AN ASYMPTOTIC UNIVERSAL FOCAL DECOMPOSITION
FOR NON-ISOCRHEONOUS POTENTIALS

C. A. A. DE CARVALHO, M. M. PEIXOTO, D. PINHEIRO, AND A. A. PINTO

Abstract. Galileo, in the seventeenth century, observed that the small oscillations of a pendulum seem to have constant period. In fact, the Taylor expansion of the period map of the pendulum is constant up to second order in the initial angular velocity around the stable equilibrium. It is well known that, for small oscillations of the pendulum and small intervals of time, the dynamics of the pendulum can be approximated by the dynamics of the harmonic oscillator. We study the dynamics of a family of mechanical systems that includes the pendulum at small neighbourhoods of the equilibrium but after long intervals of time so that the second order term of the period map can no longer be neglected. We analyze such dynamical behaviour through a renormalization scheme acting on the dynamics of this family of mechanical systems. The main theorem states that the asymptotic limit of this renormalization scheme is universal: it is the same for all the elements in the considered class of mechanical systems. As a consequence, we obtain a universal asymptotic focal decomposition for this family of mechanical systems. This paper is intended to be the first in a series of articles aiming at a semiclassical quantization of systems of the pendulum type as a natural application of the focal decomposition associated to the two-point boundary value problem.

1. INTRODUCTION

1.1. Focal decomposition. The concept of focal decomposition was introduced by Peixoto in [55] (under the name of \( \sigma \)-decomposition) and further developed by Peixoto and Thom in [59]. The starting point was the two-point boundary value problem for ordinary differential equations of the second order,

\[
\ddot{x} = f(t, x, \dot{x}), \quad t, x, \dot{x} \in \mathbb{R},
\]

\[x(t_1) = x_1, \quad x(t_2) = x_2,
\]

which was formulated for the equation of Euler at the beginning of the calculus of variations in the first half of the eighteenth century. This is the simplest and oldest of all boundary value problems. Accordingly, there is vast literature about it, mostly in the context of applied mathematics where one frequently uses the methods of functional analysis. Here we adopt a different point of view: we look for the number \( i \in \{0, 1, 2, \ldots, \infty\} = \mathbb{N} \) of solutions of problem (1.1), and how this number varies with the endpoints \((t_1, x_1)\) and \((t_2, x_2)\).
Let \( \mathbb{R}^4 = \mathbb{R}^2(t_1, x_1) \times \mathbb{R}^2(t_2, x_2) \) be the set of all pairs of points of the \((t, x)\)-plane and to each point \((t_1, x_1, t_2, x_2) \in \mathbb{R}^4\) associate the number of solutions \(i\) of the boundary value problem (1.1). When \(t_1 = t_2\), the index \(i\) is defined as being 0 if \(x_1 \neq x_2\), and \(\infty\) if \(x_1 = x_2\). The diagonals \(\delta, \Delta\) are defined as

\[
\delta = \{(t_1, x_1, t_2, x_2) \in \mathbb{R}^4 : t_1 = t_2\},
\]

\[
\Delta = \{(t_1, x_1, t_2, x_2) \in \mathbb{R}^4 : t_1 = t_2, x_1 = x_2\}.
\]

Let \(\Sigma_i \subset \mathbb{R}^4\) be the set of points to which the index \(i\) has been assigned. Clearly \(\mathbb{R}^4\) is the disjoint union of all the sets \(\Sigma_i\), that is, these \(\Sigma_i\) define a partition of \(\mathbb{R}^4\). This partition is called the focal decomposition of \(\mathbb{R}^4\) associated with the boundary value problem (1.1):

\[
\mathbb{R}^4 = \Sigma_0 \cup \Sigma_1 \cup \ldots \cup \Sigma_\infty.
\]

The decomposition above is the central object of study for a given \(f\).

From the above definition, the diagonal \(\delta\) is 3-dimensional and is naturally decomposed into the 2-dimensional connected component \(\Delta\) belonging to \(\Sigma_\infty\) and two 3-dimensional connected components belonging to \(\Sigma_0\). This is called the natural stratification of \(\delta\).

If one of the endpoints in (1.1) is kept fixed, say \((t_1, x_1) = (0, 0)\), then the sets \(\Sigma_i\) induce a decomposition of \(\mathbb{R}^2(t_2, x_2)\) by the sets \(\sigma_i = \Sigma_i \cap \{(0, 0) \times \mathbb{R}^2(t_2, x_2)\}\). The restricted problem with base point \((0, 0)\) consists of finding the corresponding focal decomposition of \(\mathbb{R}^2\) by the sets \(\sigma_i\):

\[
\mathbb{R}^2 = \sigma_0 \cup \sigma_1 \cup \ldots \cup \sigma_\infty.
\]

Throughout this paper, when we speak of the focal decomposition associated to problem (1.1), we will be referring to the focal decomposition of \(\mathbb{R}^2\) defined above.

One has some clarification of what the sets \(\sigma_i, \Sigma_i\) are if one relates them to a certain surface associated to each point \((t, x)\) called the star of this point. It is the lifted manifold introduced in [54] for other purposes. Let

\[
\frac{dt}{dt} = 1, \quad \frac{dx}{dt} = u, \quad \frac{du}{dt} = f(t, x, u)
\]

be the first system in \(\mathbb{R}^3(t, x, u)\) equivalent to (1.1). One calls the star associated to the base point \((t, x)\) the surface \(S(t, x) \subset \mathbb{R}^3(t, x, u)\) obtained by the union of the trajectories of (1.2) passing through the points of the line \((t, x, u)\), \(-\infty < u < \infty\). One immediately verifies that if \(\pi\) is the projection \(\pi(t, x, u) = (t, x)\), then: a) with respect to base point \((t_1, x_1)\), \((t_2, x_2) \in \sigma_i\) if and only if \([\pi | S(t_1, x_1)]^{-1}(t_2, x_2)\) consists of \(i\) points; also b) \((t_1, x_1, t_2, x_2) \in \Sigma_i\) if and only if \(S(t_1, x_1) \cap S(t_2, x_2)\) consists of \(i\) solutions of (1.2). See [58, 59] for more details on the relationship between the focal decomposition and stratification theory.

Concerning the star, a final comment seems appropriate. To each curve \(\gamma\) solution of a two-point boundary value problem, we associate two stars, one at each endpoint, which do intersect. When they do so transversally, we say that \(\gamma\) is transversal or “structurally stable” because \(\gamma\) cannot be perturbed away by small variations of the equation. Besides, this situation is generic [54], i.e. any equation can be approximated by one for which the solutions of the two-point boundary problem are “structurally stable”.
A notable example of a focal decomposition due to Peixoto and Thom \cite{59} is provided by the focal decomposition of the pendulum equation $\ddot{x} + \sin x = 0$ with base point $(0,0)$ (see Figure 1). This focal decomposition contains non-empty sets $\sigma_i$ with all finite indices. Every set $\sigma_{2k-1}$, $k = 1, 2, \ldots$, consists of a 2-dimensional open set plus the cusp-point $(\pm k\pi, 0)$; they all have two connected components. All four connected components of the even-indexed sets $\sigma_{2k}$ are open-arcs, asymptotic to one of the lines $x = \pm \pi$ and incident to the cusp-points $(\pm k\pi, 0)$; the lines $x = \pm \pi$ are part of $\sigma_1$, except for the points $(0, \pm \pi)$ which belong to $\sigma_0$.

![Figure 1. The pendulum's focal decomposition.](image)

There are many examples in the calculus of variations where rules are given for determining the number of solutions of the boundary value problem (1.1) and for which the corresponding focal decomposition can be derived. The first such example seems to be a theorem of Goldschmidt (1831) concerning the number of catenaries passing through two points in the plane \cite[p. 340]{65}. See also the papers by MacNeish \cite{49} and Peixoto \cite{54}, and especially the textbook by Collatz \cite[pp. 170ff.]{16}, where several focal decompositions are given. However, the number of distinct indices obtained for the focal decompositions in these examples is very small and the concept itself is not recognized.

After its recognition and early development in \cite{55,59}, the concept of focal decomposition was further developed by Kupka and Peixoto in \cite{41} in the context of geodesics and by Peixoto and Silva in \cite{58} regarding the Euler equation associated with a given Lagrangian function $L$. In \cite{42}, Kupka, Peixoto and Pugh proved generic results pointing out the existence of residual sets in the space of complete Riemannian metrics on a finite-dimensional complete Riemannian manifold $M$, for which bounds upon the number of geodesics of prescribed length on $M$ were obtained. Relationships with the arithmetic of positive definite quadratic
forms are considered by Kupka and Peixoto in [41] and by Peixoto in [56,57], and in [41,56] attention is also drawn to the relationship between focal decomposition and the Brillouin zones of solid state physics, a view that is further developed later by Veerman et al. in [66]. In [57], it is pointed out that the focal decomposition associated to (1.1) is relevant for the computation of the semiclassical quantization of this equation via the Feynman path integral method, and in [13,14] de Carvalho et al. exhibit further relations with quantum statistical mechanics.

The concept of focal decomposition might also be relevant to the study of caustic formation by focusing wavefronts. This appears in optics (see the paper by Berry and Upstill [7]), tsunami formation (see the papers by Berry [5,6]) or general relativity (see the papers by Friedrich and Stewart [32], Hasse et al. [35], Ellis et al. [28] and Ehlers and Newman [27]). It is not so surprising that focal decomposition, i.e. the two-point boundary value problem, is intimately connected to the very formulation of a variational problem, as Euler taught us.

As described above, the focal decomposition is a geometrical object assigning to each pair of end points in $\mathbb{R}^2$, the number of solutions of the two-point boundary value problem (1.1). To determine the number of solutions of each of these boundary value problems one can resort to a number of techniques, including variational methods based on the mountain pass lemma due to Ambrosetti and Rabinowitz [3] and more recent extensions. Indeed, even though we restrict ourselves to the 2-dimensional case here, we should remark that the concept of focal decomposition is still well defined when the phase space of the ordinary differential equations under consideration has higher dimension (see, e.g., [42] for focal decomposition on manifolds of arbitrary dimension).

1.2. Renormalization. Besides focal decomposition, the other key topic in this paper is renormalization — the study of asymptotic self-similarity. The main idea behind renormalization is the introduction of an operator — the renormalization operator — on a space of systems whose action on each system is to remove its small scale behaviour and to rescale the remaining variables to preserve some normalization. If a system converges to some limiting behaviour under iteration of the renormalization operator, then we say that such behaviour is universal. Since the renormalization operator relates different scales, such universal behaviour is self-similar.

In this paper we introduce a new kind of renormalization scheme acting on a family of mechanical systems which includes the pendulum. This scheme has the distinguishing feature that time is not rescaled, but rather translated, while the initial velocities and space are appropriately scaled.

Renormalization is extremely relevant in several areas of physics. The first ideas concerning renormalization were introduced in the 1940s in the context of quantum electrodynamics by Bethe, Feynman, Schwinger and Dyson. These ideas were later developed by Stueckelberg and Petermann in [61] and by Gell-Mann and Low in [33], with the introduction of the Renormalization Group as a tool to improve approximate solutions to quantum field theory equations. Later developments in the subject were due to Callan [11], Symanzik [64] and Weinberg [67], among many others. Kadanoff in [40] and Wilson in [69] introduced renormalization techniques into statistical mechanics in order to improve the understanding of critical phenomena.

The relevance of renormalization in dynamical systems was first noticed by Feigenbaum in [29,30] and, independently, by Coullet and Tresser in [17], with
their discovery of period doubling universality in unimodal maps of the interval. They introduced a renormalization operator — the period doubling operator — to show that period doubling sequences for this class of maps are asymptotically self-similar and that these sequences have an identical form for a large open set of such a class of maps. This renormalization operator consists of time and space asymptotic normalizations chosen to preserve the dynamical characteristics of the maps under iteration. Infinitely renormalizable maps separate the regions with regular (zero entropy) and chaotic (positive entropy) dynamics and, moreover, their invariant set is universal, i.e. it is the same for all infinitely renormalizable unimodal maps. It was soon realized that the period doubling operator was just a restriction of another operator acting on the space of unimodal maps — the renormalization operator — whose dynamical behaviour is much richer.

Renormalization has been a very active research area in dynamical systems in the past two decades: Sullivan in [62–63], McMullen in [52], de Melo and Pinto in [25], Lyubich in [45] and Faria, de Melo and Pinto in [23] studied the renormalization operator acting in families of unimodal maps. The renormalization operator naturally appears in several other families of maps, such as families of critical circle maps or families of annulus maps (see, for instance, MacKay [46, 47], Ostlund et al. [53], Lanford [43,44], Cawley [15], Jiang [36–38], de Melo [24], Martens [51], de Faria and de Melo [21,22], Yampolsky [70,71] and Pinto et al. [60]).

The main subject of this paper is a renormalization scheme acting on the dynamics of a family of mechanical systems that includes the pendulum. Our motivation for the introduction of such a scheme comes from the restricted focal decomposition with base point \((0,0)\) of the pendulum equation \(\ddot{x} + \sin x = 0\) in Figure 1. It turns out that the sequence formed by the even-indexed sets in the pendulum’s focal decomposition is approximately self-similar. The renormalization scheme we introduce can then be justified in the following way: for a large integer \(n\), we consider the even-indexed set \(\sigma_{2n}\) and, contrary to previous renormalizations, we do not rescale time but just shift it so that its origin is at \(t = n\pi\). We then restrict the initial velocities to a small interval so that the index corresponding to the shifted even-indexed set is equal to one; we complete the procedure by normalizing space in such a way that the shifted even-indexed set is asymptotic to the lines \(x = \pm 1\). Under iteration of this renormalization scheme, we obtain asymptotic trajectories and focal decomposition are universal and self-similar.

To be more precise, our renormalization scheme acts on the dynamics of a family of mechanical systems (see the books [1,11,15,50] for details on mechanical systems) defined by a Lagrangian function \(L(x, \dot{x}) = \dot{x}^2/2 - V(x)\), where \(V(x)\) is a given non-isochronous potential, i.e. not all the periodic solutions of the corresponding Euler–Lagrange equation have the same period. For more information on isochronous potentials, see the paper [9] by Bolotin and MacKay and the references therein, or Calogero’s book [12]. We obtain that the asymptotic limit of this renormalization scheme is universal, and as a consequence we obtain a universal asymptotic focal decomposition for this family of mechanical systems.

The main strength of our method relies on the fact that the asymptotic limit of the renormalization is explicitly known, facilitating the analysis of the corresponding focal decomposition. As it stands, it applies only to one degree of freedom
mechanical systems. However, we plan to extend our results to finite and infinite degrees of freedom in the future.

This paper is the first step towards a broader research program, proposed by Peixoto and Pinto and described in the survey articles [19, 20], connecting renormalization techniques, focal decomposition of differential equations and semiclassical physics. The next step of this research program is based on an extension of this renormalization scheme to obtain a 4-dimensional universal asymptotic focal decomposition, i.e. with no restrictions on the base point of the boundary value problem. The ultimate goal of such a program would be to deal with applications of this renormalization procedure to semiclassical physics.

The present paper is structured in the following way: in Section 2, we introduce the main concepts we deal with throughout the paper and rigorously state our main results, while Section 3 is devoted to providing the reader with a global picture of the strategy of the proof of the main results. The remaining sections are more technical: Section 4 deals with mechanical systems defined by some quartic potentials – we compute explicit solutions of the Euler–Lagrange equation associated with such mechanical systems in terms of Jacobian elliptic functions (see the textbooks [2, 10, 68] for further details on Jacobian elliptic functions) and use known estimates for such functions to obtain estimates for our explicit solutions; Section 5 is devoted to the study of mechanical systems defined by perturbations of the quartic potentials – we use an implicit function argument inspired in a paper by Bishnani and MacKay [8] to obtain estimates for the solutions of this second class of mechanical systems as a continuation of the estimates obtained previously for the explicit solutions associated with the quartic potentials; Section 6 is devoted to the study of the period map for the periodic orbits of this family of mechanical systems; Sections 7 and 8 are devoted to the renormalization scheme introduced in this paper and its convergence. In Section 9, we show why we believe that our program is relevant for semiclassical physics. We summarize in Section 10.

2. Main theorems and definitions

We start this section by fixing notation and introducing basic definitions which will be used throughout the paper. We also state the main results to be proved in the following sections.

2.1. Setting. We consider mechanical systems defined by a Lagrangian function \( \mathcal{L} : \mathbb{R}^2 \to \mathbb{R} \) of the form

\[
\mathcal{L} \left( q, \frac{dq}{d\tau} \right) = \frac{1}{2} m \left( \frac{dq}{d\tau} \right)^2 - V(q),
\]

where the potential function \( V : \mathbb{R} \to \mathbb{R} \) is a non–isochronous potential. Furthermore, we assume that the potential \( V \) is a \( C^\kappa \) map (\( \kappa \geq 5 \)) with a Taylor expansion at a point \( q^* \in \mathbb{R} \) given by

\[
V(q) = V(q^*) + \frac{V''(q^*)}{2} (q - q^*)^2 + \frac{V^{(4)}(q^*)}{4!} (q - q^*)^4 + O \left( |q - q^*|^5 \right),
\]

where \( V''(q^*) > 0 \) and \( V^{(4)}(q^*) \neq 0 \). The Euler–Lagrange equation associated with (2.1) is

\[
m \frac{d^2 q}{d\tau^2} = - \frac{dV}{dq}(q),
\]
and the corresponding Hamilton equations are given by
\[
\frac{dq}{d\tau} = \frac{p}{m}, \\
\frac{dp}{d\tau} = -\frac{d\mathcal{V}}{dq}(q).
\]
(2.3)
Therefore, the point \( q^* \) is an elliptic equilibrium of (2.2) (or equivalently, \((q^*, 0)\) is an elliptic equilibrium of (2.3)) and thus, there is a 1-parameter family of periodic orbits covering a neighbourhood of the equilibrium point.

We believe that it is worth remarking that the assumption \( \mathcal{V}^{(3)}(q^*) = 0 \) is used to provide low order symmetry to the potential function \( \mathcal{V} \) and cannot be removed through the use of a canonical transformation preserving (2.1) as a natural mechanical system.

2.2. Asymptotic universal behaviour for the trajectories. Since \( q^* \) is an elliptic equilibrium of (2.2) there is \( \phi > 0 \) such that for all initial velocity \( \nu \in [-\phi, \phi] \) the solutions \( q(\nu; \tau) \) of the Euler–Lagrange equation (2.2) with initial conditions \( q(\nu; 0) = q^* \) and \( dq/d\tau(\nu; 0) = \nu \) are periodic. Thus, the trajectories \( q : [-\phi, \phi] \times \mathbb{R} \to \mathbb{R} \) of (2.2) are well defined by \( q(\nu; \tau) \) for all \( \tau \in \mathbb{R} \) and \( \nu \in [-\phi, \phi] \). Furthermore, there exist \( \alpha > 0 \) small enough and \( N \geq 1 \) large enough such that, for every \( n \geq N \), the \( n \)-renormalized trajectories \( x_n : [-1, 1] \times [0, \alpha n] \to \mathbb{R} \) are well defined by
\[
(2.4) \quad x_n(v; t) = (-1)^n \Gamma_{n,t}^{-1} \left[ q \left( \Gamma_{n,t} \mu \omega v; \frac{n\pi - \ell t}{\omega} \right) - q^* \right],
\]
where \( \Gamma_{n,t} \) is the \((n, t)\)-scaling parameter
\[
(2.5) \quad \Gamma_{n,t} = \left( \frac{8t}{3\pi n} \right)^{1/2},
\]
\( \ell \) is \( \pm 1 \) depending on the sign of \( \mathcal{V}^{(4)}(q^*) \) and \( \omega \) and \( \mu \) are given by
\[
(2.6) \quad \omega = \left( \frac{V''(q^*)}{m} \right)^{1/2}, \quad \mu = \left( \frac{3!V'''(q^*)}{|\mathcal{V}^{(4)}(q^*)|} \right)^{1/2}.
\]
Note that \( \omega^{-1} \) and \( \mu \) are the natural time and length scales for the dynamical system defined by (2.2). Furthermore, the variables \( v \) and \( t \) are dimensionless, as well as the \((n, t)\)-scaling parameter \( \Gamma_{n,t} \). Therefore, the \( n \)-renormalized trajectories \( x_n(v; t) \) are dimensionless.

Definition 2.1. The asymptotic trajectories \( X_\ell : [-1, 1] \times \mathbb{R}^+_0 \to \mathbb{R} \) are defined by
\[
X_\ell(v; t) = v \sin \left( \ell t \left( v^2 - 1 \right) \right),
\]
where \( \ell = \pm 1 \) depending on the sign of \( \mathcal{V}^{(4)}(q^*) \).

Let \( C^{2,0}(A \times B, \mathbb{R}) \) denote the set of real valued continuous functions on \( A \times B \) which are two times continuously differentiable with respect to its first variable. We endow such space with the norm
\[
\|f\|_{C^{2,0}(A \times B, \mathbb{R})} = \sup\{\|f\|_\infty, \|\partial_1 f\|_\infty, \|\partial^2_1 f\|_\infty\},
\]
where \( \partial^i_1 f \) denotes the \( i \)th derivative with respect to the first variable of \( f \in C^{2,0}(A \times B, \mathbb{R}) \) and \( \|\|_\infty \) denotes the uniform norm
\[
\|f\|_\infty = \sup\{|f(x, y)| : (x, y) \in A \times B\}.
\]
Theorem 2.2. Let $\ell$ be the sign of $V^{(4)}(q^+)$. For every $0 < \epsilon < 1/4$ there exists $C > 0$ such that

$$
\|x_n(v; t) - X_\ell(v; t)\|_{C^{2,0}([-1, 1] \times [0, n^{1/4 - \epsilon}], \mathbb{R})} \leq C n^{-4\epsilon}.
$$

We remark that the $n$-renormalized trajectories $x_n(v; t)$ may not be defined for small values of $n \in \mathbb{N}$. However, there exists $N_0 > 0$ such that the $n$-renormalized trajectories $x_n(v; t)$ are well defined for every $n \geq N_0$. The value $N_0$ depends on the potential function $V$ introduced in Section 2.1 and $\epsilon \in (0, 1/4)$. It must be chosen in such a way that for every $\epsilon \in (0, 1/4)$ and every $n \geq N_0$, the solutions $q(\tau)$ of the Euler-Lagrange equation (2.2) with initial conditions $q(0) = q^*$ and $dq/d\tau(0) = \Gamma_{n,\ell} \mu \omega v$ are periodic for every $v \in [-1, 1]$ and every $t \in [0, n^{1/4 - \epsilon}]$, i.e.

$$
\Gamma_{n,\ell} \mu \omega \leq \left( \frac{8 n^{1/4 - \epsilon}}{3 \pi n} \right)^{1/2} \mu \omega \leq \left( \frac{8}{3\pi} \right)^{1/2} \mu \omega N_0^{-4/3 + 4\epsilon}/2 \leq \phi.
$$

Hence, it is enough to take $N_0 = ((8/3\pi)(\mu \omega / \phi)^2)^{(4)/(3 + 4\epsilon)}$.

The connection between Theorem 2.2 and the concept of focal decomposition is discussed in detail in Sections 2.3 and 2.4. Section 2.3 describes the focal decomposition induced by the asymptotic trajectories $X_\ell(v; t)$, while Section 2.4 is concerned with the focal decomposition induced by the renormalized trajectories $x_n(v; t)$ and its convergence to the asymptotic focal decomposition. This is the content of Corollaries 2.4, 2.6 and 2.7.

2.3. Asymptotic universal focal decomposition. The renormalization scheme (2.4) induces a renormalization scheme on the associated focal decompositions, as we pass to describe.

The asymptotic trajectories $X_\ell(v; t)$ induce an asymptotic focal decomposition of the half-cylinder $C = \mathbb{R}_0^+ \times [-1, 1]$ given by

$$
C = \sigma_0^* \cup \sigma_1^* \cup \ldots \cup \sigma_\infty^*,
$$

where the sets $\sigma_i^*$ have as elements pairs $(t, x) \in C$ such that $X_\ell(v; t) = x$ has exactly $i$ solutions $v(t, x) \in [-1, 1]$, with each distinct solution corresponding to an asymptotic trajectory connecting the points $(0, 0) \in C$ and $(t, x) \in C$. Therefore, for each $i \in \{0, 1, \ldots, \infty\}$, the set $\sigma_i^* \subset C$ contains all points in $(t, x) \in C$ such that there exist exactly $i$ asymptotic trajectories connecting $(0, 0) \in C$ and $(t, x) \in C$. The following result is then a consequence of Theorem 2.2 and the reasoning above. Indeed, it should be remarked that all that is needed to obtain the next corollary is the $C^0$ convergence of the $n$-renormalized trajectories to the asymptotic trajectories.

Lemma 2.3. There exists an asymptotic universal focal decomposition of $C \subset \mathbb{R}^2$ induced by the asymptotic trajectories $X_\ell(v; t)$. Furthermore, we have that the odd-indexed sets of the asymptotic universal focal decomposition are given by

$$
\bigcup_{k=0}^{\infty} \sigma_{2k+1}^* = \{(t, X_\ell(v; t)) \in C : f_\ell(t, v) = 0, v \in [-1, 1]\} \cup \{(0) \times \mathbb{R}^+\},
$$

where

$$
f_\ell(t, v) = \sin(\ell t(v^2 - 1)) + 2\ell t v^2 \cos(\ell t(v^2 - 1)).
$$
Proof. The existence of an asymptotic universal focal decomposition induced by the asymptotic trajectories $X_\ell(v; t)$ follows from the discussion above. To complete the proof we note that

$$\bigcup_{k=0}^{\infty} \sigma_{2k+1}^* = \left\{ (t, X_\ell(v; t)) \in C : \frac{\partial}{\partial v} X_\ell(v; t) = 0, v \in [-1, 1] \right\} \cup \{0\} \times \mathbb{R}^+.$$

Differentiating the expression for the asymptotic trajectories $X_\ell(v; t)$ in Definition 2.1, we obtain the required result. □

The asymptotic universal focal decomposition is shown in Figure 2. As in the case of the focal decomposition with base point $(0, 0)$ of the pendulum equation $\ddot{x} + \sin x = 0$ (see Figure 1), the asymptotic universal focal decomposition also exhibits non-empty sets $\sigma_i$ with all finite indices. However, contrary to that focal decomposition, which gives a stratification for the whole $\mathbb{R}^2$, our asymptotic universal focal decomposition gives only a stratification of the half-cylinder $C = \mathbb{R}_0^+ \times [-1, 1]$.

There are two main reasons for this to happen, which we pass to explain. First, our renormalization scheme acts only on periodic orbits, neglecting the high-energy non-periodic orbits, which restrains $X_\ell(v; t)$ to the interval $[-1, 1]$. Second, we have defined the renormalization operator only for positive times. Noticing that the mechanical systems we renormalize have time-reversal symmetry, one can extend the asymptotic universal focal decomposition to the cylinder $\mathbb{R} \times [-1, 1]$ by a symmetry on the $x$ axis.

![Figure 2. The asymptotic universal focal decomposition.](image)

We now discuss the construction of Figure 2. We start by noting that $X_\ell(\pm 1; t)$ is identically zero for every $t \in \mathbb{R}^+$, and similarly for $X_\ell(v; 0)$ for every $v \in [-1, 1]$. If $t = 0$ we obtain that $(0, 0) \in \sigma_\infty^*$ and $\{0\} \times ([-1, 1] - \{0\}) \subset \sigma_0^*$. For each fixed
value of \( t \in (0, \pi) \) it is easy to check that \( X_\ell(v;t) \) has a maximum to one side of \( v = 0 \) and a minimum to the opposite side. For such values of \( t \): for \( x \) large enough or \( x \) small enough, \((x,t) \in \sigma^*_0\); for \( x \) between zero and these extreme values, \((x,t) \in \sigma^*_2\); and finally, for \( x = 0 \) we must have \((0,t) \in \sigma^*_3\). At \( t = \pi \) the graph of \( X_\ell(v,\pi) \) is tangent to the horizontal axis at \( v = 0 \). Thus, \((0,\pi) \in \sigma^*_3\) and for each fixed \( t \in (\pi, 2\pi) \) the map \( X_\ell(v;t) \) has one more extreme to each side of \( v \), i.e. one maximum and one minimum to each side of \( v = 0 \), the new extremes being smaller in absolute value and closer to the origin than the existing ones. A reasoning analogous to the one used for the case \( t \in (0, \pi) \) provides us with the indices of the asymptotic focal decomposition for points \((x,t)\) with \( t \in (\pi, 2\pi) \). The full picture is then obtained by induction of the previous argument.

The construction above provides the following characterization for the asymptotic focal decomposition in Figure 2. For every \( k \in \mathbb{N} \), the set \( \sigma^*_2 \) is a 2-dimensional open set with two connected components. The odd-indexed sets \( \sigma^*_2 \) are the union of two open arcs, asymptotic to one of the lines \( x = \pm1 \) and incident to the cusp-point \(((k-1)\pi,0)\), and a line segment joining the cusp-points \(((k-1)\pi,0)\) and \((k\pi,0)\); the lines \( x = \pm1 \) are part of \( \sigma^*_0 \). Thus, the even-indexed sets \( \sigma^*_2 \) are 2-dimensional manifolds, while the odd-indexed sets \( \sigma^*_2 \) are not manifolds because they contain the cusp-points \((k\pi,0)\). On the other hand, if we decompose the odd-indexed sets \( \sigma^*_2 \) into a cusp-point plus three 1-dimensional manifolds (two open arcs and one line segment), then we get a decomposition of the whole plane into a collection of disjoint connected manifolds. The above decomposition of \( \mathcal{C} \) is an example of what is called a stratification of \( \mathcal{C} \), the strata being the disjoint connected manifolds into which \( \mathcal{C} \) was decomposed. Hence \( \sigma^*_1 \) consists of two 1-dimensional strata, \( \sigma^*_2 \) consists of two 2-dimensional strata, \( \sigma^*_3 \) consists of three 1-dimensional strata and one 0-dimensional strata, and so on. To complete the picture, \( \sigma^*_0 \) consists of two 2-dimensional strata plus the \( x \)-axis minus the origin which belongs to \( \sigma^*_\infty \).

We now provide some remarks concerning the differences between the asymptotic universal focal decomposition in Figure 2 and the pendulum focal decomposition of Figure 4. Let \( z(v;t) \) denote the solutions to the pendulum equation with the same initial conditions as the asymptotic trajectories \( X_\ell(v,t) \). The asymptotic trajectories \( X_\ell(v,t) \) inherit some distinctive properties from the definition of the \( n \)-renormalized trajectories in (2.4) and the \((n,t)\)-scaling parameter \( \Gamma_{n,t} \) in (2.5) (also in Section 7), which we now list:

(i) \( X_\ell(\pm1;t) = 0 \) for every \( t \in \mathbb{R}^+ \), while \( \lim_{v \to \pm2} |z_\ell(v;t)| = \pi \) for every \( t \in \mathbb{R}^+ \);
(ii) for every fixed \( t > 0 \) the image of \([-1,1]\) by \( X_\ell(v;t) \) is a proper subset of the interval \([-1,1]\), while the image of \((-2,2)\) by \( z(v;t) \) is the full interval \((-\pi,\pi)\);
(iii) for every fixed \( t \in (0, \pi) \) the map \( X_\ell(v;t) \) has two extreme points (a maximum and a minimum), while the map \( z(v;t) \) has no interior extremes for \( v \in (-2,2) \).

2.4. Convergence to the asymptotic universal focal decomposition. A justification for the labeling of the focal decomposition of Figure 2 as “asymptotic” is also in order. Let \( 0 < \epsilon < 1/4 \) be as in Theorem 2.2 and for each \( n \in \mathbb{N} \) sufficiently
large, denote by $\mathcal{C}_{n,\epsilon}$ the subset of $\mathcal{C}$ given by

$$
\mathcal{C}_{n,\epsilon} = [0, n^{1/4-\epsilon}] \times [-1, 1].
$$

Note that each $n$-renormalized trajectory $x_n$ induces a focal decomposition of the set $\mathcal{C}_{n,\epsilon} \subset \mathcal{C}$ by the sets $\sigma^n_i$ whose elements are points $(t, x) \in \mathcal{C}$ such that there exists a number $i$ of $n$-renormalized trajectories connecting $(0, 0) \in \mathcal{C}_{n,\epsilon}$ to $(t, x) \in \mathcal{C}_{n,\epsilon}$.

Therefore, the $n$-renormalized focal decomposition is given by

$$
(2.8) \quad \mathcal{C}_{n,\epsilon} = \bigcup_{k=0}^{\infty} \sigma^n_k.
$$

Let $d_H(X, Y)$ denote the Hausdorff distance of two subsets $X, Y$ of $\mathbb{R}^2$ with the standard metric. The following result is a consequence of the $\mathcal{C}^2$ convergence in Theorem 2.2. It states that the sequence of $n$-renormalized focal decompositions converges to the asymptotic universal focal decomposition with respect to the Hausdorff distance.

**Corollary 2.4.** For every $0 < \epsilon < 1/4$ and every $i \in \{0, 1, \ldots, \infty\}$ there exists $C > 0$ such that

$$
d_H(\sigma^n_i, \sigma^n_i \cap \mathcal{C}_{n,\epsilon}) \leq C n^{-4\epsilon}.
$$

Corollary 2.4 is a consequence of the $\mathcal{C}^2$ convergence in Theorem 2.2 and the following two results. Before proceeding to the statement and proof of such results, we need to introduce some notation.

**Definition 2.5.** Consider the asymptotic universal focal decomposition of $\mathcal{C}$ in (2.7) and the $n$-renormalized focal decompositions of $\mathcal{C}_{n,\epsilon}$ in (2.8):

- for each $z = (t, x) \in \sigma^n_i$, we define the index $i(z)$ of $z$ equal to $i$;
- for each $z = (t, x) \in \sigma^n_i$ we define the $n$-renormalized index $i_n(z)$ of $z$ as the integer $i$.

**Corollary 2.6.** For every $k \in \mathbb{N} \cup \{0\}$, every compact set $K \subset \sigma^n_{2k}$ and every $z \in K$, the $n$-renormalized index $i_n(z)$ of $z$ converges to the index $i(z)$.

**Proof.** Let $k \in \mathbb{N} \cup \{0\}$ and pick an arbitrary compact set $K \subset \sigma^n_{2k} \subset \mathcal{C}$. Denote by $z = (t, x) \in K$ an arbitrary element of $K$ and choose $n \in \mathbb{N}$ large enough so that for fixed $0 < \epsilon < 1/4$, we have $t < n^{1/4-\epsilon}$.

The $n$-renormalized index of $z = (t, x)$ is the number of solutions $\bar{v}(t, x)$ of the equation

$$
(2.9) \quad x^n_q(\bar{v}, t) = x.
$$

Then, Theorem 2.2 ensures that for such fixed $t$, the $n$-renormalized trajectories $x_n(v; t)$ are $\mathcal{C}^2$ close to the asymptotic trajectories $X_\ell(v; t)$. As a consequence, we obtain that the number of solutions of (2.9) converges to the number of solutions $\bar{v}(t, x)$ of

$$
X_\ell(\bar{v}, t) = x,
$$

where $X_\ell(v, t)$ are the asymptotic trajectories of Definition 2.1. The result then follows by noting that the number of solutions of the previous equation is exactly the index $i(z)$ of $z = (t, x) \in \mathcal{C}$, concluding the proof.
Recall that the odd-index sets of the asymptotic universal focal decomposition \(\sigma_{2k+1}^*\) are the union of three 1-dimensional curves for every \(k \in \mathbb{N}\) and the union of two 1-dimensional curves for \(k = 0\). We will denote by \(\sigma_{2k+1}^{*,0}\), \(k \in \mathbb{N}\), the horizontal line segment given by

\[
\sigma_{2k+1}^{*,0} = \{(t, 0) \in \mathcal{C} : (k-1)\pi < t \leq k\pi\}
\]

and by \(\sigma_{2k+1}^*\) and \(\sigma_{2k+1}^{*-}\), \(k \in \mathbb{N} \cup \{0\}\), the subsets of \(\sigma_{2k+1}^*\) contained in \(\mathbb{R}^+ \times (0,1)\) and \(\mathbb{R}^+ \times (-1,0)\), respectively. We remark that for every \(k \in \mathbb{N}\), \(\sigma_{2k+1}^{*,0}\) is the graph of the function \(s_{2k+1}^{*,0} : ((k-1)\pi, k\pi] \to \mathbb{R}\), which is identically zero. Similarly, for every \(k \in \mathbb{N} \cup \{0\}\) there exist functions \(s_{2k+1}^* : (k\pi, +\infty) \to \mathbb{R}\) and \(s_{2k+1}^{*-} : (k\pi, +\infty) \to \mathbb{R}\) whose graphs are, respectively, the sets \(\sigma_{2k+1}^*\) and \(\sigma_{2k+1}^{*-}\).

A similar reasoning applies to the \(n\)-renormalized focal decompositions of \(C_{n,\epsilon}\). Fix \(\epsilon \in (0, 1/4)\) and \(n \in \mathbb{N}\). For every \(k \in \mathbb{N}\) such that \(k\pi < n^{1/4-\epsilon}\) we have that the odd-index sets of the \(n\)-renormalized focal decompositions \(\sigma_{2k+1}^n\) are the union of three 1-dimensional curves, except for \(\sigma_1^n\) which is the union of only two 1-dimensional curves. As in the case of the asymptotic focal decomposition, these 1-dimensional curves are graphs of functions \(s_{2k+1}^{n,0} : ((k-1)\pi, k\pi] \to \mathbb{R}\) for \(k \in \mathbb{N}\) such that \(k\pi < n^{1/4-\epsilon}\), and \(s_{2k+1}^{n,+} : (k\pi, n^{1/4-\epsilon}) \to \mathbb{R}\) and \(s_{2k+1}^{n,-} : (k\pi, n^{1/4-\epsilon}) \to \mathbb{R}\) for \(k \in \mathbb{N} \cup \{0\}\) such that \(k\pi < n^{1/4-\epsilon}\).

**Corollary 2.7.** For every \(\epsilon \in (0, 1/4)\) and every \(k \in \mathbb{N} \cup \{0\}\) such that \(k\pi < n^{1/4-\epsilon}\) there exists \(C > 0\) such that

\[
\| s_{2k+1}^{n,\pm}(t) - s_{2k+1}^{n,\pm}(t) \|_{C^0((k\pi, n^{1/4-\epsilon}), \mathbb{R})} \leq C n^{-4\epsilon}
\]

and

\[
\| s_{2k+1}^{*,0}(t) - s_{2k+1}^{*,0}(t) \|_{C^0((k-1)\pi, k\pi], \mathbb{R})} \leq C n^{-4\epsilon}.
\]

**Proof.** We note that the odd-indexed sets of the \(n\)-renormalized focal decomposition are contained in the union of line segments of the form \(\{(t, x) \in C_{n,\epsilon} : x = 0\}\) with the set

\[
\left\{(t, x_n(v; t)) \in C_{n,\epsilon} : \frac{\partial}{\partial v} x_n(v; t) = 0, v \in [-1, 1]\right\}.
\]

The result then follows from Theorem 2.2.

\(\square\)

### 3. The renormalization procedure

In this section we introduce an affine change of coordinates of space and time that enables us to map trajectories of the Lagrangian system \(\mathbf{2.1}\) to trajectories of a dimensionless Lagrangian system, simplifying the study of the asymptotic properties of such trajectories.

#### 3.1. An affine change of coordinates and the perturbed quartic potentials.

In order to simplify the proofs throughout the paper, and without loss of generality, we apply a change of coordinates of space and time to the Lagrangian \(\mathbf{2.1}\) and corresponding Euler-Lagrange equation \(\mathbf{2.2}\). The new dimensionless coordinates \((x, t)\) are defined by

\[
(3.1) \quad x = \mu^{-1} (q - q^*), \quad t = \omega t,
\]
where $\omega$ and $\mu$ are as given in (2.6). We obtain a normalized Lagrangian function
\begin{equation}
L(x, \dot{x}) = \frac{1}{2} \dot{x}^2 - V(x),
\end{equation}
with potential function $V(x)$ of the form
\begin{equation}
V(x) = V(0) + \frac{1}{2}x^2 + \frac{\ell}{4}x^4 + f(x),
\end{equation}
where

(a) $f(x) = O(|x|^5)$ is a $C^\kappa$ map ($\kappa \geq 5$);

(b) $\ell = 1$ if $V^{(4)}(q^*) > 0$ and $\ell = -1$ if $V^{(4)}(q^*) < 0$.

Since adding a constant to the potential does not change the form of the associated Euler–Lagrange equation, for simplicity of notation and without loss of generality we rescale the total energy associated with the Lagrangian system (3.2) so that $V(0) = 0$. In the next definition, we introduce the quartic potentials and the perturbed quartic potentials.

**Definition 3.1.** The *quartic potentials* $V_\ell : \mathbb{R} \to \mathbb{R}$ are defined by
\begin{equation}
V_\ell(x) = \frac{1}{2}x^2 + \frac{\ell}{4}x^4, \quad \ell \in \{-1, 1\}.
\end{equation}
The $C^5$ *perturbed quartic potentials* $V : \mathbb{R} \to \mathbb{R}$ are defined by
\begin{equation}
V(x) = \frac{1}{2}x^2 + \frac{\ell}{4}x^4 + f(x),
\end{equation}
where $\ell \in \{-1, 1\}$ and $|f^{(i)}(x)| = O(|x|^{5-i})$ for $i \in \{0, \ldots, 5\}$.

![Figure 3. The quartic potentials $V_\ell(x)$.](image-url)

In the new coordinates $(x, t)$, the Euler-Lagrange equation associated with the normalized Lagrangian (3.2) reduces to
\begin{equation}
\ddot{x} = -(x + \ell x^3 + f'(x)),
\end{equation}
where the dots denote differentiation with respect to the variable $t$. The next lemma describes the relation between the solutions associated with the Lagrangian (2.1) and the normalized Lagrangian (3.2). The result follows easily from the change of coordinates (3.1).
Lemma 3.2. The solution \( q(\nu; \tau) \) of Euler-Lagrange equation (2.2) with initial conditions \( q(\nu; 0) = q^* \) and \( \frac{dq}{dt}(\nu; 0) = \nu \) is related to the solution \( x(\nu; t) \) of the Euler–Lagrange equation (3.6) with initial conditions \( x(\nu; 0) = 0 \) and \( \dot{x}(\nu; 0) = \nu \) by

\[
q(\nu; \tau) = \mu x \left( \mu^{-1} \omega^{-1} \nu; \omega \tau \right) + q^*,
\]

where \( \omega \) and \( \mu \) are as given in (2.6).

3.2. The renormalization procedure. Since \( x = 0 \) is a local minimum of the potential \( V(x) \) given by (3.3) there is \( v_{\text{Max}} > 0 \) such that for all \( |v| < v_{\text{Max}} \) the solutions \( x(\nu; t) \) of the Euler–Lagrange equation (3.6) with initial conditions \( x(\nu; 0) = 0 \) and \( \dot{x}(\nu; 0) = \nu \) are periodic. Furthermore, there exist \( \alpha > 0 \) small enough and \( N \geq 1 \) large enough such that, for every \( n \geq N \), the \( n \)-renormalized trajectories \( x_n : [-1, 1] \times [0, \alpha n] \to \mathbb{R} \) are well defined by

\[
x_n(v; t) = (-1)^n \Gamma_{n,t}^{-1} x(\Gamma_{n,t} v; n \pi - \ell t),
\]

where \( \Gamma_{n,t} \) is the \( (n, t) \)-scaling parameter defined in (2.5). Using Lemma 3.2 we obtain that Theorem 2.2 follows as a consequence of the following result.

Proposition 3.3. For every \( 0 < \epsilon < 1/4 \) there exists \( C > 0 \) such that

\[
\|x_n(v; t) - X_\ell(v; t)\|_{C^{2,0}([-1,1] \times [0, n^{1/4 - \epsilon}], \mathbb{R})} \leq C n^{-4\epsilon},
\]

where \( \ell \in \{-1, 1\} \) is equal to the sign of \( V^{(4)}(0) \).

Sections 4 to 7 are devoted to preparing the proof of Proposition 3.3 which will be completed in Section 8.

4. Quartic potentials

In this section we will make use of Jacobian elliptic functions (see [2][10][68]) to compute explicit solutions, with suitable initial conditions, for trajectories of a given class of Lagrangian systems. Furthermore, we will compute the first terms of the Taylor expansion of such solutions in terms of their initial velocities. In the next section, we will use these Taylor expansions to study the trajectories associated to the more general potentials described in Section 3.

Throughout the next sections, we will make extensive use of the big \( O \) notation. Let \( f(x) \) and \( g(x) \) be two functions defined on some subset of the real numbers. We write \( f(x) \leq O(g(x)) \) for \( x \to +\infty \) if and only if there exist \( M, N \in \mathbb{R} \) such that \( |f(x)| \leq M|g(x)| \) for all \( x \geq N \). We write \( f(x) = O(g(x)) \) for \( x \to +\infty \) if and only if \( f(x) \leq O(g(x)) \) for \( x \to +\infty \) and \( g(x) \leq O(f(x)) \) for \( x \to +\infty \). Similarly, if \( x_0 \in \mathbb{R} \), we write that \( f(x) \leq O(g(x)) \) for \( x \to x_0 \) if and only if there exist \( M, \delta \in \mathbb{R} \) such that \( |f(x)| \leq M|g(x)| \) for all \( x \) with \( |x - x_0| \leq \delta \). We write \( f(x) = O(g(x)) \) for \( x \to x_0 \) if and only if \( f(x) \leq O(g(x)) \) for \( x \to x_0 \) and \( g(x) \leq O(f(x)) \) for \( x \to x_0 \). If no confusion can arise, we simply write \( f(x) = O(g(x)) \), dropping the limit in \( x \).

We will consider Lagrangian functions of the form

\[
L_\ell(x, \dot{x}) = \frac{1}{2} \dot{x}^2 - V_\ell(x),
\]
where $V_\ell(x)$ is a quartic potential introduced in (3.4), and compute explicit solutions of the corresponding Euler–Lagrange equation

$$\dot{x} = -x - \ell x^3,$$

(4.1)

$x(0) = 0$, \hspace{1cm} $\dot{x}(0) = v$.

We denote such solutions by $x_\ell(v; t)$, where $\ell \in \{-1, 1\}$.

Recall that $x = 0$ is a local minimum of $V_\ell(x)$. Thus, there exists $v_{Max} > 0$ such that for all initial velocities $v$ satisfying $|v| < v_{Max}$, the solutions of (4.1) are periodic.

Let $\ell \in \{-1, 1\}$ and let $v > 0$ be small enough so that the solutions of (4.1) are periodic. We define the $(\ell, v)$-coefficients $a_\ell$ and $b_\ell$ by

$$a_\ell := a_\ell(v) = \left(1 + \left(1 + 2\ell v^2\right)^{1/2}\right)^{1/2},$$

$$b_\ell := b_\ell(v) = \left(-\ell + \ell \left(1 + 2\ell v^2\right)^{1/2}\right)^{1/2}.$$  

(4.2)

**Lemma 4.1.** Let $\ell = 1$. The solution $x_1(v; t)$ of (4.1) is given by

$$x_1(v; t) = A_1 \text{sd} \left(\lambda_1 t; m_1\right),$$

where:

(i) $\text{sd}(t; m)$ is a Jacobian elliptic function.

(ii) the amplitude $A_1$, the frequency $\lambda_1$, and the parameter $m_1$ are given by

$$A_1 := A_1(v) = \text{sign}(v)a_1 b_1 / \left(a_1^2 + b_1^2\right)^{1/2},$$

$$\lambda_1 := \lambda_1(v) = \left((a_1^2 + b_1^2)/2\right)^{1/2},$$

$$m_1 := m_1(v) = b_1^2 / \left(a_1^2 + b_1^2\right).$$

(4.3)

(iii) $a_1$ and $b_1$ are as given in (4.2).

**Proof.** Noting that the solution of (4.1) is also the solution of the first order differential equation

$$\frac{1}{2} \ddot{x}^2 + V_1(x) = \frac{1}{2} v^2$$

with initial condition $x_1(0) = 0$, integrating the above equation we obtain that the solution $x_1(v; t)$ is implicitly determined by

$$\int_0^x \left(\frac{1}{2} v^2 - V_1(z)\right)^{-1/2} \, dz = \text{sign}(v)t/\sqrt{2}.$$ 

Using formula 17.4.51 from [2], we get

$$\frac{1}{\lambda_1} \left(\text{sd}^{-1} \left(x_1/A_1; m_1\right)\right) = \text{sign}(v)t/\sqrt{2},$$

where $A_1$, $\lambda_1$ and $m_1$ are defined in (4.3) in terms of $a_1$ and $b_1$ given in (4.2). The lemma follows by solving the above equality with respect to $x_1$ to obtain $x_1(v; t)$.  

The proof of the next lemma is analogous to the proof of Lemma 4.1 (using formula 17.4.45 from [2]). We skip its proof.
Lemma 4.2. Let $\ell = -1$. The solution $x_{-1}(v; t)$ of (4.1) with $|v| < v_{\text{Max}}$ is given by
\[
x_{-1}(v; t) = A_{-1} \text{sn}(\lambda_{-1} t; m_{-1}),
\]
where:

(i) $\text{sn}(t; m)$ is the Jacobian elliptic sinus.
(ii) the amplitude $A_{-1}$, the frequency $\lambda_{-1}$ and the parameter $m_{-1}$ are given by
\[
A_{-1} := A_{-1}(v) = \text{sign}(v) b_{-1},
\]
\[
\lambda_{-1} := \lambda_{-1}(v) = a_{-1}/\sqrt{2},
\]
\[
m_{-1} := m_{-1}(v) = (b_{-1}/a_{-1})^2.
\]
(iii) $a_{-1}$ and $b_{-1}$ are as given in (4.2).

In the next two lemmas we provide estimates for the amplitude $A_\ell$, frequency $\lambda_\ell$ and parameter $m_\ell$. The estimates are obtained by finding the values of the successive derivatives of $A_\ell$, $\lambda_\ell$ and $m_\ell$ when $v = 0$.

Lemma 4.3. Let $\ell = 1$. There exists $v_{\text{Max}} > 0$ such that for all $v \in (-v_{\text{Max}}, v_{\text{max}})$, we have that
\[
\begin{align*}
|A_1 - \left(v - \frac{1}{2} v^3\right)| &= O\left(|v|^5\right), \\
|\lambda_1 - \left(1 + \frac{1}{2} v^2\right)| &= O\left(v^4\right), \\
|m_1 - \frac{1}{2} v^2| &= O\left(v^4\right).
\end{align*}
\]

Lemma 4.4. Let $\ell = -1$. There exists $v_{\text{Max}} > 0$ such that for all $v \in (-v_{\text{Max}}, v_{\text{Max}})$, we have that
\[
\begin{align*}
|A_{-1} - \left(v + \frac{1}{4} v^3\right)| &= O\left(|v|^5\right), \\
|\lambda_{-1} - \left(1 - \frac{1}{4} v^2\right)| &= O\left(v^4\right), \\
|m_{-1} - \frac{1}{2} v^2| &= O\left(v^4\right).
\end{align*}
\]

Lemma 4.5. Let $j_1 = \frac{1}{2}$ and $j_{-1} = -\frac{1}{4}$ and define $e_\ell(v) = 1 + j_\ell v^2$, so that we have
\[
|\lambda_\ell - e_\ell(v)| = O(v^4).
\]
There is $v_0 > 0$ such that for all $|v| < v_0$, the following estimates hold:
\[
\begin{align*}
|\sin(\lambda_\ell(v) t) - \sin(e_\ell(v) t)| &= O(tv^4), \\
(e_\ell(v) t) - \sin(t) &= O(tv^2).
\end{align*}
\]

Proof. We use the estimates for $\lambda_\ell$ in Lemmas 4.3 and 4.4, the definition of $e_\ell$ and the trigonometry formulas for the sine and cosine of the sum of two angles. $\square$

The next lemma gives us two classical estimates on Jacobian elliptic functions (formulas 16.13.1 and 16.13.3 of [2]).
Lemma 4.6. For every $T > 0$ there is $m_0 > 0$ small enough such that for all $|m| < m_0$, the estimates
\[
\left| \text{sn}(t; m) - \left( \sin(t) - \frac{1}{4} m (t - \sin(t) \cos(t)) \cos(t) \right) \right| = O \left( m^2 \right),
\]
\[
\left| \text{dn}(t; m) - \left( 1 - \frac{1}{2} m \sin^2(t) \right) \right| = O \left( m^2 \right)
\]
hold for all $t \in [0,T]$.

The constant $m_0$ in the statement of the previous lemma is a function of $T$. Indeed, for fixed $t \in \mathbb{R}$, we have that $\text{sn}(t; m)$ and $\text{dn}(t; m)$ are meromorphic functions of $m$. Lemma 4.6 provides the first two terms of the MacLaurin series of these functions with respect to $m$. The radius of convergence of such a series is then equal to the distance $d_m(t)$ between the origin and the nearest pole in the complex $m$-plane. Moreover, $d_m(t)$ is strictly positive and depends continuously on $t$. Thus, we obtain that $m_0 = \min_{t \in [0,T]} d_m(t) > 0$.

In the next lemma we rewrite the classical estimates of Lemma 4.6 in a form that will be useful for the proof of Lemma 4.8.

Lemma 4.7. Let $g_\ell(t)$ be given by
\[
g_\ell(t) = t \cos(t) + \ell \left( \sin(t) \cos^2(t) + 2 \sin(t) \right).
\]
For every $T > 0$ there is $m_0 > 0$ such that for all $|m| < m_0$, the estimates
\[
\left| \text{sd}(t; m) - \left( \sin(t) - \frac{1}{4} m (g_1(t) - 4 \sin(t)) \right) \right| = O \left( m^2 \right),
\]
\[
\left| \text{sn}(t; m) - \left( \sin(t) - \frac{1}{4} m (g_{-1}(t) + 2 \sin(t)) \right) \right| = O \left( m^2 \right)
\]
hold for every $t \in [0,T]$.

Proof. Fix $T > 0$ and let $t \in [0,T]$. For the first item, using the definition of the Jacobian elliptic function $\text{sd}(t; m)$ (see [2]), we obtain
\[
\text{sd}(t; m) = \frac{\text{sn}(t; m)}{\text{dn}(t; m)}.
\]
Using the second identity in Lemma 4.6 and the Taylor series, we have that there is $m_0 > 0$ such that for all $|m| < m_0$ we have that
\[
\left| \frac{1}{\text{dn}(t; m)} - \left( 1 + \frac{1}{2} m \sin^2(t) \right) \right| = O \left( m^2 \right)
\]
holds for every $t \in [0,T]$. Combining identity (4.6) and identity (4.7) with the first item of Lemma 4.6 we obtain that for every $|m| < m_0$ we have that
\[
\left| \text{sd}(t; m) - \left( \sin(t) - \frac{1}{4} m (g_1(t) - 4 \sin(t)) \right) \right| = O \left( m^2 \right),
\]
where $g_1(t)$ is as given in the statement, completing the proof of (4.4).

The second equality follows from Lemma 4.6 and the definition of $g_{-1}(t)$. We obtain that there is $m_0 > 0$ such that for all $|m| < m_0$, we get
\[
\left| \text{sn}(t; m) - \left( \sin(t) - \frac{1}{4} m (g_{-1}(t) + 2 \sin(t)) \right) \right| = O \left( m^2 \right),
\]
as required. \qed
Lemma 4.8. Let $\ell \in \{-1, 1\}$. For every $T > 0$ there exists $v_0 > 0$ such that for every $|v| < v_0$ and every $t \in [0, T]$, the solution $x_\ell(v; t)$ of (4.1) satisfies the estimate

$$
|x_\ell(v; t) - \left(v \sin(e_\ell(v)t) - \frac{1}{8} v^3 g_\ell(e_\ell(v)t)\right)| = O(t|v|^5 + |v|^3),
$$

where $e_\ell$ is defined in Lemma 4.5 and $g_\ell$ is defined in Lemma 4.7.

In particular, we get

$$
|x_\ell(v; t) - v \sin(t)| = O(t|v|^3 + |v|^3).
$$

Proof. Fix $T > 0$ and let $t \in [0, T]$. We will split the proof into the following two cases: $\ell = 1$ and $\ell = -1$.

Case $\ell = 1$. For simplicity of notation, we will use $A = A_1$, $\lambda = \lambda_1$, $m = m_1$, $g = g_1$ and $e = e_1$.

From Lemma 4.1, we have that

$$
x_1(v; t) = A \sin(\lambda t; m).
$$

Combining (4.4) in Lemma 4.7 with (4.8), we obtain

$$
x_1(v; t) - A \left(\sin(\lambda t) - \frac{1}{4} m (g(\lambda t) - 4 \sin(\lambda t))\right) = O(A m^2).
$$

Combining the estimate for $m$ in Lemma 4.3 with (4.9), we obtain

$$
x_1(v; t) - A \left(\sin(\lambda t) - \frac{1}{8} v^2 (g(\lambda t) - 4 \sin(\lambda t))\right) = O(A v^4).
$$

Putting together the estimate for $A$ in Lemma 4.3 and (4.10), we get

$$
x_1(v; t) - \left(v \sin(\lambda t) - \frac{1}{8} v^3 g(\lambda t)\right) = O\left(|v|^5\right).
$$

Applying Lemma 4.5 to (4.11), we get that

$$
x_1(v; t) - \left(v \sin(e(t)) - \frac{1}{8} v^3 g(e(t))\right) = O\left(t|v|^5 + |v|^5\right),
$$

for every $t \in [0, T]$.

Case $\ell = -1$. For simplicity of notation, we will use $A = A_{-1}$, $\lambda = \lambda_{-1}$, $m = m_{-1}$, $g = g_{-1}$ and $e = e_{-1}$.

From Lemma 4.2, we have that

$$
x_{-1}(v; t) = A \sin(\lambda t; m).
$$

Combining (4.5) in Lemma 4.7 with (4.13), we get

$$
x_{-1}(v; t) - A \left(\sin(\lambda t) - \frac{1}{4} m (g(\lambda t) + 2 \sin(\lambda t))\right) = O(A m^2).
$$

Combining the estimate for $m$ in Lemma 4.4 and (4.14), we obtain

$$
x_{-1}(v; t) - A \left(\sin(\lambda t) - \frac{1}{8} v^2 (g(\lambda t) + 2 \sin(\lambda t))\right) = O(A v^4).
$$

Putting together the estimate for $A$ in Lemma 4.4 and (4.15), we get

$$
x_{-1}(v; t) - \left(v \sin(\lambda t) - \frac{1}{8} v^3 g(\lambda t)\right) = O\left(|v|^5\right).\]
Using Lemma 4.5 in (4.16), we obtain that

\[|x_{-1}(v; t) - \left( v \sin(\ell t) - \frac{1}{8} v^3 g(\ell t) \right)| = O(t|v|^5 + |v|^5)\]

holds for every \( t \in [0, T] \).

Finally, again applying Lemma 4.5 to identities (4.12) and (4.17), for \( \ell \in \{-1, 1\} \), we get

\[|x_\ell(v; t) - v \sin(t)| = O(t|v|^3 + |v|^3),\]

which completes the proof. \( \square \)

The constant \( v_0 \) in the statement of Lemma 4.8 depends on the constant \( m_0 \) of Lemmas 4.6 and 4.7 through the relations detailed in Lemmas 4.1 and 4.2.

5. Perturbed quartic potentials \( V \)

In this section, we will consider perturbed quartic potentials \( V(x) \) that are perturbations of the quartic potentials \( V_\ell(x) \) near the elliptic fixed-point \((0, 0)\) of the first order differential system associated to the Euler–Lagrange equation (3.6). The results below can be seen as applications of the theorem of existence and differentiability of solutions of ODEs with respect to parameters and initial conditions, thus providing estimates for the solutions \( x(v; t) \) of the Euler–Lagrange equation (3.6) associated to \( V(x) \) using the solutions \( x_\ell(v; t) \) associated to the quartic potentials \( V_\ell(x) \).

Throughout this paper, we will denote by \( x(v; t) \) the solutions of the Euler–Lagrange equation (3.6) with initial conditions \( x(v; 0) = 0 \) and \( \dot{x}(v; 0) = v \).

Let \( C^2_0([a, b], \mathbb{R}) \) be the set of all maps \( \psi \in C^2([a, b], \mathbb{R}) \) such that \( \psi(0) = 0 \). Let \( U : C^2_0([0, T], \mathbb{R}) \to C^0([0, T], \mathbb{R}) \times \mathbb{R} \) be the linear operator given by

\[U \psi = \left( \psi + \psi, \dot{\psi}(0) \right).

**Lemma 5.1.** The linear operator \( U : C^2_0([0, T], \mathbb{R}) \to C^0([0, T], \mathbb{R}) \times \mathbb{R} \) is invertible and its inverse \( U^{-1} : C^0([0, T], \mathbb{R}) \times \mathbb{R} \to C^2_0([0, T], \mathbb{R}) \) is a linear operator with bounded norm.

**Proof.** The existence and uniqueness of \( U^{-1}(f(t), v) \) is guaranteed by the theorem of existence and uniqueness of solutions of differential equations. For every \( f \in C^0([0, T], \mathbb{R}) \) and every \( v \in \mathbb{R} \), let us show that

\[\psi_v(t) = v \sin(t) - \left( \int_0^t f(s) \sin(s)ds \right) \cos(t) + \left( \int_0^t f(s) \cos(s)ds \right) \sin(t)\]

is equal to \( U^{-1}(f, v) \). The first derivative of \( \psi_v(t) \) is given by

\[\dot{\psi}_v(t) = v \cos(t) + \left( \int_0^t f(s) \sin(s)ds \right) \sin(t) + \left( \int_0^t f(s) \cos(s)ds \right) \cos(t)\,.

The second derivative of \( \psi_v(t) \) is given by

\[\ddot{\psi}_v(t) = -v \sin(t) + f(t) + \left( \int_0^t f(s) \sin(s)ds \right) \cos(t) - \left( \int_0^t f(s) \cos(s)ds \right) \sin(t)\,.


Hence, \( \dot{\psi}_v(t) + \psi_v(t) = f(t) \), \( \psi_v(0) = 0 \) and \( \dot{\psi}_v(0) = v \), which implies that \( U(\psi_v(t)) = (f(t), v) \). Therefore, the inverse \( U^{-1} : C^0([0,T], \mathbb{R}) \times \mathbb{R} \to C^0([0,T], \mathbb{R}) \) of \( U \) is a linear operator, and

\[
\|U^{-1}(f(t), v)\|_{C^2([0,T], \mathbb{R})} \leq (1 + 2T) \|f\|_{C^0([0,T], \mathbb{R})} + |v| .
\]

Therefore, the linear operator \( U^{-1} \) has bounded norm. \( \square \)

**Theorem 5.2.** For every \( T > 0 \), there exist \( \delta_1 > 0 \), \( \delta_2 > 0 \), \( \delta_3 > 0 \) and \( K > 0 \) with the following properties: For every \( |v| < \delta_1 \) and every \( \epsilon \in C^2([-\delta_2, \delta_2], \mathbb{R}) \), with \( \|\epsilon\|_{C^2([-\delta_2, \delta_2], \mathbb{R})} < \delta_3 \), the ordinary differential equation

\[
\ddot{x} + x + \ell x^3 + \epsilon(x) = 0
\]

has a unique solution \( x_\epsilon(v; t) \in C^2([0, T], \mathbb{R}) \) with \( x_\epsilon(v; 0) = 0 \) and \( \dot{x}_\epsilon(v; 0) = v \), where \( \ell = \pm 1 \). Furthermore, we have that

\[
(5.1) \quad \left\| \frac{\delta x_\epsilon}{\delta \epsilon} (x_\epsilon(v; t), x_\epsilon(v; t)) \right\|_{C^2([0,T], \mathbb{R})} \leq K \|\epsilon(x_\epsilon(v; t))\|_{C^0([0,T], \mathbb{R})} .
\]

**Proof.** Let \( \Theta : C^2_0([0,T], \mathbb{R}) \times C^2_0([-\delta, \delta], \mathbb{R}) \to C^0([0,T], \mathbb{R}) \times \mathbb{R} \), be given by

\[
\Theta(x(t), \epsilon(x)) = (\dot{x} + x + \ell x^3 + \epsilon(x), \dot{x}(0)) .
\]

The non-linear operator \( \Theta \) is \( C^1 \), with partial derivative \( \frac{\delta \Theta}{\delta x} (x(t), \epsilon(x)) : C^2_0([0,T], \mathbb{R}) \to C^0([0,T], \mathbb{R}) \times \mathbb{R} \) given by

\[
\frac{\delta \Theta}{\delta x} (x(t), \epsilon(x)) \psi = \left( \psi + \psi + 3\ell x^2 \psi + D\epsilon(x) \psi, \psi(0) \right) ,
\]

and with partial derivative \( \frac{\delta \Theta}{\delta \epsilon} (x(t), \epsilon(x)) : C^2_0([-\delta, \delta], \mathbb{R}) \to C^0([0,T], \mathbb{R}) \times \mathbb{R} \) given by

\[
\frac{\delta \Theta}{\delta \epsilon} (x(t), \epsilon(x)) \alpha = (\alpha, 0) .
\]

By Lemma 5.1 the linear operator \( U : C^2_0([0,T], \mathbb{R}) \to C^0([0,T], \mathbb{R}) \times \mathbb{R} \) given by

\[
U\psi = \left( \dot{\psi} + \psi, \dot{\psi}(0) \right)
\]

is invertible and its inverse has bounded norm. Furthermore, we have that

\[
\left\| U\psi - \frac{\delta \Theta}{\delta x} (x(t), \epsilon(x)) \psi \right\|_{C^0([0,T], \mathbb{R}) \times \mathbb{R}} \leq \left\| (3\ell x^2 \psi + D\epsilon(x) \psi, 0) \right\|_{C^0([0,T], \mathbb{R}) \times \mathbb{R}} \leq (3C\delta_2^2 + \delta_3) \|\psi\|_{C^2([0,T], \mathbb{R})} .
\]

Hence, there is \( \delta_2 > 0 \) and \( \delta_3 > 0 \) small enough so that

\[
\left\| U - \frac{\delta \Theta}{\delta x} (x(t), \epsilon(x)) \right\| \leq \|U^{-1}\|^{-1} .
\]

Therefore, the linear operator \( \frac{\delta \Theta}{\delta x} : C^2_0([0,T], \mathbb{R}) \to C^0([0,T], \mathbb{R}) \times \mathbb{R} \) is invertible and its inverse has bounded norm

\[
\left\| \left[ \frac{\delta \Theta}{\delta x} (x(t), \epsilon(x)) \right]^{-1} \right\|^{-1} \geq \|U^{-1}\|^{-1} - \left\| U - \frac{\delta \Theta}{\delta x} (x(t), \epsilon(x)) \right\| .
\]

By Lemmas 4.1 and 4.2 there exists \( x_\epsilon(v; t) \in C^2_0([0,T], \mathbb{R}) \) such that \( \Theta(x_t(v; t), 0) = 0 \). Hence, by the implicit function theorem and by the invertibility of the operator \( \frac{\delta \Theta}{\delta x} \), there exist \( \delta_1 > 0 \) and \( \delta_3 > 0 \) small enough with the following property:
for all $|v| < \delta_1$ and for all $\|\epsilon\|_{C^2([\delta_2, \delta_2], \mathbb{R})} < \delta_3$, there exists a unique solution $x_\epsilon(v; t) \in C^3_0([0, T], \mathbb{R})$ such that

$$(5.2) \quad \Theta(x_\epsilon(v; t), \epsilon(x)) = 0.$$  

Differentiating (5.2) with respect to $\epsilon$, we get

$$\frac{\delta \Theta}{\delta \epsilon(x)} \frac{\delta x}{\delta \epsilon(x)} + \frac{\delta \Theta}{\delta \epsilon(x)} = 0.$$  

Hence, the operator $\frac{\delta x}{\delta \epsilon(x)} : C^2([\delta, \delta], \mathbb{R}) \to C^2([0, T], \mathbb{R})$ is given by

$$\frac{\delta x}{\delta \epsilon(x)} = -\left[\frac{\delta \Theta}{\delta x}(x_\epsilon(v; t), \epsilon(x))\right]^{-1}\frac{\delta \Theta}{\delta \epsilon(x)}(x_\epsilon(v; t), \epsilon(x)).$$  

Therefore, there is $K > 0$ such that

$$\left\|\frac{\delta x}{\delta \epsilon(x)}\right\|_{C^2([0, T], \mathbb{R})} \leq K \left\|\frac{\delta \Theta}{\delta \epsilon(x)}\right\|_{C^0([0, T], \mathbb{R})}.$$  

completing the proof. 

\\[ \Box \]

**Lemma 5.3.** Let $V(x)$ be a perturbed quartic potential. For every $T > 0$, there exists $v_0 > 0$ such that, for all $0 \leq t \leq T$ and for all $|v| < v_0$, the solution $x_\ell(v; t)$ of the Euler–Lagrange equation associated to the potential $V(x)$ satisfies the following estimate:

$$\|x(v; t) - x_\ell(v; t)\|_{C^2([0, T], \mathbb{R})} < O(v^4),$$

where $x_\ell(v; t)$ are the solutions of the Euler–Lagrange equation associated to the potential $V_\ell(x)$ and $\ell$ is equal to the sign of $V^{(4)}(0)$.

**Proof.** Let us consider the equation

$$(5.3) \quad \ddot{x} + x + \ell x^3 + \epsilon(x) = 0.$$  

Let $f'(x)$ be as given by (3.3). Using Theorem 5.2 there exists $v_0 > 0$ small enough such that, for every $0 \leq k \leq 1$ and every $|v| < v_0$, the equation (5.3), with $\epsilon(x) = k f'(x)$, has a unique solution $x_k(v; t)$ with initial conditions $x_k(v; 0) = 0$ and $\dot{x}_k(v; 0) = v$. Furthermore, the solution $x_k(v; t)$ is periodic and there exists $C_1 > 0$ such that $|x_k(v; t)| < C_1 |v|$ for all $0 \leq t \leq T$ (see Lemma 6.2). By (3.5), there exists $C_2 > 0$ such that

$$\|f'|_{C^0([-C_1 |v|, C_1 |v|], \mathbb{R})} \leq C_2 v^4.$$  

Hence, using (5.1), there exists $K > 0$ such that

$$\|x(v; t) - x_\ell(v; t)\|_{C^2([0, T], \mathbb{R})} \leq K \|\epsilon(x)\|_{C^0([-C_1 |v|, C_1 |v|], \mathbb{R})} \leq K C_2 v^4,$$

concluding the proof. 

\\[ \Box \]
**Proposition 5.4.** Let \( V(x) \) be a perturbed quartic potential. For every \( T > 0 \), there exists \( v_0 > 0 \) such that, for all \( 0 \leq t \leq T \) and for all \( |v| < v_0 \), the solution \( x(v; t) \) of the Euler–Lagrange equation associated to the potential \( V(x) \) satisfies the following estimate:

\[
\left| x(v; t) - \left( v \sin(e_\ell(v)t) - \frac{1}{8} v^3 g_\ell(e_\ell(v)t) \right) \right| = O(v^4),
\]

where \( x_\ell(v; t) \) are the solutions of the Euler–Lagrange equation associated to the potential \( V_\ell(x) \) and \( \ell \) is equal to the sign of \( V^{(4)}(0) \). In particular, we get

\[
|x(v; t) - v \sin(t)| = O\left(|v|^3|t| + v^4\right).
\]

**Proof.** The proof follows from combining Lemmas 4.8 and 5.3. \( \square \)

6. Period map \( T \)

In this section, we are going to compute the first three terms of Taylor’s expansion of the period map \( T(v) \) which associates to each initial velocity \( v \neq 0 \) of a solution \( x(v; t) \), with \( x(v; 0) = 0 \), its smallest period \( T(v) \) in a small neighborhood of 0.

Since \( x = 0 \) is a local minimum of the perturbed quartic potential, the solutions \( x(v; t) \) of the Euler–Lagrange equation (3.6), with initial conditions \( x(v; 0) = 0 \) and \( \dot{x}(v; 0) = v \), are periodic for all small values of \( v \). Hence, there exists \( v_0 > 0 \), small enough, such that, for all \( |v| < v_0 \), the extreme points \( x^m(v) \leq 0 \leq x^M(v) \) of a solution \( x(v; t) \) are well defined. Since the energy of the system is conserved along its orbits and the initial energy is equal to \( \frac{1}{2} v^2 \), we obtain that \( x^m(v) \) and \( x^M(v) \) are implicitly determined by

\[
V(x^M(v)) = V(x^m(v)) = \frac{1}{2} v^2.
\]

Similarly, we will denote by \( x^{m}_\ell(v) \leq 0 \leq x^{M}_\ell(v) \) the extreme points of the solution \( x_\ell(v; t) \) of (5.1).

**Definition 6.1.** The period map \( T : (-v_0, v_0) \to \mathbb{R} \) is defined by

\[
T(v) = \sqrt{2} \int_{x^m(v)}^{x^M(v)} \left( \frac{1}{2} v^2 - V(x) \right)^{-1/2} dx.
\]

We note that the period of the solution \( x(0, t) \) is equal to 0, but we keep \( T(0) = 2\pi \) for the period map \( T \) to be continuous and smooth.

**Lemma 6.2.** Let \( V(x) \) be a perturbed quartic potential and let \( V_\ell(x) \) be a quartic potential. There exists \( v_0 > 0 \) small enough such that, for all \( |v| < v_0 \), we have that

(i) \( x^m_\ell(v) = -x^M_\ell(v) \),

(ii) there exist constants \( C_1 > C_2 > 0 \) such that

\[
-C_1 |v| < x^m_\ell(v) < -C_2 |v| \quad \text{and} \quad C_2 |v| < x^M_\ell(v) < C_1 |v|,
\]

(iii) the following estimates hold:

\[
|x^M_\ell(v) - x^M_\ell(v)| = O(|v|^5) \quad \text{and} \quad |x^m_\ell(v) - x^m_\ell(v)| = O(|v|^5).
\]

**Proof.** Take \( v_0 > 0 \) small enough so that the solutions \( x(v; t) \) of the Euler–Lagrange equation (3.3), with initial conditions \( x(v; 0) = 0 \) and \( \dot{x}(v; 0) = v \), are periodic for all \( v \) such that \( |v| < v_0 \).
Item (i) is a consequence of the symmetry of the perturbed quartic potentials \( V_\ell(x) \). Item (ii) follows from the fact that

\[
\left| V(x) - \frac{1}{2} x^2 \right| = O(x^4)
\]

and \( V(x^M(v)) = V(x^m(v)) = v^2/2 \).

For the proof of item (iii), recall that

\[
V_\ell(x) = \frac{1}{2} x^2 + \frac{\ell}{4} x^4
\]

and that the solutions of \( V_\ell(x_\ell(v)) = v^2/2 \) are given by

\[
(6.1) \quad x^M_\ell(v) = -x^m_\ell(v) = \left(-\ell + \ell (1 + 2\ell v^2)^{1/2}\right)^{1/2}.
\]

By the definition of the perturbed quartic potentials in (3.5), we have that

\[
|V(x) - V_\ell(x)| = O(|x|^5) .
\]

Recalling that \( V(x^M(v)) = v^2/2 \) and combining item (ii) with the previous identity, we get

\[
(6.2) \quad \left| V_\ell(x^M(v)) - \frac{1}{2} v^2 \right| = O(|v|^5) .
\]

Putting together (6.1) and (6.2), and using the Taylor series, we obtain

\[
|x^M(v) - x^M_\ell(v)| = O(|v|^5) .
\]

Similarly, we get

\[
|x^m(v) - x^m_\ell(v)| = O(|v|^5) ,
\]

as required. \( \square \)

**Lemma 6.3.** Let \( V_\ell(x) \) be a quartic potential. The period map \( T_\ell: (-v_0, v_0) \to \mathbb{R} \) satisfies the following estimate:

\[
\left| T_\ell(v) - \left( 2\pi - \frac{3\pi}{4} \ell v^2 \right) \right| = O(v^4) ,
\]

where \( \ell \) is equal to the sign of \( V^{(4)}(0) \).

**Proof.** By Lemma 6.2, we have that

\[
T_\ell(v) = \sqrt{2} \int_{x^M_\ell(v)}^{x^M(v)} \left( \frac{1}{2} v^2 - V_\ell(x) \right)^{-1/2} \, dx
\]

\[
= 2\sqrt{2} \int_0^{x^M(v)} \left( \frac{1}{2} v^2 - V_\ell(x) \right)^{-1/2} \, dx.
\]

By the change of coordinates \( z = x/x^M_\ell(v) \), we obtain that

\[
T_\ell(v) = 2\sqrt{2} x^M_\ell(v) \int_0^1 \left( \frac{1}{2} v^2 - V(x^M_\ell(v)z) \right)^{-1/2} \, dz.
\]

Hence, we get

\[
T_\ell(0) = 2\pi, \quad \frac{dT_\ell}{dv}(0) = 0, \quad \frac{d^2T_\ell}{dv^2}(0) = -\frac{3\pi \ell}{2}, \quad \frac{d^3T_\ell}{dv^3}(0) = 0 .
\]
Therefore, we have that
\[ |T_\ell(v) - \left(2\pi - \frac{3\pi}{4} \ell v^2\right)| = O(v^4). \]

**Proposition 6.4.** Let \( V(x) \) be a a perturbed quartic potential. The period map \( T : (-v_0, v_0) \to \mathbb{R} \) satisfies the following estimate:
\[ |T(v) - \left(2\pi - \frac{3\pi}{4} \ell v^2\right)| = O(v^4). \]

**Proof.** By Lemma 5.3, there is \( v_0 > 0 \) small enough, \( T > 0 \) and \( C > 0 \) with the following properties: for all \( |v| < v_0 \), we have that \( 2T_\ell(v) < T \), and for all \( 0 \leq t \leq T \), we have that
\[ (6.3) \quad |x(v; t) - x_\ell(v; t)| < Cv^4. \]

By Lemma 4.8, there is \( K > 0 \) such that, for all \( |v| < v_0 \), we have that
\[ (6.4) \quad x_\ell(v; 2Cv^3) \geq 2Cv^4 - K|v|^6, \]
\[ (6.5) \quad x_\ell(v; -2Cv^3) \leq -2Cv^4 + K|v|^6. \]

Combining (6.3) and (6.4), we obtain that
\[ x(v; T_\ell(v) + 2Cv^3) \geq x_\ell(v; T_\ell(v) + 2Cv^3) - Cv^4 \geq x_\ell(v; +2Cv^3) - Cv^4 \geq Cv^4 - K|v|^6. \]

(6.6)

Combining (6.3) and (6.5), we get
\[ x(v; T_\ell(v) - 2Cv^3) \leq x_\ell(v; T_\ell(v) - 2Cv^3) + Cv^4 \leq x_\ell(v; -2Cv^3) + Cv^4 \leq -Cv^4 + K|v|^6. \]

(6.7)

Take \( v_1 < v_0 \) such that \( Cv^4 > K|v|^6 \). Combining (6.6) and (6.7), for all \( |v| < v_1 \) we obtain that
\[ x(v; T_\ell(v) - 2Cv^3) < 0 < x(v; T_\ell(v) + 2Cv^3). \]

Therefore, by continuity of \( x(v; t) \), we have that
\[ T_\ell(v) - 2Cv^3 < T(v) < T_\ell(v) + 2Cv^3, \]
that is,
\[ |T(v) - T_\ell(v)| = O(|v|^3). \]

Since by definition \( T(v) \) is a function of \( v^2 \), the cubic term of its Taylor expansion must vanish and thus, combining the identity above with Lemma 6.3, we get
\[ |T(v) - \left(2\pi - \frac{3\pi}{4} \ell v^2\right)| = O(v^4). \]

□

The estimate for the period map provided in the lemma above can be obtained by a variety of different methods \[34,39\].
7. THE SCALING PARAMETER

In this section, we will use the period map \( T(v) \) to determine estimates for the \((n, t)\)-scaling velocity \( \gamma_{n,t} \) with relevance for the determination of the \((n, t)\)-scaling parameter \( \Gamma_{n,t} \) introduced in \((2.5)\). Furthermore, we will compute some estimates for \( x(\gamma_{n,t}; v; t) \) which will be useful for the study of the \( n \)-renormalized trajectories \( x_n(v; t) = (-1)^n \Gamma_{n,t}^{-1} x(\Gamma_{n,t} v; n\pi - \ell t) \).

The \((n, t)\)-scaling velocity \( \gamma_{n,t} \geq 0 \) has the property that there exists \( K_1 > 0 \) such that

\[ x(\gamma_{n,t}; n\pi - \ell t) = 0 \]

for every \( 0 < t < K_1 n \). Furthermore, we have that \( x(\gamma_{n-1,t-\ell\pi}; n\pi - \ell t) = 0 \).

For simplicity of notation, we will denote \( \gamma_{n,t} \) by \( \gamma \) and \( \Gamma_{n,t} \) by \( \Gamma \) for the remainder of the paper.

**Proposition 7.1.** Let \( V(x) \) be a perturbed quartic potential and \( \ell \in \{-1, 1\} \) be equal to the sign of \( V^{(4)}(0) \). There exists \( K_1 > 0 \) such that for every \( n \geq 1 \) and for every \( t \in [0, K_1 n] \) the \((n, t)\)-scaling velocity \( \gamma = \gamma_{n,t} \geq 0 \) given by

\[
T(\gamma) = 2\pi - \frac{2\ell t}{n}
\]

is well defined. Furthermore, for every \( t \in [0, K_1 n] \) and every \( v \in [-1, 1] \) the following estimates are satisfied:

\[
|n/2 T(\gamma v) - (n\pi - \ell t - R_{\ell,v,t})| = O \left( \frac{t^2}{n} \right),
\]

where \( R_{\ell,v,t} = -\ell t \left( 1 - v^2 \right) \). In particular, we have that \( \gamma = O \left( \left( |t|/n \right)^{1/2} \right) \).

**Proof.** By Proposition \((6.4)\) there is \( v_0 = v_0(V) > 0 \) such that, for all \( |v| < v_0 \), the period map \( T(v) \) satisfies

\[
|T(v) - \left( 2\pi - \frac{3\pi}{4} \ell v^2 \right)| = O \left( v^4 \right).
\]

Hence, we have that

\[
|T(\gamma) - \left( 2\pi - \frac{3\pi}{4} \ell \gamma^2 \right)| = O \left( \gamma^4 \right).
\]

By \((7.1)\), we get

\[
T(\gamma) - 2\pi = -\frac{2\ell t}{n}.
\]

Combining \((7.5)\) and \((7.4)\), we obtain

\[
\left| \frac{2t}{n} - \frac{3\pi}{4} \gamma^2 \right| = O \left( \gamma^4 \right).
\]

Therefore, there is \( K_1 = K_1(V) > 0 \) such that \( \gamma = \gamma_{n,t} \) is well defined for every \( 0 < t < K_1 n \), and \( |\gamma^2| < O \left( t/n \right) \). Furthermore,

\[
|\gamma^2 - \frac{8t}{3\pi n}| = O \left( \left( \frac{t}{n} \right)^2 \right).
\]
Now let \( v \in [-1, 1] \). From identity (7.3), we get

\[
T(\gamma) - T(\gamma v) + \frac{3\pi}{4} \ell (\gamma^2 - \gamma^2 v^2) = O \left( \left( \frac{t}{n} \right)^2 \right).
\]

Combining (7.6) with (7.7), we obtain

\[
T(\gamma) - T(\gamma v) + 2\ell t \left( 1 - v^2 \right) = O \left( \left( \frac{t}{n} \right)^2 \right).
\]

Note that (7.1) is equivalent to

\[
n\pi - \ell t = \frac{n}{2} T(\gamma).
\]

Putting together (7.9) and (7.8), we obtain

\[
\left| n\pi - \ell t - \left( \frac{n}{2} T(\gamma v) - \ell t \left( 1 - v^2 \right) \right) \right| = O \left( \frac{t^2}{n^2} \right).
\]

Rearranging terms in the identity above, we obtain (7.3). □

**Remark 7.2.** Note that (7.2) can be restated as

\[
\left| \gamma^2 - \Gamma^2 \right| = O \left( \frac{t^2}{n^2} \right).
\]

**Lemma 7.3.** Let \( V(x) \) be a perturbed quartic potential and \( \ell \in \{-1, 1\} \) be equal to the sign of \( V^{(4)}(0) \). There exists \( K_1 > 0 \) such that for every \( n \geq 1 \), every \( t \in [0, K_1n] \) and every \( v \in [-1, 1] \), the \((n, t)\)-scaling parameter \( \Gamma = \Gamma_{n,t} \) satisfies

\[
\left| n\pi - \ell t - \left( \frac{n}{2} T(\gamma v) - \ell t \left( 1 - v^2 \right) \right) \right| = O \left( \frac{t^2 \gamma^4}{n^2} \right),
\]

where \( R_{\ell,v,t} = -\ell t \left( 1 - v^2 \right) \).

**Proof.** By Proposition 7.1, there exists \( K_1 > 0 \) such that for every \( n \geq 1 \), every \( t \in [0, K_1n] \) and every \( v \in [-1, 1] \), the following estimate holds:

\[
\left| n\pi - \ell t - \left( \frac{n}{2} T(\gamma v) - \ell t \left( 1 - v^2 \right) \right) \right| = O \left( \frac{t^2 \gamma^4}{n^2} \right).
\]

From Proposition 6.4, we obtain that

\[
\left| T(\gamma v) - \left( 2\pi - \frac{3\pi}{4} \ell \gamma^2 v^2 \right) \right| = O \left( \gamma^4 v^4 \right).
\]

Putting together identities (7.10) and (7.12), we get

\[
\left| T(\gamma v) - \left( 2\pi - \frac{3\pi}{4} \ell \gamma^2 v^2 \right) \right| = O \left( \Gamma^4 v^4 \right).
\]

From the previous identity and Proposition 6.4, we obtain

\[
\left| T(\gamma v) - T(\Gamma v) \right| = O \left( \Gamma^4 v^4 \right).
\]

Combining (7.11) and (7.13) and recalling that \( \Gamma < O \left( \left| t/n \right|^{1/2} \right) \), we obtain

\[
\left| n\pi - \ell t - \left( \frac{n}{2} T(\gamma v) - \ell t \left( 1 - v^2 \right) \right) \right| = O \left( \frac{t^2 \gamma^4}{n} \right),
\]

completing the proof. □
Lemma 7.4. Let $V(x)$ be a perturbed quartic potential and $\ell \in \{-1,1\}$ be equal to the sign of $V^{(4)}(0)$. There exists $K_1 > 0$ such that for every $n \geq 1$, every $t \in [0, K_1 n]$ and every $v \in [-1,1]$, we have that

$$x (\Gamma v; n\pi - \ell t) \in \left\{ x (\Gamma v; \hat{t}_{\ell,n}(v,t)) : \hat{t}_{\ell,n}(v,t) \in A_{\ell,n} \right\},$$

where $A_{\ell,n}$ is the set of $C^2$ functions $\hat{t}_{\ell,n} : [-1,1] \times [0, K_1 n] \to \mathbb{R}$ with the property that

$$\left| \hat{t}_{\ell,n}(v,t) - \left( \frac{n}{2} \right) T(\Gamma v) + R_{\ell,v,t} \right| = O \left( \Gamma^4 v^4 n \right).$$

Proof. Let $K_1 > 0$ be as in Proposition 7.1. By Lemma 7.3, we have that

$$\left| n\pi - \ell t - \left( \frac{n}{2} \right) T(\Gamma v) + R_{\ell,v,t} \right| = O \left( \Gamma^4 v^4 n \right).$$

Thus, we get that

$$x (\Gamma v; n\pi - \ell t) \in \left\{ x (\Gamma v; \hat{t}_{\ell,n}(v,t)) : \hat{t}_{\ell,n}(v,t) \in B_{\ell,n} \right\},$$

where $B_{\ell,n}$ is the set of $C^2$ functions $\hat{t}_{\ell,n} : [-1,1] \times [0, K_1 n] \to \mathbb{R}$ with the property that

$$\left| \hat{t}_{\ell,n}(v,t) - \left( \frac{n}{2} \right) T(\Gamma v) + R_{\ell,v,t} \right| = O \left( \Gamma^4 v^4 n \right).$$

Since $T(\Gamma v)$ is the period of $x (\Gamma v; t)$, we get

$$x (\Gamma v; n\pi - \ell t) \in \left\{ x (\Gamma v; \hat{t}_{\ell,n}(v,t)) : \hat{t}_{\ell,n}(v,t) \in A_{\ell,n} \right\},$$

where $[n/2]$ denotes the fractional part of $n/2$ and $A_{\ell,n}$ is as given in the statement. \(\square\)

Lemma 7.5. Let $V(x)$ be a perturbed quartic potential and $\ell \in \{-1,1\}$ be equal to the sign of $V^{(4)}(0)$. There exists $K_1 > 0$ such that for every $n \geq 1$, every $t \in [0, K_1 n]$ and every $v \in [-1,1]$, the following estimate is satisfied:

$$\left| \left[ \frac{n}{2} \right] T(\Gamma v) + R_{\ell,v,t} - \left( 2\pi \left[ \frac{n}{2} \right] + R_{\ell,v,t} \right) \right| = O \left( \Gamma^2 \right).$$

Proof. Consequence of Proposition 6.4 \(\square\)

Lemma 7.6. Let $V(x)$ be a perturbed quartic potential and $\ell \in \{-1,1\}$ be equal to the sign of $V^{(4)}(0)$. There exists $K_2 > 0$ such that for every $n \geq 1$, every $t \in [0, K_2 n^{1/2}]$ and every $v \in [-1,1]$, if $\hat{t}_{\ell,n} : [-1,1] \times [0, K_2 n^{1/2}] \to \mathbb{R}$ is such that

$$\hat{t}_{\ell,n}(v,t) - \left( 2\pi \left[ \frac{n}{2} \right] + R_{\ell,v,t} \right) = O \left( \Gamma^2 + \Gamma^4 n \right),$$

then

$$\left| \cos \left( \hat{t}_{\ell,n}(v,t) \right) - (-1)^n \cos \left( R_{\ell,v,t} \right) \right| = O \left( \Gamma^2 + \Gamma^4 n \right),$$

$$\left| \sin \left( \hat{t}_{\ell,n}(v,t) \right) - (-1)^n \sin \left( R_{\ell,v,t} \right) \right| = O \left( \Gamma^2 + \Gamma^4 n \right).$$

Proof. Recall that $\Gamma = O \left( (t/n)^{1/2} \right)$ and let $C > 0$ be such that

$$\left| \hat{t}_{\ell,n}(v,t) - \left( 2\pi \left[ \frac{n}{2} \right] + R_{\ell,v,t} \right) \right| \leq C \left( \Gamma^2 + \Gamma^4 n \right).$$

Fix $K_2 > 0$ sufficiently small so that $C \left( \Gamma^4 n + \Gamma^2 \right) < 1$ for every $t \in [0, K_2 n^{1/2}]$. The result then follows through the use of the trigonometric formulas. \(\square\)
8. ASYMPTOTIC UNIVERSALITY OF TRAJECTORIES

In this section we combine the results of the previous sections to prove Proposition [3.3] We break its proof into the following three lemmas.

**Lemma 8.1.** There exists $N \geq 1$ such that for every $n \geq N$, every $0 < \epsilon < 1/2$ and every $t \in [0, n^{1/2-\epsilon}]$, the estimate

$$|x_n(v; t) - X_\ell(v; t)| = O\left(n^{-2\epsilon}\right)$$

holds for all $v \in [-1, 1]$.

**Proof.** Let $K_1 > 0$ be as in Proposition [7.1], $K_2 > 0$ be as in Lemma [7.6] and $K = \min\{K_1, K_2\}$. Combining Lemmas [7.3] and [7.5] with Proposition [5.4] we obtain

$$|x(Gv; n\pi - \ell t) - \Gamma v \sin(\tilde{\ell}_v,n(v,t))| = O\left(\Gamma^3 t + \Gamma^4\right),$$

for some $\tilde{\ell}_v,n : [-1, 1] \times [0, Kn] \to \mathbb{R}$ such that

$$|\tilde{\ell}_v,n(v,t) - \left(2\pi \left[\frac{n}{2}\right] + R_{\ell_v,t}\right)| = O\left(\Gamma^2 + \Gamma^4n\right).$$

Combining (8.1) with Lemma [7.6] we obtain that for every $n \geq 1$, every $t \in [0, Kn^{1/2}]$ and every $v \in [-1, 1]$, the following estimate holds:

$$|x(Gv; n\pi - \ell t) - (-1)^n \Gamma v \sin(R_{\ell_v,t})| = O\left(\Gamma^3 + \Gamma^3 + \Gamma^4 + \Gamma^5n\right).$$

Clearly, the last identity is equivalent to

$$\big(-1\big)^n \Gamma^{-1} x(Gv; n\pi - \ell t) - v \sin(\ell t(v^2 - 1)) = O\left(\Gamma^2 + \Gamma^2 t + \Gamma^3 + \Gamma^4n\right).$$

Hence, we obtain that

$$|x_n(v; t) - X_\ell(v; t)| = O\left(\frac{t}{n} + \frac{t^2}{n} + \frac{t^3/2}{n^{3/2}}\right),$$

which yields the estimate in the statement for every $n$ sufficiently large and every $0 < \epsilon < 1/2$ and $t \in [0, n^{1/2-\epsilon}]$. □

**Lemma 8.2.** There exists $N \geq 1$ such that for every $n \geq N$, every $0 < \epsilon < 1/3$ and every fixed $t \in [0, Kn^{1/3-\epsilon}]$, the estimate

$$\left|\frac{\partial x_n}{\partial v}(v; t) - \frac{\partial X_\ell}{\partial v}(v; t)\right| = O\left(n^{-3\epsilon}\right)$$

holds for all $v \in [-1, 1]$.

**Proof.** Let $K_1 > 0$ be as in Proposition [7.1], $K_2 > 0$ be as in Lemma [7.6] and $K = \min\{K_1, K_2\}$. Using Lemma [7.4] we have that

$$x_n(v; t) = (-1)^n \Gamma^{-1} x(Gv; \tilde{\ell}_v,n(v,t))$$

for some $\tilde{\ell}_v,n : [-1, 1] \times [0, Kn] \to \mathbb{R}$ such that

$$|\tilde{\ell}_v,n(v,t) - \left(\left[\frac{n}{2}\right] T(Gv) + R_{\ell_v,t}\right)| = O(\Gamma^4 n).$$

Differentiating the previous identity with respect to $v$, we obtain

(8.2) \[ \frac{\partial x_n}{\partial v}(v; t) = (-1)^n \Gamma^{-1} (X_1 + X_2), \]
where \(X_i, i \in \{1, 2\}\), are given by
\[
X_1 = \Gamma \frac{\partial x}{\partial v} \left( \Gamma_v; t_{\ell,n}(v,t) \right), \\
X_2 = \frac{\partial t_{\ell,n}}{\partial v} (v,t) \frac{\partial x}{\partial t} \left( \Gamma_v; t_{\ell,n}(v,t) \right).
\]

Noting that
\[
(8.3) \left| \frac{\partial t_{\ell,n}}{\partial v} (v,t) - \left[ \frac{n}{2} \right] \Gamma T' (\Gamma_v) + \frac{\partial R_{\ell,v,t}}{\partial v} \right| = O \left( \Gamma^4 v^3 n \right)
\]
and using Proposition 5.4 and Lemma 7.5, we obtain that
\[
(8.4) \left| X_2 - \Gamma v \frac{\partial R_{\ell,v,t}}{\partial v} \cos (s_{\ell,n}(v,t)) \right| = O \left( \Gamma^3 + \Gamma^4 t \right),
\]
where \(s_{\ell,n} : [-1, 1] \times [0, Kn] \rightarrow \mathbb{R}\) is such that
\[
\left| s_{\ell,n}(v,t) - \left( 2\pi \left[ \frac{n}{2} \right] + R_{\ell,v,t} \right) \right| = O \left( \Gamma^2 + \Gamma^4 n \right).
\]

Combining (8.4) and Lemma 7.6 with (8.2), we obtain that for every \(n \geq 1\), every \(t \in [0, Kn^{1/2}]\) and every \(v \in [-1, 1]\), the following estimate holds:
\[
\left| \frac{\partial x_n}{\partial v} (v,t) - \left( \sin (R_{\ell,v,t}) + v \frac{\partial R_{\ell,v,t}}{\partial v} \cos (R_{\ell,v,t}) \right) \right| = O \left( \Gamma^2 + \Gamma^4 t \right).
\]

Using Lemma 7.3 we get
\[
\left| \frac{\partial x_n}{\partial v} (v,t) - \left( \sin \left( \ell t (v^2 - 1) \right) + 2ltv^2 \cos \left( \ell t (v^2 - 1) \right) \right) \right| = O \left( \Gamma^2 + \Gamma^4 t \right).
\]

Therefore, we obtain that
\[
\left| \frac{\partial x_n}{\partial v} (v,t) - \frac{\partial X_{\ell}}{\partial v} (v,t) \right| = O \left( \frac{t}{n} + \frac{t^2}{n} + \frac{t^3}{n^{3/2}} + \frac{t^5/2}{n^{3/2}} \right),
\]
yielding the estimate in the statement for every \(n \in \mathbb{N}\) sufficiently large and every \(0 < \epsilon < 1/3\) and \(t \in [0, n^{1/3-\epsilon}]\).

\[ \square \]

**Lemma 8.3.** There exists \(N \geq 1\) such that for every \(n \geq N\), every \(0 < \epsilon < 1/4\) and every fixed \(t \in [0, Kn^{1/4-\epsilon}]\), the estimate
\[
\left| \frac{\partial^2 x_n}{\partial v^2} (v,t) - \frac{\partial^2 X_{\ell}}{\partial v^2} (v,t) \right| = O \left( n^{-4\epsilon} \right)
\]
holds for all \(v \in [-1, 1]\).

**Proof.** Let \(K_1 > 0\) be as in Proposition 7.1 \(K_2 > 0\) be as in Lemma 7.6 and \(K = \min\{K_1, K_2\}\). Recall from (8.2) that
\[
(8.5) \frac{\partial x_n}{\partial v} (v,t) = (-1)^n \Gamma^{-1} (X_1 + X_2),
\]
where $X_i, i \in \{1, 2\}$, satisfy
\[
X_1 = \Gamma \frac{\partial x}{\partial v} (\Gamma v; \tilde{t}_{\ell,n}(v,t)), \\
X_2 = \frac{\partial \tilde{t}_{\ell,n}(v,t)}{\partial v} \frac{\partial x}{\partial t} (\Gamma v; \tilde{t}_{\ell,n}(v,t))
\]
and $\tilde{t}_{\ell,n} : [-1, 1] \times [0, Kn] \to \mathbb{R}$ is such that
\[
\left| \tilde{t}_{\ell,n}(v,t) - \left[ \frac{n}{2} \right] T (\Gamma v) + R_{\ell,v,t} \right| = O \left( (\Gamma v)^4 n \right).
\]
Differentiating identity (8.5) with respect to $v$, we obtain that
\[
\frac{\partial^2 x_n}{\partial v^2} (v; t) = (-1)^n \Gamma^{-1} (Y_1 + 2Y_2 + Y_3 + Y_4),
\]
where $Y_i, i \in \{1, ..., 4\}$, are given by
\[
Y_1 = \Gamma^2 \frac{\partial^2 x}{\partial v^2} (\Gamma v; \tilde{t}_{\ell,n}(v,t)), \\
Y_2 = \Gamma \frac{\partial \tilde{t}_{\ell,n}(v,t)}{\partial v} \frac{\partial^2 x}{\partial t \partial v} (\Gamma v; \tilde{t}_{\ell,n}(v,t)), \\
Y_3 = \frac{\partial^2 \tilde{t}_{\ell,n}(v,t)}{\partial v^2} \frac{\partial x}{\partial t} (\Gamma v; \tilde{t}_{\ell,n}(v,t)), \\
Y_4 = \left( \frac{\partial \tilde{t}_{\ell,n}(v,t)}{\partial v} \right)^2 \frac{\partial^2 x}{\partial t^2} (\Gamma v; \tilde{t}_{\ell,n}(v,t)).
\]
Recalling (8.3), noting that
\[
\left| \frac{\partial^2 \tilde{t}_{\ell,n}}{\partial v^2} (v,t) - \left[ \frac{n}{2} \right] \Gamma^2 T'' (\Gamma v) + \frac{\partial^2 R_{\ell,v,t}}{\partial v^2} (\Gamma v; \tilde{t}_{\ell,n}(v,t)) \right| = O \left( (\Gamma v)^2 n \right),
\]
and using Proposition 5.4 and Lemma 7.5, we obtain that
\[
Y_1 = O(\Gamma^3 t + \Gamma^4), \\
Y_2 - \Gamma \frac{\partial R_{\ell,v,t}}{\partial v} \cos (\tilde{s}_{\ell,n}(v,t)) = O \left( \Gamma^3 + \Gamma^3 t^2 + \Gamma^4 t + \Gamma^5 n \right), \\
Y_3 - \Gamma v \frac{\partial^2 R_{\ell,v,t}}{\partial v^2} \cos (\tilde{s}_{\ell,n}(v,t)) = O \left( \Gamma^3 + \Gamma^3 t^2 + \Gamma^4 t + \Gamma^5 n \right), \\
Y_4 + \Gamma v \left( \frac{\partial R_{\ell,v,t}}{\partial v} \right)^2 \sin (\tilde{s}_{\ell,n}(v,t)) = O \left( \Gamma^3 t + \Gamma^3 t^3 + \Gamma^4 t^2 + \Gamma^5 + \Gamma^5 nt + \Gamma^6 n + \Gamma^9 n^2 \right),
\]
where $\tilde{s}_{\ell,n} : [-1, 1] \times [0, Kn] \to \mathbb{R}$ is such that
\[
\left| \tilde{s}_{\ell,n}(v,t) - \left( 2\pi \left[ \frac{n}{2} \right] + R_{\ell,v,t} \right) \right| = O \left( \Gamma^2 + \Gamma^4 n \right).
\]
Combining (8.8) and Lemma 7.6 with (8.6), we obtain that for every $n \geq 1$, every $t \in [0, Kn^{1/2}]$ and every $v \in [-1, 1]$, the following estimate holds:
\[
\left| \frac{\partial^2 x_n}{\partial v^2} (v; t) - \left( 2 \frac{\partial R_{\ell,v,t}}{\partial v} + \frac{\partial^2 R_{\ell,v,t}}{\partial v^2} \right) \cos (R_{\ell,v,t}) - v \left( \frac{\partial R_{\ell,v,t}}{\partial v} \right)^2 \sin (R_{\ell,v,t}) \right| = O \left( \Gamma^2 + \Gamma^2 t + \Gamma^2 t^2 + \Gamma^2 t^3 + \Gamma^3 + \Gamma^3 t + \Gamma^3 t^2 + \Gamma^4 n + \Gamma^4 nt + \Gamma^4 nt^2 + \Gamma^5 n \right).
\]
Using Lemma 7.3, we get
\[
\left| \frac{\partial^2 x_n}{\partial v^2}(v; t) - \left( 6t v \cos \left( t (v^2 - 1) \right) - 4t^2 v^3 \sin \left( t (v^2 - 1) \right) \right) \right| = O \left( \Gamma^2 + \Gamma^2 t + \Gamma^2 t^2 + \Gamma^3 + \Gamma^3 t + \Gamma^4 t^2 + \Gamma^4 n + \Gamma^4 n t + \Gamma^5 n \right).
\]
Therefore, we obtain that
\[
\left| \frac{\partial^2 x_n}{\partial v^2}(v; t) - \left( \frac{\partial^2 X_t}{\partial v^2}(v; t) \right) \right| = O \left( \frac{t}{n} + \frac{t^2}{n} + \frac{t^3}{n} + \frac{t^4}{n^3/2} + \frac{t^{5/2}}{n^{3/2}} + \frac{t^{7/2}}{n^{3/2}} \right),
\]
yielding the estimate in the statement for every \( n \in \mathbb{N} \) sufficiently large and every \( 0 < \epsilon < 1/4 \) and \( t \in [0, n^{1/4-\epsilon}] \).

9. Application to semiclassical physics

Focal decomposition is in fact a first step towards semiclassical quantization. This was already recognized in the semiclassical calculation of partition functions for quantum mechanical systems, where the need to consider a varying number of classical solutions in different temperature regimes became evident [18].

We shall illustrate this for two physical quantities which can be expressed in terms of Feynman path integrals: in quantum mechanics, we shall consider the 1-dimensional propagator between an initial point \( x_1 \) at time \( t_1 \) and a final point \( x_2 \) at time \( t_2 \); in quantum statistical mechanics, we shall study the 1-dimensional thermal density matrix element in position representation. For further details on semiclassical quantization, see the textbook [31] by Feynman and Hibbs or the review paper [26] by DeWitt-Morette.

In the quantum mechanical case, the propagator is just the time evolution operator \( \exp[-iH(t_2 - t_1)/\hbar] \) computed between position eigenstates \( |x_1\rangle \) and \( |x_2\rangle \):
\[
G(x_1, t_1; x_2, t_2) = \langle x_2 | \exp[-iH(t_2 - t_1)/\hbar] | x_1 \rangle,
\]
where \( H \) is the Hamiltonian
\[
H = \frac{p^2}{2m} + V(x).
\]
The propagator may be written as a path integral
\[
G(x_1, t_1; x_2, t_2) = \int_{z(t_1) = x_1}^{z(t_2) = x_2} [Dz(t)] \exp \left( \frac{iS[z]}{\hbar} \right),
\]
where \( S \) is the classical action of the mechanical problem
\[
S[z] = \int_{t_1}^{t_2} \left[ \frac{1}{2} m z'^2 - V(z) \right] dt.
\]
The integral over \( z \) in (9.1) stands for a sum over all trajectories that connect \( x_1 \) at \( t_1 \) to \( x_2 \) at \( t_2 \), an object whose mathematical characterization has led to much investigation over the years.

We shall only be interested in the leading semiclassical expression for the propagator, which can be formally derived from (9.1). The resulting expression is
\[
G_{sc}(x_1, t_1; x_2, t_2) = \sum_{n=1}^{i} \left\{ \det[S_n^{(2)}] \right\}^{-1/2} \exp \left( \frac{iS_n}{\hbar} \right).
\]
The \( i \) in the upper limit is the same as in the focal decomposition, and indicates that the approximation is restricted to all classical trajectories, i.e. all solutions of
the classical equation of motion which satisfy the boundary conditions \(z(t_1) = x_1\) and \(z(t_2) = x_2\). The \(S_n\) in the exponential stands for the value of the action of the \(n\)th classical trajectory, whereas \(S^{(2)}_n\) denotes the second functional derivative of the action with respect to \(z(t)\), computed at the \(n\)th classical trajectory. The inverse square root of the determinant of that operator, the so-called van Vleck determinant, accounts for the first quantum corrections in a semiclassical expansion.

In the case of quantum statistical mechanics, the thermal density matrix element is just the Boltzmann operator \(\exp[-\beta H]\) computed between position eigenstates \(|x_1\rangle\) and \(|x_2\rangle\), i.e. \(\rho(x_2, x_1) = \langle x_2 | \exp[-\beta H | x_1 \rangle\). The path integral for this quantity is given by

\[
\rho(x_2, x_1) = \int_{z(0) = x_1}^{z(\beta \hbar) = x_2} [Dz(\tau)] \exp \left(-\frac{S[z]}{\hbar} \right),
\]

with the so-called euclidean action defined as

\[
S[z] = \int_0^{\beta \hbar} \left[ \frac{1}{2} m \dot{z}^2 + V(z) \right] d\tau.
\]

There is a crucial difference between this formula and that for the classical action of the mechanical problem: the sign in front of the potential. In fact, in the path integral formulation of quantum statistical mechanics one is led to investigate the mechanical problem defined by minus the potential. With this in mind, one may proceed along the same lines as in quantum mechanics to obtain a semiclassical approximation to the thermal density matrix element. It is given by

\[
\rho_{sc}(x_2, x_1) = \sum_{n=1}^{i'} \left\{ \det[S^{(2)}_n] \right\}^{-1/2} \exp \left(-\frac{S_n}{\hbar} \right).
\]

Another important difference with respect to quantum mechanics is that the sum runs only over those solutions of the (euclidean) equation of motion satisfying the boundary conditions which are local minima of the euclidean action \(S[z]\), whereas in quantum mechanics all solutions, i.e. any extremum, must be taken into account. That is why we use \(i'\) as an upper limit of the sum, defined as the number of extrema that are minima. Clearly, \(i' \leq i\). In reality, we need a refinement of the focal decomposition to tell us which of the solutions are local minima.

Either in quantum mechanics or in quantum statistical mechanics, the semiclassical approximation has to sum over all, or part of, the classical paths satisfying fixed point boundary conditions. Given the pairs \((x_1, t_1)\) and \((x_2, t_2)\), or \((x_1, 0)\) and \((x_2, \beta \hbar)\), the number and type of classical trajectories are the very ingredients which lead to a focal decomposition. It should, therefore, be no surprise that the focal decomposition can be viewed as the starting point for a semiclassical calculation.

As for the renormalization procedure, it was introduced to study the behaviour of classical trajectories for very short space and very long time separations of the fixed endpoints. It maps those trajectories into \(n\)-renormalized ones, whose time separations are shifted by \(n\) half-periods, and whose space separations are scaled up to values of order one. As will be shown in the sequel, this procedure converges to an asymptotic universal family of trajectories that have a well-defined and simple functional form, and which define an asymptotic universal focal decomposition self-similar to the original one.
The natural question to pose is whether the combination of focal decomposition and renormalization can be used to calculate semiclassical expansions for propagators in the short space, long time separation of the endpoints, or analogously, for thermal density matrices for short space separation and low temperatures (long euclidean time $\beta \hbar$ is equivalent to low temperatures $T = 1/(kB\beta)$) by using the simple asymptotic forms alluded to in the previous paragraph.

The conjecture to be investigated in a forthcoming article is that this can be done in a relatively simple way, thanks to the simple form of the asymptotes. This will bypass a much more difficult (if not impossible) calculation involving Jacobi's elliptic functions. Should our expectation be realized, we would obtain semiclassical estimates for both propagators and thermal density matrices in the short space/long time or short space/low temperature limits. Expressing those quantities in terms of energy eigenfunctions and energy eigenvalues, we have

$$G(x_1,t_1;x_2,t_2) = \sum_{m=1}^{\infty} \psi^*_m(x_2)\psi_m(x_1) \exp\left(-iE_m(t_2-t_1)/\hbar\right)$$

in the case of quantum mechanics, or

$$(9.2) \quad \rho(x_2,x_1) = \sum_{m=1}^{\infty} \psi^*_m(x_2)\psi_m(x_1) \exp(-\beta E_m)$$

in the case of quantum statistical mechanics.

Clearly, in the latter case, if we take $\beta$ large, only the lowest energy $E_0$ (the ground state) will contribute. Furthermore, the points $x_1$ and $x_2$ are to be taken in the limit of short space separation. If we choose one of them to be the origin, then (9.2) becomes

$$\rho(x_2,x_1) \approx \rho(0,0) \approx |\psi_0|^2 \exp(-\beta E_0).$$

This means that the combination of focal decomposition and renormalization may lead us to a direct estimate of the ground state energy for a quantum mechanical system from the asymptotic forms obtained in this article.

10. Conclusions

We have studied the dynamics of a family of mechanical systems that includes the pendulum at small neighbourhoods of an elliptic equilibrium and characterized such dynamical behaviour through a renormalization scheme. We have proved that the asymptotic limit of the renormalization scheme introduced in this paper is universal: it is the same for all the elements in the considered class of mechanical systems. As a consequence we have obtained a universal asymptotic focal decomposition for this family of mechanical systems.

We believe that the existence of a universal asymptotic focal decomposition might be useful not only in the theory of boundary value problems of ordinary differential equations but also in several distinct fields of the physical sciences such as quantum statistical mechanics, optics, general relativity and even tsunami formation. Our belief in such applications is based on the relevance that focal decomposition may have on the study of caustic formation by focusing wavefronts of such significance to those fields.
ACKNOWLEDGMENTS

The authors thank Michael Berry, Edson de Faria, Welington de Melo and Robert MacKay for helpful discussions and we thank the anonymous referee for his useful comments. The authors also thank the financial support of LIAAD–INESC TEC through program PEst, USP-UP project, Faculty of Sciences, University of Porto, Calouste Gulbenkian Foundation, FEDER and COMPETE Programmes, and Fundação para a Ciência e a Tecnologia (FCT) through the project “Dynamics and Applications” (PTDC/MAT/121107/2010). The authors thank the financial support of CEMAPRE and Fundação para a Ciência e a Tecnologia (FCT) through the Program POCI 2010 and the project “Randomness in Deterministic Dynamical Systems and Applications” (PTDC/MAT/105448/2008). The authors also thank CNPq, FAPERJ, and FUJB for financial support. Part of this research was done during visits by the authors to IMPA (Brazil), The University of Warwick (United Kingdom), IHES (France), Université Toulouse III - Paul Sabatier (France), CUNY (USA), SUNY (USA) and MSRI (USA), and we thank them for their hospitality. The third author’s research was supported by FCT - Fundação para a Ciência e Tecnologia through the grant with reference SFRH/BPD/27151/2006 and the program “Ciência 2007”.

REFERENCES

[1] Ralph Abraham and Jerrold E. Marsden, Foundations of mechanics, Benjamin/Cummings Publishing Co. Inc. Advanced Book Program, Reading, Mass., 1978. Second edition, revised and enlarged; With the assistance of Tudor Rațiu and Richard Cushman. MR515141 (81e:58025)

[2] M. Abramowitz and I. A. Stegun, Handbook of mathematical functions, Dover Publications (New York), 1965.

[3] Antonio Ambrosetti and Paul H. Rabinowitz, Dual variational methods in critical point theory and applications, J. Functional Analysis 14 (1973), 349–381. MR0370183 (51 #6412)

[4] V. I. Arnol’d, Mathematical methods of classical mechanics, 2nd ed., Graduate Texts in Mathematics, vol. 60, Springer-Verlag, New York, 1989. Translated from the Russian by K. Vogtmann and A. Weinstein. MR997295 (90c:58046)

[5] M. V. Berry, Tsunami asymptotics, New J. Phys. 7 (2005), 129.

[6] M. V. Berry, Focused tsunami waves, Proc. R. Soc. Lond. Ser. A Math. Phys. Eng. Sci. 463 (2007), no. 2087, 3055–3071, DOI 10.1080/rspa.2007.0051. MR2360189 (2009b:86002)

[7] M. V. Berry and C. Upstill, Catastrophe optics: morphologies of caustics and their diffraction patterns, Prog. Optics 18 (1980), 257–346.

[8] Z. Bishnani and R. S. Mackay, Safety criteria for aperiodically forced systems, Dyn. Syst. 18 (2003), no. 2, 107–129, DOI 10.1080/146893603003087. MR1994783 (2004f:37037)

[9] S. V. Bolotin and R. S. Mackay, Isochronous potentials, Localization and energy transfer in nonlinear systems (L. Vazquez, R. S. MacKay, and M. P. Zorzano, eds.), World Sci, 2003, pp. 217–224.

[10] F. Bowman, Introduction to elliptic functions with applications, English Universities Press, Ltd., London, 1953. MR0058760 (15,420d)

[11] C. G. Callan, Broken scale invariance in scalar field theory, Phys. Rev. D 2 (1970), 1541–1547.

[12] Francesco Calogero, Isochronous systems, Oxford University Press, Oxford, 2008. MR2383111 (2009b:37114)

[13] C. A. A. de Carvalho, R. M. Cavalcanti, E. S. Fraga, and S. E. Jorás, Semiclassical series at finite temperature, Ann. Physics 273 (1999), no. 1, 146–170, DOI 10.1006/aphy.1998.5900. MR1690497 (2000a:81040)

[14] C. A. A. de Carvalho, R. M. Cavalcanti, E. S. Fraga, and S. E. Jorás, Improved semiclassical density matrix: Taming caustics, Physical Review E 65 (2002), no. 5, 56112–56221.
[15] Elise E. Cawley, *The Teichmüller space of an Anosov diffeomorphism of T^2*, Invent. Math. **112** (1993), no. 2, 351–376, DOI 10.1007/BF01232439. MR1213107 (94c:58141)

[16] Lothar Collatz, *Differential equations*, A Wiley-Interscience Publication, John Wiley & Sons Ltd., Chichester, 1986. An introduction with applications; Translated from the German by E. R. Dawson. MR858128 (87j:34001)

[17] P. Coullet and C. Tresser, *Itérations d'endomorphismes et groupe de renormalisation*, C. R. Acad. Sci. Paris Sér. A-B **287** (1978), no. 7, A577–A580.

[18] C. A. A. de Carvalho, Mauricio M. Peixoto, Diogo Pinheiro, and Alberto A. Pinto, *Renormalization and focal decomposition*, J. Difference Equ. Appl. **17** (2011), no. 7, 1019–1029, DOI 10.1080/10236191003685916. MR2822433 (2012g:37099)

[19] Carlos A. A. de Carvalho, Mauricio M. Peixoto, Diogo Pinheiro, and Alberto A. Pinto, *Renormalization and focal decomposition*, Dynamics, games and science. II, Springer Proc. Math., vol. 2, Springer, Heidelberg, 2011, pp. 25–40, DOI 10.1007/978-3-642-14788-3_2. MR2883266 (2012j:34036)

[20] W. de Melo, *Rigidity and renormalization in one-dimensional dynamical systems*, Proceedings of the International Congress of Mathematicians, Vol. II (Berlin, 1998), 1998, pp. 765–778. MR1648124 (99m:58082)

[21] G. F. R. Ellis, B. A. C. C. Bassett, and P. K. S. Dunsby, *Lensing and caustic effects on cosmological distances*, Classical Quantum Gravity **15** (1998), no. 8, 2345–2361, DOI 10.1088/0264-9381/15/8/015. MR1645577 (99k:83064)

[22] Mitchell J. Feigenbaum, *Quantitative universality for a class of nonlinear transformations*, J. Statist. Phys. **19** (1978), no. 1, 25–52. MR0501179 (58 #18601)

[23] Mitchell J. Feigenbaum, *The universal metric properties of nonlinear transformations*, J. Statist. Phys. **21** (1979), no. 6, 669–706, DOI 10.1007/BF01107909. MR55919 (82e:58072)

[24] R. P. Feynman and A. R. Hibbs, *Quantum mechanics and path integrals*, McGraw-Hill, New York, 1965.

[25] H. Friedrich and J. M. Stewart, *Characteristic initial data and wavefront singularities in general relativity*, Proc. Roy. Soc. London Ser. A **385** (1983), no. 1789, 345–371, DOI 10.1098/rspa.1983.0018. MR692204 (84d:83006)

[26] M. Gell-Mann and F. E. Low, *Quantum electrodynamics at small distances*, Physical Rev. (2) **95** (1954), 1300–1312. MR0064652 (16,315e)

[27] Jack K. Hale, *Ordinary differential equations*, 2nd ed., Robert E. Krieger Publishing Co. Inc., Huntington, N.Y., 1980. MR0587488 (82e:34001)

[28] Wolfgang Hasse, Marcus Kriele, and Volker Perlick, *Caustics of wavefronts in general relativity*, Classical Quantum Gravity **13** (1996), no. 5, 1161–1182, DOI 10.1088/0264-9381/13/5/027. MR1390106 (97g:53087)
[36] Yun Ping Jiang, *Asymptotic differentiable structure on Cantor set*, Comm. Math. Phys. **155** (1993), no. 3, 503–509. MR1231640 (94e:58100)

[37] Yunping Jiang, *Renormalization and geometry in one-dimensional and complex dynamics*, Advanced Series in Nonlinear Dynamics, vol. 10, World Scientific Publishing Co. Inc., River Edge, NJ, 1996. MR1442953 (98e:58070)

[38] Yunping Jiang, *Smooth classification of geometrically finite one-dimensional maps*, Trans. Amer. Math. Soc. **348** (1996), no. 6, 2391–2412, DOI 10.1090/S0002-9947-96-01487-0. MR1321579 (96j:58054)

[39] D. W. Jordan and P. Smith, *Nonlinear ordinary differential equations*, 3rd ed., Oxford Texts in Applied and Engineering Mathematics, vol. 2, Oxford University Press, Oxford, 1999. An introduction to dynamical systems. MR1743361 (2000j:34001)

[40] L. P. Kadanoff, *Scaling laws for Ising models near $T_c$*, Physics **2** (1966), 263–272.

[41] Ivan A. K. Kupka and M. M. Peixoto, *On the enumerative geometry of geodesics*, From Topology to Computation: Proceedings of the Smalefest (Berkeley, CA, 1990), Springer, New York, 1993, pp. 243–253. MR1246123 (94i:58153)

[42] Ivan Kupka, Mauricio Peixoto, and Charles Pugh, *Focal stability of Riemann metrics*, J. Reine Angew. Math. **593** (2006), 31–72, DOI 10.1515/CRELLE.2006.029. MR2271396 (2007c:53050)

[43] Oscar E. Lanford III, *Renormalization group methods for critical circle mappings with general rotation number*, VIIIth international congress on mathematical physics (Marseille, 1986), World Sci. Publishing, Singapore, 1987, pp. 532–536. MR0915597

[44] R. S. MacKay, *A renormalisation approach to invariant circles in area-preserving maps*, Phys. D **7** (1983), no. 1-3, 283–300, DOI 10.1016/0167-2789(83)90131-8. Order in chaos (Los Alamos, N.M., 1982). MR719057 (85d:58066)

[45] M. MacNeish, *On determination of a catenary with given directrix and passing through two given points*, Ann. of Math. **7** (1906), 65–80.

[46] Curtis T. McMullen, *Complex dynamics and renormalization*, Annals of Mathematics Studies, vol. 135, Princeton University Press, Princeton, NJ, 1994. MR1312365 (96b:58097)

[47] Stellan Östlund, David Rand, James Sethna, and Eric Siggia, *Universal properties of the transition from quasiperiodicity to chaos in dissipative systems*, Phys. D **8** (1983), no. 3, 303–342, DOI 10.1016/0167-2789(83)90229-4. MR719057 (85d:58066)

[48] Jerrold E. Marsden and Tudor S. Ratiu, *Introduction to mechanics and symmetry*, 2nd ed., Texts in Applied Mathematics, vol. 17, Springer-Verlag, New York, 1999. A basic exposition of classical mechanical systems. MR1723696 (2000i:70002)

[49] Marco Martens, *The periodic points of renormalization*, Ann. of Math. (2) **147** (1998), no. 3, 543–584, DOI 10.2307/120959. MR1637651 (99f:58057)

[50] Curtis T. McMullen, *Complex dynamics and renormalization*, Annals of Mathematics Studies, vol. 135, Princeton University Press, Princeton, NJ, 1994. MR1312365 (96b:58097)

[51] Stellan Östlund, David Rand, James Sethna, and Eric Siggia, *Universal properties of the transition from quasiperiodicity to chaos in dissipative systems*, Phys. D **8** (1983), no. 3, 303–342, DOI 10.1016/0167-2789(83)90229-4. MR719057 (85d:58066)

[52] M. M. Peixoto, *On a generic theory of end point boundary value problems*, An. Acad. Brasil. Ci. **41** (1969), 1–6. MR0252747 (40 #5964)

[53] M. M. Peixoto, *On end-point boundary value problems*, J. Differential Equations **44** (1982), no. 2, 273–280, DOI 10.1016/0022-0396(82)90017-1. Special issue dedicated to J. P. LaSalle. MR657782 (83h:34016)

[54] M. M. Peixoto, *Sigma décomposition et arithmétique de quelques formes quadratiques définies positives*, R. Thom Festschrift volume: Passion des Formes (M. Porte, ed.), ENS Editions (Paris), 1994, pp. 455–479.

[55] M. Peixoto, *Geometry, topology and physics*, Walter de Gruyter & Co., Berlin, 1997. Edited by Boris N. Apanasov, Steven B. Bradlow, Waldyr A. Rodrigues, Jr. and Karen K. Uhlenbeck. MR1605264 (98k:00014)
AN ASYMPTOTIC UNIVERSAL FOCAL DECOMPOSITION

[58] M. M. Peixoto and A. R. da Silva, Focal decomposition and some results of S. Bernstein on the 2-point boundary value problem, J. London Math. Soc. (2) 60 (1999), no. 2, 517–547, DOI 10.1112/S0024610799007929. MR1718712 (2000i:34040)

[59] M. M. Peixoto and R. Thom, Le point de vue énumératif dans les problèmes aux limites pour les équations différentielles ordinaires, C. R. Acad. Sci., Paris, Sér. I 303 (1986), 629–633; erratum, 307 (1988) 197–198; II, 303 (1986) 693–698.

[60] Alberto A. Pinto, David A. Rand, and Flávio Ferreira, Fine structures of hyperbolic diffeomorphisms, Springer Monographs in Mathematics, Springer-Verlag, Berlin, 2009. With a preface by Jacob Palis and Enrique R. Pujals. MR2464147 (2010e:37036)

[61] E. C. G. Stueckelberg and A. Petermann, La normalisation des constantes dans la théorie des quanta, Helv. Phys. Acta 26 (1953), 499–520 (French). MR0087504 (19,363d)

[62] Dennis Sullivan, Bounds, quadratic differentials, and renormalization conjectures, (Providence, RI, 1988), Amer. Math. Soc., Providence, RI, 1992, pp. 417–466. MR1184622

[63] Dennis Sullivan, Linking the universalities of Milnor-Thurston, Feigenbaum and Ahlfors-Bers, Topological methods in modern mathematics (Stony Brook, NY, 1991) (Houston, Texas) (L. R. Goldberg and A. V. Phillips, eds.), Publish or Perish, 1993, pp. 543–563.

[64] K. Symanzik, Small distance behaviour in field theory and power counting, Comm. Math. Phys. 18 (1970), no. 3, 227–246, DOI 10.1007/BF01649434. MR1552571

[65] I. Todhunter, A history of the calculus of variations, Chelsea Publishing (New York), 1861 (reprint).

[66] J. J. P. Veerman, Mauricio M. Peixoto, André C. Rocha, and Scott Sutherland, On Brillouin zones, Comm. Math. Phys. 212 (2000), no. 3, 725–744. MR1779166 (2002a:52021)

[67] S. Weinberg, New approach to the renormalization group, Phys. Rev. D 8 (1973), 3497–3509.

[68] E. T. Whittaker and G. N. Watson, A course of modern analysis, Cambridge Mathematical Library, Cambridge University Press, Cambridge, 1996. An introduction to the general theory of infinite processes and of analytic functions; with an account of the principal transcendental functions; Reprint of the fourth (1927) edition. MR1424469 (97k:01072)

[69] Kenneth G. Wilson, The renormalization group: critical phenomena and the Kondo problem, Rev. Modern Phys. 47 (1975), no. 4, 773–840. MR0438986 (55 #11887)

[70] Michael Yampolsky, Complex bounds for renormalization of critical circle maps, Ergodic Theory Dynam. Systems 19 (1999), no. 1, 227–257, DOI 10.1017/S0143385799120947. MR1677153 (2000d:37003)

[71] Michael Yampolsky, Renormalization horseshoe for critical circle maps, Comm. Math. Phys. 240 (2003), no. 1-2, 75–96, DOI 10.1007/s00220-003-0891-8. MR2004980 (2004h:37055)

Instituto de Física, Universidade Federal do Rio de Janeiro, Rio de Janeiro, Brazil
E-mail address: aragao@if.ufrj.br

Instituto de Matemática Pura e Aplicada, Rio de Janeiro, Brazil
E-mail address: peixoto@impa.br

Department of Mathematics, Brooklyn College of the City University of New York, Brooklyn, New York 11210
E-mail address: dpinheiro@brooklyn.cuny.edu

LIAAD - INESC TEC and Department of Mathematics, Faculty of Science, University of Porto, Portugal
E-mail address: aapinto@fc.up.pt