid<\pi / 2\}$ such that (4.3) holds. We use the function $\psi_{0}(t)$ instead of the more natural $\dot{q}_{0}(t)$ because the closest singularities to the real line of $\psi_{0}(z)$ are located in the complex plain at $\pm \mathrm{i} \pi$, while the closest singularities of $\dot{q}_{0}(t)$ are located at $\pm \mathrm{i}(\pi / 2)$. Hence near the boundary of the complex strip $S$ the analytic function $\psi_{0}$ is regular and this improves the bounds of Lemmas 4.6 and 4.7.

Arguing as in Lemma 2.1 we can prove, for $\mu$ small, the existence of " 1 bump pseudohomoclinics" $Q_{A, \theta}^{\mu} \in q_{\theta}+\widetilde{E}_{\theta}$ which is near $q_{\theta}$.

LEMmA 4.1. - There exist $\mu_{1}>0$ (independent of $\omega$ ) and $C_{4}>0$ such that for all $0<\mu<\mu_{1}$, for all $\theta \in \mathbf{R}$, for all $A \in \mathbf{R}^{n}$ there exists a unique function $Q_{A, \theta}^{\mu}(t): \mathbf{R} \rightarrow \mathbf{R}$, and a constant $\alpha_{A, \theta}^{\mu}$ smooth in $(A, \theta, \mu)$, such that
(i) $\quad-\ddot{Q}_{A, \theta}^{\mu}(t)+\sin Q_{A, \theta}^{\mu}(t)=\mu \sin Q_{A, \theta}^{\mu}(t) f(\omega t+A)+\alpha_{A, \theta}^{\mu} \psi_{\theta}(t)$;
(ii) $\quad \int_{\mathbf{R}}\left(Q_{A, \theta}^{\mu}(t)-q_{\theta}(t)\right) \psi_{\theta}(t) \mathrm{d} t=0$;
(iii) $\max \left(\left|Q_{A, \theta}^{\mu}(t)-q_{\theta}(t)\right|,\left|\dot{Q}_{A, \theta}^{\mu}(t)-\dot{q}_{\theta}(t)\right|\right) \leqslant C_{4} \mu \exp \left(-\frac{|t-\theta|}{2}\right)$;
(iv) $\max \left(\left|\partial_{A} Q_{A, \theta}^{\mu}(t)\right|,\left|\partial_{A} \dot{Q}_{A, \theta}^{\mu}(t)\right|,\left|\omega \cdot \partial_{A} Q_{A, \theta}^{\mu}(t)\right|\right.$,

$$
\left.\left|\omega \cdot \partial_{A} \dot{Q}_{A, \theta}^{\mu}(t)\right|\right) \leqslant C_{4} \mu \exp \left(-\frac{|t-\theta|}{2}\right)
$$

(v) $\quad Q_{A, \theta}^{\mu}(t)=Q_{A+k 2 \pi, \theta}^{\mu}(t), \quad \forall k \in \mathbf{Z}^{n}$ and $Q_{A, \theta+\eta}^{\mu}(t)=Q_{A+\omega \eta, \theta}^{\mu}(t-\eta), \forall \theta, \eta \in \mathbf{R}$.

We define the function $\widetilde{F}_{\mu}(A, \theta): \mathbf{T}^{n} \times \mathbf{R} \rightarrow \mathbf{R}$ as the action functional of Lagrangian 1.6 evaluated on the "1-bump pseudo-homoclinic solutions" $Q_{A, \theta}^{\mu}(t)$ obtained in Lemma 4.1, namely

$$
\begin{equation*}
\widetilde{F}_{\mu}(A, \theta):=\Phi_{\mu, A}\left(Q_{A, \theta}^{\mu}\right)=\int_{\mathbf{R}} \mathcal{L}_{\mu, A}\left(Q_{A, \theta}^{\mu}(t), \dot{Q}_{A, \theta}^{\mu}(t), t\right) \mathrm{d} t \tag{4.4}
\end{equation*}
$$

and $\widetilde{G}_{\mu}(A): \mathbf{T}^{n} \rightarrow \mathbf{R}$ as $\widetilde{G}_{\mu}(A)=\widetilde{F}_{\mu}(A, 0)$. By Lemma (4.1)(v) the following invariance property holds $\widetilde{F}_{\mu}(A, \theta+\eta)=\widetilde{F}_{\mu}(A+\omega \eta, \theta), \forall \theta, \eta \in \mathbf{R}$; in particular $\widetilde{F}_{\mu}(A, \theta)=$ $\widetilde{G}_{\mu}(A+\omega \theta), \forall \theta \in \mathbf{R}$.

Lemma 4.1(i)-(ii) suggests that $\left\{Q_{A, \theta}^{\mu} \mid \theta \in \mathbf{R}\right\}$ is "natural constraint" for $\Phi_{\mu, A}$. For completeness we prove it in detail (it is the analogue of (2.6)), together with useful bounds of the derivatives of $\widetilde{F}_{\mu}$ (and $\widetilde{G}_{\mu}$ ) up to order 2.

LEMMA 4.2. - Let $0<\mu<\min \left\{\mu_{1},\left(\bar{K} / 8 C_{4}\right)\right\}$, where $\bar{K}$ is defined in (4.3). If $\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)=0$ then $Q_{A, \theta}^{\mu}$ is a true solution of (1.5). More precisely

$$
\begin{equation*}
\left|\alpha_{A, \theta}^{\mu}\right| \leqslant \frac{2}{|\bar{K}|}\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right| . \tag{4.5}
\end{equation*}
$$

Moreover we also have

$$
\begin{equation*}
\left|\nabla^{s} \widetilde{G}_{\mu}(A)\right|=\mathrm{O}(\mu), \quad\left|\nabla^{s} \widetilde{F}_{\mu}(A, \theta)\right|=\mathrm{O}(\mu), \quad s=1,2 \tag{4.6}
\end{equation*}
$$

Proof. - We have

$$
\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)=\int_{\mathbf{R}} \dot{Q}_{A, \theta}^{\mu}(t) \partial_{\theta} \dot{Q}_{A, \theta}^{\mu}(t)
$$

$$
+\left(\sin Q_{A, \theta}^{\mu}(t)-\mu \sin Q_{A, \theta}^{\mu}(t) f(\omega t+A)\right) \partial_{\theta} Q_{A, \theta}^{\mu}(t) \mathrm{d} t
$$

Integrating by parts and using also that $Q_{A, \theta}^{\mu}(t)$ solves (i) in Lemma 4.1 we obtain

$$
\begin{equation*}
\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)=\left[\partial_{\theta} Q_{A, \theta}^{\mu}(t) \dot{Q}_{A, \theta}^{\mu}(t)\right]_{-\infty}^{+\infty}+\int_{\mathbf{R}} \partial_{\theta} Q_{A, \theta}^{\mu}(t) \alpha_{A, \theta}^{\mu} \psi_{\theta}(t) \mathrm{d} t \tag{4.7}
\end{equation*}
$$

By Lemma 4.1(iii)-(iv) the boundary term in (4.7) vanishes. Moreover, deriving w.r.t. $\eta$ the second equality in lemma 4.1(v), we get

$$
\begin{equation*}
\partial_{\theta} Q_{A, \theta}^{\mu}(t)=\omega \cdot \partial_{A} Q_{A, \theta}^{\mu}(t)-\dot{Q}_{A, \theta}^{\mu}(t) \tag{4.8}
\end{equation*}
$$

Since, by Lemma 4.1(ii), $\int_{\mathbf{R}} Q_{A, \theta}^{\mu} \psi_{\theta} \mathrm{d} t=\int_{\mathbf{R}} q_{\theta} \psi_{\theta} \mathrm{d} t$ is independent of $A, \int_{\mathbf{R}} \partial_{A} Q_{A, \theta}^{\mu} \psi_{\theta}(t)$ $=0$. Hence

$$
\begin{align*}
\partial_{\theta} \widetilde{F}_{\mu}(A, \theta) & =\alpha_{A, \theta}^{\mu} \int_{\mathbf{R}}\left(\omega \cdot \partial_{A} Q_{A, \theta}^{\mu}(t)-\dot{Q}_{A, \theta}^{\mu}(t)\right) \psi_{\theta}(t) \mathrm{d} t \\
& =-\alpha_{A, \theta}^{\mu} \int_{\mathbf{R}} \dot{Q}_{A, \theta}^{\mu}(t) \psi_{\theta}(t) \mathrm{d} t \\
& =-\alpha_{A, \theta}^{\mu}\left(\bar{K}+\int_{\mathbf{R}}\left(\dot{Q}_{A, \theta}^{\mu}(t)-\dot{q}_{\theta}(t)\right) \psi_{\theta} \mathrm{d} t\right) \tag{4.9}
\end{align*}
$$

By Lemma 4.1(iii), noting that $\left|\psi_{\theta}\right| \leqslant 1$, we have

$$
\begin{equation*}
\left|\int_{\mathbf{R}}\left(\dot{Q}_{A, \theta}^{\mu}(t)-\dot{q}_{\theta}(t)\right) \psi_{\theta}(t) \mathrm{d} t\right| \leqslant 4 C_{4} \mu \tag{4.10}
\end{equation*}
$$

Finally by (4.10) and (4.9), if $0<\mu 4 C_{4} \leqslant|\bar{K}| / 2$, then

$$
\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right| \geqslant\left|\alpha_{A, \theta}^{\mu}\right|\left(|\bar{K}|-4 \mu C_{4}\right) \geqslant\left|\alpha_{A, \theta}^{\mu}\right||\bar{K}| / 2
$$

namely (4.5).
The proof of (4.6) can be performed with similar computations. Let us for instance justify that $\partial_{\theta}^{2} \widetilde{F}_{\mu}(A, \theta)=\mathrm{O}(\mu)$. We have

$$
\begin{align*}
\partial_{\theta} \widetilde{F}_{\mu}(A, \theta) & =-\alpha_{A, \theta}^{\mu} \int_{\mathbf{R}} \dot{Q}_{A, \theta}^{\mu}(t) \psi_{\theta}(t) \mathrm{d} t \\
& =\int_{\mathbf{R}} \ddot{Q}_{A, \theta}^{\mu} \dot{Q}_{A, \theta}^{\mu}-\sin Q_{A, \theta}^{\mu} \dot{Q}_{A, \theta}^{\mu}+\mu f(\omega t+A) \sin Q_{A, \theta}^{\mu} \dot{Q}_{A, \theta}^{\mu} \mathrm{d} t \\
& =\mu \int_{\mathbf{R}} f(\omega t+A) \sin Q_{A, \theta}^{\mu} \dot{Q}_{A, \theta}^{\mu} \mathrm{d} t \tag{4.11}
\end{align*}
$$

Hence

$$
\partial_{\theta}^{2} \widetilde{F}_{\mu}(A, \theta)=\mu \int_{\mathbf{R}}\left(\cos Q_{A, \theta}^{\mu} \partial_{\theta} Q_{A, \theta}^{\mu} \dot{Q}_{A, \theta}^{\mu}+\sin Q_{A, \theta}^{\mu} \partial_{\theta} \dot{Q}_{A, \theta}^{\mu}\right) f(\omega t+A) \mathrm{d} t
$$

The bound of $\partial_{\theta}^{2} \widetilde{F}_{\mu}(A, \theta)$ (independent of $\omega$ ) is then a straightforward consequence of (4.8) and Lemma 4.1(i)-(iii)-(iv).

The relation between the two reduced action functionals $F_{\mu}(A, \theta)=\Phi_{\mu, A}\left(q_{A, \theta}^{\mu}\right)$ and $\widetilde{F}_{\mu}(A, \theta)=\Phi_{\mu, A}\left(Q_{A, \theta}^{\mu}\right)$ is given below: they are the same "up to a change of variables near the identity". The next theorem is formulated to handle also non-analytic perturbations $f$.

THEOREM 4.1. - Assume that $f$ is smooth. There exist $\mu_{2}>0$ and $C_{5}>0$ (independent of $\omega$ ) such that, for $0<\mu<\mu_{2}$, there exist a Lipschitz continuous function $h_{\mu}: \mathbf{T}^{n} \times \mathbf{R} \rightarrow \mathbf{R}$, with $\left|h_{\mu}(A, \theta)\right| \leqslant C_{5} \mu,\left|h_{\mu}\left(A^{\prime}, \theta^{\prime}\right)-h_{\mu}(A, \theta)\right| \leqslant C_{5} \mu\left(\left|A^{\prime}-A\right|+\right.$ $\left.\left|\theta^{\prime}-\theta\right|\right), h_{\mu}(A, \theta+\eta)=h_{\mu}(A+\eta \omega, \theta)$, such that

$$
\begin{equation*}
F_{\mu}(A, \theta)=\widetilde{F}_{\mu}\left(A, \theta+h_{\mu}(A, \theta)\right) \tag{4.12}
\end{equation*}
$$

In particular, setting $g_{\mu}(A)=h_{\mu}(A, 0), \psi_{\mu}: \mathbf{T}^{n} \rightarrow \mathbf{T}^{n}$ defined by $\psi_{\mu}(A)=A+g_{\mu}(A) \omega$ is a homeomorphism provided that $C_{5} \mu \leqslant 1 / 2$ and satisfies

$$
G_{\mu}=\widetilde{G}_{\mu} \circ \psi_{\mu}
$$

The inverse of $\psi_{\mu}$ writes $\psi_{\mu}^{-1}(A)=A+k_{\mu}(A) \omega$ where $k_{\mu}$ is $2 C_{5} \mu$-Lipschitz continuous. Moreover, if $f$ is analytic and satisfies (4.1) with $r_{i}>0$ for all $i=1, \ldots, n$, then the homeomorphism $\psi_{\mu}$ is a real analytic diffeomorphism.

In order to prove Theorem 4.1 we need the next two lemmas, proved in Appendix A. The next lemma states that $Q_{A, \theta}^{\mu}$ belongs to the supplementary space $E_{\theta+l_{\mu}(A, \theta)}$ for some $l_{\mu}(A, \theta)=\mathrm{O}(\mu)$.

Lemma 4.3. - For $\mu$ small enough (independently of $\omega$ ) there exists a smooth function $l_{\mu}(A, \theta)$ with $l_{\mu}(A, \theta)=\mathrm{O}(\mu), \nabla l_{\mu}(A, \theta)=\mathrm{O}(\mu), l_{\mu}(A, \theta+\eta)=l_{\mu}(A+$ $\eta \omega, \theta)$ such that $Q_{A, \theta}^{\mu}\left(\theta+l_{\mu}(A, \theta)\right)=\pi$.

Define $V_{\mu}(A, \theta):=F_{\mu}\left(A, \theta+l_{\mu}(A, \theta)\right)=\Phi_{\mu, A}\left(q_{A, \theta+l_{\mu}(A, \theta)}^{\mu}\right)$.
LEMMA 4.4. - There exists a positive constant $C_{6}$ such that, for all $(A, \theta) \in \mathbf{T}^{n} \times \mathbf{R}$, there holds

$$
\begin{equation*}
\left|\widetilde{F}_{\mu}(A, \theta)-V_{\mu}(A, \theta)\right| \leqslant C_{6}\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|^{2} \tag{4.13}
\end{equation*}
$$

In particular if $\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)=0$ then $\widetilde{F}_{\mu}(A, \theta)=V_{\mu}(A, \theta)$. In addition

$$
\begin{equation*}
\left|\nabla \widetilde{F}_{\mu}(A, \theta)-\nabla V_{\mu}(A, \theta)\right| \leqslant C_{6} \mu\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right| \tag{4.14}
\end{equation*}
$$

Proof of Theorem 4.1. - By Lemma 4.3, there is a smooth function $\bar{l}_{\mu}$ such that $\bar{l}_{\mu}(A, \theta)=\mathrm{O}(\mu), \nabla \bar{l}_{\mu}(A, \theta)=\mathrm{O}(\mu), \bar{l}_{\mu}(A, \theta+\eta)=\bar{l}_{\mu}(A+\eta \omega, \theta)$ and $F_{\mu}(A, \theta)=$ $V_{\mu}\left(A, \theta+\bar{l}_{\mu}(A, \theta)\right) . \bar{l}_{\mu}(A, \theta)$ is the unique solution of the equation $x=-l_{\mu}(A, \theta+x)$. It is enough to find, for all $\theta, y=\bar{h}_{\mu}(A, \theta)$ such that

$$
\begin{equation*}
V_{\mu}(A, \theta)=\widetilde{F}_{\mu}(A, \theta+y) \tag{4.15}
\end{equation*}
$$

Indeed, then

$$
F_{\mu}(A, \theta)=V_{\mu}\left(A, \theta+\bar{l}_{\mu}(A, \theta)\right)=\widetilde{F}_{\mu}\left(A,\left(\theta+\bar{l}_{\mu}(A, \theta)\right)+\bar{h}_{\mu}\left(A, \theta+\bar{l}_{\mu}(A, \theta)\right)\right)
$$

and $h_{\mu}(A, \theta)$ will be defined by

$$
\begin{equation*}
h_{\mu}(A, \theta)=\bar{l}_{\mu}(A, \theta)+\bar{h}_{\mu}\left(A, \theta+\bar{l}_{\mu}(A, \theta)\right) \tag{4.16}
\end{equation*}
$$

Note that if $\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)=0$ then, by Lemma 4.4, Eq. (4.15) is solved by $y=0$. In general we look for $y$ of the form $y=\partial_{\theta} \widetilde{F}_{\mu}(A, \theta) z$. Then we can write

$$
\begin{align*}
\widetilde{F}_{\mu}(A, \theta+y) & =\widetilde{F}_{\mu}(A, \theta)+\partial_{\theta} \widetilde{F}_{\mu}(A, \theta) y+R_{\mu}(A, \theta, y) y^{2}  \tag{4.17}\\
& =\widetilde{F}_{\mu}(A, \theta)+\left(\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right)^{2} z+R_{\mu}\left(A, \theta, \partial_{\theta} \widetilde{F}_{\mu}(A, \theta) z\right)\left(\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right)^{2} z^{2}
\end{align*}
$$

where

$$
R_{\mu}(A, \theta, y)=\frac{1}{y^{2}}\left[\widetilde{F}_{\mu}(A, \theta+y)-\widetilde{F}_{\mu}(A, \theta)-\partial_{\theta} \widetilde{F}_{\mu}(A, \theta) y\right]
$$

is smooth and, by the estimates (4.6) on the derivatives of $\widetilde{F}_{\mu}$, satisfies $R_{\mu}(A, \theta, y)=$ $\mathrm{O}(\mu), \partial_{y} R_{\mu}(A, \theta, y)=\mathrm{O}(\mu /|y|)$. By (4.17) Eq. (4.15) is then equivalent to

$$
\frac{V_{\mu}(A, \theta)-\widetilde{F}_{\mu}(A, \theta)}{\left(\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right)^{2}}=z+R_{\mu}\left(A, \theta, \partial_{\theta} \widetilde{F}_{\mu}(A, \theta) z\right) z^{2}
$$

We have $R_{\mu}\left(A, \theta, \partial_{\theta} \widetilde{F}_{\mu}(A, \theta) z\right) z^{2}=\mathrm{O}\left(\mu z^{2}\right)$ and $\partial_{z}\left(R_{\mu}\left(A, \theta, \partial_{\theta} \widetilde{F}_{\mu}(A, \theta) z\right) z^{2}\right)=\mathrm{O}(\mu z)$. By the contraction mapping theorem, for $\mu$ small enough, for all $u \in \mathbf{R}$ such that $|u|<2 C_{6}$, there exists a unique solution $z=\varphi(\mu, A, \theta, u)$ of the equation

$$
\begin{equation*}
u=z+R_{\mu}\left(A, \theta, \partial_{\theta} \widetilde{F}_{\mu}(A, \theta) z\right) z^{2} \tag{4.18}
\end{equation*}
$$

such that $|z|<3 C_{6}$. Moreover, the function $\varphi$ defined in this way is smooth and is real analytic if $\widetilde{F}_{\mu}$ is real analytic. Setting

$$
\begin{equation*}
\bar{h}_{\mu}(A, \theta):=\varphi\left(\mu, A, \theta, \frac{V_{\mu}(A, \theta)-\widetilde{F}_{\mu}(A, \theta)}{\left(\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right)^{2}}\right) \partial_{\theta} \widetilde{F}_{\mu}(A, \theta) \tag{4.19}
\end{equation*}
$$

if $\partial_{\theta} \widetilde{F}_{\mu}(A, \theta) \neq 0$ and $\bar{h}_{\mu}(A, \theta)=0$ if $\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)=0$, we get a continuous function $\bar{h}_{\mu}$ which satisfies (4.15) and $\left|\bar{h}_{\mu}(A, \theta)\right| \leqslant 3 C_{6}\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|$, which implies $\left|\bar{h}_{\mu}\right|=\mathrm{O}(\mu)$. Moreover $\bar{h}_{\mu}$ is the unique function that enjoys these properties. By (4.19) the restriction of $\bar{h}_{\mu}$ to

$$
U_{\mu}:=\left\{(A, \theta) \in \mathbf{T}^{n} \times \mathbf{R}: \partial_{\theta} \widetilde{F}_{\mu}(A, \theta) \neq 0\right\}
$$

is smooth. Deriving the identity $V_{\mu}(A, \theta)=\widetilde{F}_{\mu}\left(A, \theta+\bar{h}_{\mu}(A, \theta)\right)$ we obtain

$$
\left(\partial_{\theta} \widetilde{F}_{\mu}\right)\left(A, \theta+\bar{h}_{\mu}(A, \theta)\right) \nabla \bar{h}_{\mu}(A, \theta)=\nabla V_{\mu}(A, \theta)-\nabla \widetilde{F}_{\mu}\left(A, \theta+\bar{h}_{\mu}(A, \theta)\right)
$$

for $(A, \theta) \in U_{\mu}$. By (4.6)

$$
\begin{aligned}
\nabla \widetilde{F}_{\mu}\left(A, \theta+\bar{h}_{\mu}(A, \theta)\right) & =\nabla \widetilde{F}_{\mu}(A, \theta)+\mathrm{O}\left(\mu\left|\bar{h}_{\mu}(A, \theta)\right|\right) \\
& =\nabla \widetilde{F}_{\mu}(A, \theta)+\mathrm{O}\left(\mu\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right)
\end{aligned}
$$

Hence, if $\partial_{\theta} \widetilde{F}_{\mu}(A, \theta) \neq 0$ then, by (4.14),

$$
\nabla \bar{h}_{\mu}=\frac{\nabla V_{\mu}(A, \theta)-\nabla \widetilde{F}_{\mu}(A, \theta)+\mathrm{O}\left(\mu\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right)}{\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)(1+\mathrm{O}(\mu))}=\mathrm{O}(\mu)
$$

uniformly in $U_{\mu}$. Since $\bar{h}_{\mu}$ is continuous and $\bar{h}_{\mu}(A, \theta)=0$, if $(A, \theta) \notin U_{\mu}$, the Lipschitz continuity of $\bar{h}_{\mu}$ follows.

Regarding the properties of $h_{\mu}$, we observe that $\bar{h}_{\mu}(A, \theta+\eta)=\bar{h}_{\mu}(A+\eta \omega, \theta)$, which is a consequence of uniqueness.

Hence, by (4.16) and since $\bar{l}_{\mu}(A, \theta+\eta)=\bar{l}_{\mu}(A+\eta \omega, \theta), h_{\mu}(A, \theta+\eta)=h_{\mu}(A+$ $\eta \omega, \theta)$. Moreover, since $\left|\nabla \bar{l}_{\mu}(A, \theta)\right|=\mathrm{O}(\mu)$, by the Lipschitz continuity of $\bar{h}_{\mu}$, there is a constant $C_{5}$ such that $\left|h_{\mu}(A, \theta)\right| \leqslant C_{5} \mu$ and $h_{\mu}$ is Lipschitz continuous with ratio $C_{5} \mu$.

To complete the proof, we remark that for a given $A \in \mathbf{T}^{n}$, the equation in $x$, $\psi_{\mu}(A+x \omega)=A$ is equivalent to

$$
\begin{equation*}
x=-g_{\mu}(A+\omega x)=-h_{\mu}(A, x) \tag{4.20}
\end{equation*}
$$

where $h_{\mu}$ is $C_{5} \mu$-Lipschitz continuous. By the contracting mapping theorem, if $C_{5} \mu \leqslant$ $1 / 2$ then (4.20) has a unique solution $k_{\mu}(A) \in \mathbf{R}$. Moreover, one can prove without much effort that $\left|k_{\mu}(A)\right| \leqslant C_{5} \mu$ and $k_{\mu}$ is Lipschitz continuous of ratio $2 C_{5} \mu$. As a consequence $\psi_{\mu}$ is a homeomorphism and its inverse is defined by $\psi_{\mu}^{-1}(A)=A+$ $k_{\mu}(A) \omega$.

The proof of the analyticity of $\psi_{\mu}$ when $f$ is analytic, is given in the next Remark 4.2 and in Appendix A.

By Theorem 4.1 we deduce the following lemma which enables us to apply the shadowing Theorem 2.3 once a "splitting condition" for $\widetilde{G}_{\mu}$ is verified.

LEMMA 4.5. - Assume that $\widetilde{G}_{\mu}$ satisfies the "splitting condition" 2.1 for some bounded open set $U, A_{0} \in U$, with constants $\delta$ and $\alpha$. Then, for $\mu|\omega| C_{5}<1 / 2, G_{\mu}$ too satisfies the "splitting condition" with $A_{0}^{\prime}=\psi_{\mu}^{-1}\left(A_{0}\right), U^{\prime}=\psi_{\mu}^{-1}(U)$ and the constants $\delta^{\prime}=\delta$ and $\alpha^{\prime}=\alpha / 2$.

Proof. - In Appendix A.
Remark 4.2. - Assume that $r_{i}>0$ for all $i$ (i.e. that the perturbation $f$ is analytic). Then we can prove, using the arguments of the next subsection, that the homoclinic function $G_{\mu}(\cdot)=F_{\mu}(\cdot, 0)$ can be extended to a complex analytic function over the interior of $D$. Hence $F_{\mu}(A, \theta)=G_{\mu}(A+\omega \theta)$ can be defined in an open neighbourhood of $\mathbf{T}^{n} \times \mathbf{R}$ in $\left(\mathbf{T}^{n}+\mathrm{i} \mathbf{R}^{n}\right) \times \mathbf{C}$, so that the extension is analytic. One could check that $l_{\mu}$ and $V_{\mu}$, defined in Lemma 4.3 have analytic extensions too, and that the
inequality of Lemma 4.4 still holds in the new set of definition. Moreover in the next Lemma 4.6 it is proved that $\widetilde{F}_{\mu}$ is analytic w.r.t. $(A, \theta)$. As a consequence, $\left(V_{\mu}(A, \theta)-\widetilde{F}_{\mu}(A, \theta)\right) /\left(\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right)^{2}$ is real analytic, and so is the function $h_{\mu}$ defined in the proof of Theorem 4.1. Therefore if $r_{i}>0$ for all $i$, then the homeomorphism $\psi_{\mu}$ defined in Theorem 4.1 is a real analytic diffeomorphism. For completeness, the proof of the claims included in this remark is spelled out in Appendix A.

### 4.2. Analytic extension

The unpertubed homoclinic $q_{0}(t)=4 \arctan \mathrm{e}^{t}$ can be extended to a holomorphic function over the strip $S:=\mathbf{R}+\mathrm{i}(-\pi / 2, \pi / 2)$. Moreover Eq. (1.5) may be considered also for complex values of $q$ and, for $\mu=0, q_{\theta}$ is a solution of (1.5) for all $\theta \in S$. The unperturbed manifold of critical points $Z=\left\{q_{\theta} \mid \theta \in S\right\}$ can be seen as a holomorphic curve in the complex affine space $q_{\theta}+H^{1}(\mathbf{R}, \mathbf{C})$.

We want to perform the Lyapunov-Schmidt reduction on the complex strip

$$
S_{\sigma}:=\mathbf{R}+\mathrm{i}\left(-\left(\frac{\pi}{2}-\sigma\right), \frac{\pi}{2}-\sigma\right)
$$

for $\sigma \in(0, \pi / 2)$. We have

$$
\begin{aligned}
& \dot{q}_{\theta}(z)=\frac{2}{\cosh (z-\theta)}, \quad \ddot{q}_{\theta}(z)=\sin q_{\theta}(z)=-2 \frac{\sinh (z-\theta)}{\cosh ^{2}(z-\theta)} \\
& \quad\left(1-\cos q_{\theta}(z)\right)=\frac{2}{\cosh ^{2}(z-\theta)}
\end{aligned}
$$

Assume that $\theta \in S_{\sigma}, \operatorname{Re}(\theta)=0$. The following estimates hold, where $t \in \mathbf{R}$

$$
\begin{align*}
\left|\dot{q}_{\theta}(t)\right| & \leqslant \frac{C}{\min \{(|t|+\sigma), 1\}} \exp (-|t|)  \tag{4.21}\\
\left|\sin q_{\theta}(t)\right| & \leqslant \frac{C}{\min \left\{(|t|+\sigma)^{2}, 1\right\}} \exp (-|t|)  \tag{4.22}\\
\left|\cos q_{\theta}(t)\right| & \leqslant \frac{C}{\min \left\{(|t|+\sigma)^{2}, 1\right\}} ;  \tag{4.23}\\
\frac{1}{\left|\dot{q}_{\theta}(t)\right|} & \leqslant C \exp (|t|) \min \{(|t|+\sigma), 1\} \tag{4.24}
\end{align*}
$$

In what follows we consider the Banach spaces

$$
X=\left\{w \in C^{2}(\mathbf{R}, \mathbf{C}) \mid \sup _{t} \exp (|t| / 2)(|w(t)|+|\dot{w}(t)|+|\ddot{w}(t)|)<+\infty\right\}
$$

and

$$
\bar{X}=\{w \in X \mid w(0)=0\}
$$

endowed with norm

$$
\|w\|_{2, \sigma}=\sup _{|t|>1}(|w(t)|+|\dot{w}(t)|+|\ddot{w}(t)|) \exp \left(\frac{|t|}{2}\right)
$$

$$
+\sup _{|t|<1}\left(\frac{|w(t)|}{(|t|+\sigma)^{2}}+\frac{|\dot{w}(t)|}{(|t|+\sigma)}+|\ddot{w}(t)|\right) .
$$

Note that the function $\psi_{0}(t)$ can be extended to a holomorphic function on $\mathbf{R}+$ $\mathrm{i}(-\pi, \pi)$. With the definition for $\|f\|$ given in (4.2), we have

Lemma 4.6. - There exist positive constants $\eta, C_{7}$ such that for all $\sigma \in(0, \pi / 2)$, $\forall 0<\mu \leqslant\left(\eta \sigma^{3}\right) /\|f\|$, for all $\omega$, for all $A \in D$, for all $\theta \in S_{\sigma}$ there exist a unique $Q_{A, \theta}^{\mu}: \mathbf{R} \rightarrow \mathbf{C}$ and a unique $\alpha_{A, \theta}^{\mu} \in \mathbf{C}$ such that

- $Q_{A, \theta}^{\mu}=q_{\theta+v_{A, \theta}^{\mu}}^{\mu}+w_{A, \theta}^{\mu}$, where $v_{A, \theta}^{\mu} \in \mathbf{C}, w \in \bar{X}$ and $\left\|w_{A, \theta}^{\mu}\right\|_{2, \sigma}+\left|v_{A, \theta}^{\mu}\right|+\left|\alpha_{A, \theta}^{\mu}\right| \leqslant$ $C_{7} \mu\|f\| / \sigma^{2}$;
- $-\ddot{Q}_{A, \theta}^{\mu}(t)+\sin Q_{A, \theta}^{\mu}(t)=\mu \sin Q_{A, \theta}^{\mu}(t) f(\omega t+A)+\alpha_{A, \theta}^{\mu} \psi_{\theta}(t)$;
- $\int_{\mathbf{R}}\left(Q_{A, \theta}^{\mu}(t)-q_{\theta}(t)\right) \psi_{\theta}(t) \mathrm{d} t=0$.

Moreover $Q_{A, \theta}^{\mu}$ and $\alpha_{A, \theta}^{\mu}$ depend analytically on $\theta$ and on the $A_{i}$ for which $r_{i}>0$.
Proof. - 1st step. Let us consider the Banach space

$$
Y=\left\{v \in C(\mathbf{R}, \mathbf{C})\left|\sup _{t}\right| v(t) \left\lvert\, \exp \left(\frac{|t|}{2}\right)<+\infty\right.\right\}
$$

endowed with norm $\|v\|_{-1, \sigma}=\sup _{|t|>1}|v(t)| \exp \left(\frac{|t|}{2}\right)+\sup _{|t|<1}(|t|+\sigma)|v(t)|$. Let $\theta \in S_{\sigma}$ be given once for all. We may assume without loss of generality that $\operatorname{Re}(\theta)=0$.

For $\theta^{\prime} \in S_{\sigma / 2}$ such that $\left|\theta^{\prime}-\theta\right| \leqslant \sigma / 2$ we introduce the linear operator $L_{\theta^{\prime}}: \bar{X} \times \mathbf{C} \rightarrow Y$ defined by

$$
L_{\theta^{\prime}}(w, \alpha)=-\ddot{w}+\left(\cos q_{\theta^{\prime}}\right) w-\alpha \psi_{\theta} .
$$

Using that $\dot{q}_{\theta^{\prime}}$ is a solution of $-\ddot{y}+\cos q_{\theta^{\prime}} y=0$ and that $\dot{q}_{\theta^{\prime}}$ does not vanish anywhere, we can compute the inverse of $L_{\theta^{\prime}}$. It is given by $L_{\theta^{\prime}}^{-1}(g)=(w, \alpha)$ with

$$
\begin{gather*}
\alpha=-\frac{\int_{\mathbf{R}} g(t) \dot{q}_{\theta^{\prime}}(t) \mathrm{d} t}{\int_{\mathbf{R}} \psi_{\theta}(t) \dot{q}_{\theta^{\prime}}(t) \mathrm{d} t},  \tag{4.25}\\
w(t)=\dot{q}_{\theta^{\prime}}(t)\left[\int_{0}^{t}-\frac{1}{\dot{q}_{\theta^{\prime}}^{2}(s)}\left(\int_{-\infty}^{s}\left(g(\sigma)+\alpha \psi_{\theta}(\sigma)\right) \dot{q}_{\theta^{\prime}}(\sigma) d \sigma\right) \mathrm{d} s\right]  \tag{4.26}\\
=\dot{q}_{\theta^{\prime}}(t)\left[\int_{0}^{t} \frac{1}{\dot{q}_{\theta^{\prime}}^{2}(s)}\left(\int_{s}^{+\infty}\left(g(\sigma)+\alpha \psi_{\theta}(\sigma)\right) \dot{q}_{\theta^{\prime}}(\sigma) d \sigma\right) \mathrm{d} s\right] \tag{4.27}
\end{gather*}
$$

Note that since $\left|\theta-\theta^{\prime}\right| \leqslant \sigma / 2, \operatorname{Re}\left(\theta^{\prime}\right) \leqslant \sigma / 2$. Therefore estimates (4.21)-(4.24) hold as well (with perhaps different constants) when $\theta$ is replaced by $\theta^{\prime}$. We derive from (4.25)-(4.27) that

$$
\begin{equation*}
|\alpha|+\|w\|_{2, \sigma} \leqslant \frac{C}{\sigma}\|g\|_{-1, \sigma} \tag{4.28}
\end{equation*}
$$

2nd step. We shall search $Q$ as $Q=q_{\theta+v}+w$ with $|\nu|<\sigma / 2, w \in \bar{X}$. Let $B$ denote the open ball of radius $\sigma / 2$ in $\mathbf{C}$ centered at 0 . Let $J_{\mu}: B \times \bar{X} \times \mathbf{C} \rightarrow Y \times \mathbf{C}$ be defined by

$$
\begin{aligned}
J_{\mu}(v, w, \alpha)= & \left(-\ddot{q}_{\theta+v}-\ddot{w}+\sin \left(q_{\theta+v}+w\right)-\mu \sin \left(q_{\theta+v}+w\right) f(\varphi)-\alpha \psi_{\theta}\right. \\
& \left.\int_{\mathbf{R}}\left(q_{\theta+v}+w-q_{\theta}\right) \psi_{\theta}(t)\right)
\end{aligned}
$$

From now we shall use the norms $\|(v, w, \alpha)\|_{2}=|v|+\|w\|_{2, \sigma}+|\alpha|$ on $B \times \bar{X} \times \mathbf{C}$ and $\|(g, \beta)\|_{-1}=\|g\|_{-1, \sigma}+|\beta|$ on $Y \times \mathbf{C}$. $J_{\mu}$ is of class $C^{1}$ and

$$
\begin{aligned}
& D J_{\mu}(v, w, \alpha)[z, W, a] \\
& \quad=\left(-z \ddot{q}_{\theta+v}-\ddot{W}+\cos \left(q_{\theta+v}+w\right)\left(z \dot{q}_{\theta+v}+W\right)\right. \\
& \left.\quad-\mu \cos \left(q_{\theta+\nu}+w\right)\left(z \dot{q}_{\theta+v}+W\right) f(\varphi)-a \psi_{\theta}, \int_{\mathbf{R}}\left(z \dot{q}_{\theta+v}+W\right) \psi_{\theta}(t)\right) .
\end{aligned}
$$

We shall prove that, provided $\|(\nu, w, \alpha)\|_{2} / \sigma$ and $\mu\|f\| / \sigma^{3}$ are small enough $D J_{\mu}(v, w, \alpha)$ is invertible. We first consider the case when $w=0$ and $\mu=0$. Let $T_{v}=D J_{0}(v, 0, \alpha)$ (independent of $\alpha$ ). Observing that $-\dddot{q}_{\theta+v}+\cos \left(q_{\theta+v}\right) \dot{q}_{\theta+v}=0$, we obtain

$$
T_{\nu}[z, W, a]=\left(-\ddot{W}+\cos q_{\theta+\nu} W-a \psi_{\theta}, \int_{\mathbf{R}}\left(z \dot{q}_{\theta+v}+W\right) \psi_{\theta}(t)\right)
$$

Using the first step we derive that $T_{v}$ is invertible and that, for a suitable positive constant $C$,

$$
\begin{equation*}
\left\|T_{v}^{-1}(g, \beta)\right\|_{2} \leqslant \frac{C}{\sigma}\|(g, \beta)\|_{-1} \tag{4.29}
\end{equation*}
$$

Now we estimate $\left\|\left(D J_{\mu}(\nu, w, \alpha)-T_{\nu}\right)[z, W, a]\right\|_{-1}$. We have

$$
\begin{aligned}
\left(D J_{\mu}(v, w, \alpha)-T_{v}\right)[z, W, a]= & \left(\left(\cos \left(q_{\theta+v}+w\right)-\left(\cos q_{\theta+v}\right)\right)\left(z \dot{q}_{\theta+v}+W\right)\right. \\
& \left.-\mu \cos \left(q_{\theta+v}+w\right)\left(z \dot{q}_{\theta+v}+W\right) f(\varphi), 0\right)
\end{aligned}
$$

We easily get

$$
\begin{aligned}
& \left\|\left(D J_{\mu}(v, w, \alpha)-T_{v}\right)[z, W, a]\right\|_{-1} \\
& \quad \leqslant C\|w\|_{2, \sigma}\left(\|W\|_{2, \sigma}+|z|\right)+\frac{\mu\|f\|}{\sigma^{2}}|z|+|\mu|\|f\|\|W\|_{2, \sigma} \\
& \quad \leqslant C\left(\|w\|_{2, \sigma}+\frac{\mu\|f\|}{\sigma^{2}}\right)\|(z, W, a)\|_{2}
\end{aligned}
$$

As a consequence, by (4.29), if $\mu\|f\| / \sigma^{3} \leqslant K_{0}$ and $\|w\|_{2, \sigma} / \sigma \leqslant K_{0}$, for $K_{0}$ small enough, then $D J_{\mu}(\mu, w, \alpha)$ is invertible and

$$
\left\|\left(D J_{\mu}(v, w, \alpha)\right)^{-1}\right\| \leqslant \frac{K_{1}}{\sigma}
$$

for a suitable positive constant $K_{1}$.
$3 r d$ step. We now prove the existence of a constant $K_{2}$ such that $(0,0,0)$ is the unique solution of the equation $J_{0}(v, w, \alpha)=0$ in $B\left(K_{2} \sigma\right)$, ball centered at the origin and of radius $K_{2} \sigma$ for the norm $\left\|\|_{2}\right.$. First we observe that, since $\ddot{q}_{\theta+\nu}=\sin \left(q_{\theta+v}\right)$, there holds

$$
\begin{aligned}
J_{0}(\nu, w, \alpha)=T_{\nu}[v, w, \alpha]+ & \left(\sin \left(q_{\theta+\nu}+w\right)-\sin \left(q_{\theta+v}\right)-\cos \left(q_{\theta+v}\right) w,\right. \\
& \left.\int_{\mathbf{R}}\left(q_{\theta+v}-q_{\theta}-v \dot{q}_{\theta+\nu}\right) \psi_{\theta}\right) .
\end{aligned}
$$

Moreover, by the analyticity of $q_{0}, \dot{q}_{0}, \psi_{0}$ over $S$,

$$
\int_{\mathbf{R}}\left(q_{\theta+v}(t)-q_{\theta}(t)-v \dot{q}_{\theta+v}(t)\right) \psi_{\theta}(t) \mathrm{d} t=\int_{\mathbf{R}}\left(q_{v}(t)-q_{0}(t)-v \dot{q}_{v}(t)\right) \psi_{0}(t) \mathrm{d} t
$$

Hence there is a constant $C^{\prime}$ such that

$$
\begin{aligned}
& \left\|\left(\sin \left(q_{\theta+\nu}+w\right)-\sin \left(q_{\theta+v}\right)-\cos \left(q_{\theta+v}\right) w, \int_{\mathbf{R}}\left(q_{\theta+v}-q_{\theta}-v \dot{q}_{\theta+v}\right) \psi_{\theta}\right)\right\|_{-1} \\
& \quad \leqslant C^{\prime}\left(\|w\|_{2, \sigma}^{2}+|v|^{2}\right)
\end{aligned}
$$

So, if $J_{0}(\nu, w, \alpha)=0$ then, by (4.29)
$\|(v, w, \alpha)\|_{2}$

$$
\begin{aligned}
& =\left\|-T_{v}^{-1}\left(\sin \left(q_{\theta+v}+w\right)-\sin \left(q_{\theta+v}\right)-\cos \left(q_{\theta+v}\right) w, \int_{\mathbf{R}}\left(q_{\theta+v}-q_{\theta}-v \dot{q}_{\theta+v}\right) \psi_{\theta}\right)\right\|_{2} \\
& \leqslant \frac{C C^{\prime}}{\sigma}\|(v, w, \alpha)\|_{2}^{2}
\end{aligned}
$$

Let $K_{2}<1 /\left(C C^{\prime}\right)$. By the latter inequality, if $J_{0}(\nu, w, \alpha)=0$ and $\|(\nu, w, \alpha)\|_{2} \leqslant K_{2} \sigma$, then $v=0, w=0, \alpha=0$.

4th step. By the previous steps we know that there exist positive constants $K_{0}, K_{1}$ and $K_{2}$ such that
(i) $\quad\left(J_{0}(\nu, w, \alpha)=0 \quad\right.$ and $\left.\quad\|(\nu, w, \alpha)\|_{2} \leqslant K_{2} \sigma\right) \Leftrightarrow \nu=w=\alpha=0$;
(ii) If $|\nu|<\sigma / 2,\|w\|_{2, \sigma} \leqslant K_{0} \sigma, \mu\|f\| \leqslant K_{0} \sigma^{3}$ then $D J_{\mu}(\nu, w, \alpha)$
is invertible and $\left\|\left(D J_{\mu}(\nu, w, \alpha)\right)^{-1}\right\| \leqslant K_{1} / \sigma$.
Moreover there exists a constant $K_{3}>0$ such that
(iii) $\quad\left\|\partial_{\mu} J_{\mu}(\nu, w, \alpha)\right\|_{-1}=\left\|\left(\sin \left(q_{\theta+\nu}+w\right) f(\varphi), 0\right)\right\|_{-1} \leqslant\|f\| K_{3} / \sigma$.

We say that (i), (ii), (iii) imply that there is $\eta$ such that, for all $0<\mu<\eta \sigma^{3} /\|f\|$, the equation $J_{\mu}(v, w, \alpha)=0$ has a unique solution such that $\|(v, w, \alpha)\|_{2}<K_{2} \sigma / 2$. In addition $\|(\nu, w, \alpha)\|_{2}=O\left(\mu\|f\| / \sigma^{2}\right)$. To prove existence, we can proceed as follows. Let $\mathcal{S}$ denote the set of all $\mu \in\left[0, K_{0} \sigma^{3} /\|f\|\right]$ such that there exists a $C^{1}$ function
$S_{\mu}:[0, \mu] \rightarrow\left\{(\nu, w, \alpha):\|(\nu, w, \alpha)\|_{2}<K_{2} \sigma / 2\right\}$ such that $S_{\mu}(0)=0, J_{t}\left(S_{\mu}(t)\right)=0$ for all $t \in[0, \mu]$. $\mathcal{S}$ is a bounded interval. Let us call $\bar{\mu}$ its supremum. By (ii) and the Implicit Function Theorem, $\bar{\mu}>0$. In addition, for $\mu \in \mathcal{S}$, there is a unique function $S_{\mu}$ with the required properties. As a consequence, for $0<\mu<\mu^{\prime}, S_{\mu}=S_{\mu^{\prime}[0, \mu]}$ and we can define a $C^{1}$ function $S:[0, \bar{\mu}) \rightarrow\left\{(\nu, w, \alpha):\|(\nu, w, \alpha)\|_{2}<K_{2} \sigma / 2\right\}$ such that $S(t)=S_{\mu}(t)$ for all $\mu \in(0, \bar{\mu})$. By (ii) and (iii), we can write, for all $t \in(0, \bar{\mu})$,

$$
\left\|S^{\prime}(t)\right\|_{2}=\left\|\left[D J_{t}(S(t))\right]^{-1} \cdot\left(\frac{\partial J_{t}}{\partial t}(S(t))\right)\right\|_{2} \leqslant \frac{K_{1} K_{3}\|f\|}{\sigma^{2}}
$$

Hence

$$
\begin{equation*}
\|S(t)\|_{2} \leqslant \frac{K_{1} K_{3}}{\sigma^{2}}\|f\||t| \tag{4.30}
\end{equation*}
$$

Now, since $S^{\prime}(t)$ is bounded, $S(t)$ converges to some $\bar{S}$ as $t \rightarrow \bar{\mu}$. Either $\bar{\mu}=K_{0} \sigma^{3} /\|f\|$ or $\|\bar{S}\|_{2}=K_{2} \sigma / 2$ (if not, by the Implicit Function Theorem, we could extend the solution $S$ to an interval $[0, \bar{\mu}+\xi), \xi>0$, contradicting the definition of $\bar{\mu})$. In the latter case, by (4.30),

$$
\|\bar{S}\|_{2}=\frac{K_{2} \sigma}{2} \leqslant \frac{K_{1} K_{3}}{\sigma^{2}} \bar{\mu}\|f\|
$$

So the existence assertion holds for $0<\mu<\eta \sigma^{3} /\|f\|$, where $\eta=\min \left(K_{0}\right.$, $\left.K_{2} /\left(2 K_{1} K_{3}\right)\right)$.

In order to prove uniqueness, we assume that there are $b_{1}, b_{2}$ such that $\left\|b_{i}\right\|_{2}<K_{2} \sigma / 2$, $J_{\mu}\left(b_{i}\right)=0$. Then, by the same argument as previously, we can prove the existence of two functions of class $C^{1} S_{1}, S_{2}:[0, \mu] \rightarrow\left\{b:\|b\|_{2}<K_{2} \sigma\right\}$ such that $S_{i}(\mu)=b_{i}$, $J_{t}\left(S_{i}(t)\right)=0$. Moreover, by (ii) and the Implicit Function Theorem, $S_{1}(\mu) \neq S_{2}(\mu)$ implies that $S_{1}(t) \neq S_{2}(t)$ for all $t \in[0, \mu]$, which contradicts (i), proving uniqueness.

The bound of $\left\|w_{A, \theta}^{\mu}\right\|_{2, \sigma}+\left|v_{A, \theta}^{\mu}\right|+\left|\alpha_{A, \theta}^{\mu}\right|$ given in the statement is a direct consequence of (4.30).

To complete the proof, we point out that $J_{\mu}$ is complex differentiable w.r.t all its variables. Therefore, as a consequence of the Implicit Function Theorem (see for example [2]), $Q_{A, \theta}^{\mu}=q_{\theta+\nu^{\mu}(A, \theta)}+w^{\mu}(A, \theta)$ depends analytically on $\theta$ and on $A_{i}$ if $r_{i}>0$.

We now consider the analytic extension of the function $\widetilde{F}_{\mu}$ defined for $(A, \theta) \in D \times S_{\sigma}$ by

$$
\begin{aligned}
\widetilde{F}_{\mu}(A, \theta) & =\Phi_{\mu, A}\left(Q_{A, \theta}^{\mu}\right) \\
& =\int_{\mathbf{R}} \frac{\left(\dot{Q}_{A, \theta}^{\mu}\right)^{2}(t)}{2}+\left(1-\cos Q_{A, \theta}^{\mu}(t)\right)+\mu\left(\cos Q_{A, \theta}^{\mu}(t)-1\right) f(\omega t+A) \mathrm{d} t
\end{aligned}
$$

Let us consider also the analytic extension for $(A, \theta) \in D \times S_{\sigma}$ of the Melnikov function

$$
M(A, \theta)=\int_{\mathbf{R}}\left(\cos q_{\theta}(t)-1\right) f(\omega t+A) \mathrm{d} t
$$

We have $M(A, \theta)=\Gamma(A+\omega \theta)$, where $\Gamma$ is defined in (2.36). We now prove

Lemma 4.7. - For $\mu\|f\| \sigma^{-3}$ small enough, we have

$$
\begin{equation*}
\widetilde{F}_{\mu}(A, \theta)=b+\mu M(A, \theta)+\mathrm{O}\left(\frac{\mu^{2}\|f\|^{2}}{\sigma^{4}}\right) \tag{4.31}
\end{equation*}
$$

uniformly for $(A, \theta) \in D \times S_{\sigma}$, where $b=\Phi_{0}\left(q_{0}\right)$.
Proof. - We have $Q_{A, \theta}^{\mu}=q_{\theta+\nu_{A, \theta}^{\prime}}^{\mu}+w_{A, \theta}^{\mu}$ which we denote for brevity $Q_{A, \theta}^{\mu}=q_{\theta+\nu}+w$.

$$
\begin{aligned}
& \widetilde{F}_{\mu}(A, \theta) \\
& =\int_{\mathbf{R}} \frac{\left(\dot{q}_{\theta+v}+\dot{w}\right)^{2}}{2}+\left(1-\cos \left(q_{\theta+v}+w\right)\right)+\mu\left(\cos \left(q_{\theta+v}+w\right)-1\right) f(\omega t+A) \mathrm{d} t \\
& =b+\int_{\mathbf{R}}\left[-\ddot{q}_{\theta+\nu} w+\frac{1}{2} \dot{w}^{2}+\left(\cos q_{\theta+v}-\cos \left(q_{\theta+v}+w\right)\right)\right. \\
& \left.\quad+\mu\left(\cos \left(q_{\theta+v}\right)-1\right) f(\omega t+A)+\mu\left(\cos \left(q_{\theta+v}+w\right)-\cos q_{\theta+v}\right) f(\omega t+A)\right] \mathrm{d} t \\
& =b+\mu M(\theta+v, A)+\int_{\mathbf{R}} \frac{1}{2} \dot{w}^{2}+\left(-\cos \left(q_{\theta+v}+w\right)+\cos q_{\theta+v}-\sin q_{\theta+v} w\right) \mathrm{d} t \\
& \quad+\mu \int_{\mathbf{R}}\left(\cos \left(q_{\theta+v}+w\right)-\cos q_{\theta+v}\right) f(\omega t+A) \mathrm{d} t .
\end{aligned}
$$

By the estimate $\|w\|_{2, \sigma} \leqslant C \mu\|f\| / \sigma^{2}$, it follows easily

$$
\widetilde{F}_{\mu}(A, \theta)=b+\mu M(A, \theta+v)+\mathrm{O}\left(\frac{\mu^{2}\|f\|^{2}}{\sigma^{4}}\right)
$$

For example we can get that $\int_{\mathbf{R}} \cos q_{\theta+v}-\cos \left(q_{\theta+v}+w\right)-\left(\sin q_{\theta+v}\right) w=\mathrm{O}\left(\mu^{2}\|f\|^{2} / \sigma^{4}\right)$ by writing $\cos q_{\theta+v}-\cos \left(q_{\theta+\nu}+w\right)-\left(\sin q_{\theta+v}\right) w=w^{2} \int_{0}^{1}(1-s) \cos \left(q_{\theta+v}+s w\right) \mathrm{d} s$ and using (4.22)-(4.23) together with $\|w\|_{2, \sigma} \leqslant \mu\|f\| / \sigma^{2}$. Moreover

$$
|M(A, \theta+v)-M(A, \theta)|=\mathrm{O}\left(\frac{|v|}{\sigma^{2}}\right)=\mathrm{O}\left(\frac{\mu\|f\|}{\sigma^{4}}\right)
$$

which completes the proof of the lemma.
The Fourier coefficients of the Melnikov function $\Gamma(A)=\sum_{k} \Gamma_{k} \exp (\mathrm{i} k A)$ are explicitely given by

$$
\begin{equation*}
\Gamma_{k}=f_{k} \frac{2 \pi(k \cdot \omega)}{\sinh \left((k \cdot \omega) \frac{\pi}{2}\right)} \quad(k \neq 0), \Gamma_{0}=4 f_{0} \tag{4.32}
\end{equation*}
$$

Since $\widetilde{F}_{\mu}(A, \theta)=\widetilde{G}_{\mu}(A+\omega \theta)$ and $M(A, \theta)=\Gamma(A+\omega \theta)$, we obtain, via a standard argument on Fourier coefficients of analytic functions, the following result

THEOREM 4.2. - There exists a positive constant $C_{8}$ such that, if $\mu\|f\| \sigma^{-3}$ is small enough, then $\forall k \in \mathbf{Z}^{n} \backslash\{0\}$, for all $\sigma \in\left(0, \frac{\pi}{2}\right)$, for all $\omega$,

$$
\begin{equation*}
\left|\widetilde{G}_{k}-\mu \Gamma_{k}\right| \leqslant \frac{C_{8} \mu^{2}\|f\|^{2}}{\sigma^{4}} \exp \left(-\sum_{i=1}^{n} r_{i}\left|k_{i}\right|\right) \exp \left(-|k \cdot \omega|\left(\frac{\pi}{2}-\sigma\right)\right) \tag{4.33}
\end{equation*}
$$

Proof. $\widetilde{F}_{\mu}(A, \theta)-b-\mu M(A, \theta)=\widetilde{G}_{\mu}(A+\omega \theta)-b-\mu \Gamma(A+\omega \theta)$ is analytic in $D \times S_{\sigma}$ and, by (4.31), its modulus is bounded by $C_{8} \mu^{2} /\left(\|f\|^{2} \sigma^{4}\right)$ where $C_{8}$ is a positive contant. By Lemma 3 in [21], (4.33) follows.

Remark 4.3. - Compare with Theorem 3.4.5 in [36] which gives the exponentially small estimate of the Fourier coefficients.

## 5. Three time scales

We consider in this section three-time-scales systems, introduced and investigated in connection with the problem of Arnold Diffusion in [15] (see also [25] and [38]), as

$$
\mathcal{H}=\frac{I_{1}}{\sqrt{\varepsilon}}+\varepsilon^{a} \beta \cdot I_{2}+\frac{p^{2}}{2}+(\cos q-1)+\mu(1-\cos q) f\left(\varphi_{1}, \varphi_{2}\right), \quad \varepsilon>0, a \geqslant 0
$$

with $n \geqslant 2, \varphi_{1} \in \mathbf{T}^{1}, \varphi_{2} \in \mathbf{T}^{n-1}, I_{1} \in \mathbf{R}^{1}, I_{2} \in \mathbf{R}^{n-1}, \beta \in \mathbf{R}^{n-1}$ and $\varepsilon$ is a positive small parameter. The frequency vector is $\omega=\left(1 / \sqrt{\varepsilon}, \varepsilon^{a} \beta\right)$, where $\beta=\left(\beta_{2}, \ldots, \beta_{n}\right) \in \mathbf{R}^{n-1}$ is fixed.

We assume through this section that $\mu\|f\| \varepsilon^{-3 / 2}$ and $\varepsilon$ are small.
Given $\kappa_{2}=\left(k_{2}, \ldots, k_{n}\right) \in \mathbf{Z}^{n-1}$, we shall use the notation $\kappa_{2}^{+}:=\left(\left|k_{2}\right|, \ldots,\left|k_{n}\right|\right)$. Moreover we shall use the abbreviation $\rho_{2}:=\left(r_{2}, \ldots, r_{n}\right)$, so that $\kappa_{2}^{+} \cdot \rho_{2}:=\sum_{i=2}^{n} r_{i}\left|k_{i}\right|$. We recall that $r_{1}, \ldots, r_{n}$ are defined in formula (4.1). Writing

$$
f\left(\varphi_{1}, \varphi_{2}\right)=\sum_{\left(k_{1}, \kappa_{2}\right) \in \mathbf{Z} \times \mathbf{Z}^{n-1}} f_{k_{1}, \kappa_{2}} \exp \left(\mathrm{i}\left(k_{1} \varphi_{1}+\kappa_{2} \cdot \varphi_{2}\right)\right),
$$

we assume that $f$ is analytic w.r.t $\varphi_{2}$. More precisely for $i \geqslant 2, r_{i}>0$, whereas $r_{1}$ may be zero in (4.1). If $a=0$, we impose in addition that $r_{i}>\left|\beta_{i}\right| \pi / 2$ for $i \geqslant 2$.

We shall use (4.33) in order to give an expansion for the "homoclinic function"

$$
\widetilde{G}_{\mu}(A)=\sum_{\left(k_{1}, \kappa_{2}\right) \in \mathbf{Z} \times \mathbf{Z}^{n-1}} \widetilde{G}_{k_{1}, \kappa_{2}} \exp \left(\mathrm{i}\left(k_{1} A_{1}+\kappa_{2} \cdot A_{2}\right)\right)=\sum_{k_{1} \in \mathbf{Z}} \widetilde{g}_{k_{1}}\left(A_{2}\right) \exp \left(\mathrm{i} k_{1} A_{1}\right)
$$

where

$$
\widetilde{g}_{k_{1}}\left(A_{2}\right)=\frac{1}{2 \pi} \int_{0}^{2 \pi} \widetilde{G}_{\mu}\left(A_{1}, A_{2}\right) \exp \left(-\mathrm{i} k_{1} A_{1}\right) \mathrm{d} A_{1}=\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \widetilde{G}_{k_{1}, \kappa_{2}} \exp \left(\mathrm{i} \kappa_{2} \cdot A_{2}\right)
$$

We start with

LEMMA 5.1. - There exists a positive constant $C_{9}$ such that, for $\mu\|f\| \varepsilon^{-3 / 2}$ small enough,

$$
\begin{equation*}
\left|\sum_{\left|k_{1}\right| \geqslant 2} \widetilde{g}_{k_{1}}\left(A_{2}\right) \exp \left(i k_{1} A_{1}\right)\right| \leqslant \sum_{\kappa_{2} \in \mathbf{Z}^{n-1},\left|k_{1}\right| \geqslant 2}\left|\widetilde{G}_{k_{1}, \kappa_{2}}\right| \leqslant C_{9} \frac{\mu\|f\|}{\sqrt{\varepsilon}} \exp \left(-\frac{\pi}{\sqrt{\varepsilon}}\right) \tag{5.1}
\end{equation*}
$$

Proof. - Choosing $\sigma=\sqrt{\varepsilon}$ and recalling that by (4.1) $\left|f_{k}\right| \leqslant\|f\| \mathrm{e}^{-\kappa_{2}^{+} \cdot \rho_{2}}$, we get from (4.32) and (4.33) that

$$
\begin{aligned}
\left|\widetilde{G}_{k_{1}, \kappa_{2}}\right| \leqslant & \mu\left|\Gamma_{k_{1}, \kappa_{2}}\right|+\left|\widetilde{G}_{k_{1}, \kappa_{2}}-\mu \Gamma_{k_{1}, \kappa_{2}}\right| \\
\leqslant & C \mu\|f\| \mathrm{e}^{-\kappa_{2}^{+} \cdot \rho_{2}}\left(\left|\frac{k_{1}}{\sqrt{\varepsilon}}+\kappa_{2} \cdot \beta \varepsilon^{a}\right|+1\right) \mathrm{e}^{-\left|\frac{k_{1}}{\sqrt{\varepsilon}}+\kappa_{2} \cdot \beta \varepsilon^{a}\right| \pi / 2} \\
& +C \frac{\mu^{2}}{\varepsilon^{2}}\|f\|^{2} \mathrm{e}^{-\kappa_{2}^{+} \cdot \rho_{2}} \mathrm{e}^{-\left|\frac{k_{1}}{\sqrt{\varepsilon}}+\kappa_{2} \cdot \beta \varepsilon^{a}\right|(\pi / 2-\sqrt{\varepsilon})}
\end{aligned}
$$

Since $\left|\frac{k_{1}}{\sqrt{\varepsilon}}+\kappa_{2} \cdot \beta \varepsilon^{a}\right| \geqslant\left|\frac{k_{1}}{\sqrt{\varepsilon}}\right|-\left|\kappa_{2} \cdot \beta \varepsilon^{a}\right|$ we get,

$$
\begin{aligned}
& \left|\widetilde{G}_{k_{1}, \kappa_{2}}\right| \\
& \leqslant \\
& \leqslant C \frac{\mu\|f\|}{\sqrt{\varepsilon}}\left(\left|k_{1}\right|+\left|\kappa_{2}\right|\right) \mathrm{e}^{-\kappa_{2}^{+} \cdot \rho_{2}+\left|\kappa_{2} \cdot \beta\right| \varepsilon^{a} \pi / 2} \mathrm{e}^{-\frac{\left|k_{1}\right|}{\sqrt{\varepsilon}} \pi / 2} \\
& \\
& \quad+C \frac{\mu^{2}}{\varepsilon^{2}}\|f\|^{2} \mathrm{e}^{-\kappa_{2}^{+} \cdot \rho_{2}+\left|\kappa_{2} \cdot \beta\right| \varepsilon^{a}(\pi / 2-\sqrt{\varepsilon})} \mathrm{e}^{-\left|k_{1}\right|\left(\frac{\pi}{2 \sqrt{\varepsilon}}-1\right)} \\
& \leqslant
\end{aligned} \begin{aligned}
& \frac{\mu\|f\|}{\sqrt{\varepsilon}}\left(\left|k_{1}\right|+\left|\kappa_{2}\right|\right) \exp \left(-\sum_{j=2}^{n}\left|k_{j}\right|\left(r_{j}-\left|\beta_{j}\right| \varepsilon^{a} \pi / 2\right)\right) \exp \left(-\left|k_{1}\right|\left(\frac{\pi}{2 \sqrt{\varepsilon}}-1\right)\right)
\end{aligned}
$$

We have used in the last line that $\mu\|f\| / \varepsilon^{3 / 2}=\mathrm{O}(1)$. Now $r_{j}-\left|\beta_{j}\right| \varepsilon^{a} \pi / 2>0$ for $\varepsilon$ small enough both if $a=0$ or if $a>0$. Summing in $\left|k_{1}\right|>2$ and in $\kappa_{2} \in \mathbf{Z}^{n-1}$ we obtain (5.1).

The Poincaré-Melnikov primitive defined in (2.36) can be written as

$$
\Gamma(\varepsilon, A)=\sum_{\left(k_{1}, \kappa_{2}\right) \in \mathbf{Z} \times \mathbf{Z}^{n-1}} \Gamma_{k_{1}, \kappa_{2}} \exp \left(\mathrm{i}\left(k_{1} A_{1}+\kappa_{2} \cdot A_{2}\right)\right)=\sum_{k_{1} \in \mathbf{Z}} \Gamma_{k_{1}}\left(\varepsilon, A_{2}\right) \exp \left(\mathrm{i} k_{1} A_{1}\right)
$$

Lemma 5.2. - Define

$$
R_{0}\left(\varepsilon, \mu, A_{2}\right)=\widetilde{g}_{0}\left(A_{2}\right)-b-\mu \Gamma_{0}\left(\varepsilon, A_{2}\right)
$$

and

$$
R_{1}\left(\varepsilon, \mu, A_{2}\right)=\widetilde{g}_{ \pm 1}\left(A_{2}\right)-\mu \Gamma_{ \pm 1}\left(\varepsilon, A_{2}\right)
$$

The following estimates hold

$$
R_{0}\left(\varepsilon, \mu, A_{2}\right)=\mathrm{O}\left(\mu^{2}\|f\|^{2}\right) \quad \text { and } \quad R_{1}\left(\varepsilon, \mu, A_{2}\right)=\mathrm{O}\left(\frac{\mu^{2}\|f\|^{2}}{\varepsilon^{2}} \exp \left(-\frac{\pi}{2 \sqrt{\varepsilon}}\right)\right)
$$

Proof. - We get immediately the estimate for $R_{0}\left(\varepsilon, \mu, A_{2}\right)$ by taking the mean value w.r.t. $A_{1} \in \mathbf{T}$ in formula (4.31) where $\theta=0, A_{2} \in \mathbf{T}^{n-1}$ and $\sigma=\pi / 4$.

We now prove the estimate for $R_{1}\left(\varepsilon, \mu, A_{2}\right)$. By (4.33) (where we choose $\sigma=\sqrt{\varepsilon}$ and $k_{1}= \pm 1$ ), we can obtain as in the proof of Lemma 5.1

$$
\begin{aligned}
& \left|\widetilde{g}_{ \pm 1}\left(A_{2}\right)-\mu \Gamma_{ \pm 1}\left(\varepsilon, A_{2}\right)\right| \\
& \quad \leqslant C \frac{\mu^{2}}{\varepsilon^{2}}\|f\|^{2} \sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \exp \left(-\sum_{j=2}^{n}\left|k_{j}\right|\left(r_{j}-\left|\beta_{j}\right| \varepsilon^{a} \pi / 2\right)\right) \exp \left(-\left(\frac{\pi}{2 \sqrt{\varepsilon}}-1\right)\right) \\
& \quad \leqslant C \frac{\mu^{2}}{\varepsilon^{2}}\|f\|^{2} \mathrm{e}^{-\frac{\pi}{2 \sqrt{\varepsilon}}},
\end{aligned}
$$

Since $\Gamma(A)$ and $G_{\mu}(A)$ are real valued functions we have that $\widetilde{g}_{-1}\left(A_{2}\right)=\overline{\widetilde{g}}_{1}\left(A_{2}\right)$ and $\Gamma_{-1}\left(A_{2}\right)=\bar{\Gamma}_{1}\left(A_{2}\right)$, where $\bar{z}$ denotes the complex conjugate of the complex number $z$. We deduce from the previous two lemmas the following result.

THEOREM 5.1. - For $\mu\|f\| \varepsilon^{-3 / 2}$ small there holds

$$
\begin{aligned}
\widetilde{G}_{\mu}\left(A_{1}, A_{2}\right)= & \widetilde{g}_{0}\left(A_{2}\right)+2 \operatorname{Re}\left[\widetilde{g}_{1}\left(A_{2}\right) \mathrm{e}^{\mathrm{i} A_{1}}\right]+\sum_{\left|k_{1}\right| \geqslant 2} \widetilde{g}_{k_{1}}\left(A_{2}\right) \exp \left(\mathrm{i} k_{1} A_{1}\right) \\
= & b+\left(\mu \Gamma_{0}\left(\varepsilon, A_{2}\right)+R_{0}\left(\varepsilon, \mu, A_{2}\right)\right) \\
& +2 \operatorname{Re}\left[\left(\mu \Gamma_{1}\left(\varepsilon, A_{2}\right)+R_{1}\left(\varepsilon, \mu, A_{2}\right)\right) \mathrm{e}^{i A_{1}}\right] \\
& +\mathrm{O}\left(\mu \varepsilon^{-1 / 2}\|f\| \exp \left(-\frac{\pi}{\sqrt{\varepsilon}}\right)\right)
\end{aligned}
$$

where

$$
R_{0}\left(\varepsilon, \mu, A_{2}\right)=\mathrm{O}\left(\mu^{2}\|f\|^{2}\right) \quad \text { and } \quad R_{1}\left(\varepsilon, \mu, A_{2}\right)=\mathrm{O}\left(\frac{\mu^{2}\|f\|^{2}}{\varepsilon^{2}} \exp \left(-\frac{\pi}{2 \sqrt{\varepsilon}}\right)\right)
$$

Remark 5.1. - (i) This improves the results in [38] which require $\mu=\varepsilon^{p}$ with $p>2+a$.
(ii) Theorem 5.1 certainly holds in any dimension, while the results of [25], which hold for more general systems, are proved for 2 rotators only.
(iii) In order to prove a splitting condition using Theorem 5.1 it is necessary, accordingly with [25] and [38], that $\exists m, l \in \mathbf{Z}^{n-1}$ such that $f_{0, l}, f_{1, m} \neq 0$. If not, recalling (4.32), there results that $\Gamma_{0}\left(\varepsilon, A_{2}\right)=\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \Gamma_{0, \kappa_{2}} \exp \mathbf{i} \kappa_{2} \cdot A_{2}=0$ and also $\Gamma_{1}\left(\varepsilon, A_{2}\right)=\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \Gamma_{1, \kappa_{2}} \exp \mathrm{i} \kappa_{2} \cdot A_{2}=0$.

Theorem 5.1 enables us to provide conditions implying the existence of diffusion orbits. For instance we obtain the following result.

Lemma 5.3. - Assume that there are $\bar{A}_{2} \in \mathbf{R}^{n-1}$ and $d_{0}, c_{0}>0$ such that, for all small $\varepsilon>0$,
(i) $\left|\Gamma_{1}\left(\varepsilon, A_{2}\right)\right|>\left(c_{0} / \sqrt{\varepsilon}\right) \mathrm{e}^{-\pi /(2 \sqrt{\varepsilon})}, \quad \forall A_{2} \in \mathbf{R}^{n-1}$ such that $\left|A_{2}-\bar{A}_{2}\right| \leqslant d_{0}$;
(ii) $\quad \Gamma_{0}\left(\varepsilon, A_{2}\right)>\Gamma_{0}\left(\varepsilon, \bar{A}_{2}\right)+c_{0}, \quad \forall A_{2} \in \mathbf{R}^{n-1}$ such that $\left|A_{2}-\bar{A}_{2}\right|=d_{0}$.

Then there is $v\left(c_{0}, d_{0}\right)>0$ and a constant $K>0$ such that, for all $\varepsilon, \mu$ with $\varepsilon, \mu\|f\| \varepsilon^{-3 / 2}<\nu\left(c_{0}, d_{0}\right)$, condition 2.1 is satisfied by $\widetilde{G}_{\mu}$ for some bounded open set $U \in \mathbf{R}^{n}$ and $\delta=\left(c_{0} \mu / 2 \sqrt{\varepsilon}\right) \mathrm{e}^{-\pi /(2 \sqrt{\varepsilon})}, \alpha=K c_{0} \mathrm{e}^{-\pi /(2 \sqrt{\varepsilon})} / 2 \sqrt{\varepsilon}, \rho_{U} \leqslant 3 d_{0}$.

Proof. - First we can derive from (4.32) and (4.33) in the same way as in the proof of Lemmas 5.1 and 5.2 that

$$
\begin{equation*}
\left|\widetilde{g}_{1}\left(A_{2}\right)\right|+\left|\nabla \widetilde{g}_{1}\left(A_{2}\right)\right| \leqslant \sum_{\kappa_{2} \in \mathbf{Z}^{n-1}}\left(1+\left|\kappa_{2}\right|\right)\left|\widetilde{G}_{1, \kappa_{2}}\right| \leqslant K_{0} \frac{\mu\|f\|}{\sqrt{\varepsilon}} \mathrm{e}^{-\pi /(2 \sqrt{\varepsilon})} \tag{5.2}
\end{equation*}
$$

for some constant $K_{0}$. By the bounds of $R_{0}$ and $R_{1}$ in Theorem 5.1, for $\varepsilon$ and $\mu\|f\| \varepsilon^{-3 / 2}$ small enough, we have
(a) $\quad\left|\widetilde{g}_{1}\left(A_{2}\right)\right|=\left|\left(\mu \Gamma_{1}+R_{1}\right)\left(A_{2}\right)\right|>\left(\mu c_{0} /(2 \sqrt{\varepsilon})\right) \mathrm{e}^{-\pi /(2 \sqrt{\varepsilon})} \quad \forall A_{2} \in B_{d_{0}}$,
(b) $\quad \tilde{g}_{0}\left(A_{2}\right)=\left(\mu \Gamma_{0}+R_{0}\right)\left(A_{2}\right)>\left(\mu \Gamma_{0}+R_{0}\right)\left(\bar{A}_{2}\right)+\left(\mu c_{0}\right) / 2 \quad \forall A_{2} \in \partial B_{d_{0}}$,
where $B_{d_{0}}$ is the open ball centered at $\bar{A}_{2}$ of radius $d_{0}$.
So we can write $\widetilde{g}_{1}\left(A_{2}\right)=\left|\widetilde{g}_{1}\left(A_{2}\right)\right| \mathrm{e}^{\mathrm{i} \phi\left(A_{2}\right)}$, where $\phi$ is a smooth real function defined in $B_{d_{0}}$ : 5.2 ) and the previous lower bound of $\left|\widetilde{g}_{1}\left(A_{2}\right)\right|$ provide a bound of $\nabla \phi\left(A_{2}\right)$. Precisely,

$$
\begin{equation*}
\left|\nabla \phi\left(A_{2}\right)\right| \leqslant \frac{\left|\nabla \widetilde{g}_{1}\left(A_{2}\right)\right|}{\left|\widetilde{g}_{1}\left(A_{2}\right)\right|} \leqslant 2\|f\| \frac{K_{0}}{c_{0}} . \tag{5.3}
\end{equation*}
$$

For $A_{2} \in B_{d_{0}}$, by Theorem 5.1 we have

$$
\begin{aligned}
\widetilde{G}_{\mu}\left(A_{1}, A_{2}\right)= & b+\left(\mu \Gamma_{0}+R_{0}\right)\left(\varepsilon, \mu, A_{2}\right)+2\left|\left(\mu \Gamma_{1}+R_{1}\right)\left(\varepsilon, \mu, A_{2}\right)\right| \cos \left(A_{1}+\phi\left(A_{2}\right)\right) \\
& +\mathrm{O}\left(\mu \varepsilon^{-1 / 2}\|f\| \exp \left(-\frac{\pi}{\sqrt{\varepsilon}}\right)\right) .
\end{aligned}
$$

Let

$$
U=\left\{\left(A_{1}, A_{2}\right) \in \mathbf{R} \times \mathbf{R}^{n-1}\left|A_{2} \in B_{d_{0}},\left|A_{1}+\phi\left(A_{2}\right)-\pi\right|<\frac{\pi}{2}\right\}\right.
$$

We now prove that $\widetilde{G}_{\mu}$ satisfies point (i) of condition 2.1 , with $\delta=c_{0} \mu \mathrm{e}^{-\pi / 2 \sqrt{\varepsilon}} / 2 \sqrt{\varepsilon}$ and $A_{0}=\left(\bar{A}_{1}, \bar{A}_{2}\right)=\left(\pi-\phi\left(\bar{A}_{2}\right), \bar{A}_{2}\right)$. Assume that $\left(A_{1}, A_{2}\right) \in \partial U$. Either $A_{2} \in \partial B_{d_{0}}$ and then

$$
\begin{equation*}
\widetilde{G}_{\mu}\left(A_{1}, A_{2}\right)-\widetilde{G}_{\mu}\left(\pi-\phi\left(\bar{A}_{2}\right), \bar{A}_{2}\right) \geqslant \frac{\mu c_{0}}{2}+\mathrm{O}\left(\mu \varepsilon^{-1 / 2}\|f\| \exp \left(-\frac{\pi}{2 \sqrt{\varepsilon}}\right)\right) \tag{5.4}
\end{equation*}
$$

or $\left|A_{1}+\phi\left(A_{2}\right)-\pi\right|=\pi / 2$ and then

$$
\begin{align*}
& \widetilde{G}_{\mu}\left(A_{1}, A_{2}\right)-\widetilde{G}_{\mu}\left(\pi-\phi\left(A_{2}\right), A_{2}\right) \\
& \quad=2\left|\left(\mu \Gamma_{1}+R_{1}\right)\left(A_{2}\right)\right|+\mathrm{O}\left(\mu \varepsilon^{-1 / 2}\|f\| \exp \left(-\frac{\pi}{\sqrt{\varepsilon}}\right)\right) \\
& \quad \geqslant \frac{\mu c_{0}}{\sqrt{\varepsilon}} \mathrm{e}^{-\pi /(2 \sqrt{\varepsilon})}+\mathrm{O}\left(\mu \varepsilon^{-1 / 2}\|f\| \exp \left(-\frac{\pi}{\sqrt{\varepsilon}}\right)\right) . \tag{5.5}
\end{align*}
$$

Hence, by (5.4) and (5.5), for $\varepsilon$ and $\mu\|f\| \varepsilon^{-3 / 2}$ small enough,

$$
\begin{equation*}
\inf _{\partial U} \widetilde{G}_{\mu} \geqslant \inf _{U} \widetilde{G}_{\mu}+\frac{c_{0} \mu \mathrm{e}^{-\pi / 2 \sqrt{\varepsilon}}}{2 \sqrt{\varepsilon}} \tag{5.6}
\end{equation*}
$$

Thus $\left(\widetilde{G}_{\mu}, U\right)$ satisfies point (i) of condition 2.1 , with $\delta=c_{0} \mu \mathrm{e}^{-\pi / 2 \sqrt{\varepsilon}} / 2 \sqrt{\varepsilon}$. Now $\left|\nabla \widetilde{G}_{\mu}\right|=O(\mu)$, hence there exists a positive constant $K_{1}$ such that, for all $A, B \in \mathbf{R}^{n}$, $\left|\widetilde{G}_{\mu}(A)-\widetilde{G}_{\mu}(B)\right| \leqslant K_{1} \mu|A-B|$. Let $\alpha=K \delta / \mu=K c_{0} \mathrm{e}^{-\pi / 2 \sqrt{\varepsilon}} / 2 \sqrt{\varepsilon}$, where $K=$ $1 / 8 K_{1}$. Then

$$
\begin{equation*}
\forall(A, B) \in \mathbf{R}^{n} \times \mathbf{R}^{n} \quad|A-B| \leqslant 2 \alpha \Rightarrow\left|\widetilde{G}_{\mu}(A)-\widetilde{G}_{\mu}(B)\right| \leqslant \frac{\delta}{4} \tag{5.7}
\end{equation*}
$$

Let $A^{0}=\left(A_{1}^{0}, A_{2}^{0}\right) \in U$ be such that $\inf _{U} \widetilde{G}_{\mu}=\widetilde{G}_{\mu}\left(A^{0}\right)$. Then by (5.7) and (5.6), $\sup _{B_{\alpha}\left(A^{0}\right)} \widetilde{G}_{\mu} \leqslant \inf _{U} \widetilde{G}_{\mu}+\delta / 4$, and the ball $B_{\alpha}\left(A^{0}\right)$ is included in $U$. So (ii) is satisfied with the above choice of $\alpha$.

That (iii) is satisfied is a straightforward consequence of (5.7).
To complete the proof, we observe that, by (5.3), $\left|\phi\left(A_{2}\right)-\phi\left(\bar{A}_{2}\right)\right| \leqslant 2 d_{0} K_{0}\|f\| / c_{0}$ for all $A_{2} \in B_{d_{0}}$. Hence any $A \in U$ can be written as

$$
A=\left(\pi-\phi\left(\bar{A}_{2}\right), \bar{A}_{2}\right)+\left(l_{1}, l_{2}\right), \quad\left|l_{1}\right| \leqslant \pi+\frac{2 d_{0} K_{0}}{c_{0}}\|f\|, \quad\left|l_{2}\right| \leqslant d_{0}
$$

Now, since $\left(l_{1}, l_{2}\right)=\sqrt{\varepsilon} l_{1} \omega+\left(0, l_{2}-\varepsilon^{a+1 / 2} l_{1} \beta\right)$,

$$
\left|\Pi_{\omega}\left(l_{1}, l_{2}\right)\right|=\left|\Pi_{\omega}\left(0, l_{2}-\varepsilon^{a+1 / 2} l_{1} \beta\right)\right| \leqslant\left|l_{2}-\varepsilon^{a+1 / 2} l_{1} \beta\right| \leqslant d_{0}+\mathrm{O}(\sqrt{\varepsilon})
$$

Hence $\rho_{U} \leqslant 2 d_{0}+\mathrm{O}(\sqrt{\varepsilon})$, and for $\varepsilon$ small enough, $\rho_{U} \leqslant 3 d_{0}$.
The condition given in the previous lemma is not easily handable. We now want to provide simpler conditions, involving properties of the perturbation $f$. For $A=$ $\left(A_{1}, A_{2}\right) \in \mathbf{T}^{1} \times \mathbf{T}^{n-1}$, let

$$
f\left(A_{1}, A_{2}\right)=\sum_{\left(k_{1}, \kappa_{2}\right) \in \mathbf{Z} \times \mathbf{Z}^{n-1}} f_{k_{1}, \kappa_{2}} \exp \left(\mathrm{i}\left(k_{1} A_{1}+\kappa_{2} \cdot A_{2}\right)\right)=\sum_{k_{1} \in \mathbf{Z}} f_{k_{1}}\left(A_{2}\right) \exp \left(\mathrm{i} k_{1} A_{1}\right)
$$

Since $f$ is analytic w.r.t $A_{2}, f_{k_{1}}\left(A_{2}\right)=(1 / 2 \pi) \int_{0}^{2 \pi} f\left(\sigma, A_{2}\right) \mathrm{e}^{-\mathrm{i} k_{1} \sigma} \mathrm{~d} \sigma$ depends analytically of $A_{2}$.

THEOREM 5.2. - Assume that $f$ satisfies one of the following conditions:
(i) $a>0, f_{0}\left(A_{2}\right)$ admits a strict local minimum at the point $\bar{A}_{2}$ and $f_{1}\left(\bar{A}_{2}\right) \neq 0$.
(ii) $a=0, f_{0}\left(A_{2}\right)$ admits a strict local minimum at the point $\bar{A}_{2}$ and $f_{1}\left(\bar{A}_{2}+\right.$ $\mathrm{i}(\pi / 2) \beta) \neq 0$.

Then, for all small $\varepsilon$ such that $\omega_{\varepsilon}=\left(1 / \sqrt{\varepsilon}, \beta \varepsilon^{a}\right)$ satisfies

$$
\left|\omega_{\varepsilon} \cdot k\right| \geqslant \frac{\gamma_{\varepsilon}}{|k|^{\tau}}, \quad \forall k \in \mathbf{Z}^{n}, k \neq 0
$$

for all $I_{0}, I_{0}^{\prime}$ with $\omega_{\varepsilon} \cdot I_{0}=\omega_{\varepsilon} \cdot I_{0}^{\prime}$, for all $\mu\|f\| \varepsilon^{-3 / 2}$ small enough, there is a heteroclinic orbit of the perturbed three-time-scales system, connecting the invariant tori $\mathcal{T}_{I_{0}}$ and $\mathcal{T}_{I_{0}^{\prime}}$. In addition, for all $\eta>0$ small enough the "diffusion time" $T_{d}$ needed to go from a $\eta$-neighbourhood of $\mathcal{T}_{I_{0}}$ to a $\eta$-neighbourhood of $\mathcal{T}_{I_{0}^{\prime}}$ is $\mathrm{O}\left(\mid I_{0}-\right.$ $I_{0}^{\prime}\left|(\sqrt{\varepsilon} / \mu) \mathrm{e}^{\pi /(2 \sqrt{\varepsilon})}\left[\left(\gamma_{\varepsilon}\right)^{-1}\left(\sqrt{\varepsilon} \mathrm{e}^{\pi /(2 \sqrt{\varepsilon})}\right)^{\tau}+|\log \mu|\right]+|\log (\eta)|\right)$.

Proof. - It is enough to prove that, if (i) or (ii) is satisfied, then the condition given in Lemma 5.3 holds for some $c_{0}, d_{0}>0$. The statement is then a direct consequence of Theorem 2.3 and Lemma 4.5.

We first assume that condition (i) is satisfied. In what follows, the notation $u=\mathrm{O}(v)$ means that $|u| \leqslant C|v|$, where $C$ is a universal constant. We have

$$
\begin{aligned}
\Gamma_{0}\left(\varepsilon, A_{2}\right) & =\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \frac{2 \pi \kappa_{2} \cdot \beta \varepsilon^{a}}{\sinh \left(\pi \kappa_{2} \cdot \beta \varepsilon^{a} / 2\right)} f_{0, \kappa_{2}} \mathrm{e}^{\mathrm{i} \kappa_{2} \cdot A_{2}} \\
& =\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}}\left(4+\mathrm{O}\left(\varepsilon^{2 a}\left|\kappa_{2}\right|^{2}\right)\right) f_{0, \kappa_{2}} \mathrm{e}^{\mathrm{i} \kappa_{2} \cdot A_{2}} \\
& =4 f_{0}\left(A_{2}\right)+\mathrm{O}\left(\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \varepsilon^{2 a}\left|\kappa_{2}\right|^{2} \mathrm{e}^{-\kappa_{2}^{+} \cdot \rho_{2}}\right) \\
& =4 f_{0}\left(A_{2}\right)+\mathrm{O}\left(\varepsilon^{2 a}\right)
\end{aligned}
$$

Moreover

$$
\begin{aligned}
\Gamma_{1}\left(\varepsilon, A_{2}\right) & =\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \frac{2 \pi\left(\varepsilon^{-1 / 2}+\kappa_{2} \cdot \beta \varepsilon^{a}\right)}{\sinh \left((\pi / 2)\left(\varepsilon^{-1 / 2}+\kappa_{2} \cdot \beta \varepsilon^{a}\right)\right)} f_{1, \kappa_{2}} \mathrm{e}^{\mathrm{i} \kappa_{2} \cdot A_{2}} \\
& =\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \frac{4 \pi}{\sqrt{\varepsilon}} \mathrm{e}^{-(\pi / 2)\left(\varepsilon^{-1 / 2}+\kappa_{2} \cdot \beta \varepsilon^{a}\right)}\left(1+\mathrm{O}\left(\left(\left|\kappa_{2}\right|+1\right) \varepsilon^{a+(1 / 2)}\right)\right) f_{1, \kappa_{2}} \mathrm{e}^{\mathrm{i} \kappa_{2} \cdot A_{2}}
\end{aligned}
$$

Since, by (4.1), $\left|f_{1, \kappa_{2}}\right|=O\left(\mathrm{e}^{-\kappa_{2}^{+} \cdot \rho_{2}}\right)$, we have

$$
\begin{aligned}
\Gamma_{1}\left(\varepsilon, A_{2}\right)= & \frac{4 \pi}{\sqrt{\varepsilon}} \mathrm{e}^{-(\pi / 2) \varepsilon^{-1 / 2}}\left[f_{1}\left(A_{2}\right)\right. \\
& \left.+\mathrm{O}\left(\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \mathrm{e}^{-\kappa_{2}^{+} \cdot \rho_{2}}\left(\left|\mathrm{e}^{-(\pi / 2) \kappa_{2} \cdot \beta \varepsilon^{a}}-1\right|+\left(\left|\kappa_{2}\right|+1\right) \varepsilon^{a+(1 / 2)} \mathrm{e}^{\left|\kappa_{2} \cdot \beta\right| \varepsilon^{a}}\right)\right)\right] \\
= & \frac{4 \pi}{\sqrt{\varepsilon}} \mathrm{e}^{-(\pi / 2) \varepsilon^{-1 / 2}}\left[f_{1}\left(A_{2}\right)\right. \\
& \left.+\mathrm{O}\left(\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}}\left(\left|\kappa_{2}\right|+1\right) \varepsilon^{a} \exp \left(-\sum_{j=2}^{n}\left|k_{j}\right|\left(r_{j}-(\pi / 2)\left|\beta_{j}\right| \varepsilon^{a}\right)\right)\right)\right] \\
= & \frac{4 \pi}{\sqrt{\varepsilon}} \mathrm{e}^{-(\pi / 2) \varepsilon^{-1 / 2}}\left[f_{1}\left(A_{2}\right)+\mathrm{O}\left(\varepsilon^{a}\right)\right]
\end{aligned}
$$

provided that $\varepsilon$ is small enough. Now, since $f_{1}\left(\bar{A}_{2}\right) \neq 0$ and $f_{0}$ admits a strict local minimum at $\bar{A}_{2}$, by the continuity of $f_{1}$, there is $d_{0}>0$ such that $\left|f_{1}\left(A_{2}\right)\right|>\left|f_{1}\left(\bar{A}_{2}\right)\right| / 2$
for all $\left|A_{2}-\bar{A}_{2}\right| \leqslant d_{0}$ and $\Delta\left(d_{0}\right):=\inf _{\left|A_{2}-\bar{A}_{2}\right|=d_{0}} f_{0}\left(A_{2}\right)-f_{0}\left(\bar{A}_{2}\right)>0$. Hence, for $\varepsilon$ small enough, the assumptions of lemma 5.3 hold with $c_{0}=\min \left(\pi\left|f_{1}\left(\bar{A}_{2}\right)\right|, 2 \Delta\left(d_{0}\right)\right)$.

We now assume that condition (ii) is satisfied. As previously, we have

$$
\begin{aligned}
\Gamma_{1}\left(\varepsilon, A_{2}\right)= & \sum_{\kappa_{2} \in \mathbf{Z}^{n-1}} \frac{4 \pi}{\sqrt{\varepsilon}} \mathrm{e}^{-(\pi / 2)\left(\varepsilon^{-1 / 2}+\kappa_{2} \cdot \beta\right)}\left(1+\mathrm{O}\left(\left(\left|\kappa_{2}\right|+1\right) \varepsilon^{1 / 2}\right)\right) f_{1, \kappa_{2}} \mathrm{e}^{\mathrm{i} \kappa_{2} \cdot A_{2}} \\
= & \frac{4 \pi}{\sqrt{\varepsilon}} \mathrm{e}^{-(\pi / 2) \varepsilon^{-1 / 2}}\left[f_{1}\left(A_{2}+\mathrm{i}(\pi / 2) \beta\right)\right. \\
& \left.+\mathrm{O}\left(\sum_{\kappa_{2} \in \mathbf{Z}^{n-1}}\left(\left|\kappa_{2}\right|+1\right) \varepsilon^{1 / 2} \exp \left(-\sum_{j=2}^{n}\left|k_{j}\right|\left(r_{j}-(\pi / 2)\left|\beta_{j}\right|\right)\right)\right)\right] \\
= & \frac{4 \pi}{\sqrt{\varepsilon}} \mathrm{e}^{-(\pi / 2) \varepsilon^{-1 / 2}}\left[f_{1}\left(A_{2}+\mathrm{i}(\pi / 2) \beta\right)+\mathrm{O}(\sqrt{\varepsilon})\right]
\end{aligned}
$$

We observe also that, if $a=0$, then $\Gamma_{0}\left(\varepsilon, A_{2}\right)=4 f_{0}\left(A_{2}\right)$ is independent of $\varepsilon$. It follows easily that condition (ii) as well implies that the assumption of lemma 5.3 holds true.

Remark 5.2. - In the time estimate obtained in Theorem 5.2, $\tau$ is the exponent of an exponetially large number, since, in order to apply the shadowing Theorem 2.3, which relies on condition 2.1, the diffusion orbit must approach to the homoclinic point at an exponentially small distance (namely $\alpha=\mathrm{O}(\exp -(\pi / 2 \sqrt{\varepsilon})$ ). This result can be improved via a shadowing theorem suited to the case when the splitting is exponentially small in one direction only. This allows to find diffusion orbits which approach to the homoclinic point just at a polynomially small distance in the orthogonal directions to $\omega$. This is proved in [7]. Moreover, for this reason, we are able to obtain diffusion orbits which drift in polynomial time w.r.t $1 / \varepsilon$ along the fast directions $\left(I_{2}, \ldots, I_{n}\right)$ (no diffusion in $I_{1}$ ), see [7].

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## Appendix A

In the proof of the following lemmas we will closely follow the arguments developed in the papers $[4,5]$ to which we refer for further details. In the sequel the notation $u=\mathrm{O}(v)$ (resp. $u=\mathrm{o}(v)$ ) will mean that there is a constant $C$ (resp. a function $\varepsilon(v)$ ) independent of anything except $f$ such that $|u| \leqslant C|v|$ (resp. $|u| \leqslant \varepsilon(v)|v|$ and $\left.\lim _{v \rightarrow 0} \varepsilon(v)=0\right)$.

Proof of Lemma 2.1. - We first assume that $\theta=0$ and give the existence proof in $[0,+\infty)$. We are looking for a solution of (1.5) in the form of $q=q_{0}+w$ with $w(0)=0$
and $\lim _{t \rightarrow+\infty} w(t)=0$. The function $w$ must satisfy the equation

$$
-\ddot{w}+w=-\left(\sin \left(q_{0}+w\right)-\sin q_{0}-w\right)+\mu \sin \left(q_{0}+w\right) f(\omega t+A)
$$

Let

$$
\mathbf{X}=\left\{w(\cdot) \in W^{1, \infty}([0,+\infty)) \mid\|w\|_{1}:=\sup _{t \in \mathbf{R}} \max (|w(t)|,|\dot{w}(t)|) \exp \left(\frac{|t|}{2}\right)<+\infty\right\}
$$

and

$$
\mathbf{X}^{\prime}=\left\{w(\cdot) \in L^{\infty}([0,+\infty))\left|\|w\|_{0}:=\sup _{t \in \mathbf{R}}\right| w(t) \left\lvert\, \exp \left(\frac{|t|}{2}\right)<+\infty\right.\right\}
$$

$\mathbf{X}$ and $\mathbf{X}^{\prime}$, endowed respectively with norms $\left\|\|_{1}\right.$ and $\| \|_{0}$, are Banach spaces. Let $\mathcal{L}_{0}$ be the linear operator which assigns to $h \in \mathbf{X}^{\prime}$ the unique solution $u=\mathcal{L}_{0} h$ of the problem:

$$
\left\{\begin{array}{l}
-\ddot{u}+u=h \\
u(0)=0, \lim _{t \rightarrow+\infty} u(t)=0
\end{array}\right.
$$

An explicit computation shows that, for $t \in[0,+\infty)$,

$$
\begin{equation*}
u(t)=\left(\mathcal{L}_{0} h\right)(t)=\frac{1}{2} \int_{0}^{+\infty}\left(\mathrm{e}^{-|t-s|}-\mathrm{e}^{-(t+s)}\right) h(s) \mathrm{d} s \tag{A.1}
\end{equation*}
$$

As an easy consequence $\mathcal{L}_{0}$ sends $\mathbf{X}^{\prime}$ into $\mathbf{X}$ continuously.
We define the non-linear operator $H: \mathbf{R} \times \mathbf{R}^{n} \times \mathbf{X} \rightarrow \mathbf{X}$ by

$$
\begin{equation*}
H(\mu, A, w):=w-\mathcal{L}_{0}\left(-\left(\sin \left(q_{0}+w\right)-\sin q_{0}-w\right)+\mu \sin \left(q_{0}+w\right) f(\omega t+A)\right) \tag{A.2}
\end{equation*}
$$

$H$ is smooth, $2 \pi \mathbf{Z}^{n}$-periodic w.r.t. $A$ and we have $H(0, A, 0)=0$. The unknown $w$ must solve the equation $H(\mu, A, w)=0$. We can apply the Implicit Function Theorem. In fact, let us check that

$$
\partial_{w} H(0, A, 0): W \rightarrow W-\mathcal{L}_{0}\left[\left(1-\cos q_{0}\right) W\right]
$$

is invertible. Since $\lim _{t \rightarrow \infty}\left(1-\cos q_{0}(t)\right)=0, \partial_{w} H(0, A, 0)$ is of the type "Identity + Compact" and therefore it is sufficient to show that it is injective. $W$ is in the kernel of $\partial_{w} H(0, A, 0)$ iff $W(0)=0$ and W satisfies in $(0,+\infty)$ the equation

$$
\begin{equation*}
-\ddot{W}+\cos q_{0} W=0 \tag{A.3}
\end{equation*}
$$

Multiplying by $\dot{q}_{0}$ in (A.3) and integrating over $[0,+\infty)$ by parts twice we obtain that $\dot{W}(0) \dot{q}_{0}(0)=0$. Since $\dot{q}_{0}(0) \neq 0$ we get also $\dot{W}(0)=0$ and as a consequence $W=0$. Thus the kernel of $\partial_{w} H(0, A, 0)$ is reduced to 0 , and this operator is invertible. We derive by the Implicit Function Theorem that there are $\rho_{0}>0$ and $\mu_{0}>0$ such that, for all $|\mu|<\mu_{0}$, for all $A \in \mathbf{R}^{n}$, the equation $H(\mu, A, w)=0$ has a unique solution $w_{A}^{\mu}$ in $\mathbf{X}$ such that $\left\|w_{A}^{\mu}\right\|<\rho_{0}$.

Note that $\mu_{0}$ and $\rho_{0}$ may be chosen independent of $A$ (and of $\omega$ too) because $\partial_{w} H(0, A, 0)$ is independent of $A$ and $\omega, \partial_{\mu} H(0, A, 0)$ is uniformly bounded, and $\partial_{w} H(\mu, A, w)\left(\right.$ resp. $\left.\partial_{\mu} H(\mu, A, w)\right)$ tend to $\partial_{w} H(0, A, 0)$ (resp. $\left.\partial_{\mu} H(0, A, 0)\right)$ as $(\mu, w) \rightarrow(0,0)$ uniformly in $(A, \omega)$.

Since $H$ is smooth $w_{A}^{\mu}$ depends smoothly on $\mu$ and $A$ and $w_{A+2 \pi k}^{\mu}=w_{A}^{\mu}$ by the $2 \pi \mathbf{Z}^{n}-$ periodicity of $H$ w.r.t. $A$. By the properties of $\partial_{\mu} H$ mentioned above, $\left\|w_{A}^{\mu}\right\|_{1}=\mathrm{O}(\mu)$.

In a similar way we can prove the existence and uniqueness of ${w_{A}^{\prime \mu}}_{A}:(-\infty, 0] \rightarrow \mathbf{R}$ which satisfies analogous properties over the interval $(-\infty, 0]$. We can define $q_{A, 0}^{\mu}$ by $q_{A, 0}^{\mu}(t)=q_{0}(t)+w_{A}^{\mu}(t)$ if $t \geqslant 0, q_{A, 0}^{\mu}=q_{0}(t)+w_{A}^{\prime \mu}(t)$ if $t<0$. This is the unique function for which (i), (ii) (with $\theta=0$ ) hold.

If $\theta \neq 0$, we observe that $q$ satisfies (i) iff

$$
\left\{\begin{array}{l}
-\left(T_{-\theta} q\right)^{\prime \prime}+\sin \left(T_{-\theta} q\right)=\mu \sin \left(T_{-\theta} q\right) f(\omega t+A+\omega \theta) \\
\left(T_{-\theta} q\right)(0)=\pi
\end{array}\right.
$$

where $T_{-\theta} q(t)=q(t+\theta)$. Hence there is a unique $q_{A, \theta}^{\mu}$ which satisfies (i), (ii), defined by $q_{A, \theta}^{\mu}=T_{\theta} q_{A+\omega \theta, 0}^{\mu}$, i.e. $q_{A, \theta}^{\mu}(t)=q_{A+\omega \theta, 0}^{\mu}(t-\theta)$; (iii) and (iv) clearly hold. The regularity of $q_{A, \theta}^{\mu}$ w.r.t. $A, \mu$ is a consequence of the regularity of $w_{A}^{\mu}$ and $w_{A}^{\prime \mu}$ w.r.t. $A$ and $\mu$. (v) follows from

$$
\partial_{A} w_{A}^{\mu}=-\left[\partial_{w} H\left(\mu, A, w_{A}^{\mu}\right)\right]^{-1}\left(\partial_{A} H\right)\left(\mu, A, w_{\mu}^{A}\right)
$$

provided we can justify that $\left\|\left(\partial_{A} H\right)\left(\mu, A, w_{A}^{\mu}\right)\right\|_{1}=\mathrm{O}(\mu), \quad\left\|\omega \cdot \partial_{A} H\left(\mu, A, w_{A}^{\mu}\right)\right\|_{1}=$ $\mathrm{O}(\mu)$. The second bound (uniform in $\omega$ ) is less obvious. We just point out that

$$
\omega \cdot \partial_{A} H\left(\mu, A, w_{A}^{\mu}\right)=-\mathcal{L}_{0}\left(\mu \sin \left(q_{0}+w_{A}^{\mu}\right) \frac{\mathrm{d}}{\mathrm{~d} t} f(\omega t+A)\right)
$$

and that we can use the "regularizing" properties of $\mathcal{L}_{0}$.
Proof of Lemma 2.4. - We give the proof in the interval $\left[\theta_{1}, \theta_{2}\right]$. We may assume without loss of generality that $\theta_{1}=0$ since, by the remark at the end of the proof of Lemma 2.1, a translation of the time by $-\theta_{1}$ amounts to adding $\omega \theta_{1}$ to $A$. For simplicity of notations, we shall write $\theta_{2}=\theta$.

We are looking for a solution $q=q_{0, \theta}^{*}+w$ of (1.5) over $[0, \theta]$ with $w(0)=w(\theta)=0$, where $q_{0, \theta}^{*}$ is the following smooth "approximate solution"

$$
q_{0, \theta}^{*}(t)= \begin{cases}q_{A, 0}^{\mu}(t) & \text { if } t \in(0, \theta / 2-1) \\ r_{\theta}^{*}(t) & \text { if } t \in[\theta / 2-1, \theta / 2+1] \\ 2 \pi+q_{A, \theta}^{\mu}(t) & \text { if } t \in(\theta / 2+1, \theta)\end{cases}
$$

where

$$
r_{\theta}^{*}(t)=(1-R(t-\theta / 2)) q_{A, 0}^{\mu}(t)+R(t-\theta / 2)\left(q_{A, \theta}^{\mu}(t)+2 \pi\right)
$$

and $R: \mathbf{R} \rightarrow[0,1]$ is a $C^{\infty}$ function such that $R(s)=0$ if $s \leqslant-1, R(s)=1$ if $s \geqslant 1$. Let $\mathcal{L}_{0, \theta}$ be the linear operator which assigns to $h \in L^{\infty}([0, \theta])$ the unique solution $u=\mathcal{L}_{0, \theta} h$
of the problem:

$$
\left\{\begin{array}{l}
-\ddot{u}+u=h  \tag{A.4}\\
u(0)=0, u(\theta)=0 .
\end{array}\right.
$$

An explicit computation shows that for $t \in[0, \theta]$ the solution $u$ of (A.4) is given by

$$
u(t)=\frac{1}{\sinh (\theta)}\left[\int_{0}^{t} h(s) \sinh (s) \sinh (\theta-t) \mathrm{d} s+\int_{t}^{\theta} h(s) \sinh (\theta-s) \sinh (t) \mathrm{d} s\right]
$$

Note that $\mathcal{L}_{0, \theta}$ sends $L^{\infty}([0, \theta])$ into $W^{1, \infty}([0, \theta])\left(W^{2, \infty}([0, \theta])\right.$ in fact $)$ and that there is a constant $K$ independent of $\theta$ such that $\left\|\mathcal{L}_{0, \theta} W\right\|_{1, \infty} \leqslant K\|W\|_{\infty}$, where $\left\|\|_{\infty}\right.$ denotes the infty norm in $[0, \theta]$ and $\|W\|_{1, \infty}:=\|W\|_{\infty}+\|\dot{W}\|_{\infty}$.

We define the smooth non-linear operator $H^{\theta}: \mathbf{R} \times \mathbf{R}^{n} \times W^{1, \infty}([0, \theta]) \rightarrow W^{1, \infty}([0, \theta])$ by

$$
H^{\theta}(\mu, A, w):=w-\mathcal{L}_{0, \theta}\left(-\left(\sin \left(q_{0, \theta}^{*}+w\right)-\ddot{q}_{0, \theta}^{*}-w\right)+\mu \sin \left(q_{0, \theta}^{*}+w\right) f(\omega t+A)\right)
$$

We immediately remark for further purpose that

$$
\begin{equation*}
\left\|\partial_{w}^{2} H^{\theta}(\mu, A, w)[W, W]\right\|=\mathrm{O}\left(\|W\|_{\infty}^{2}\right) \tag{A.5}
\end{equation*}
$$

Moreover, by Lemma 2.1(i) and the definition of $q_{0, \theta}^{*}, \|-\sin q_{0, \theta}^{*}+\ddot{q}_{0, \theta}^{*}+\mu \sin \left(q_{0, \theta}^{*}\right) \times$ $f(\omega t+A) \|_{\infty}=\mathrm{O}(\exp (-\theta / 2))$ hence

$$
\begin{equation*}
\left\|H^{\theta}(\mu, A, 0)\right\|_{1, \infty}=\mathrm{O}(\exp (-\theta / 2)) \tag{A.6}
\end{equation*}
$$

$q_{0, \theta}^{*}+w$ is a solution of (1.5) with the appropriate boundary conditions iff $H^{\theta}(\mu, A, w)$ $=0$. We shall show that there exist $\bar{C}, \bar{L}, \bar{\mu}>0$ such that $\forall \theta>\bar{L}$, for all $|\mu|<\bar{\mu}$, for all $A$ and $\omega, \partial_{w} H^{\theta}(\mu, A, 0)$ is invertible and

$$
\begin{equation*}
\left\|\left(\partial_{w} H^{\theta}(\mu, A, 0)\right)^{-1}\right\| \leqslant \bar{C} \tag{A.7}
\end{equation*}
$$

Since $\partial_{w} H^{\theta}(\mu, A, 0)$ is of the type "Id + Compact", it is enough to prove that

$$
\forall W \in W^{1, \infty}([0, \theta]) \quad\left\|\partial_{w} H^{\theta}(\mu, A, 0) W\right\|_{1, \infty} \geqslant \frac{1}{\bar{C}}\|W\|_{1, \infty}
$$

We shall just sketch the proof of this assertion (see also Lemma 2 of [4]). Arguing by contradiction, we assume that there are sequences $\left(\mu_{n}\right) \rightarrow 0,\left(\theta_{n}\right) \rightarrow \infty,\left(A_{n}\right),\left(\omega_{n}\right)$, ( $W_{n}$ ) such that $W_{n} \in W^{1, \infty}\left(\left[0, \theta_{n}\right]\right),\left\|W_{n}\right\|_{1, \infty}=1$,

$$
\begin{equation*}
\left\|\partial_{w} H^{\theta_{n}}\left(\mu_{n}, A_{n}, 0\right) W_{n}\right\|_{1, \infty} \rightarrow 0 \tag{A.8}
\end{equation*}
$$

Let $\xi_{n} \in\left[0, \theta_{n}\right]$ be such that $m_{n}:=\max _{t \in\left[0, \theta_{n}\right]}\left|W_{n}(t)\right|=W_{n}\left(\xi_{n}\right)$. By (A.8) and the properties of $\mathcal{L}_{0, \theta},\left\|W_{n}\right\|_{1, \infty}=\mathrm{O}\left(m_{n}\right)$. Hence $\liminf \left(m_{n}\right)>0$. Taking a subsequence,
we may assume that $\left(\xi_{n}\right)$ is bounded or $\left(\theta_{n}-\xi_{n}\right)$ is bounded or $\left(\left(\xi_{n}\right) \rightarrow \infty\right.$ and $\left.\left(\theta_{n}-\xi_{n}\right) \rightarrow \infty\right)$.

In the first case, still up to a subsequence $W_{n} \rightarrow W \neq 0$ uniformly in compact subsets of $[0, \infty)$. Taking limits in (A.8) we obtain that $W(0)=0,-\ddot{W}+\cos q_{0} W=0$, which contradicts $W \neq 0$. The second case can be dealt with similarly. In the third case, up to a subsequence, $W_{n}\left(\cdot+\xi_{n}\right) \rightarrow W \neq 0$ uniformly in compact subsets of $\mathbf{R}$, with $|W(t)| \leqslant|W(0)|$ for all $t \in \mathbf{R}$. Taking limits in (A.8), we obtain that $-\ddot{W}+W=0$ over $\mathbf{R}$, which contradicts $W \neq 0$ bounded.

From now we shall assume that $|\mu|<\mu_{0} \leqslant \bar{\mu}, \theta>\bar{L}$. Let

$$
R^{\theta}(\mu, A, w)=H^{\theta}(\mu, A, w)-H^{\theta}(\mu, A, 0)-\partial_{w} H^{\theta}(\mu, A, 0) w
$$

By the previous assertion,

$$
\begin{aligned}
H^{\theta}(\mu, A, w)=0 \Leftrightarrow w= & -\left(\partial_{w} H^{\theta}(\mu, A, 0)\right)^{-1} H^{\theta}(\mu, A, 0) \\
& -\left(\partial_{w} H^{\theta}(\mu, A, 0)\right)^{-1} R^{\theta}(\mu, A, w):=F_{\mu, A}^{\theta}(w)
\end{aligned}
$$

We just have to show that $F_{\mu, A}^{\theta}$ is a contraction in some ball $B(0, \rho) \subset W^{1, \infty}([0, \theta])$. For this, we derive from (A.5) and (A.6) in a standard way that, for all $\|w\|_{1, \infty},\left\|w^{\prime}\right\|_{1, \infty} \leqslant \rho$, $|\mu|<\bar{\mu}, \theta>\bar{L}$ there holds

$$
\begin{equation*}
\left\|F_{\mu, A}^{\theta}(w)\right\|_{1, \infty}=\mathrm{O}\left(\exp (-\theta / 2)+\rho^{2}\right) ; \quad\left\|F_{\mu, A}^{\theta}(w)-F_{\mu, A}^{\theta}\left(w^{\prime}\right)\right\|=\mathrm{O}\left(\rho\left\|w^{\prime}-w\right\|\right) \tag{A.9}
\end{equation*}
$$

We can deduce that $F_{\mu, A}^{\theta}$ is a contraction in $\bar{B}(0, \rho)$, with $\rho=C \exp (-\theta / 2)$, for some constant $C$, provided that $\theta>\bar{L}, \bar{L}$ large enough. Applying the Contraction Mapping Theorem we conclude that there is a unique solution $\left\|w_{\mu}^{L}(A, \theta)\right\|_{1, \infty} \leqslant C \exp (-\theta / 2)$ of the equation $H_{\mu, A}^{\theta}(w)=0$. Note that by (A.9) uniqueness holds in $B\left(0, \rho_{0}\right)$ for some $\rho_{0}>0$ independent of $\theta$. The regularity of the solutions in $(A, \theta, \mu)$ follows as in [4].

Proof of Remark 2.3. - Assume with no loss of generality that $A_{0}=0$ and that $D^{2} G_{\mu}(0)=\operatorname{diag}\left\{\lambda_{1}, \ldots, \lambda_{n}\right\}$ (this can be always obtained by an orthogonal change of variables). Assume also that $\lambda_{1}=\min \lambda_{i}$ and that $\lambda_{n}=\max \lambda_{i}$. Let $Q(A)=\frac{1}{2}\left(\lambda_{1} A_{1}^{2}+\right.$ $\cdots+\lambda_{n} A_{n}^{2}$. By the Taylor expansion

$$
\begin{equation*}
G_{\mu}(A)=Q(A)+\frac{1}{6} D^{3} G_{\mu}(\bar{A}) A^{3} \tag{A.10}
\end{equation*}
$$

for a suitable $\bar{A}$ with $|\bar{A}|<|A|$. For $A \in B_{\rho}(0)$, with $\rho=\lambda_{1} /\left(3 \sup _{A \in \mathbf{T}^{n}}\left|D^{3} G_{\mu}\right|\right)$, there holds $\left|(1 / 6) D^{3} G_{\mu}(\bar{A}) A^{3}\right| \leqslant(1 / 9) Q(A)$, hence, by (A.10)

$$
\begin{equation*}
\frac{8}{9} Q(A) \leqslant G_{\mu}(A) \leqslant \frac{10}{9} Q(A) \tag{A.11}
\end{equation*}
$$

If we choose $U=B_{\rho}(0)$ condition (i) is then satisfied with $\delta=4 \lambda_{1} \rho^{2} / 9$ (we have $\inf _{B_{\rho}(0)} G_{\mu}=0$ ). Moreover, by (A.11) we have that

$$
\left\{G_{\mu} \leqslant \frac{\delta}{2}\right\} \subset\left\{Q(A) \leqslant \frac{9 \delta}{16}\right\} \quad \text { and } \quad\left\{G_{\mu} \geqslant \frac{3 \delta}{4}\right\} \subset\left\{Q(A) \geqslant \frac{27 \delta}{40}\right\}
$$

One can easily compute that

$$
\operatorname{dist}\left(\left\{Q(A) \leqslant \frac{9 \delta}{16}\right\},\left\{Q(A) \geqslant \frac{27 \delta}{40}\right\}\right)=((\sqrt{6}-\sqrt{5}) / \sqrt{10}) \rho \sqrt{\lambda_{1} / \lambda_{n}}=2 \alpha
$$

hence $\operatorname{dist}\left(\left\{G_{\mu} \leqslant \frac{\delta}{2}\right\},\left\{G_{\mu} \geqslant \frac{3 \delta}{4}\right\}\right) \geqslant 2 \alpha$.
Condition (ii) holds as well since, by (A.11), $\left\{G_{\mu}(A) \leqslant \delta / 4\right\}$ contains $\{Q(A) \leqslant \delta / 5\}$, and $B_{\alpha}(0) \subset\{Q(A) \leqslant \delta / 5\}$.

Proof of Lemma 4.3. - Let us consider the function $H: \mathbf{R} \times \mathbf{T}^{n} \times \mathbf{R} \times \mathbf{R} \rightarrow \mathbf{R}$ defined by

$$
H(\mu, A, \theta, l)=Q_{A, \theta}^{\mu}(\theta+l)-\pi
$$

The unknown $l_{\mu}(A, \theta)$ can be implicitely defined by the equation $H(\mu, A, \theta, l)=0$. We have $H(0, A, \theta, 0)=0$ and

$$
\partial_{l} H(0, A, \theta, 0)=\dot{q}_{\theta}(\theta)=\dot{q}_{0}(0) \neq 0
$$

Hence by the Implicit function theorem, for $\mu$ small enough (independently of $A, \theta, \omega$ because $\partial_{l} H$ and $\partial_{\mu} H$ are continuous uniformly in $A, \theta$ with a modulus of continuity independent of $\omega$ ), there exists a unique smooth solution $l_{\mu}(A, \theta)=\mathrm{O}(\mu)$ of $H(\mu, A, \theta, l)=0$. Moreover, by the uniform estimates in $A$ and $\omega$ that we can obtain for $\partial_{A} Q_{A, \theta}^{\mu}, \omega \cdot \partial_{A} Q_{A, \theta}^{\mu}$ and by (4.8), $\left|\partial_{A} H\right|+\left|\partial_{\theta} H\right|=\mathrm{O}(\mu)$. Hence there holds $\left|\nabla l_{\mu}(A, \theta)\right|=\mathrm{O}(\mu)$.

Proof of Lemma 4.4. - The first step is to prove that

$$
\begin{align*}
& \max \left(\left|q_{A, \theta+l_{\mu}(A, \theta)}^{\mu}(t)-Q_{A, \theta}^{\mu}(t)\right|,\left|\dot{q}_{A, \theta+l_{\mu}(A, \theta)}^{\mu}(t)-\dot{Q}_{A, \theta}^{\mu}(t)\right|\right) \\
& \quad \leqslant K_{0}\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right| \exp \left(-\frac{|t-\theta|}{2}\right), \quad \forall t \in \mathbf{R} \tag{A.12}
\end{align*}
$$

By Lemma 2.1(iv) we have $q_{A, \theta+l_{\mu}(A, \theta)}^{\mu}=T_{\theta+l_{\mu}(A, \theta)} q_{A^{\prime}, 0}^{\mu} ; Q_{A, \theta}^{\mu}=T_{\theta+l_{\mu}(A, \theta)} Q_{A^{\prime}, \theta^{\prime}}^{\mu}$, where $A^{\prime}=A+\omega\left(\theta+l_{\mu}(A, \theta)\right)$ and $\theta^{\prime}=-l_{\mu}(A, \theta)$. So it is enough to prove the estimate for $w:=Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{A^{\prime}, 0}^{\mu}$.

Note that $Q_{A^{\prime}, \theta^{\prime}}^{\mu}(0)=q_{A^{\prime}, 0}^{\mu}(0)=\pi$. So $Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{0}, q_{A^{\prime}, 0}^{\mu}-q_{0}$ belong to $\mathbf{X}$ and satisfy

$$
H\left(\mu, A^{\prime}, q_{A^{\prime}, 0}^{\mu}-q_{0}\right)=0 ; \quad H\left(\mu, A^{\prime}, Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{0}\right)=\alpha_{A^{\prime}, \theta^{\prime}}^{\mu} \mathcal{L}_{0}\left(\psi_{\theta^{\prime}}\right)
$$

where $\mathbf{X}, H$ and $\mathcal{L}_{0}$ are defined in the proof of Lemma 2.1. Therefore

$$
\alpha_{A^{\prime}, \theta^{\prime}}^{\mu} \mathcal{L}_{0}\left(\psi_{\theta^{\prime}}\right)=\partial_{w} H\left(\mu, A^{\prime}, q_{A^{\prime}, 0}^{\mu}-q_{0}\right)\left(Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{A^{\prime}, 0}^{\mu}\right)+\mathrm{o}\left(\left\|Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{A^{\prime}, 0}^{\mu}\right\|_{1}\right)
$$

Moreover $\left\|Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{0}\right\|_{1}+\left\|q_{A^{\prime}, 0}^{\mu}-q_{0}\right\|_{1}=\mathrm{O}(\mu)$. Hence, by the properties of $H$ mentioned in the proof of Lemma 2.1 (in particular the fact that $\partial_{w} H\left(0, A^{\prime}, 0\right)$ is invertible) we obtain, for $\mu$ small enough the following bound:

$$
\left\|Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{A^{\prime}, 0}^{\mu}\right\|_{1}=\mathrm{O}\left(\left|\alpha_{A^{\prime}, \theta^{\prime}}^{\mu}\right|\right)
$$

Since, by Lemma 4.2, $\left|\alpha_{A^{\prime}, \theta^{\prime}}^{\mu}\right|=\mathrm{O}\left(\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right)$ we deduce estimate (A.12).
We can now estimate $\widetilde{F}_{\mu}(A, \theta)-V_{\mu}(A, \theta)$. Consider the action functional $\Phi_{\mu, A^{\prime}}$ defined in (1.7). For $|\mu|<\mu_{0}$,

$$
D^{2} \Phi_{\mu, A^{\prime}}(q)[w, w]=\int_{\mathbf{R}} \dot{w}^{2}+\cos (q) w^{2}-\mu \cos (q) w^{2} f\left(A^{\prime}+\omega t\right) \mathrm{d} t=\mathrm{O}\left(\|w\|_{1}^{2}\right)
$$

By the definition of $q_{A^{\prime}, 0}^{\mu}$ (i) in Lemma 2.1), we easily obtain with an integration by parts that $D \Phi_{\mu, A^{\prime}}\left(q_{A^{\prime}, 0}^{\mu}\right) w=0$ for all $w \in \mathbf{X}$ such that $w(0)=0$. Therefore

$$
\Phi_{\mu, A^{\prime}}\left(q_{A^{\prime}, 0}^{\mu}+w\right)=\Phi_{\mu, A^{\prime}}\left(q_{A^{\prime}, 0}^{\mu}\right)+\mathrm{O}\left(\|w\|_{1}^{2}\right)
$$

for all $w \in \mathbf{X}$ such that $w(0)=0$. Hence since $\left(Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{A^{\prime}, 0}^{\mu}\right)(0)=0$,

$$
\begin{aligned}
\widetilde{F}_{\mu}(A, \theta)-V_{\mu}(A, \theta) & =\widetilde{F}_{\mu}\left(A, \theta+l_{\mu}(A, \theta)+\theta^{\prime}\right)-F_{\mu}\left(A, \theta+l_{\mu}(A, \theta)\right) \\
& =\widetilde{F}_{\mu}\left(A^{\prime}, \theta^{\prime}\right)-F_{\mu}\left(A^{\prime}, 0\right) \\
& =\Phi_{\mu, A^{\prime}}\left(Q_{A^{\prime}, \theta^{\prime}}^{\mu}\right)-\Phi_{\mu, A^{\prime}}\left(q_{A^{\prime}, 0}^{\mu}\right)=\mathrm{O}\left(\left\|Q_{A^{\prime}, \theta^{\prime}}^{\mu}-q_{A^{\prime}, 0}^{\mu}\right\|_{1}^{2}\right)
\end{aligned}
$$

We obtain by (A.12) that

$$
\left|\widetilde{F}_{\mu}(A, \theta)-V_{\mu}(A, \theta)\right|=\mathrm{O}\left(\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)^{2}\right)
$$

namely (4.13). We now prove (4.14). As a consequence of (4.11), (2.6) and (A.12), we have

$$
\begin{equation*}
\left|\left(\partial_{\theta} F_{\mu}\right)\left(A, \theta+l_{\mu}(A, \theta)\right)-\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|=\mathrm{O}\left(\mu\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right) \tag{A.13}
\end{equation*}
$$

In the same way, using that

$$
\nabla_{A} F_{\mu}(A, \theta)=\mu \int_{\mathbf{R}} \partial_{\varphi} f(\omega t+A)\left(\cos \left(q_{A, \theta}^{\mu}(t)\right)-1\right) \mathrm{d} t
$$

and an analogous expression for $\nabla_{A} \widetilde{F}_{\mu}(A, \theta)$, we get

$$
\begin{equation*}
\left|\left(\nabla_{A} F_{\mu}\right)\left(A, \theta+l_{\mu}(A, \theta)\right)-\nabla_{A} \widetilde{F}_{\mu}(A, \theta)\right|=\mathrm{O}\left(\mu\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right) \tag{A.14}
\end{equation*}
$$

We compute

$$
\nabla_{A} V_{\mu}(A, \theta)-\nabla_{A} \widetilde{F}_{\mu}(A, \theta)=\left(\nabla_{A} F_{\mu}\right)\left(A, \theta+l_{\mu}(A, \theta)\right)-\nabla_{A} \widetilde{F}_{\mu}(A, \theta)
$$

$$
\begin{aligned}
& +\left(\partial_{\theta} F_{\mu}\right)\left(A, \theta+l_{\mu}(A, \theta)\right) \nabla_{A} l_{\mu}(A, \theta) \\
= & \mathrm{O}\left(\mu\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right)+\mathrm{O}\left(\mu \mid\left(\partial_{\theta} F\right)\left(A, \theta+l_{\mu}(A, \theta)\right)\right. \\
& \left.-\partial_{\theta} \widetilde{F}_{\mu}(A, \theta) \mid\right)
\end{aligned}
$$

by (A.14) and since $\left|\nabla_{A} l_{\mu}(A, \theta)\right|=\mathrm{O}(\mu)$. By (A.13) we get

$$
\left|\nabla_{A} V_{\mu}(A, \theta)-\nabla_{A} \widetilde{F}_{\mu}(A, \theta)\right|=\mathrm{O}\left(\mu\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right)
$$

The estimate

$$
\left|\partial_{\theta} V_{\mu}(A, \theta)-\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|=\mathrm{O}\left(\mu\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right)
$$

is obtained in a similar way.
Proof of Lemma 4.5. - By Theorem 4.1, for $C_{5} \mu|\omega| \leqslant 1 / 2$, for all $A, A^{\prime} \in \mathbf{R}^{n}$ there holds

$$
\begin{aligned}
\left|\left(\psi_{\mu}\right)(A)-\left(\psi_{\mu}\right)\left(A^{\prime}\right)\right| & =\left|A-A^{\prime}+\left(g_{\mu}(A)-g_{\mu}\left(A^{\prime}\right)\right) \omega\right| \\
& \leqslant\left|A-A^{\prime}\right|+\left|g_{\mu}(A)-g_{\mu}\left(A^{\prime}\right)\right||\omega| \leqslant\left|A-A^{\prime}\right|+C_{5} \mu|\omega|\left|A-A^{\prime}\right| \\
& \leqslant 3\left|A-A^{\prime}\right| / 2
\end{aligned}
$$

Hence, for all $A, A^{\prime} \in \mathbf{R}^{n}$,

$$
\begin{equation*}
\left|\psi_{\mu}^{-1}(A)-\psi_{\mu}^{-1}\left(A^{\prime}\right)\right| \geqslant \frac{2}{3}\left|A-A^{\prime}\right| \geqslant \frac{1}{2}\left|A-A^{\prime}\right| \tag{A.15}
\end{equation*}
$$

We now verify that the image of $B_{\alpha}\left(A_{0}\right)$ under the homeomorphism $\psi_{\mu}^{-1}$ contains the ball $B_{\alpha / 2}\left(\psi_{\mu}^{-1}\left(A_{0}\right)\right)$ of radius $\alpha^{\prime}=\alpha / 2$ centered at $\psi_{\mu}^{-1}\left(A_{0}\right)$. Since $\psi_{\mu}^{-1}$ is a homeomorphism, $\partial \psi_{\mu}^{-1}\left(B_{\alpha}\left(A_{0}\right)\right)=\psi_{\mu}^{-1}\left(\partial B_{\alpha}\left(A_{0}\right)\right)$. Let $A \in \partial B_{\alpha}\left(A_{0}\right)$, namely $\left|A-A_{0}\right|=\alpha$. Then by (A.15) $\left|\psi_{\mu}^{-1}(A)-\psi_{\mu}^{-1}\left(A_{0}\right)\right| \geqslant \alpha / 2$ and then $\operatorname{dist}\left\{\partial \psi_{\mu}^{-1}\left(B_{\alpha}\left(A_{0}\right)\right), \psi_{\mu}^{-1}\left(A_{0}\right)\right\}>\alpha / 2$. This means that $B_{\alpha / 2}\left(\psi_{\mu}^{-1}\left(A_{0}\right)\right) \subset \psi_{\mu}^{-1}\left(B_{\alpha}\right)$. Recalling that $\widetilde{G}_{\mu}=G_{\mu} \circ \psi_{\mu}^{-1}$ we have that $\inf _{\psi_{\mu}^{-1}(U)} G_{\mu}=\inf _{U} \widetilde{G}_{\mu}, \inf _{\partial \psi_{\mu}^{-1}(U)} G_{\mu}=\inf _{\partial U} \widetilde{G}_{\mu}, \sup _{\psi_{\mu}^{-1}\left(B_{\alpha}\right)} G_{\mu}=\sup _{B_{\alpha}} \widetilde{G}_{\mu}$ and then we easily deduce (i)-(ii). Properties (iii) follows as well since by (A.15) the homeomorphism $\psi_{\mu}^{-1}$ reduces the distance between sets at most by a factor $1 / 2$.

Proof of remark 4.2. - We shall use in this proof the results of Section 4.2.
(i) We first prove that $F_{\mu}(A, \theta)$ has a complex analytic extension. Let

$$
Y=\left\{w \in W^{1, \infty}(\mathbf{R}, \mathbf{C}) \mid\|w\|:=\sup _{t \in \mathbf{R}}(|w(t)|+|\dot{w}(t)|) \mathrm{e}^{|t| / 2}<\infty\right\}
$$

and $\bar{Y}=\{w \in Y \mid w(0)=0\}$. We assume that $f$ has an analytic extension defined in an open neighbourhood of

$$
D_{r}:=\mathbf{T}^{n}+\mathrm{i}[-r, r]^{n}
$$

for some $r>0$ and we introduce $|f|_{r}:=\sup _{A \in D_{r}}|f(A)|<\infty$. Let $\Phi_{\mu}: D_{r} \times Y \rightarrow \mathbf{C}$ be defined by

$$
\begin{aligned}
\Phi_{\mu}(A, w)= & \int_{-\infty}^{\infty} \frac{1}{2}\left(\dot{q}_{0}(t)+\dot{w}(t)\right)^{2}+\left(1-\cos \left(q_{0}(t)+w(t)\right)\right) \\
& -\mu\left(1-\cos \left(q_{0}(t)+w(t)\right)\right) f(A+\omega t) \mathrm{d} t
\end{aligned}
$$

$\Phi_{\mu}$ is a smooth functional. Furthermore, since $\cos : \mathbf{C} \rightarrow \mathbf{C}$ is analytic, $\Phi_{\mu}$ is complex differentiable, which means that its differential $D \Phi_{\mu}$ is $\mathbf{C}$-linear everywhere. We shall denote by $\bar{\Phi}_{\mu}$ the restriction of $\Phi_{\mu}$ to $D_{r} \times \bar{Y} ; w \in \bar{Y}$ is a critical point of $\bar{\Phi}_{\mu}(A, \cdot)$ iff the restriction of $q:=q_{0}+w$ to each interval $(-\infty, 0)$ and $(0, \infty)$ is smooth and satisfies $-\ddot{q}(t)+\sin q(t)=\mu \sin q(t) f(A+\omega t)$.

Note that, in analogy with the real case, if $w \in \bar{Y}$ is a solution of the linear equation $-\ddot{w}+\left(\cos q_{0}\right) w=0$ in each interval $(-\infty, 0)$ and $(0, \infty)$, then $w=0$.

As a consequence we can apply the Implicit Function Theorem and obtain in the same way as in the proof of Lemma 2.1 that there is $\rho>0$ such that, for $\mu|f|_{r}$ small enough, for all $A \in D_{r}$ there is a unique $q_{A, 0}^{\mu}=q_{0}+w_{A}^{\mu}$ in $q_{0}+\bar{Y}$ such that $\left\|w_{A}^{\mu}\right\|<\rho$ and

$$
-\ddot{q}_{A, 0}^{\mu}(t)+\sin q_{A, 0}^{\mu}(t)=\mu \sin q_{A, 0}^{\mu} f(A+\omega t) \quad \text { in }(-\infty, 0) \text { and }(0, \infty), \quad q_{A}^{\mu}(0)=\pi
$$

Moreover $q_{A, 0}^{\mu}$ depends smoothly on $A$.
For later purpose, we remark that, in connection with the Implicit Function Theorem, there are two constants $\rho^{\prime}, C>0$ such that, if $Q \in q_{0}+\bar{Y}$ is a solution of $-\ddot{q}+\sin q=$ $\mu \sin q f(A+\omega t)=g(t)$ in each interval $(-\infty, 0)$ and $(0,+\infty)$ with $g \in Y,\|g\|<\rho^{\prime}$ then

$$
\begin{equation*}
\left\|Q-q_{A, 0}^{\mu}\right\| \leqslant C\|g\| \tag{A.16}
\end{equation*}
$$

From (A.16) we can derive as in the proof of lemma 4.4 that

$$
\begin{align*}
& \bar{\Phi}_{\mu}\left(A, Q-q_{0}\right)-\bar{\Phi}_{\mu}\left(A, q_{A, 0}^{\mu}-q_{0}\right) \\
& \quad=D_{w} \bar{\Phi}_{\mu}\left(A, q_{A, 0}^{\mu}-q_{0}\right) \cdot\left(Q-q_{A, 0}^{\mu}\right)+\mathrm{O}\left(\left\|Q-q_{A, 0}^{\mu}\right\|^{2}\right) \\
& \quad=\mathrm{O}\left(\left\|Q-q_{A, 0}^{\mu}\right\|^{2}\right) \tag{A.17}
\end{align*}
$$

By the complex differentiability of $\Phi_{\mu}$, the differential of the map $\left(A \mapsto q_{A, 0}^{\mu}\right.$ ) is $\mathbf{C}$ linear at all point, hence $q_{A, 0}^{\mu}$ depends analytically on $A \in D_{r}$. As a result, $G_{\mu}: \mathbf{T}^{n} \rightarrow \mathbf{R}$ has an analytic extension which we shall still denote by $G_{\mu}$, defined by $G_{\mu}(A)=$ $\Phi_{\mu}\left(w_{A}^{\mu}\right)$. Now let

$$
U_{r, \omega}=\{\theta \in \mathbf{C}:|\omega \| \operatorname{Im} \theta|<r / 2\} .
$$

Since $F_{\mu}(A, \theta)=G_{\mu}(A+\omega \theta)$ for all $A \in \mathbf{T}^{n}, \theta \in \mathbf{R}, F_{\mu}$ too has an analytic extension, defined in $D_{r / 2} \times U_{r, \omega}$ by $F_{\mu}(A, \theta)=G_{\mu}(A+\theta \omega)$.
(ii) Now we justify that $l_{\mu}$ too has an analytic extension and that the estimate $\left|V_{\mu}-\widetilde{F}_{\mu}\right|=\mathrm{O}\left(\left|\partial_{\theta} \widetilde{F}_{\mu}\right|^{2}\right)$ may be extended to complex values of $(A, \theta)$. Let

$$
S^{\prime}=\left\{\theta \in \mathbf{C}: \operatorname{Im} \theta \in\left(-\frac{\pi}{4}, \frac{\pi}{4}\right)\right\}
$$

As a straightforward consequence of Lemma 4.6 (where $\sigma$ is not small but for instance equal to $\pi / 6$ ) for $\mu|f|_{r}$ small enough, for all $A \in D_{r}, \theta \in S^{\prime}$, there is a unique
$Q_{A, \theta}^{\mu} \in q_{0}+Y=q_{\theta}+Y$, there is a unique $\alpha_{A, \theta}^{\mu} \in \mathbf{C}$ such that $\left\|Q_{A, \theta}^{\mu}-q_{\theta}\right\|+\left|\alpha_{A, \theta}^{\mu}\right| \leqslant \rho$ and

$$
\left\{\begin{array}{l}
-\ddot{Q}_{A, \theta}^{\mu}+\sin Q_{A, \theta}^{\mu}=\mu \sin Q_{A, \theta}^{\mu} f(A+\omega t)+\alpha_{A, \theta}^{\mu} \psi_{\theta} \\
\int_{-\infty}^{\infty}\left(Q_{A, \theta}^{\mu}-q_{\theta}\right) \psi_{\theta} \mathrm{d} t=0 .
\end{array}\right.
$$

Moreover $\left\|Q_{A, \theta}^{\mu}-q_{\theta}\right\|+\left|\alpha_{A, \theta}^{\mu}\right|=\mathrm{O}\left(\mu|f|_{r}\right)$,

$$
\begin{equation*}
\left\|\partial_{A} Q_{A, \theta}^{\mu}\right\|=\mathrm{O}\left(\mu|f|_{r}\right) \tag{A.18}
\end{equation*}
$$

and $Q_{A, \theta}^{\mu}, \alpha_{A, \theta}^{\mu}$ depend analytically on $(A, \theta) \in D_{r} \times S^{\prime}$.
We have $\widetilde{F}_{\mu}(A, \theta)=\Phi_{\mu}\left(A, Q_{A, \theta}^{\mu}-q_{\theta}\right)$. In a similar way as in Lemma 4.2 we obtain

$$
\begin{equation*}
\left|\alpha_{A, \theta}^{\mu}\right|=\mathrm{O}\left(\left|\partial_{\theta} \widetilde{F}_{\mu}(A, \theta)\right|\right) \tag{A.19}
\end{equation*}
$$

Now by Lemma 4.1,

$$
\begin{equation*}
Q_{A, 0}^{\mu}(t+\eta)=Q_{A+\omega \eta,-\eta}^{\mu}(t) \tag{A.20}
\end{equation*}
$$

for all $A \in \mathbf{T}^{n}, \eta \in \mathbf{R}$. Hence, by analyticity, (A.20) holds for all $A \in D_{r}, \eta \in \mathbf{R}$. In particular $Q_{A, 0}^{\mu}(\eta)=Q_{A+\omega \eta,-\eta}^{\mu}(0)$. This implies that for $A \in D_{r / 2}, Q_{A, 0}^{\mu}$ has an analytic extension defined in $U_{r, \omega}$ (provided $\left.r /(2|\omega|) \leqslant \pi / 4\right)$. We can write for $A \in D_{r / 2}$, $\eta \in U_{r, \omega}$,

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} \eta} Q_{A+\omega \eta,-\eta}^{\mu}(0)={\dot{Q^{\mu}}}_{A+\omega \eta,-\eta}(0)=\dot{q}_{0}(\eta)+\mathrm{O}(\mu)=\dot{q}_{0}(0)+\mathrm{O}(|\eta|)+\mathrm{O}(\mu) \tag{A.21}
\end{equation*}
$$

We already know (Lemma 4.3) that, for $A \in \mathbf{T}^{n}$, there is $\widetilde{l}(A)=l_{\mu}(A, 0)=\mathrm{O}(\mu) \in \mathbf{R}$ such that $\left.Q_{A, 0}^{\mu} \widetilde{l}(A)\right)=\pi$, which writes also $Q_{A+\omega \bar{l}(A),-\widetilde{l}(A)}^{\mu}(0)=\pi$. Hence by (A.18), for $A^{\prime} \in[-r, r]^{n}$ and $B=A+\mathrm{i} A^{\prime}+\omega \widetilde{l}(A)$,

$$
Q_{B,-\widetilde{l}(A)}^{\mu}(0)=\pi+\mathrm{O}\left(\mu\left|A^{\prime}\right|\right)
$$

Hence, since $\dot{q}_{0}(0) \neq 0$, by (A.21) and by the Implicit Function Theorem, if $\mu|\omega|\left|A^{\prime}\right|$ is small enough, there is $\widetilde{l}\left(A+\mathrm{i} A^{\prime}\right)=\widetilde{l}(A)+\mathrm{O}\left(\mu\left|A^{\prime}\right|\right)$ such that

$$
\left.Q_{A+\mathrm{i} A^{\prime}+\omega l}^{\mu} \widetilde{l}\left(A+\mathrm{i} A^{\prime}\right),-\widetilde{l}\left(A+\mathrm{i} A^{\prime}\right) 0\right) ~(0)=\pi
$$

Moreover $\tilde{l}$ thus defined is analytic in $D_{r^{\prime}}$ for $r^{\prime}$ small enough (how small depending on $|\omega|)$. For a given $A \in D_{r^{\prime}}$, let $q=Q_{A+\omega \bar{l}(A),-\widetilde{l}(A)}^{\mu}$. Then, with the abbreviation $\alpha=\alpha_{A+\omega \bar{l}(A),-\widetilde{l}(A)}^{\mu}$,

$$
\left\{\begin{array}{l}
-\ddot{q}+\sin q=\mu \sin q f(A+\omega \widetilde{l}(A)+\omega t)+\alpha \psi_{-\widetilde{l}(A)}, \\
q(0)=\pi
\end{array}\right.
$$

Therefore, by (A.17),

$$
\begin{aligned}
G_{\mu}(A+\widetilde{l}(A) \omega)-\widetilde{G}_{\mu}(A) & =G_{\mu}(A+\widetilde{l}(A) \omega)-\widetilde{F}_{\mu}(A+\omega \widetilde{l}(A),-\widetilde{l}(A)) \\
& =\Phi_{\mu}\left(q_{A+\widetilde{l}(A) \omega}^{\mu}-q_{0}\right)-\bar{\Phi}_{\mu}\left(Q_{A+\widetilde{l}(A) \omega,-\widetilde{l}(A)}^{\mu}-q_{0}\right) \\
& =\mathrm{O}\left(\|\alpha \psi-\widetilde{l}(A)\|^{2}\right) \\
& =\mathrm{O}\left(\left|\partial_{\theta} \widetilde{F}_{\mu}(A+\omega \widetilde{l}(A),-\widetilde{l}(A))\right|^{2}\right)=\mathrm{O}\left(\left|\partial_{\theta} \widetilde{F}_{\mu}(A, 0)\right|^{2}\right)
\end{aligned}
$$

Let $l_{\mu}$ be the analytic map defined in $D_{r^{\prime} / 2} \times U_{r^{\prime} / 2, \omega}$ by $l_{\mu}(A, \theta)=\widetilde{l}(A+\omega \theta)$. Let

$$
V_{\mu}(A, \theta)=F_{\mu}\left(A, \theta+l_{\mu}(A, \theta)\right) .
$$

This is an analytic extension of the map defined in Lemma 4.3. By the previous estimates,

$$
\begin{equation*}
\left|V_{\mu}-\widetilde{F}_{\mu}\right|=\mathrm{O}\left(\left|\partial_{\theta} \widetilde{F}_{\mu}\right|^{2}\right) \tag{A.22}
\end{equation*}
$$

uniformly for $(A, \theta) \in D_{r^{\prime} / 2} \times U_{r^{\prime} / 2, \omega}$.
(iii) To complete the proof, we remark that, if $f, g: U \rightarrow \mathbf{C}$ are analytic in $U$, open subset of $\mathbf{C}^{m}$ and $g \not \equiv 0, f=\mathrm{O}(g)$ locally in $U$, then $f / g$ has an analytic extension defined in the whole set $U$ (see for instance [40], p. 32, Theorem 3.4). Hence $\left(V_{\mu}-\widetilde{F}_{\mu}\right) /\left(\partial_{\theta} \widetilde{F}_{\mu}\right)^{2}$ is real analytic in $\mathbf{T}^{n} \times \mathbf{R}$. As a consequence, $\bar{h}_{\mu}(A, \theta)$ given by (4.19) is real analytic w.r.t. $(A, \theta)$. Moreover, as $l_{\mu}$ is an analytic function, so is $\bar{l}_{\mu}$. Finally $h_{\mu}(A, \theta)=\bar{l}_{\mu}(A, \theta)+\bar{h}_{\mu}\left(A, \theta+\bar{l}_{\mu}(A, \theta)\right)$ is real analytic.

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