ASYMPTOTIC EXPANSION OF THE HOMOCLINIC SPLITTING MATRIX FOR THE RAPIDLY, QUASIPERIODICALLY, FORCED PENDULUM.

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ABSTRACT. We study a Hamiltonian describing a pendulum coupled with several anisochronous oscillators, devising an asymptotic expansion for the splitting (matrix) associated with a homoclinic point. This expansion consists of contributions that are manifestly exponentially small in the limit of vanishing hyperbolicity, by a shift-of-contour argument. Hence, we infer a similar upper bound on the splitting itself.

1. Main Concepts and Results

1.1. **Background and history.** The study of "separatrix splitting" dates back to Poincaré's classic *Les Méthodes Nouvelles de la Mécanique Céleste* [Poi93].

Starting with Kolmogorov's 1954 note [Kol54], it was proved in a series of papers over a period of twenty years that quasiperiodic motions (invariant tori) are typical for nearly integrable Hamiltonians [Mos62, Arn63, Mos66a, Mos66b, Mos67], and that motions which become quasiperiodic asymptotically in time (stable/unstable manifolds) are stable under small perturbations [Mos67, Gra74].

Arnold [Arn64] described a mechanism how a chain of such "whiskered" tori could provide a way of escape for special trajectories, resulting in instability in the system. (A trajectory would typically lie on a torus and therefore stay eternally within a bounded region in phase space.) The latter is often called Arnold mechanism and the general idea of instability goes by the name Arnold diffusion. It is conjectured in [AA68] that Arnold diffusion due to Arnold mechanism is present quite generically, for instance in the three body problem.

Arnold mechanism is based on Poincaré's concept of biasymptotic solutions, discussed in the last chapter of [Poi93], that are formed at intersections of whiskers of tori. Following such intersections, a trajectory can "diffuse" in a finite time from a neighbourhood of one torus to a neighbourhood of another, and so on.

Chirikov's report [Chi79] is a very nice physical account on Arnold diffusion, while Lochak's compendium [Loc99] discusses more recent developments in a readable fashion and is a good point to start learning about diffusion. Gelfreich's introduction [Gel01] to splitting of separatrices is excellent, and we recommend it to anyone intending to study the topic. From there one should advance to [GL01], which covers more topics with more details. The extensive memoir [LMS03] by Lochak, Marco, and Sauzin is written from the geometric point of view. It has a historical flavor, making it interesting and accessible to virtually anyone.

For the separatrix splitting of the *periodically* forced pendulum, see [HMS88,SMH91,DS92, EKS93]. In the case of the standard map, an asymptotic expression of the splitting has been obtained in [Gel99], as the culmination of a series of works starting with Lazutkin's [Laz03].

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For the quasiperiodically forced pendulum, several studies exist. We mention [DGJS97a, DGJS97b, DJSG99, DG03]. Especially [Tre97, LMS03, RW00, Sau01, RT04] emphasize geometrical aspects of separatrix splitting in the quasiperiodic setting. A lower bound on the splitting has been obtained in [GGM99b], where the pendulum is coupled to two rotators—one speeding up and one slowing down. A generalization to the case of several slow rotators can be found in [Pro03]. These results extend [CG94, Gal94, Gen95a]. See also [GGM99a].

1.2. **The model.** We consider the Hamiltonian

$$\mathcal{H}(\phi, \psi, I, A) = \frac{1}{2}I^2 + g^2 \cos \phi + \frac{1}{2}A^2 - \lambda f(\phi, \psi)$$
 (1.1)

of a pendulum coupled to d rotators, with $\phi \in \mathbb{S}^1 := \mathbb{R}/2\pi\mathbb{Z}$ and $I \in \mathbb{R}$ the coordinate and momentum of the pendulum, and $\psi \in \mathbb{T}^d := (\mathbb{S}^1)^d$ and $A \in \mathbb{R}^d$ the angles and actions of the rotators, respectively. The perturbation f is assumed to be real-valued and real-analytic in its arguments, and λ is a (small) real number, whereas g > 0. This Hamiltonian is sometimes called the generalized Arnold model or the Thirring model. It is the prototype of a nearly integrable Hamiltonian system close to a simple resonance, as is explained in the introduction of [Gen95b]. A review of applications can be found in [Chi79].

The equations of motion are

$$\dot{\phi} = I, \quad \dot{\psi} = A, \quad \dot{I} = g^2 \sin \phi + \lambda \, \partial_{\phi} f, \quad \dot{A} = \lambda \, \partial_{\psi} f.$$
 (1.2)

For the parameter value $\lambda=0$, which is addressed as the unperturbed case, the pendulum and the rotators decouple. The former then has the separatrix flow $\phi:\mathbb{R}\to\mathbb{S}^1$ given by

$$\phi(t) = \Phi^0(e^{gt})$$
 with $\Phi^0(z) = 4 \arctan z$.

By elementary trigonometry, the odd function Φ^0 possesses the symmetry property

$$\Phi^0(z) = 2\pi - \Phi^0(z^{-1}). \tag{1.3}$$

In the phase space of the pendulum, the separatrix—given by Φ^0 —separates closed trajectories (libration) from open ones (rotation).

On the other hand, $\psi: \mathbb{R} \to \mathbb{T}^d$ is quasiperiodic: the vector $\omega := A(0) \equiv A(t)$ in

$$\psi(t) = \psi(0) + \omega t \pmod{2\pi}$$

is assumed to satisfy for some positive numbers a and ν the Diophantine condition

$$|\omega \cdot q| > a |q|^{-\nu} \quad \text{for} \quad q \in \mathbb{Z}^d \setminus \{0\}.$$
 (1.4)

Thus, at the instability point of the pendulum, the flow possesses the invariant tori

$$\mathcal{T}_0 := \left\{ (\phi, \psi, I, A) = (0, \theta, 0, \omega) \mid \theta \in T^d \right\}$$

indexed by ω , with stable and unstable manifolds (\mathcal{W}_0^s and \mathcal{W}_0^u , respectively) coinciding:

$$\mathcal{W}_0^{s,u} = \left\{ (\phi, \psi, I, A) = \left(\Phi^0(z), \theta, gz \partial_z \Phi^0(z), \omega \right) \mid z \in [-\infty, \infty], \ \theta \in T^d \right\}. \tag{1.5}$$

Remark 1.1. The constant g is the Lyapunov exponent for the unstable fixed point of the pendulum motion; in the limit $s \to -\infty$ two nearby initial angles $\phi(s)$ and $\phi(s+\delta s)$ separate at the exponential rate e^{gs} . As $\phi(t) = \Phi^0(e^{t/g^{-1}})$, this fixes a natural time scale of g^{-1} units, characteristic of the pendulum motion in the unperturbed Hamiltonian system (1.1).

When the perturbation is switched on $(\lambda \neq 0)$, it is known that most of the invariant tori survive and have stable and unstable manifolds—or "whiskers" as Arnold has called them—that may not coincide anymore. We prove a bound on their splitting.

1.3. Main theorem. In [Ste07] we construct the perturbed manifolds in a form similar to (1.5) as graphs of analytic functions over a piece of $[-\infty, \infty] \times \mathbb{T}^d$. To that end, we look for solutions of the form

$$(\phi(t), \psi(t)) = (\Phi(e^{\gamma t}, \omega t), \omega t + \Psi(e^{\gamma t}, \omega t)) = (0, \omega t) + (\Phi, \Psi)(e^{\gamma t}, \omega t)$$

$$(1.6)$$

with quasiperiodic behavior in one of the two limits $t \to \pm \infty$. Note especially that the Lyapunov exponent $\gamma > 0$ depends on λ , with $\gamma|_{\lambda=0} = g$.

Remark 1.2. One should not expect asymptotic quasiperiodicity in both limits $t \to \pm \infty$, as the unstable and stable manifolds, \mathcal{W}^u_{λ} and \mathcal{W}^s_{λ} , are generically expected to depart for nonzero values of the perturbation parameter λ . Therefore, either the past or future asymptotic of a trajectory will evolve so as to ultimately reach the deformed invariant torus \mathcal{T}_{λ} .

Let us denote the total derivative d/dt by ∂_t and the complete angular gradient $(\partial_{\phi}, \partial_{\psi})$ by ∂ for short. Substituting (1.6) into the equations of motion, we get the equation

$$(\omega \cdot \partial_{\theta} + \gamma e^{\gamma t} \partial_{z})^{2} X(e^{\gamma t}, \omega t) = [(g^{2} \sin \Phi, 0) + \lambda \partial_{t} f(X + (0, \theta))](e^{\gamma t}, \omega t)$$

for $X := (\Phi, \Psi)$, where θ stands for the canonical projection $[-\infty, \infty] \times \mathbb{T}^d \to \mathbb{T}^d$. Notice that the partial differential operator

$$\mathcal{L} := \omega \cdot \partial_{\theta} + \gamma z \partial_{z}$$

satisfies the characteristic identity

$$\mathcal{L}F(ze^{\gamma t}, \theta + \omega t) = \partial_t F(ze^{\gamma t}, \theta + \omega t), \tag{1.7}$$

which reflects the time derivative nature of \mathcal{L} . In fact, if T is the "time-reversal map"

$$T(z,\theta) \equiv (z^{-1}, -\theta), \tag{1.8}$$

then, by the chain rule,

$$\mathcal{L}(F \circ T) = -(\mathcal{L}F) \circ T. \tag{1.9}$$

Let us abbreviate

$$\Omega(X) := (g^2 \sin \Phi, 0) + \lambda \widetilde{\Omega}(X) \quad \text{with} \quad \widetilde{\Omega}(X) := \partial f(X + (0, \theta)). \tag{1.10}$$

We have then encoded the equations of motion into the PDE

$$\mathcal{L}^2 X = \Omega(X). \tag{1.11}$$

The action variables, or momenta, trivially follow from the knowledge of $X(z,\theta)$:

$$(I(t), A(t)) = (0, \omega) + Y(e^{\gamma t}, \omega t), \quad Y := \mathcal{L}X.$$

The solutions X provide a parametrization of the deformed tori and their stable and unstable manifolds. As hinted below (1.6), we find two kinds of solutions, $X^u(z,\theta)$ defined for $z \in$ $[-z_0, z_0]$ and $X^s(z, \theta)$ defined for $z \in [-\infty, -z_0^{-1}] \cup [z_0^{-1}, \infty]$. Here, $z_0 > 1$. We will consider λ small in a g-dependent fashion, taking

$$\epsilon := \lambda g^{-2} \tag{1.12}$$

small. Such a choice is needed for studying the limit $g \to \infty$, which corresponds to rapid forcing; see Remark 1.1. The domain we restrict ourselves to is

$$D := \{ (\epsilon, g) \in \mathbb{C} \times \mathbb{R} \mid |\epsilon| < \epsilon_0, \ 0 < g < g_0 \}, \tag{1.13}$$

for some positive values of ϵ_0 and g_0 .

The following theorem is from [Ste07]. It is a version of a classical result, and by no means new; earlier treatments include [Mel63, Mos67, Gra74, Eli94, Gal94, Gen95a, Gen95b].

Theorem 1 (Tori and their whiskers). Let f be real-analytic and even, i.e.,

$$f(\phi, \psi) = f(-\phi, -\psi).$$

Also, suppose ω satisfies the Diophantine condition (1.4), and fix $g_0 > 0$. Then there exist a positive number ϵ_0 and a function $\gamma(\epsilon, g)$ on D, analytic in ϵ with $|\gamma - g| < Cg|\epsilon|$, such that equation (1.11) has a solution X^u which is analytic in ϵ as well as in (z, θ) in a neighbourhood of $[-1, 1] \times \mathbb{T}^d$ and which satisfies $(X^0 := (\Phi^0, 0))$

$$X^{u}(1,0) = (\pi,0), \quad X^{u}(z,\theta) = X^{0}(z) + \mathcal{O}(\epsilon).$$
 (1.14)

Corresponding to the same γ , there exists a solution $X^s(z,\theta) = X^0(z) + \mathcal{O}(\epsilon)$ which is an analytic function of $(z^{-1}, -\theta)$ in a neighbourhood of $[-1, 1] \times \mathbb{T}^d$. The maps

$$W^{s,u}(z,\theta) = (X^{s,u}, Y^{s,u})(z,\theta) + ((0,\theta), (0,\omega)), \quad Y^{s,u} := \mathcal{L}X^{s,u}, \tag{1.15}$$

provide analytic parametrizations of the stable and unstable manifolds $\mathcal{W}_{\lambda}^{s,u}$ of the torus \mathcal{T}_{λ} .

Remark 1.3. The number ϵ_0 above depends on the Diophantine exponent ν and on f. The perturbation $(\phi, \psi) \mapsto f(\phi, \psi)$ is analytic on the compact set $\mathbb{S}^1 \times \mathbb{T}^d$. By Abel's Lemma (multivariate power series converge on polydisks), it extends to an analytic map on a "strip" $|\Im \mathfrak{m} \phi|, |\Im \mathfrak{m} \psi| \leq \eta \ (\eta > 0)$ around $\mathbb{S}^1 \times \mathbb{T}^d$. By Theorem 1, there exists some $0 < \sigma < \eta$ such that each $\theta \mapsto X^{s,u}(\cdot, \theta)$ is analytic on $|\Im \mathfrak{m} \theta| \leq \sigma$.

An important part of Theorem 1 is that the domains of X^u and X^s overlap. Namely, if X solves equation (1.11), then so does $(2\pi,0) - X \circ T$. This is due to (1.9) and the parity of f. Consequently, by time-reversal, the stable and unstable manifolds are related through

$$X^{s} = (2\pi, 0) - X^{u} \circ T. \tag{1.16}$$

In particular, as T(1,0) = (1,0), $X^s(1,0) = X^u(1,0)$. The actions $Y^{s,u} = \mathcal{L}X^{s,u}$ satisfy $Y^s = Y^u \circ T$, yielding $Y^s(1,0) = Y^u(1,0)$. In other words, a homoclinic intersection of the stable and the unstable manifolds $\mathcal{W}^{s,u}_{\lambda}$ occurs at $(z,\theta) = (1,0)$, as their parametrizations (1.15) coincide at this homoclinic point. Since the manifolds $\mathcal{W}^{s,u}_{\lambda}$ are invariant, there in fact exists a homoclinic trajectory on which the parametrizations agree:

$$W^{s}(e^{\gamma t}, \omega t) \equiv W^{u}(e^{\gamma t}, \omega t). \tag{1.17}$$

Coming to the second result of [Ste07], let us expand $X^u = \sum_{\ell=0}^{\infty} \epsilon^{\ell} X^{u,\ell}$. It turns out that the common analyticity domain of each $X^{u,\ell}$ in the z-variable is in fact much larger than the (small) neighbourhood of [-1,1]—the corresponding domain of X^u according to Theorem 1. Namely, it includes the wedgelike region

$$\mathbb{U}_{\tau,\vartheta} := \left\{ |z| \le \tau \right\} \bigcup \left\{ \arg z \in [-\vartheta,\vartheta] \cup [\pi - \vartheta,\pi + \vartheta] \right\} \subset \mathbb{C}$$
 (1.18)

(with some positive τ and ϑ). We repeat parts of the argument in Section 2 for convenience.

Theorem 2 (Analytic continuation). Each order $X^{u,\ell}$ of the solution extends analytically to a common region $\mathbb{U}_{\tau,\vartheta} \times \{|\Im \mathfrak{m}\,\theta| \leq \sigma\}$. Moreover, if $\psi \mapsto f(\cdot,\psi)$ is a trigonometric polynomial of degree N, i.e., N is the smallest integer such that $\hat{f}(\cdot,q) = 0$ whenever |q| > N, then $\theta \mapsto X^{u,\ell}(\cdot,\theta)$ is a trigonometric polynomial of degree ℓN , at most.

Remark 1.4. With η and σ as in Remark 1.3, the numbers τ and ϑ are specified by the following observation: $\Phi^0(z) = 4 \arctan z$ implies that $|\Im \mathfrak{m} \Phi^0(z)| \leq \eta$ in $\mathbb{U}_{\tau,\vartheta}$ with τ and ϑ sufficiently small. By Remark 1.3, $(z,\theta) \mapsto f(\Phi^0(z),\theta)$ is analytic on $\mathbb{U}_{\tau,\vartheta} \times \{|\Im \mathfrak{m} \theta| \leq \sigma\}$, which is the basis of the proof.

In spite of Theorem 2, (a straightforward upper bound on) $X^{u,\ell}$ grows without a limit as $|\Re \mathfrak{e} \, z| \to \infty$, such that there is no reason whatsoever to expect absolute convergence of the series $\sum_{\ell=0}^{\infty} \epsilon^{\ell} X^{u,\ell}$ in an unbounded z-domain with a fixed ϵ . In fact, it is known that the behavior of the unstable manifold gets extremely complicated for large values of z even with innocent looking Hamiltonian systems. Still, it seems to us that the possibility of a uniform analytic extension of the coefficients $X^{u,\ell}$ has not been appreciated in the literature.

Due to (1.16), an analog of Theorem 2 is seen to hold for X^s , with z replaced by z^{-1} .

Theorem 2 allows one (at each order in ϵ) to track trajectories $t \mapsto W^{s,u}(e^{\gamma t}, \theta + \omega t)$ on the invariant manifolds $\mathcal{W}^{s,u}_{\lambda}$ for arbitrarily long times in a uniform complex neighbourhood $|\Im t| \leq g^{-1}\vartheta$ of the real line, for arbitrary $\theta \in \mathbb{T}^d$. The motivation for doing this stems from studying the splitting of the manifolds $\mathcal{W}^{s,u}_{\lambda}$ in the vicinity of the homoclinic trajectory (1.17). The general ideology that, being able to extend "splitting related functions" to a large complex domain yields good estimates, is due to Lazutkin [Laz03], as is emphasized in [LMS03].

In order to study the intersection more closely, we express the actions as functions of the original angle variables $(\phi, \psi) = X^{s,u}(z, \theta) + (0, \theta)$ appearing in the Hamiltonian (1.1). To this end, let $F^{s,u}: (z,\theta) \mapsto (\phi,\psi)$ be the above coordinate transformations, and write $Y^{s,u} = \bar{Y}^{s,u} \circ F^{s,u}$. We set $\partial := (\partial_{\phi}, \partial_{\psi})$ and $D := (\partial_{z}, \partial_{\theta})$. By the chain rule,

$$DY^{s,u} = (\partial \bar{Y}^{s,u} \circ F^{s,u}) DF^{s,u}. \tag{1.19}$$

The invertibility of the matrix $DF^{s,u}$ is a consequence of Theorem 1. At the homoclinic point $(z,\theta)=(1,0)$, equation (1.17) implies that $F^s(e^{\gamma t},\omega t)\equiv F^u(e^{\gamma t},\omega t)$ and $F^{s,u}(1,0)=(\pi,0)$. Casting in the obvious manner $\bar{Y}^{u,s}=(\bar{Y}_{\Phi}^{u,s},\bar{Y}_{\Psi}^{u,s})$, we define the *splitting vector*:

$$\Delta(\phi, \psi) := (A^u - A^s)^T = (\bar{Y}_{\Psi}^u - \bar{Y}_{\Psi}^s)^T (\phi, \psi).$$

Transposition indicates that we consider Δ a column vector. We define the *splitting matrices*:

$$\Upsilon(t) := \partial_{\theta} (Y_{\Psi}^{u} - Y_{\Psi}^{s})^{T} \left(e^{\gamma t}, \omega t \right), \tag{1.20}$$

$$\widetilde{\Upsilon}(t) := \partial_{\psi} (\bar{Y}_{\Psi}^{u} - \bar{Y}_{\Psi}^{s})^{T} (F^{u,s}(e^{\gamma t}, \omega t)). \tag{1.21}$$

Notice that if $(\phi, \psi)(t)$ stands for the moment for the angles on the homoclinic trajectory at time t, with the particular initial condition $(\phi, \psi)(0) = (\pi, 0)$, then $\widetilde{\Upsilon} = \partial_{\psi} \Delta(\phi, \psi)$.

It is nontrivial but straightforward to establish that Υ and Υ differ by a close-to-identity factor, if ϵ is taken to be small proportionally to g:

Theorem 3. Fix $t \in \mathbb{R}$ and define $\tilde{\epsilon} := g^{-1}e^{g|t|}\epsilon$. Then, as $g \to 0$,

$$\Upsilon(t) = \widetilde{\Upsilon}(t) (\mathbb{1} + \mathcal{O}(\widetilde{\epsilon})),$$

if $\tilde{\epsilon}$ is sufficiently small (independently of g and t).

Theorem 3 is our first result. Its proof in Appendix A is based on energy conservation and the fact that the actions are given by scalar potentials. We can now state the main theorem:

Theorem 4 (Homoclinic splitting). Let us assume that $\psi \mapsto f(\cdot, \psi)$ is a trigonometric polynomial in addition to the assumptions of Theorem 1. Then for each $t \in \mathbb{R}$ there exist positive constants C and c, such that the exponentially small upper bound

$$|\Upsilon_{ij}(t)| \le C|\epsilon|e^{-cg^{-1/(\nu+1)}}$$

where ν is the Diophantine exponent, holds, provided ϵg^{-4} is sufficiently small.

Theorem 4 is derived from Proposition 3.12—the centerpiece of this work—stated in Subsection 3.3. Technicalities aside, we believe that Subsection 3.3 can be understood without any further introduction, and we urge the reader to go through it before moving on to Section 2.

Why just study a $d \times d$ submatrix of the full $(d+1) \times (d+1)$ Jacobi matrix of $(I^u - I^s, A^u - A^s)$ when measuring transversality of the homoclinic intersection $\mathcal{W}^s_{\lambda} \cap \mathcal{W}^u_{\lambda}$? Because the latter is singular. As a matter of fact, (A.2) simply means that the splitting vanishes in the direction of the homoclinic trajectory. The author is grateful to Dr Mischa Rudnev for pointing this out. For further motivation, see item R1 of Appendix R and Section 8 in [GGM99b].

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2. Analytic Continuation of the Solution

Here we present the proof of Theorem 2. We drop the superscript u from the notation and single out the uncoupled part X^0 of the complete solution X;

$$X = X^0 + \widetilde{X}$$
 with $\widetilde{X}|_{\epsilon=0} \equiv 0$.

Equation (1.11) and $\mathcal{L}^2 X^0 = (\gamma^2 \sin \Phi^0, 0)$ imply $\mathcal{L}^2 \widetilde{X} = -(\gamma^2 \sin \Phi^0, 0) + \Omega(X^0 + \widetilde{X})$. Thus,

$$K\widetilde{X} = \widetilde{W}(\widetilde{X}), \tag{2.1}$$

when we define the linear operator

$$\mathcal{K} := \begin{pmatrix} L & 0 \\ 0 & \mathcal{L}^2 \end{pmatrix} \quad \text{with} \quad L := \mathcal{L}^2 - \gamma^2 \cos \Phi^0$$
 (2.2)

and the nonlinear operator

$$\widetilde{W}(\widetilde{X}) := (-\gamma^2 \sin \Phi^0 - \gamma^2 (\cos \Phi^0) \widetilde{\Phi}, 0) + \Omega(X^0 + \widetilde{X}). \tag{2.3}$$

Throughout the rest of the work, we shall refer to different parts of the Taylor expansion of a suitable function $h(z, \theta)$ around z = 0 using the notation

$$h_k(\theta) := \frac{\partial_z^k h(0,\theta)}{k!}, \quad h_{\leq k}(z,\theta) := \sum_{i=0}^k z^k h_k(\theta) \quad \text{and} \quad \delta_k h := h - h_{\leq k-1}.$$

By (2.1), $\delta_2 \widetilde{X}$ satisfies

$$\mathcal{K}\delta_2\widetilde{X} = W(\delta_2\widetilde{X}),\tag{2.4}$$

where

$$W(Z) := \delta_2 \left[\widetilde{W}(\widetilde{X}_{\leq 1} + Z) + \begin{pmatrix} \gamma^2 (\cos \Phi^0) \widetilde{\Phi}_{\leq 1} \\ 0 \end{pmatrix} \right]. \tag{2.5}$$

Let us consider analytic functions Z on a small, compact and complex, neighbourhood Π of the set $[-1,1] \times \mathbb{T}^d$. It is a complex Banach space \mathcal{A} , once equipped with the supremum norm, and has the closed subspace

$$\mathcal{A}_1 := \{ Z \in \mathcal{A} \mid Z_{\leq 1} = 0 \}. \tag{2.6}$$

We have shown in [Ste07] that K maps A_1 to itself and has a bounded, $\mathcal{O}(\gamma^{-2})$, inverse.

By virtue of (2.4) and \widetilde{W} 's analyticity, $\delta_2\widetilde{X}$ admits the representation

$$\delta_{2}\widetilde{X} = \mathcal{K}^{-1}\delta_{2} \left[\begin{pmatrix} \gamma^{2} \cos \Phi^{0} & 0 \\ 0 & 0 \end{pmatrix} \widetilde{X}_{\leq 1} + \sum_{k=0}^{\infty} w^{(k)} (\widetilde{X}_{\leq 1})^{\otimes k} \right] + \mathcal{K}^{-1} \sum_{k=1}^{\infty} \left[w^{(k)} (\widetilde{X}_{\leq 1} + \delta_{2}\widetilde{X})^{\otimes k} - w^{(k)} (\widetilde{X}_{\leq 1})^{\otimes k} \right]$$

$$(2.7)$$

on the set Π , taking ϵ small enough, and denoting

$$w^{(k)} := \frac{1}{k!} D^k \widetilde{W}(0) \tag{2.8}$$

as well as a repeated argument of such a symmetric k-linear operator by $(x)^{\otimes k} := (x, \dots, x)$, for the sake of brevity. Observe that we have omitted a δ_2 in front of the square brackets on the second line of (2.7) as redundant.

Equation (2.7) may be viewed as a recursion relation for $\delta_2 \widetilde{X}$. It is crucial that

$$w^{(0)}, w^{(1)} = \mathcal{O}(\epsilon g^2),$$
 (2.9)

when $(\epsilon, g) \in D$; see (1.13). Namely, any given order $\delta_2 \widetilde{X}^\ell$ in the convergent expansion $\delta_2 \widetilde{X} = \sum_{\ell=1}^{\infty} \epsilon^\ell \delta_2 \widetilde{X}^\ell$ is then completely determined by $\widetilde{X}_{\leq 1}$ and the lower orders $\delta_2 \widetilde{X}^\ell$ $(1 \leq l \leq \ell - 1)$ through the right-hand side of (2.7). Moreover, since $\widetilde{X}_{\leq 1} = \mathcal{O}(\epsilon)$, only finitely many terms in the sum over the index k are involved. Together these facts imply that only finitely many recursive steps using (2.7) are needed to completely describe any given order $\delta_2 \widetilde{X}^\ell$ in terms of $\widetilde{X}_{\leq 1}$ alone and that, at each such step, only finitely many terms from the k-sum contribute.

It is important to understand that $\widetilde{X}_{\leq 1}$ is a predetermined function. As we shall see, the recursion procedure will then provide the analytic continuation of each $X^{u,\ell} = \widetilde{X}_{\leq 1}^{\ell} + \delta_2 \widetilde{X}^{\ell}$ $(\ell \geq 1)$ to the large region $\mathbb{U}_{\tau,\vartheta} \times \{|\Im \mathfrak{m} \, \theta| \leq \sigma\}$ of Theorem 2.

2.1. **Tree expansion.** We next give a pictorial representation of the above recursion. It involves tree diagrams similar to those of Gallavotti, et al. (see, e.g., [Gal94, CG94]), with one difference: there will be no resummations nor cancellations, as the expansion in (2.7) contains no resonances and is instead well converging. This so-called tree expansion is needed for bookkeeping and pedagogical purposes; we simply choose to draw a tree instead of spelling out a formula.

Let us first define the auxiliary functions

$$h^{(k)} := \begin{cases} w^{(0)} + \left[\begin{pmatrix} \gamma^2 \cos \Phi^0 & 0 \\ 0 & 0 \end{pmatrix} + w^{(1)} \right] \widetilde{X}_{\leq 1} & \text{if } k = 1, \\ w^{(k)} \left(\widetilde{X}_{\leq 1} \right)^{\otimes k} & \text{if } k = 2, 3, \dots, \end{cases}$$

and make the identifications

$$\longrightarrow (k) := \mathcal{K}^{-1} \delta_2 h^{(k)} \quad \text{and} \quad \longrightarrow (2.10)$$

Furthermore, let

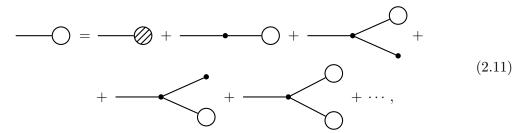
$$\longrightarrow$$
 := $\delta_2 \widetilde{X}$, \longrightarrow := $\widetilde{X}_{\leq 1}$, and \longrightarrow $k \text{ lines}$:= $\mathcal{K}^{-1} w^{(k)}$.

In the diagram representing the k-linear $w^{(k)}$, the k "free" lines to the right of the node stand for the arguments. We say that these lines *enter* the *internal* node, whereas the single line to the left of the node *leaves* it. For instance,

$$= \mathcal{K}^{-1} w^{(3)} \big(\widetilde{X}_{\leq 1}, \, \delta_2 \widetilde{X}, \, \mathcal{K}^{-1} \delta_2 h^{(4)} \big).$$

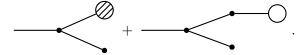
Notice that, as $w^{(k)}$ is symmetric, permuting the lines entering a node does not change the resulting function. We emphasize that all of the functions introduced above are analytic on Π and $|\epsilon| < \epsilon_0$.

In terms of such tree diagrams, or simply trees, equation (2.7) reads



using multilinearity to split the sums $\widetilde{X}_{\leq 1} + \delta_2 \widetilde{X}$ into pieces. Above, the sum after the first tree consists of all trees having one internal node and an arbitrary number of end nodes, at least one of which, however, is a white circle. This rule encodes the fact that on the second line of (2.7) the summation starts from k=1 and that the contributions with only $\widetilde{X}_{\leq 1}$ in the argument (i.e., trees with only black dots as end nodes) are cancelled.

Using (2.7) recursively now amounts to replacing each of the lines with a white-circled end node by the complete expansion of such a tree above. This is to be understood additively, so that replacing one end node, together with the line leaving it, by a sum of two trees results in a sum of two new trees. For example, such a replacement in the third tree on the right-hand side of (2.11) by the first two trees gives the sum



Before proceeding, we introduce a little bit of terminology. The leftmost line in a tree is called the *root line*, whereas the node it leaves (*i.e.*, the uniquely defined leftmost node) is called the *root*. A line leaving a node v and entering a node v' can always be interpreted as the root line of a *subtree*, the maximal tree consisting of lines and nodes in the original tree with v as its root. We call v a (not necessarily unique) successor of v', whereas v' is the unique predecessor of v.

The recursion (2.11) can be repeated on a given tree if it has at least one white circle left. Otherwise, the tree in question must satisfy

- (R1') The tree has only filled circles (\bigcirc) and black dots (\bullet) as its end nodes, together with
- (**R2**') Any internal node has an entering (line that is the root line of a) subtree containing at least one filled circle as an end node.

After all, the recursion can only stop by replacing an existing white circle with a filled one. Continuing ad infinitum yields the expansion

—
$$= \sum (\text{Trees satisfying (R1') and (R2')}) = \sum_{\text{trees } T}' T,$$
 (2.12)

where the prime restricts the summation to trees T satisfying (R1') and (R2'). We point out that each admissible tree appears precisely once in this sum, considering different two trees that can be superposed by a (nontrivial) permutation of subtrees that enter the same node.

The earlier discussion concerning the description of $\delta_2 \widetilde{X}^{\ell}$ in terms of a finite sum involving only $\widetilde{X}_{\leq 1}$ translates to the language of trees in a straightforward fashion. First, the second part of (2.9) and $\widetilde{X}_{\leq 1} = \mathcal{O}(\epsilon)$ amount pictorially to

since \mathcal{K}^{-1} produces a factor of g^{-2} . Second, $w^{(k)} = \mathcal{O}(g^2)$ and the first part of (2.9) yield

Expanding the filled end nodes

$$\longrightarrow \bigotimes = \sum_{k=1}^{\infty} \longrightarrow k ,$$
 (2.13)

according to (2.10), on the right-hand side of (2.12), we get a new version of the latter by replacing the rules (R1') and (R2'), respectively, with

(R1) The tree has only numbered circles ((k)) with arbitrary values of k) and black dots (\bullet) as its end nodes,

and

(R2) Any internal node has an entering (line that is the root line of a) subtree containing at least one numbered circle as an end node.

Let us define the *degree* of a tree as the positive integer

$$\deg T := \#(\underline{\hspace{1cm}}) + \#(\underline{\hspace{1cm}}) + \sum_{k=1}^{\infty} k \#(\underline{\hspace{1cm}})$$
 (2.14)

for any tree T satisfying (R1) and (R2). By #(G) we mean the number of occurrences of the graph G in the tree T. That is, the degree of a tree is the number of its end nodes with suitable weights plus the number of nodes with precisely one entering line. Since a tree has finitely many nodes, its degree is well-defined. Then a rearrangement of the sum arising from (2.13) being inserted into (2.12) yields formally

$$\longrightarrow \sum_{l=1}^{\infty} \sum_{\substack{\text{trees } T \\ \text{deg } T = l}}^{*} T,$$
 (2.15)

where the asterisk reminds us that the rules (R1) and (R2) are being respected.

According to the analysis above, the particular graphs appearing in the definition of deg T are the only possible single-node subgraphs of T proportional to a positive power of ϵ . Since

each tree is analytic in ϵ , writing $(\cdot)^k$ for the kth coefficient of the power series in ϵ , we have

$$T = \sum_{k=\deg T}^{\infty} \epsilon^k T^k = \epsilon^{\deg T} \sum_{k=0}^{\infty} \epsilon^k T^{k+\deg T}.$$

Hence, only trees with degree at most equal to ℓ can contribute to $\delta_2 \widetilde{X}^{\ell}$:

$$\delta_2 \widetilde{X}^{\ell} = \sum_{l=1}^{\ell} \sum_{\deg T=l}^* T^{\ell} = \left(\sum_{\deg T \le \ell}^* T\right)^{\ell} \tag{2.16}$$

or, alternatively,

$$\delta_2 \widetilde{X} = \sum_{\deg T \le \ell}^* T + \mathcal{O}(\epsilon^{\ell+1}) \qquad (\epsilon \to 0)$$
 (2.17)

for each and every $\ell=1,2,\ldots$. The expansion in (2.15) is just a compact way of writing (2.17). We emphasize that the latter can be derived completely rigorously, for each value of ℓ separately, but resorting to the use of formal series allowed us to treat all orders of $\delta_2 \widetilde{X}$ at once. We call the series (2.15) an asymptotic expansion of $\delta_2 \widetilde{X}$; the partial sums $\sum_{\deg T \leq \ell}^* T$ need not converge to $\delta_2 \widetilde{X}$ for any fixed ℓ as $\ell \to \infty$, but for a fixed ℓ the error is bounded by an ℓ -dependent constant times $|\epsilon|^{\ell+1}$ on the mutual domain of analyticity, $|\epsilon| < \epsilon_0$.

Example 2.1. The beginning of the asymptotic expansion (2.17) reads

$$\delta_2 \widetilde{X} = \underbrace{} 1 + \mathcal{O}(\epsilon^2) = \underbrace{} 1 + \underbrace{} 2 + \underbrace{} 1 + \underbrace{\phantom{A$$

2.2. Analyticity domain of trees. As already pointed out, all trees T above are analytic functions of (z, θ, ϵ) on $\Pi \times \{|\epsilon| < \epsilon_0\}$. Due to the projections δ_2 appearing in (2.10), they also satisfy $T|_{z=0} = \partial_z T|_{z=0} = 0$, i.e., are elements of the space \mathcal{A}_1 defined in (2.6). On this space, the operator \mathcal{K} in (2.2) has a bounded inverse given by a simple integral kernel [Ste07]. Consequently, the analyticity domain of a tree in the z-variable is in fact much larger than the neighbourhood of [-1,1] included in Π ; it contains the wedgelike region $\mathbb{U}_{\tau,\vartheta}$ in (1.18).

Lemma 2.2 (Analytic continuation of trees). Without affecting the analyticity domain with respect to ϵ , there exist numbers $0 < \tau < 1$, $0 < \vartheta < \pi/2$, and $0 < \sigma < \eta$ such that each tree in the sums (2.15) and (2.17) extends to an analytic function of (z, θ) on $\mathbb{U}_{\tau, \vartheta} \times \{|\Im \mathfrak{m} \theta| \leq \sigma\}$.

Since the sum in (2.16) is finite and the functions $\widetilde{X}_{\leq 1}$ and X^0 in $X = X^0 + \widetilde{X}_{\leq 1} + \delta_2 \widetilde{X}$ are analytic on $\mathbb{U}_{\tau,\vartheta} \times \{|\Im \mathfrak{m} \theta| \leq \sigma\}$, Theorem 2 follows from Lemma 2.2.

3. Size of the Homoclinic Splitting

The Lyapunov exponent $\gamma = \gamma(\epsilon, g)$ is an analytic function of ϵ in a neighbourhood of the origin $(|\epsilon| < \epsilon_0)$. In fact, $\gamma = g(1 + \mathcal{O}(\epsilon))$ for small ϵ . Further, since \mathcal{L} depends on γ , we have $\mathcal{L} = \mathcal{L}_0 + \mathcal{O}(\epsilon)$ with $\mathcal{L}_0 = \omega \cdot \partial_\theta + gz\partial_z$. In particular, as $X^0(z,\theta) \equiv (4\arctan z, 0)^T$, the expansion of the splitting matrix Υ —see (1.20)—begins linearly with respect to ϵ :

$$\Upsilon = \sum_{\ell=1}^{\infty} \epsilon^{\ell} \Upsilon^{\ell}.$$

It is customary to call the first order coefficient, Υ^1 , the Melnikov term or the Melnikov matrix.

3.1. **The Melnikov term.** Let us study Υ^1 . The symmetry (1.3) comes in very handy, since it implies, together with the even parity of f, that $f(\Phi^0(z^{-1}), -\theta) \equiv f(\Phi^0(z), \theta)$. Setting

$$F(s,\theta) \equiv \int_{-\infty}^{\infty} \partial_{\psi}^{2} f(\Phi^{0}(e^{gt+s}), \theta + \omega t) - \partial_{\psi}^{2} f(0, \theta + \omega t) dt$$
(3.1)

then yields, after some work,

$$\Upsilon^1(t) = g^2 F(0,0)$$

for all t. Here the second term in the integrand removes the quasiperiodic limit of the first one, making the integral absolutely convergent. In fact, the integrand tends to zero at an exponential rate as $|t| \to \infty$.

We now give a result, due to Lazutkin in it's original form, [Laz03], and extended by [DGJS97a], [Sau01], and [LMS03] to the quasiperiodic setting. See Appendix A for the proof.

Lemma 3.1 (Lazutkin). Suppose a function F is analytic on $S := [-i\vartheta, i\vartheta] \times \{|\Im \mathfrak{m} \theta| < \eta\}$, continuous on the closure \bar{S} , and satisfies the identity

$$F(s,\theta) \equiv F(0,\theta - g^{-1}\omega s). \tag{3.2}$$

Then F extends analytically to $\{|\Im \mathfrak{m} s| \leq \vartheta\} \times \{|\Im \mathfrak{m} \theta| < \eta\}$, with (3.2) still holding. On $\mathbb{R} \times \mathbb{T}^d$ it obeys a bound of the form

$$|F(s,\theta) - \widetilde{F}| \le CB(g)e^{-c\,g^{-1/(\nu+1)}}, \quad B(g) := \sup_{(s,\theta) \in \bar{S}} |F(s,\theta)|,$$

where \widetilde{F} stands for the θ -average $\langle F(s,\,\cdot\,) \rangle$ and is independent of s.

Remark 3.2. Notice that, for small values of g, the analyticity domain of F in Lemma 3.1 is much larger than what the right-hand side of (3.2) suggests.

The lemma applies to the function defined in (3.1), due to Remarks 1.3 and 1.4 regarding the analyticity of f. Because the ψ -derivatives in the integrand are interchangeable with total θ -derivatives, $\tilde{F} = 0$, and we get

Proposition 3.3. There exist positive constants c and C such that the Melnikov term satisfies

$$|\Upsilon^1| \le Cg^2 e^{-c g^{-1/(\nu+1)}}.$$

Remark 3.4. Identity (3.2) in Lazutkin's lemma is equivalent to

$$(\partial_s + g^{-1}\omega \cdot \partial_\theta)F = 0, (3.3)$$

or, calling $\bar{F}(z,\theta) \equiv F(g^{-1} \ln z, \theta)$, to $\mathcal{L}_0 \bar{F} = 0$. There exists a whole industry trying to push through the argument given above replacing Υ^1 with a "better measure" of the splitting, such as Υ , by searching for a coordinate system in which the "measure" satisfies (3.3). The state of this quest is best described in [LMS03].

Neither Lemma 3.1 nor Proposition 3.3 is optimal. There is a refinement in [Sau01] which, however, is not optimal. In many special cases one can derive a much stronger bound.

3.2. **Regularized integrals.** This subsection borrows heavily from Gallavotti, $et\ al.$; see for instance [CG94] and [GGM99b]. We adhere to the notation

$$X(t; z, \theta) := X(ze^{\gamma t}, \theta + \omega t),$$

which makes the dependence on the Lyapunov exponent γ implicit (Recall that γ is a function of ϵ), such that the perturbation expansion reads compactly $X(t;z,\theta) = \sum_{\ell=0}^{\infty} \epsilon^{\ell} X^{\ell}(t;z,\theta)$. With little trickery, the equations of motion retain their appearance: defining $\Omega(X)(t;z,\theta) \equiv$

 $[(g^2 \sin \Phi, 0) + \lambda \widetilde{\Omega}(X)](t; z, \theta)$ together with $\widetilde{\Omega}(X)(t; z, \theta) \equiv \partial f(X(t; z, \theta) + (0, \theta + \omega t))$, and implementing (1.7), we obtain the analogue of (1.11):

$$\partial_t Y(t; z, \theta) := \partial_t^2 X(t; z, \theta) = \Omega(X)(t; z, \theta). \tag{3.4}$$

There is a need to consider functions $h(t; z, \theta)$ that can be expanded as finite sums

$$\sum_{p=0}^{P} t^p h^p (ze^{\gamma t}, \theta + \omega t), \tag{3.5}$$

where h^p is analytic on $\{0 < |z| < \tau\} \times \mathbb{T}^d$ —a punctured neighbourhood of the origin times the d-torus—with at worst a finite pole at z = 0.

Let us set $T_N h(t;z,\theta) := \sum_{|q| \le N} \hat{h}(t;z,q) e^{iq\cdot\theta}$ and define a regularized integral

$$\int h\left(t;z,\theta\right) := \lim_{N \to \infty} \mathop{\mathrm{res}}_{R=0} \frac{1}{R} \int_{-\infty}^{t} e^{-R|\tau|} T_N h(\tau;z,\theta) \, d\tau,$$

the residue at R=0 meaning that of the analytic extension from the complex half-plane with $\Re e R$ sufficiently large. This is well-known from the theory of Laplace transforms. The truncation, T_N , is needed to insure that the origin of the R-plane is not an accumulation point for poles, since factors $(R + i\omega \cdot q)^{-1}$ appear from the integral.

By the following proposition, it is natural to employ the compelling notation

$$\int_{-\infty}^{t} h(\tau; z, \theta) d\tau := \int h(t; z, \theta),$$

and to view z and θ as mere parameters:

Proposition 3.5. The regularized integral is linear. Moreover,

- (1) f h (t; z, θ) = f h (t₀; z, θ) + ∫_{t₀}^t h(τ; z, θ) dτ for any real t₀,
 (2) if h and h' are trigonometric polynomials (in θ) and h(t; z, θ) = h'(t; z, θ) for all t, then $f h(t; z, \theta) = f h'(t; z, \theta)$ for all t,
- (3) $T_N
 otin h =
 otin T_N h$,
- (4) $\partial_t \int h(t;z,\theta) = \frac{d}{dt} \int_{-\infty}^t h(\tau;z,\theta) d\tau = h(t;z,\theta),$
- (5) $\iint h(t;z,\theta) =: \int_{-\infty}^{t} \int_{-\infty}^{\tau} h(s;z,\theta) \, ds \, d\tau = \int_{-\infty}^{t} (t-\tau) \, h(\tau;z,\theta) \, d\tau.$

Proof. For a finite t_0 , $\operatorname{res}_{R=0} \frac{1}{R} \int_{t_0}^t e^{-R|\tau|} u(\tau) d\tau$ equals the evaluation $\int_{t_0}^t e^{-R|\tau|} u(\tau) d\tau \Big|_{R=0} = 0$ $\int_{t_0}^t u(\tau) d\tau$, since the integral is entire in R, from which (1) follows. If h is a trigonometric polynomial, then $\int h(t;z,\theta) = \underset{R=0}{\text{res}} \frac{1}{R} \int_{-\infty}^{t} e^{-R|\tau|} h(\tau;z,\theta) d\tau$, yielding (2). Claim (3) is trivial and (4) follows from (1). In order to prove (5), one first checks that it holds for trigonometric polynomials and then deals with the general case.

We also define

$$\int_{t}^{\infty} h(\tau; z, \theta) d\tau := \lim_{N \to \infty} \operatorname{res}_{R=0}^{1} \frac{1}{R} \int_{t}^{\infty} e^{-R|\tau|} T_{N} h(\tau; z, \theta) d\tau$$

when it makes sense, and write

$$f_{\pm\infty}^t = -f_t^{\pm\infty}$$
 and $f_{-\infty}^t + f_t^{\infty} = f_{-\infty}^{\infty}$.

The operator \mathcal{D}^{-1} is defined by giving the expression of its Fourier kernel:

$$\mathcal{D}^{-1}(p,q) = \begin{cases} (i\omega \cdot q)^{-1} & \text{if } p = q \in \mathbb{Z}^d \setminus \{0\}, \\ 0 & \text{otherwise.} \end{cases}$$

It is the formal inverse of \mathcal{D} on functions $h(z,\theta)$ with vanishing θ -average. Finally,

$$\mathcal{I}h(z,\theta) := \int_{-\infty}^{0} h(ze^{\gamma\tau}, \theta + \omega\tau) d\tau.$$

Applying (1.7) in the case $F = \mathcal{I}h$, the formal identity $\mathcal{L}\mathcal{I} = \mathbb{1}$ follows. Recalling $\delta_1 h(z, \theta) := h(z, \theta) - h_0(\theta) = h(z, \theta) - h(0, \theta)$, one also gets $\mathcal{I}\mathcal{L}\delta_1 = \delta_1$ by a direct computation.

Proposition 3.6. Let h be analytic on $\{|z| < \tau\} \times \mathbb{T}^d$, and denote $h(t; z, \theta) \equiv h(ze^{\gamma t}, \theta + \omega t)$.

$$\int_{-\infty}^{t} h(\tau; z, \theta) d\tau = \langle h_0 \rangle t + \mathcal{D}^{-1} h_0(\theta + \omega t) + \mathcal{I} \delta_1 h(z e^{\gamma t}, \theta + \omega t)$$

and

$$\int_{-\infty}^{t} \partial_{\tau} h(\tau; z, \theta) d\tau = h(t; z, \theta) - \langle h_0 \rangle.$$

Proof. Decompose $h = \langle h_0 \rangle + [h_0 - \langle h_0 \rangle] + \delta_1 h$ and use the linearity of f. For the second identity, notice $\partial_{\tau} h(\tau; z, \theta) = \mathcal{D}h_0(\theta + \omega \tau) + \mathcal{L}\delta_1 h(ze^{\gamma \tau}, \theta + \omega \tau)$, and recall the identities $\mathcal{D}^{-1}\mathcal{D}h_0 = h_0 - \langle h_0 \rangle$ and $\mathcal{I}\mathcal{L}\delta_1 h = \delta_1 h$.

Remark 3.7. Observe that, unless the integrals converge in the traditional sense, $f_{-\infty}^t$ does generically not tend to $f_{-\infty}^{\infty}$ or 0 as t tends to ∞ or $-\infty$, respectively.

A detailed proof of the following elementary result can be found in [Ste06].

Proposition 3.8. Let S_0 and S_1 be the spaces of analytic functions h on $\{|z| < \tau\} \times \mathbb{T}^d$ satisfying $\langle h_0 \rangle = 0$ and $\langle h_1 \rangle = 0$, respectively. Denote $h(t; z, \theta) \equiv h(ze^{\gamma t}, \theta + \omega t)$. The operators $L_t = \partial_t^2 - \gamma^2 \cos \Phi^0(ze^{\gamma t})$, $L = \mathcal{L}^2 - \gamma^2 \cos \Phi^0$, and $\mathcal{L} = \mathcal{D} + \gamma z \partial_z$ satisfy

$$\partial_t h(t; z, \theta) \equiv \mathcal{L}h(ze^{\gamma t}, \theta + \omega t)$$
 and $L_t h(t; z, \theta) \equiv Lh(ze^{\gamma t}, \theta + \omega t).$

The operator f maps S_0 into itself and, in fact, $f = \partial_t^{-1}$ on S_0 . Similarly, the operator defined by $\int_{-\infty}^t K_{\Phi}(t,\tau;z) h(\tau;z,\theta) d\tau$ maps S_1 into itself and is the inverse of L_t on S_1 . In particular, the operator $\mathcal{K} = \begin{pmatrix} L & 0 \\ 0 & \mathcal{L}^2 \end{pmatrix} = \begin{pmatrix} L_t & 0 \\ 0 & \partial_t^2 \end{pmatrix}$ is invertible on the space $S_1 \times S_0$, where

$$\mathcal{K}^{-1}h(t;z,\theta) \equiv \int_{-\infty}^{t} K(t,\tau;z)h(\tau;z,\theta) d\tau.$$

Above, the kernel $K = \begin{pmatrix} K_{\Phi} & 0 \\ 0 & K_{\Psi} \end{pmatrix}$ splits into the useful sum

$$K(t,\tau;z) = \sum_{i=0}^{1} \sum_{j=1}^{2} (t-\tau)^{i} K_{ij}(t;z) \bar{K}_{ij}(\tau;z),$$
(3.6)

setting $P(t;z) := (z^2 e^{2\gamma t} + 1)^{-1} z e^{\gamma t}$, $Q(t;z) := z^{-1} e^{-\gamma t} (z^2 e^{2\gamma t} - 1)$, and

$$K_{01} = \frac{1}{2\gamma} \begin{pmatrix} Q & 0 \\ 0 & 0 \end{pmatrix}, \qquad \bar{K}_{01} = \bar{K}_{11} = \begin{pmatrix} P & 0 \\ 0 & 0 \end{pmatrix}, \qquad K_{02} = -\frac{1}{2\gamma} \begin{pmatrix} P & 0 \\ 0 & 0 \end{pmatrix},$$

$$\bar{K}_{02} = \begin{pmatrix} Q & 0 \\ 0 & 0 \end{pmatrix}, \qquad K_{11} = 2 \begin{pmatrix} P & 0 \\ 0 & 0 \end{pmatrix}, \qquad K_{12} = \bar{K}_{12} = \begin{pmatrix} 0 & 0 \\ 0 & 1 \end{pmatrix}.$$

3.3. Asymptotic expansion for the splitting matrix. Starting from (3.4) and using Proposition 3.8 with $\langle Y_0^u \rangle = \langle \mathcal{D} X_0^u \rangle = 0$, we have

$$Y^{u}(t;z,\theta) = \int_{-\infty}^{t} \Omega(X^{u})(\tau;z,\theta) d\tau$$

The identity $X^s(t;z,\theta) \equiv (2\pi,0) - X^u(-t;z^{-1},-\theta)$ yields $Y^s(t;z,\theta) \equiv Y^u(-t;z^{-1},-\theta)$, such that defining the d component column vector $\Delta(t;z,\theta) := (Y_{\Psi}^u - Y_{\Psi}^s)(t;z,\theta)$ and the functions $f^{s,u}(t;z,\theta) := f(X^{s,u}(t;z,\theta) + (0,\theta+\omega t))$ thus results in an expression for the splitting matrix Υ of (1.20) in terms of regularized integrals: we get

$$\Upsilon(t) = \partial_{\theta} \Delta(t; 1, 0) = \partial_{\theta} \Delta(0; e^{\gamma t}, \omega t)$$
(3.7)

because

$$\partial_{\theta} \Delta(t; z, \theta) = \lambda \int_{-\infty}^{t} \partial_{\theta} f_{\psi}^{u}(\tau; z, \theta) d\tau + \lambda \int_{t}^{\infty} \partial_{\theta} f_{\psi}^{s}(\tau; z, \theta) d\tau,$$

where the subindices attached to $f^{s,u}$ stand for partial derivatives (our convention is $\varphi = (\phi, \psi)$ and $\partial = \partial_{\varphi} = (\partial_{\phi}, \partial_{\psi})$.

Remark 3.9. Extracting here the first order in $\epsilon = g^{-2}\lambda$ casts the Melnikov matrix in a compact form in terms of the regularized integrals:

$$\Upsilon^1 = g^2 \int_{-\infty}^{\infty} \partial_{\psi}^2 f(\Phi^0(e^{g\tau}), \omega \tau) d\tau.$$

The reader is invited to compare this with (3.1) and the equation following it. We point out that a similar procedure, using Lemma 3.1, that resulted in the exponentially small bound on Υ^1 can be applied here, despite the unusual integrals.

Suppose that, even with the superscript u, the integral

$$\int_{1}^{\infty} \partial_{\theta} f_{\psi}^{u}(\tau; z, \theta) d\tau \tag{3.8}$$

exists. If the integrand is a trigonometric polynomial in θ , then Proposition 3.5 implies that the value of the integral only depends on the integrand's restriction to $\{(t; z, \theta) | t \in \mathbb{R}\}$, as the notation suggests ¹. Now, fixing $(z,\theta)=(1,0)$ and dropping it from the notation,

$$\Upsilon(t) = \lambda \int_{-\infty}^{\infty} \partial_{\theta} f_{\psi}^{u}(\tau) d\tau + \lambda \int_{t}^{\infty} \left[f_{\psi\varphi}^{s} \partial_{\theta} (X^{s} - X^{u}) \right](\tau) d\tau, \tag{3.9}$$

since $X^{u}(\tau; 1, 0) \equiv X^{s}(\tau; 1, 0)$.

Let us make the assumption that $\psi \mapsto f(\cdot, \psi)$ is a trigonometric polynomial, such that, by Theorem 2, each order of X in ϵ is a trigonometric polynomial in θ . Even so, equation (3.9) is formal in the sense that the integrands are defined for $\tau \gg 0$ only up to an arbitrarily high order in ϵ . However, it is asymptotic, which is to say that at each order the identity is exact. Put differently, (3.9) is a collection of exact identities, one for each order in ϵ , written in closed form. Moreover, at each order, the integrands are analytic functions by virtue of the extension result in Theorem 2. Of course, we need to check that each order of (3.8), i.e., the integral $f_t^{\infty} (\partial_{\theta} f_{\psi}^u)^{\ell}(\tau; z, \theta) d\tau$ with arbitrary $\ell \in \mathbb{N}$, exists ². The point of the formula above is twofold. First, the integral $f_{-\infty}^{\infty} \partial_{\theta} f_{\psi}^u(\tau) d\tau$ (at each order

in ϵ) turns out to be exponentially small in the limit $g \to 0$. Second, we can actually construct

¹This is the first of the two places where we need to restrict ourselves to trigonometric polynomials...

²... and this is the second.

a (formal) recursion relation for the function $\partial_{\theta}(X^s - X^u)(\tau)$, taking us to ever-increasing orders in ϵ , as follows. Differentiating both sides of (3.4) with respect to θ one obtains

$$\begin{pmatrix} L_t & 0 \\ 0 & \partial_t^2 \end{pmatrix} \partial_\theta X^{s,u} = M^{s,u} \partial_\theta X^{s,u} + \lambda f_{\varphi\psi}^{s,u}$$

with $L_t = \partial_t^2 - \gamma^2 \cos \Phi^0(ze^{\gamma t})$, as in Proposition 3.8, and

$$M^{s,u}(t;z,\theta) := \begin{pmatrix} g^2 \cos \Phi^{s,u}(t;z,\theta) - \gamma^2 \cos \Phi^0(ze^{\gamma t}) & 0 \\ 0 & 0 \end{pmatrix} + \lambda f_{\varphi\varphi}^{s,u}(t;z,\theta).$$

Owing to the θ -derivative acting on $X^{s,u}$ and Proposition 3.8, we get

$$\partial_{\theta} X^{s,u}(t;z,\theta) = \int_{+\infty}^{t} K(t,t';z) \left[M^{s,u} \partial_{\theta} X^{s,u} + \lambda f_{\varphi\psi}^{s,u} \right] (t';z,\theta) dt',$$

upon choosing ∞ in conjunction with X^s and $-\infty$ with X^u . The most convenient way to see this is to infer the identity for $\partial_{\theta}X^u$ and employing $K(-t, -\tau; z^{-1}) \equiv -K(t, \tau; z)$ and $X^s(t; z, \theta) \equiv (2\pi, 0) - X^u(-t; z^{-1}, -\theta)$.

At $(z,\theta)=(1,0),\ M^s(t)\equiv M^u(t)$ and $f^s_{\varphi\psi}(t)\equiv f^u_{\varphi\psi}(t)$, such that (2) of Proposition 3.5 validates

$$\partial_{\theta}(X^{s} - X^{u})(\tau) = \int_{-\infty}^{-\infty} K(\tau, \tau') \left[M^{u} \partial_{\theta} X^{u} + \lambda f_{\varphi\psi}^{u} \right](\tau') d\tau' + \int_{-\infty}^{\tau} K(\tau, \tau') \left[M^{s} \partial_{\theta}(X^{s} - X^{u}) \right](\tau') d\tau'.$$

$$(3.10)$$

By construction, the quantities appearing in square brackets above are θ -gradients:

$$M^{s,u}\partial_{\theta}X^{s,u} + \lambda f_{\varphi\psi}^{s,u} = \partial_{\theta}\mathfrak{A}^{s,u}$$

with, for instance ³,

$$\mathfrak{A}^{s,u}:=\begin{pmatrix} L_t & 0 \\ 0 & \partial_t^2 \end{pmatrix} X^{s,u} = \Omega(X^{s,u}) - \begin{pmatrix} \gamma^2 \cos \Phi^0(ze^{\gamma t}) & 0 \\ 0 & 0 \end{pmatrix} X^{s,u},$$

regardless of the value of (z, θ) . In addition, but only at $(z, \theta) = (1, 0)$,

$$[M^s \partial_{\theta} (X^s - X^u)](t) = [\partial_{\theta} (\mathfrak{A}^s - \mathfrak{A}^u)](t).$$

Notice that (3.10) is a fixed point equation of the form $\zeta = \mathcal{E} + BM^s\zeta$ solved formally by $\zeta(t) = \partial_{\theta}(X^s - X^u)(t)$. This is highly interesting from the point of view of asymptotic analysis, since $M^s = \mathcal{O}(\epsilon)$ multiplies ζ on the right-hand side. Indeed, truncating the Taylor series in ϵ , e.g., $\zeta^{\leq k} = \sum_{i=0}^k \epsilon^i \zeta^i$, the identities $\zeta^k = (\mathcal{E} + BM^s\zeta)^k = (\mathcal{E} + BM^s\zeta^{\leq k-1})^k$ become exact. Hence, we can recursively construct any order of $\partial_{\theta}(X^s - X^u)(t)$ from its lower orders, without the need of diverting from $\{(t; z, \theta) = (t; 1, 0) \mid t \in \mathbb{R}\}$ in order to compute the θ -derivative. Cumulatively, we have $\zeta^k = \left[\sum_{i=0}^{k-1} (BM^s)^j \mathcal{E}\right]^k$, where j is a power.

Proposition 3.10. In brief, the (presumably) divergent Neumann series $\sum_{j=0}^{\infty} (BM^s)^j \mathcal{E}$ is an asymptotic expansion of $\partial_{\theta}(X^s - X^u)(t)$ in the sense that

- (1) At each order in ϵ it terminates after finitely many well-defined terms, and
- (2) $\left[\partial_{\theta}(X^s X^u) \left[\sum_{j=0}^{\infty} (BM^s)^j \mathcal{E}\right]^{\leq k}\right](t) = \mathcal{O}(\epsilon^{k+1}) \text{ for each } k \in \mathbb{N}.$

Once inserted into (3.9), the latter series provides us with an asymptotic expansion for the splitting matrix Υ in (3.7):

³We could make $\mathfrak{A}^{s,u}$ of order ϵ by deleting its lowest order, because the latter does not depend on θ .

Corollary 3.11. In the asymptotic sense of Proposition 3.10,

$$\Upsilon(t) = \lambda \int_{-\infty}^{\infty} \partial_{\theta} f_{\psi}^{u}(\tau) d\tau + \lambda \int_{t}^{\infty} \left[f_{\psi\varphi}^{s} \sum_{i=0}^{\infty} (BM^{s})^{j} \mathcal{E} \right] (\tau) d\tau.$$

The operator B above has the expression

$$Bh(t;z,\theta) = \int_{-\infty}^{t} K(t,\tau;z)h(\tau;z,\theta) d\tau, \qquad (3.11)$$

whereas \mathcal{E} is the restriction of $\mathcal{E}(t;z,\theta) := \int_{-\infty}^{-\infty} K(t,\tau;z) \, \partial_{\theta} \mathfrak{A}^{u}(\tau;z,\theta) \, d\tau$ to $(z,\theta) = (1,0)$. Using (3.6), also \mathcal{E} splits into pieces: $\mathcal{E}(t;z,\theta) = \sum_{i=0}^{1} \sum_{p=0}^{i} \sum_{j=1}^{2} t^{p} K_{ij}(t;z) \, \mathcal{E}_{ij}^{p}(z,\theta)$, where

$$\mathcal{E}_{ij}^p(z,\theta) := -\int_{-\infty}^{\infty} (-\tau)^{\delta_{i1}-p} \bar{K}_{ij}(\tau;z) \,\partial_{\theta} \mathfrak{A}^u(\tau;z,\theta) \,d\tau.$$

We shall see that these t-independent factors are exponentially small with respect to g, as the latter tends to zero. Furthermore, denoting $K_{ij}^p(t;z) := t^p K_{ij}(t;z)$,

$$(BM^s)^l \mathcal{E} = \sum_{i=0}^1 \sum_{p=0}^i \sum_{j=1}^2 \left[(BM^s)^l K_{ij}^p \right] \mathcal{E}_{ij}^p.$$

This is so because, by Proposition 3.5, all integrals due to (3.11) factorize formally: if h is a trigonometric polynomial at each order, then

$$\int_{t}^{\infty} (h\mathcal{E})(\tau;z,\theta) d\tau = \sum_{i=0}^{1} \sum_{n=0}^{i} \sum_{j=1}^{2} \left[\int_{t}^{\infty} h(\tau;z,\theta) \tau^{p} K_{ij}(\tau;z) d\tau \right] \mathcal{E}_{ij}^{p}(z,\theta).$$

By virtue of Corollary 3.11, we infer

Proposition 3.12. On the homoclinic trajectory, i.e., setting $(z, \theta) = (1, 0)$, the following asymptotic expansion of the splitting matrix holds:

$$\Upsilon(t) = \int_{-\infty}^{\infty} \lambda \,\partial_{\theta} f_{\psi}^{u}(\tau) \,d\tau + \lambda \,c_{ij}^{p}(t) \mathcal{E}_{ij}^{p}, \tag{3.12}$$

where repeated indices are contracted ($i \in \{0,1\}$, $p \in \{0,i\}$, $j \in \{1,2\}$) and

$$c_{ij}^p(t) := \int_t^\infty \left[f_{\psi\varphi}^s \sum_{l=0}^\infty (BM^s)^l K_{ij}^p \right] (\tau) \, d\tau.$$

As already pointed out, the terms appearing in the asymptotic expansion of Proposition 3.12 will turn out exponentially small with respect to g. The factors c_{ij}^p shall pose no problems, the functions M^s and K_{ij}^p being explicit and simple. In brief, Theorem 4 begins to emerge!

For the record, writing $B^u h(t;z,\theta) := \int_{-\infty}^t K(t,\tau;z) h(\tau;z,\theta) d\tau$, we get (at $(z,\theta) = (1,0)$)

$$c_{ij}^{p}(t) = (-1)^{p+j\delta_{i0}} \int_{-\infty}^{-t} \left[f_{\psi\varphi}^{u} \sum_{l=0}^{\infty} (B^{u} M^{u})^{l} K_{ij}^{p} \right] (\tau) d\tau.$$
 (3.13)

3.4. Emergence of exponential smallness. The "asymptotic" integrals in (3.12), including $\lambda \mathcal{E}_{ij}^p$, are of the form

$$\int_{-\infty}^{\infty} \partial_{\theta} F(X^{u}; t; z, \theta) dt \equiv \sum_{\ell=1}^{\infty} \epsilon^{\ell} \sum_{0 < |q| < \ell N} iq \, e^{iq \cdot \theta} \int_{-\infty}^{\infty} \left[\hat{F}(\hat{X}^{u}; t; z, q) \right]^{\ell} dt,$$

evaluated at $(z, \theta) = (1, 0)$, where the integrand on the right-hand side is a trigonometric polynomial of degree $\leq \ell N$ by Theorem 2. We also point out that, by construction, the latter are θ -gradients, which allows us to omit the harmful q = 0 terms. Above, F depends on $X = X^u$ locally, i.e., only through $X(t; z, \theta)$, as well as analytically near X^0 , i.e., the series

$$F(X) = F(X^0 + \widetilde{X}) = \sum_{k=0}^{\infty} F^{(k)} (\widetilde{X})^{\otimes k}$$
(3.14)

converges. In fact, since F is one of the functions in

$$\left\{\lambda f_{\psi}^{u}\right\} \bigcup \left\{-\lambda (-t)^{\delta_{i1}-p} \bar{K}_{ij} \mathfrak{A}^{u} \mid 0 \le p \le i \le 1 \text{ and } 1 \le j \le 2\right\},\tag{3.15}$$

we observe that

$$F(X^{u}; t; z, \theta) \equiv t^{p} F(X^{u}; 0; ze^{\gamma t}, \theta + \omega t) \qquad (p \in \{0, 1\}),$$
 (3.16)

with $(z,\theta) \mapsto [F(X^u;0;z,\theta)]^{\ell}$ analytic on $(\mathbb{U}_{\tau,\theta} \setminus \{0\}) \times \{|\Im \mathfrak{m} \theta| \leq \sigma\}$ for all ℓ by virtue of Theorem 2. At z=0 there is a simple pole, due to \bar{K}_{02} , at worst.

We use the following lemma for analyzing such integrals. Its proof is given in Appendix A.

Lemma 3.13 (Shift of contour). Suppose that the function $h(t; z, \theta) \equiv t^p h(ze^{gt}, \theta + \omega t)$ is analytic with respect to $(z, \theta) \in (\mathbb{U}_{\tau, \theta} \setminus \{0\}) \times \{|\Im \theta| \leq \sigma\}$ and $\hat{h}(\cdot, 0) = 0$. Moreover, set $t_q := \operatorname{sgn}(\omega \cdot q)\vartheta g^{-1}$ and

$$H_q(R) := \int_0^\infty e^{(iq \cdot \omega - R)t} (t + it_q)^p \hat{h}(e^{g(t + it_q)}, q) dt$$
 (3.17)

for each $q \in \mathbb{Z}^d \setminus \{0\}$. If

$$\left| \int_{-\infty}^{0} e^{iq \cdot \omega t} (t + it_q)^p \hat{h}(e^{g(t + it_q)}, q) dt \right| \le A(g)e^{-\sigma|q|}$$

and if $H_q(R)$ admits an analytic continuation to $\{0 < |R| \le \rho\}$ with a pole of order k at R = 0, respecting the bound

$$\sup_{|R|=\rho} |H_q(R)| \le B(g)e^{-\sigma|q|},$$

then we obtain the exponentially small (c > 0) bound

$$\left| \oint_{-\infty}^{\infty} h(t;1,0) dt \right| \le C \left[A(g) + B(g) \sum_{j=0}^{k} \frac{1}{j!} \left(\frac{\rho \vartheta}{g} \right)^j \right] e^{-cg^{-1/(\nu+1)}}.$$

With the aid of Lemma 3.13, we shall prove in Section 4 the following key result:

Proposition 3.14 (Convergence vs. exponential smallness). Fix $a \ t \in \mathbb{R}$. There exist positive constants c, ϵ_1 and C, such that the estimates

$$|\Upsilon^{\ell}(t)| \leq C \begin{cases} \epsilon_1^{-\ell} \ell!^{4(\nu+1)} e^{-cg^{-1/(\nu+1)}} \\ \epsilon_1^{-\ell} \end{cases}$$

both hold true for all $\ell = 1, 2, \ldots$

This dichotomy, in which the exponential smallness competes with the usual bound due to convergence, is not new; see [Gal94,GGM99b,Pro03]. Remarkably, we get the same exponent of the factorial, $4(\nu + 1)$, as the latter articles—even though our method is quite different.

Theorem 4 follows immediately from Proposition 3.14 by an argument due to Gallavotti, et al.: For each g, let n(g) be a positive integer. If $|\tilde{\epsilon}| < \frac{1}{2}$ and $\epsilon = \tilde{\epsilon} \epsilon_1 n(g)^{-4(\nu+1)}$,

$$\begin{split} |\Upsilon(t)| &\leq C e^{-cg^{-1/(\nu+1)}} \sum_{\ell=1}^{n(g)} \left(\frac{|\epsilon|}{\epsilon_1}\right)^{\ell} \ell!^{4(\nu+1)} + C \sum_{\ell=n(g)+1}^{\infty} \left(\frac{|\epsilon|}{\epsilon_1}\right)^{\ell} \\ &\leq C e^{-cg^{-1/(\nu+1)}} \frac{|\epsilon|}{\epsilon_1} n(g)^{4(\nu+1)} + C \left(\frac{|\epsilon|}{\epsilon_1}\right)^{n(g)+1} \\ &\leq C |\tilde{\epsilon}| e^{-cg^{-1/(\nu+1)}} + C |\tilde{\epsilon}| e^{n(g)\ln|\tilde{\epsilon}|}, \end{split}$$

since $\ell! \leq \ell^{\ell} \leq n(g)^{\ell}$ in the first sum. Now, set $n(g) = cg^{-1/(\nu+1)}/\ln 2$, such that $n(g)\ln|\tilde{\epsilon}| \leq -cg^{-1/(\nu+1)}$ and $\tilde{\epsilon} = (\epsilon/\epsilon_1)(c/\ln 2)^{4(\nu+1)}g^{-4}$ (demanding in particular that ϵg^{-4} be small).

Remark 3.15. We used the exponentially small but diverging estimate to bound the partial sum $\sum_{\ell=1}^{n(g)} \epsilon^{\ell} \Upsilon^{\ell}$, whereas the remainder of the series was easily controlled by convergence. As $n(g) \to \infty$ with $g \to 0$, the important thing here is to have the exponentially small bound on Υ^{ℓ} for arbitrarily large ℓ , in addition to the ϵ -analyticity of Υ .

4. Proof of Theorem 4

We are left with proving Proposition 3.14, since Theorem 4 was already shown to be its corollary. Here things are most conveniently explained using tree diagrams. However, each tree will be treated as an *individual*, solely for bookkeeping benefits, and no cancellations nor regroupings are forced upon them.

By Proposition 3.12, we need to consider the simple factors c_{ij}^p , as well as the integrals

$$\int_{-\infty}^{\infty} \hat{F}(\hat{X}^u; t, z, q) dt \equiv \sum_{\ell=1}^{\infty} \epsilon^{\ell} \int_{-\infty}^{\infty} \left[\hat{F}(\hat{X}^u; t; z, q) \right]^{\ell} dt$$
 (4.1)

of Subsection 3.4, at each order $\ell \ge |q|/N > 0$. As F comes from the collection in (3.15), all integrals of the latter type shall be controlled with the aid of Lemma 3.13.

Due to the superscript u—referring to the unstable manifold—in the integrand above, the required bounds on the integrals over \mathbb{R}_- in $f_{-\infty}^{\infty} = f_{-\infty}^0 + f_0^{\infty}$ are straightforward, and are discussed later.

In order to deal with f_0^{∞} , we present the procedure below, which amounts to *little more than integration by parts*. First, we expand F according to (3.14) and, like in Section 2, split

$$\widetilde{X} = \widetilde{X}_{\leq 1} + \delta_2 \widetilde{X} = \longrightarrow + \longrightarrow ,$$

dropping the superscript u from the notation. We can then express F pictorially as

$$\sum_{m=0}^{\infty} \sum_{0 \le m' \le m} \binom{m}{m'} \xrightarrow{F^{(m)}} \cdots \xrightarrow{m'},$$

where the binomial coefficient comes from the combinatorics of shuffling the arguments of the symmetric $F^{(m)}$, which is attached to the root (node). Next, we replace the subtrees $\longrightarrow = \delta_2 \widetilde{X}$ with the expansion (2.15) derived from (2.11).

In a tree T_{v_0} with root v_0 , a generic subtree T_v with root v has then the expression

$$T_{v} = U_{v} \left(\mathcal{K}^{-1} T_{w_{1}}, \dots, \mathcal{K}^{-1} T_{w_{m'_{v}}}; \left(\widetilde{X}_{\leq 1} \right)^{\otimes (m_{v} - m'_{v})} \right), \tag{4.2}$$

where T_{w_j} is a subtree—with root w_j —entering v and the node function

$$U_v = \begin{cases} F^{(m_v)} \text{ coming from (3.15)}, & \text{if } v = v_0, \\ w^{(m_v)} \text{ as in (2.8)}, & \text{otherwise.} \end{cases}$$

In a pictorial representation, there are m_v lines entering the node v and precisely $m_v - m'_v$ of the latter are leaving a black dot (\bullet) end node:

$$T_v = U_v \qquad T_{w_{n'_v}} \qquad .$$

$$T_{w_{n'_v}} \qquad .$$

$$m_v - m'_v$$

The "whole" tree T_{v_0} contributes at orders ℓ satisfying

$$\ell \ge 1 + (m_{v_0} - m'_{v_0}) + \sum_{j=1}^{m'_{v_0}} \deg T_{w_j} =: d(T_{v_0}), \tag{4.3}$$

where deg T_w —defined in (2.14)—counts end nodes with suitable weights as well as nodes with exactly one entering line in the subtree T_w , and the constant term 1 counts the root $(F^{(m_{v_0})} = \mathcal{O}(\epsilon))$. By this we mean that $d(T_{v_0})$ is the largest integer such that

$$T_{v_0} = \mathcal{O}(\epsilon^{\mathrm{d}(T_{v_0})})$$
 as $\epsilon \to 0$.

4.1. Some combinatorics.

Lemma 4.1. The number of nodes, $n(T_{v_0})$, of a tree T_{v_0} contributing at order $\ell \geq 2$ obeys

$$n(T_{v_0}) \le 2(d(T_{v_0}) - 1) \le 2(\ell - 1).$$
 (4.4)

The number of end nodes is at most $d(T_{v_0}) - 1 \le \ell - 1$.

Proof. $n(T_{v_0})$ attains its maximum with respect to the "degree" $d(T_{v_0})$ as follows:

- (1) If $n(T_{v_0})$ is even, there is only one line entering the root v_0 and the rest of the tree is binary, *i.e.*, contains only end nodes and nodes with exactly two entering lines.
- (2) If $n(T_{v_0})$ is odd, the tree is binary.

Moreover, each of the end nodes is either \bullet or \bigcirc 1, which contribute the least to the degree; see (2.14). These choices minimize the number of end nodes when the $\mathcal{O}(\epsilon)$ nodes having exactly one entering line are excluded (except at v_0 which is always $\mathcal{O}(\epsilon)$). Therefore, $d(T_{v_0})$ gets minimized with respect to the number of all nodes, $n(T_{v_0})$. Since in a binary tree of j end nodes there are 2j-1 nodes, we infer (4.4).

The root and each of the end nodes increases $d(T_{v_0})$ by at least one, which implies the bound on the number of end nodes.

⁴This convention merely facilitates drawing.

Corollary 4.2. At most $2^{6\ell}$ trees contribute at order ℓ .

Proof. It is well-known that the number of (rooted) trees with k indistinguishable nodes is $N(k) := \frac{1}{k} \binom{2k-2}{k-1} \le \frac{1}{k} 4^{k-1}$, which follows from generating functions (see [Drm04]) and the bound $(2m)! \le 4^m (m!)^2$. By Lemma 4.1, the number of end nodes is less than ℓ . We label the latter arbitrarily by the labels in $\{\bullet\} \bigcup \{m \mid 1 \le m \le \ell-1\}$ in order to form an upper bound on the number of our trees; these are the only possible labels, as otherwise $d(T_{v_0})$ certainly exceeds the order ℓ —contradicting (4.3). The labeling of a tree with j end nodes can be carried out in at most $\binom{j+\ell-1}{j-1} \le 2^{j+\ell-1}$ ways. The desired bound is thus obtained from $1 + \sum_{k=2}^{\ell-1} N(k) \binom{\ell+k-2}{k-2} + \binom{2\ell-2}{\ell-2} \sum_{k=\ell}^{2(\ell-1)} N(k) \le 2^{6\ell}$, because by Lemma 4.1 the number of end nodes is at most $\ell-1$ even though the total number of nodes can be as large as $2(\ell-1)$. The term 1 on the left-hand side counts the single node tree $F^{(0)}$.

4.2. **Simplification of integrals: scalar trees.** We take a preliminary step towards bounding the values of the trees.

Let us split the kernels K of the operators \mathcal{K}^{-1} appearing in (4.2) into four pieces according to (3.6). These operators are attached to the lines between the nodes of a tree. Each of the $2^{6\ell}$ trees counted in Corollary 4.2 thus breaks into at most $4^{2\ell}$ new trees, as there are no more than 2ℓ such lines by Lemma 4.1.

At the same time, we also expand the matrix products due to the coordinate representation of (4.2) at each node: $(U_v)^i(\widetilde{T}_1,\ldots,\widetilde{T}_m) = \sum_{j_1=1}^{d+1}\cdots\sum_{j_m=1}^{d+1}(U_v)^i_{j_1\ldots j_m}\widetilde{T}_1^{j_1}\ldots\widetilde{T}_m^{j_m}$, where the \widetilde{T}_k 's represent all arguments of U_v (i.e., subtrees entering the node v)—including $\widetilde{X}_{\leq 1}$ —and the superindices specify vector components. We next separate each scalar term into its own tree, thus getting up to $(d+1)^{2\ell}$ of these scalar trees from each old tree.

There are lines carrying a factor $t-\tau$, coming from the i=1 terms of (3.6). They correspond to double integrals: $\iint_{-\infty}^{t} (d\tau)^2 = \int_{-\infty}^{t} d\tau (t-\tau)$. We remove these factors, insert a new node \tilde{v} on the line with a node function $U_{\tilde{v}} \equiv 1$ and an integral sign both on the line leaving and on the line entering \tilde{v} . This operation can be depicted as

$$\frac{1}{\int_{-\infty}^{t}} \longmapsto \frac{1}{\int_{-\infty}^{t} \int_{-\infty}^{\tau}}.$$
(4.5)

It does not affect the number of trees, but slightly simplifies the discussion below, even though the number of lines in a single tree can be as much as doubled.

Altogether, we arrive at the following conclusion:

Proposition 4.3. There are less than C_d^{ℓ} scalar trees to be considered, with $C_d = 2^{10}(d+1)^2$, at each order ℓ . Each of these trees has at most $4(\ell-1)$ lines.

Remark 4.4 (Some conventions). From now on, by a tree we will always refer to a scalar tree, where all the decompositions above have been carried out. In order to alleviate notation, we systematically omit all vector component indices.

Recall that the node w' is the unique predecessor of the node w. Let T_{v_0} be a generic (scalar) tree with root v_0 . Denote by V the set of all nodes and by V_{int} the set of integrated nodes, i.e., the nodes whose leaving line carries an integral. We consider the root v_0 an integrated node, such that V_{int} consists of all nodes of T_{v_0} except black dot (\bullet) end nodes. We can describe T_{v_0} by giving its structure recursively: if T_v is the subtree of T with root $v \in V_{\text{int}}$, then

$$T_v(t;z,\theta) = u_v(t;z,\theta) \prod_{\substack{w \in V_{\text{int}} \\ w'=v}} \int_{-\infty}^t T_w(\tau;z,\theta) d\tau.$$
 (4.6)

We call u_v the *multiplier* of the node v. It is a scalar function comprising all factors in the expression of the tree carrying the same variable of integration ("time label"), constricted in-between integral signs. In particular, it contains all subtrees ——— entering v, as these are the functions $\widetilde{X}_{\leq 1}$ involving no integrals.

To be completely explicit, the multiplier u_v of a node $v \in V_{\text{int}}$ is one of the functions below:

- (1) $\bar{K}_v \delta_2 h^{(k_v)}$, if v is the end node (k_v) ; see (2.10).
- (2) $\bar{K}_v w^{(m_v)} \Big(\prod_{\substack{w \notin V_{\text{int}} \\ w'=v}} \tilde{X}_{\leq 1} \Big) \Big(\prod_{\substack{w \in V_{\text{int}} \\ w'=v}} K_w \Big)$, if v is neither an end node nor the root v_0 , and didn't appear from splitting a double integral according to (4.5).
- (3) 1, if v appeared from splitting a double integral; see (4.5).

(4)
$$F^{(m_v)} \left(\prod_{\substack{w \notin V_{\text{int}} \\ w'=v}} \widetilde{X}_{\leq 1} \right) \left(\prod_{\substack{w \in V_{\text{int}} \\ w'=v}} K_w \right)$$
, if $v = v_0$; see (3.15).

The functions K_v and \bar{K}_v refer to diagonal elements of K_{ij} and \bar{K}_{ij} in (3.6), respectively. m'_v is the cardinality of $\{w \in V_{\text{int}} \mid w' = v\}$ and $m_v - m'_v$ the cardinality of $\{w \in V \setminus V_{\text{int}} \mid w' = v\}$. We draw the scalar (sub)tree in (4.6)—originated from (4.2)—as

$$T_v = u_v \xrightarrow{T_{w_1}} T_{w_{n'_v}}$$

This diagram is reminiscent of the one below (4.2), except that the operator U_v has changed into the multiplier u_v and the subtrees — have been absorbed into u_v . Moreover, recalling the decompositions above, the present diagram carries $\int_{-\infty}^{t}$ on its lines instead of \mathcal{K}^{-1} —compare (4.6) with (4.2)—and has possibly more nodes due to the diagram in (4.5).

We do not consider the end nodes with a black dot (\bullet) nodes anymore. Subsequently, we will refer to "integrated nodes" (elements of $V_{\rm int}$) as just "nodes".

4.3. Integration by parts: one step. Reinserting the implicit arguments (z, θ) ,

$$u_v(t;z,\theta) \equiv t^p u_v(ze^{\gamma t}, \theta + \omega t) \tag{4.7}$$

for the obviously defined $u_v(z,\theta)$. The power p can be nonzero only at the root, $v=v_0$, where it possibly assumes the value 1 due to case (4) in the list of all possible multipliers above.

We now turn our attention to the Fourier transformed integrals

$$\int_0^\infty T_{v_0}(t;z,q) dt = \int_0^\infty \sum_{q_{v_0} + \sum_{w' = v_0} \tilde{q}_w = q} u_{v_0}(t;z,q_{v_0}) \prod_{w' = v_0} \int_{-\infty}^t T_w(\tau;z,\tilde{q}_w) d\tau dt, \qquad (4.8)$$

which—recalling (4.6)—arise from (4.1), at each order $\ell \geq |q|/N > 0$. We omit the usual hat in Fourier transforms, since there is no danger of confusion. Due to Theorem 2, we may impose the finiteness condition

$$|q_{v_0}| + \sum_{w'=v_0} |\tilde{q}_w| \le \ell N,$$
 (4.9)

such that $\int_0^\infty T_{v_0}(t;z,q) dt$ agrees at order ℓ with

$$\sum_{q_{v_0} + \sum_{w' = v_0} \tilde{q}_w = q}^{\star} \operatorname{res}_{R=0}^{\star} \frac{1}{R} \int_0^{\infty} e^{-Rt} u_{v_0}(t; z, q_{v_0}) \prod_{w' = v_0} \int_{-\infty}^t T_w(\tau; z, \tilde{q}_w) d\tau dt.$$
(4.10)

Here the asterisk reminds us that (4.9) is being respected by the sum.

Remark 4.5. The task is to show that the t-integral in (4.10) extends analytically from large positive values of $\Re e R$ to a punctured neighbourhood of R = 0, such that the residue can be computed. For this, we need the specific structure of the multipliers u_v .

Let us define $\xi_{01} = \bar{\xi}_{02} = 1$, $\bar{\xi}_{01} = \xi_{02} = \xi_{11} = \bar{\xi}_{11} = -1$, and $\xi_{12} = \bar{\xi}_{12} = 0$. There exist positive, continuous, functions a_{ij} and \bar{a}_{ij} such that

$$|K_{ij}(t;z)| \le a_{ij}(z)e^{\xi_{ij}\gamma|t|} \quad \text{and} \quad |\bar{K}_{ij}(t;z)| \le \bar{a}_{ij}(z)e^{\bar{\xi}_{ij}\gamma|t|} \tag{4.11}$$

hold on $\{(t,z) \in \mathbb{R} \times \mathbb{C} \mid |z| = 1 \text{ and } ze^{\gamma t} \neq \pm i\}$. Moreover,

$$a_{ij}\bar{a}_{ij} \le \frac{A}{\gamma^{1-i}},\tag{4.12}$$

where the constant A is independent of γ . Notice that, in the Kronecker delta notation,

$$\xi_{ij} + \bar{\xi}_{ij} = -2\,\delta_{i1}\delta_{j1} \le 0. \tag{4.13}$$

Emphasizing the node in question, instead of specifying the subindices ij we write ξ_v and $\bar{\xi}_v$. For each node v, we define the numbers

$$r_v := \sup_{(z,\theta) \in B} \min \left\{ k \in \mathbb{Z} \mid \forall \, \delta > 0 : u_v(t;z,\theta) e^{-(k+\delta)\gamma t} \to 0 \text{ as } t \to \infty \right\}$$

with

$$B := \{ |z| = 1, |\arg z| \le \vartheta \} \times \{ |\Im \mathfrak{m} \theta| \le \sigma \}. \tag{4.14}$$

Thus, recalling (4.7), we are inside the analyticity domain $\mathbb{U}_{\tau,\vartheta} \times \{|\Im \mathfrak{m} \theta| \leq \sigma\}$ of Lemma 2.2. These numbers measure the divergence rate of the multipliers u_v in the limit $t \to \infty$.

Lemma 4.6. Ordered according to the list of possible u_v 's on p. 21,

$$r_{v} = \begin{cases} \bar{\xi}_{v} + k_{v}, & case (1), \\ \bar{\xi}_{v} + (m_{v} - m'_{v}) + \sum_{w'=v} \xi_{w}, & case (2), \\ 0, & case (3), \\ \bar{\xi}_{F} + (m_{v} - m'_{v}) + \sum_{w'=v} \xi_{w}, & case (4), \end{cases}$$

where $\bar{\xi}_F \in \{0,1\}$ depends on the choice of F in (3.15). Moreover $u_v(z,\theta)$, see (4.7), is analytic on $\{|\Im \mathfrak{m}\,\theta| \leq \sigma\}$ with respect to θ and on $\mathbb{U}_{\tau,\theta}\setminus\{0\}$ with respect to z. It is also analytic in the punctured neighbourhood $\{|z| \geq \tau^{-1}\}$ of $z = \infty$, at which point there is a (possible) pole of order r_v (if $r_v > 0$).

Proof. The maps $\Phi^0(z)$ and $f(\Phi^0(z), \theta)$ are analytic in these domains, without singularities at $z = 0, \infty$; see (1.3). The rest follows by staring at the expression of u_v in each case. \square

Starting from the end nodes of a tree—setting $s_v := 0$ for them—we recursively define

$$s_v := \sum_{w'=v} \max(0, n_w) \text{ and } n_v := r_v + s_v.$$
 (4.15)

These numbers measure the divergence rate of (sub)trees in the limit $t \to \infty$:

Lemma 4.7. For T_v as in (4.6), and any $\delta > 0$, the estimates

$$T_v(t; z, \theta)e^{-(n_v + \delta)\gamma t} \to 0 \quad as \quad t \to \infty,$$
 (4.16)

$$\left(\prod_{w'=v} \int_{-\infty}^{t} T_w(\tau; z, \theta) d\tau\right) e^{-(s_v + \delta)\gamma t} \to 0 \quad as \quad t \to \infty$$
(4.17)

hold true in the region B of (4.14).

Proof. Assume that (4.16) is true for each successor w of v (w' = v). By (4.6), we only need to observe that, given $\delta > 0$, $\left| \int_{-\infty}^{t} T_w(\tau; z, \theta) d\tau \right| \le C_1 + C_2 \int_0^t e^{(n_w + \delta)\gamma\tau} d\tau = \mathcal{O}(e^{[\max(0, n_w) + \delta]\gamma t})$ in the limit $t \to \infty$. For then (4.15) implies (4.17) for the node v. Now, (4.16) follows from the definition of r_v . Since (4.16) holds for end nodes ($n_v = r_v$), induction proves the claim. \square

Recalling Remark 4.5, let us now come back to the integral in (4.10), setting $v = v_0$, i.e.,

$$\int_{0}^{\infty} e^{-Rt} u_{v}(t; z, q_{v}) \prod_{w'=v} \int_{-\infty}^{t} T_{w}(\tau; z, \tilde{q}_{w}) d\tau dt.$$
 (4.18)

An obvious problem is the exponential divergence of the integrand; see (4.16) and (4.17). Our cure is the following. Since u_v is meromorphic at $z = \infty$, by Lemma 4.6, we may expand

$$u_v(z,\theta) = \sum_{k=-s}^{r_v} z^k u_{v,k}(\theta) + u_{v,<-s}(z,\theta)$$
(4.19)

for any integer $s \ge -r_v$, with $u_{v,<-s}(z,\theta) = \mathcal{O}(z^{-s-1})$ as $z \to \infty$. Extending (4.7), we write

$$u_{v,k}(t;\theta) \equiv t^p e^{k\gamma t} u_{v,k}(\theta + \omega t)$$
 and $u_{v,<-s}(t;z,\theta) \equiv t^p u_{v,<-s}(ze^{\gamma t},\theta + \omega t),$

where $p \in \{0,1\}$ depends on the choice of F in case (4) on p. 21. In particular, the integral $\int_0^\infty e^{-Rt} u_{v,<-s_v}(t;z,q_v) \prod_{w'=v} \int_{-\infty}^t T_w(\tau;z,\tilde{q}_w) d\tau dt$ is convergent for $\Re \epsilon R > -\gamma$, by virtue of (4.17), and can be estimated on a circle $|R| = \rho < \gamma$ for the purpose of Lemma 3.13.

The rest of the integral (4.18) is integrated by parts: for $-s_v \le k \le r_v$ and sufficiently large positive values of $\Re e R$,

$$\int_{0}^{\infty} e^{-Rt} z^{k} u_{v,k}(t; q_{v}) \prod_{w'=v} \int_{-\infty}^{t} T_{w}(\tau; z, \tilde{q}_{w}) d\tau dt$$

$$= \frac{z^{k} u_{v,k}(q_{v})}{R_{v}} \left\{ \frac{E_{v}(0)}{R_{v}^{p}} + \sum_{p'=0}^{p} \frac{1}{R_{v}^{p-p'}} \int_{0}^{\infty} e^{-R_{v}t} t^{p'} \frac{dE_{v}}{dt} dt \right\},$$
(4.20)

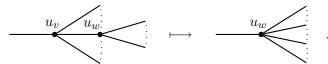
where $R_v := R - k\gamma - i\omega \cdot q_v$ and $E_v(t) := E_v(t; z, \{\tilde{q}_w \mid w' = v\}) := \prod_{w'=v} \int_{-\infty}^t T_w(\tau; z, \tilde{q}_w) d\tau$. But $\frac{dE_v}{dt} = \sum_{w'=v} T_w(t) \prod_{\bar{w}'=v, \bar{w} \neq w} \int_{-\infty}^t T_{\bar{w}}(\tau) d\tau$, so that the recursion relation (4.6) yields

$$\frac{dE_v}{dt} = \sum_{w'=v} \sum_{q_w + \sum_{\bar{w}'=w} \tilde{q}_{\bar{w}} = \tilde{q}_w} u_w(t;z,q_w) \bigg(\prod_{\bar{w}'=w} \int_{-\infty}^t T_{\bar{w}}(\tau;z,\tilde{q}_{\bar{w}}) \, d\tau \bigg) \bigg(\prod_{\substack{\bar{w}'=v \\ \bar{w} \neq w}} \int_{-\infty}^t T_{\bar{w}}(\tau;z,\tilde{q}_{\bar{w}}) \, d\tau \bigg).$$

We now fix w, and collect all the Fourier sums together so that they run over the set specified by $q_v + q_w + \sum_{\bar{w}' = w} \tilde{q}_{\bar{w}} + \sum_{\bar{w}' = v, \bar{w} \neq w} \tilde{q}_{\bar{w}} = q$. Again, without affecting the ℓ th order, we also impose the finiteness condition $|q_v| + |q_w| + \sum_{\bar{w}' = w} |\tilde{q}_{\bar{w}}| + \sum_{\bar{w}' = v, \bar{w} \neq w} |\tilde{q}_{\bar{w}}| \leq \ell N$ similar to (4.9), which allows us to bring the sum out of the t-integral. We observe that the remaining integral in (4.20) produces integrals similar to the original (4.18), with the following changes: (i) u_v changes to u_w , (ii) R changes to R_v , (iii) P changes to P, and (iv) the integral on the line leaving v, v changes to v c

Thus, a single step in the integration-by-parts scheme can be described in terms of trees as follows: Given a tree, consider one of the successors, w, of the root, v = w'. The line from w to v is "contracted" by erasing the node w, reattaching to v all subtrees originally entering w,

and replacing the multiplier of the root, v, by u_w . Finally, we rename the root w. Pictorially,



4.4. Integration by parts: exhausting the entire tree. We can proceed recursively and repeat the procedure in the previous subsection until there are no nodes left in the tree. Start with a tree T_v having root v. First, call $T_0 := T_v$ and set $v_0 := v$. Then choose a successor w of v and define $v_1 := w$, contracting the line from w to v. Next, in the new tree (the rightmost diagram above) called T_1 , choose a successor of w and call it v_2 . Contract the line from v_2 to v_1 . Repeat until the tree has been exhausted and all nodes have been numbered. We can express the sequence of trees formed as

$$T_i(t) = u_{v_i} \prod_{w \in T_i: w' = v_i} \int_{-\infty}^t T_w(\tau) d\tau,$$
 (4.21)

where in the product we consider the tree T_i with the root v_i having the successors w. Let us define the numbers, the divergence rates of the multipliers u_{v_i} ,

$$r_i := r_{v_i}$$
.

Analogously to the numbers s_v in (4.15), we set

$$s_i := \sum_{w \in T_i: w' = v_i} \max(0, n_w).$$

The numbers n_w are the ones defined in (4.15) for the original tree, T_0 . Notice that, although $s_0 = s_{v_0}$, s_i is not simply equal to s_{v_i} , but is the analogue in the tree T_i of which v_i is the root. Similarly to (4.17), s_i bounds the divergence rate of the product in (4.21): if $\delta > 0$,

$$\left(\prod_{w \in T_i: w'=v_i} \int_{-\infty}^t T_w(\tau) d\tau\right) e^{-(s_i + \delta)\gamma t} \to 0 \quad \text{as} \quad t \to \infty.$$
 (4.22)

The following, elementary, property is useful.

Lemma 4.8. For every i, $s_{i+1} + r_{i+1} \leq s_i$.

Proof. Since $r_{i+1} := r_{v_{i+1}}$, $n_{v_{i+1}} := r_{v_{i+1}} + s_{v_{i+1}}$, and $\sum_{w \in T_i: w' = v_{i+1}} \max(0, n_w) = s_{v_{i+1}}$, and

$$s_{i+1} := \sum_{w \in T_{i+1}: w' = v_{i+1}} \max\left(0, n_w\right) = \sum_{w \in T_i: w \neq v_{i+1}, w' = v_i} \max\left(0, n_w\right) + \sum_{w \in T_i: w' = v_{i+1}} \max\left(0, n_w\right),$$

we get $s_{i+1} + r_{i+1} = \sum_{w \in T_i: w \neq v_{i+1}, w' = v_i} \max(0, n_w) + n_{v_{i+1}} \leq s_i$. Here we used the fact that $v'_{i+1} = v_i$ in the tree T_i .

Let $\tilde{k}_0 := 0$ and $Q_0 := 0$. Suppose that at the *i*th step we are considering the integral

$$\int_0^\infty e^{(iQ_i \cdot \omega + \tilde{k}_i \gamma - R)t} t^{p_i} u_{v_i}(ze^{\gamma t}, q_{v_i}) e^{iq_{v_i} \cdot \omega t} \prod_{w \in T_i: w' = v_i} \int_{-\infty}^t T_w(\tau; z, \tilde{q}_w) d\tau dt;$$

cf. (4.18). Then the integration-by-parts procedure described above takes us at the (i + 1)st step to a similar integral with i replaced by i + 1, defining

$$\tilde{k}_{i+1} := \tilde{k}_i + k_i \quad \text{and} \quad Q_{i+1} := Q_i + q_{v_i},$$
(4.23)

multiplied by the factor

$$\frac{z^{k_i} u_{v_i,k_i}(q_{v_i})}{\left(R - \tilde{k}_{i+1} \gamma - i\omega \cdot Q_{i+1}\right)^{1+p_i - p_{i+1}}}.$$
(4.24)

The integer indices p_{i+1} and k_i can assume the values

$$0 \le p_{i+1} \le p_i \le 1$$
 and $k_i = -\tilde{k}_i - s_i, \dots, r_i$. (4.25)

With the aid of Lemma 4.8, it is straightforward to see that

$$0 \le \tilde{k}_i + r_i + s_i \le n_{v_0},\tag{4.26}$$

where $n_{v_0} = s_{v_0} + r_{v_0} = s_0 + r_0$ is the number describing the divergence rate of the original tree, T_0 , in the sense of (4.16). The number of possible values of k_i above is thus at most $1 + n_{v_0}$. This implies the two-sided (equivalent) bounds

$$r_i - n_{v_0} \le k_i \le r_i$$
 and $0 \le r_i - k_i \le n_{v_0}$.

We continue recursively until there are no nodes—alternatively, integrals that require regularization—left. The result is a sum with terms of two different species. We set

$$R_i := R - \tilde{k}_i \gamma - i\omega \cdot Q_i,$$

in order to make the presentation more compact, and also define (cf. E_v in (4.20))

$$E_j(t) := E_j(t; z, \{\tilde{q}_w \mid w \in T_j, w' = v_j\}) := \prod_{w \in T_j: w' = v_j} \int_{-\infty}^t T_w(\tau; z, \tilde{q}_w) d\tau,$$

with the understanding that $E_{|V_{\rm int}|-1}\equiv 1$. The first class of terms is

$$\left(\prod_{i=0}^{j-1} \frac{z^{k_i} u_{v_i, k_i}(q_{v_i})}{R_{i+1}^{1+p_i-p_{i+1}}}\right) \int_0^\infty e^{(iq_{v_j} \cdot \omega - R_j)t} t^{p_j} u_{v_j, < -\tilde{k}_j - s_j}(ze^{\gamma t}, q_{v_j}) E_j(t) dt, \tag{4.27}$$

for $0 < j < |V_{\rm int}|$. Notice that $s_{|V_{\rm int}|-1} = 0$. Second, there are the terms $(0 < j < |V_{\rm int}|)$

$$\left(\prod_{i=0}^{j-1} \frac{z^{k_i} u_{v_i, k_i}(q_{v_i})}{R_{i+1}^{1+p_i-p_{i+1}}}\right) \frac{z^{k_j} u_{v_j, k_j}(q_{v_j})}{R_{j+1}^{1+p_j}} E_j(0).$$
(4.28)

4.5. **Estimates.** By the definition of s_i , we can bound

$$|E_j(t)| \le C_{E_j} e^{(s_j + \frac{1}{4})\gamma t} e^{-\sigma \sum |\tilde{q}_w|} \qquad (t \ge 0),$$

for |z|=1 with $|\arg z|\leq \vartheta$. Deferring the proof to Appendix A, we formulate

Lemma 4.9. The coefficients $u_{v_i,k_i}(q)$ satisfy

$$|u_{v_i,k_i}(q_{v_i})| \le C_{v_i} \tau^{k_i - r_i} e^{-\sigma |q_{v_i}|},$$

where τ is as in Lemma 4.6. On $\{1 \le |z| \le \tau^{-1} \text{ with } |\arg z| \le \vartheta\} \bigcup \{|z| \ge \tau^{-1}\},$

$$|u_{v_i < -\tilde{k}_i - s_i}(z, q_{v_i})| \le C_{v_i}(\tau|z|)^{-\tilde{k}_i - s_i - 1} 2^{n_{v_0}} \tau^{-r_i} e^{-\sigma|q_{v_i}|}.$$

When |z| = 1, $|\arg z| \le \vartheta$, and $|\Im \mathfrak{m} \theta| \le \sigma$ it holds true, for suitable $\bar{r}_i \in \mathbb{Z}$, that

$$|u_{v_i}(ze^{\gamma t}, \theta + \omega t)| \le C_{v_i} \begin{cases} e^{\bar{r}_i \gamma t}, & t < 0, \\ e^{r_i \gamma t}, & t \ge 0. \end{cases}$$

Recalling (4.26), the integral in (4.27) is bounded by $(2\tau^{-1})^{n_{v_0}+1}\frac{C_{v_j}C_{E_j}}{\gamma^{1+p_j}}e^{-\sigma|q_{v_j}|-\sigma\sum|\tilde{q}_w|}$, when $\Re \epsilon R \ge -\frac{1}{2}\gamma$, |z|=1, and $|\arg z|\le \vartheta$. Thus, the bounds on (4.27) and (4.28) read

$$\left(\prod_{i=0}^{j-1} \frac{C_{v_i}}{|R_{i+1}^{1+p_i-p_{i+1}}|}\right) \frac{C_{v_j}C_{E_j}}{\gamma^{1+p_j}} (2\tau^{-1})^{n_{v_0}+1} \tau^{\sum_{i=0}^{j-1} (k_i-r_i)} e^{-\sigma \sum_{i=0}^{j} |q_{v_i}| -\sigma \sum |\tilde{q}_w|}$$
(4.29)

and

$$\left(\prod_{i=0}^{j-1} \frac{C_{v_i}}{|R_{i+1}^{1+p_i-p_{i+1}}|}\right) \frac{C_{v_j}C_{E_j}}{|R_{j+1}^{1+p_j}|} \tau^{\sum_{i=0}^{j} (k_i - r_i)} e^{-\sigma \sum_{i=0}^{j} |q_{v_i}| - \sigma \sum |\tilde{q}_w|}, \tag{4.30}$$

respectively. Here $e^{-\sigma \sum_{i=0}^{j} |q_{v_i}| - \sigma \sum_{i=0}^{j} |q_{w_i}|} \leq e^{-\sigma |q|}$, because the tree is a Fourier transform with index q, and the q_{v_i} , \tilde{q}_w come from convolutions. The factors of τ are controlled by

Lemma 4.10. For all j, $\sum_{i=0}^{j} (k_i - r_i) \ge -n_{v_0}$. Moreover, $n_{v_0} \le d(T_{v_0})$ (see (4.3)).

Proof. From (4.23) and (4.25), we clearly have $\tilde{k}_{i+1} \geq -s_i$. Then,

$$\sum_{i=0}^{j} (k_i - r_i) = \tilde{k}_{j+1} - \sum_{i=0}^{j} r_i \ge -s_j - r_j - \sum_{i=0}^{j-1} r_i \ge -s_{j-1} - \sum_{i=0}^{j-1} r_i \ge \cdots \ge -s_0 - r_0 = -n_{v_0},$$

with the aid of Lemma 4.8. The estimate $n_{v_0} \leq d(T_{v_0})$ follows easily from (4.3) and (4.13). \square

By the last bound of Lemma 4.9 it is clear that $C_{E_j} = \left(\prod_{i=j+1}^{|V_{\rm int}|-1} C_{v_i}\right) \widetilde{C}_{E_j}$. In fact, each integral in E_j produces a γ^{-1} in the upper bound, and we get $\widetilde{C}_{E_j} \leq (C\gamma^{-1})^{|V_{\rm int}|-j-1}$. Starting from the explicit expressions of the multipliers, we estimate $\prod_{i=0}^{|V_{\rm int}|-1} C_{v_i} \leq (Cg)^{|V|} (C\epsilon)^{\mathrm{d}(T_{v_0})}$ with the aid of (4.11) and (4.12). Here V, defined above (4.6), is the set of all nodes including the black dot (\bullet) end nodes. Since $|V_{\rm int}| \leq |V| \leq 2 \operatorname{n}(T_{v_0}) \leq 4 \left(\mathrm{d}(T_{v_0}) - 1\right)$,

$$\left(\prod_{i=0}^{j} C_{v_i}\right) C_{E_j} \le g^{|V|+1+j-|V_{\text{int}}|} (C\epsilon)^{\operatorname{d}(T_{v_0})} \qquad (0 < j < |V_{\text{int}}|).$$

Notice that the fraction $1/R_i$ is analytic on the domain where $|R| \neq |\tilde{k}_i \gamma + i\omega \cdot Q_i|$. Thus, all the $1/R_i$ are analytic in the punctured neighbourhood of the origin

$$0 < |R| \le \rho := \frac{1}{2} \min \{ \{ \gamma \} \cup \{ |\omega \cdot Q_i| \mid Q_i \ne 0 \} \}.$$

On the circle $|R| = \rho$, $|R_i^{1+p}|^{-1} \le \rho^{-1-p}$, provided $p \ge -1$. According to (4.25),

$$\prod_{i=0}^{j-1} \left| \frac{1}{R_{i+1}^{1+p_i - p_{i+1}}} \right| \le \rho^{-j - p_0 + p_j} \qquad (0 < j < |V_{\text{int}}|).$$

Subsequently, summing over j the bounds on (4.29) and (4.30) we get

$$\sum_{j=1}^{V_{\text{int}}|-1} (C\epsilon)^{\mathrm{d}(T_{v_0})} g^{|V|+1+j-|V_{\text{int}}|} \rho^{-j-1-p_0} e^{-\sigma|q|} \le (C\epsilon)^{\mathrm{d}(T_{v_0})} (g/\rho)^{|V|+1} g^{-1} e^{-\sigma|q|} =: B(g) e^{-\sigma|q|}.$$

At the ℓ th order, $|Q_i| \leq \ell N$ for each i, and thus $\rho \geq \frac{1}{2} \min(\gamma, a(\ell N)^{-\nu})$.

From the above it is also clear that the order of the pole at R=0 in our integrals does not exceed $|V_{\rm int}|+1$. For the purposes of Lemma 3.13, we compute (using $g, \gamma \leq C$)

$$B(g) \sum_{J=0}^{|V_{\text{int}}|+1} \frac{1}{J!} \left(\frac{\rho \vartheta}{g}\right)^{J} \le CB(g) \le (C_N \epsilon)^{\operatorname{d}(T_{v_0})} g^{-1} (\ell!)^{4\nu}.$$

There are at most $(4\ell)! \leq 4^{4\ell}(\ell!)^4$ orders in which we can exhaust all of the up to 4ℓ lines of a tree contributing at order ℓ .

4.6. Remaining integrals. The integral over \mathbb{R}_{-} in (4.1) is simple, because the integrand satisfies the identity (3.16) and has the analyticity properties stated below that equation. In other words, we can separate the (possible) pole and the constant term from the rest: writing $F(X^u; 0, z, \theta) \equiv F(z, \theta),$

$$F(X^u;t,z,\theta) = t^p z^{-1} e^{-\gamma t} F_{-1}(\theta + \omega t) + t^p F_0(\theta + \omega t) + t^p \delta_1 F(ze^{\gamma t}, \theta + \omega t).$$

Here p=0,1. Applying $\int_{-\infty}^{0}$ on each term separately, we get that $\int_{-\infty}^{0} F(X^{u};t,z,\theta)$ equals

$$(-1)^{p} \sum_{q'} \frac{z^{-1} e^{iq' \cdot \theta} \hat{F}_{-1}(q')}{(iq' \cdot \omega - \gamma)^{1+p}} + (-1)^{p} \sum_{q' \neq 0} \frac{e^{iq' \cdot \theta} \hat{F}_{0}(q')}{(iq' \cdot \omega)^{1+p}} + \int_{-\infty}^{0} t^{p} \delta_{1} F(ze^{\gamma t}, \theta + \omega t) dt.$$

These terms are small compared to to the large bounds obtained for the f_0^{∞} part above. We also have to study the coefficients c_{ij}^p appearing in (3.12). This is most conveniently done in terms of the representation (3.13). Since $M^u = \mathcal{O}(\epsilon)$, only

$$\int_{-\infty}^{-t} \left[f_{\psi\varphi}^u (B^u M^u)^l K_{ij}^p \right] (\tau) \, d\tau$$

with $l \leq \ell$ can contribute to c_{ij}^p at order ℓ . We only need to consider $t \geq 0$, as $\Upsilon(t) = \Upsilon(-t)$. The above integral consists in obvious shorthand notation, through (3.6), of 4^l terms like

$$\int_{-\infty}^{-t} (-t-\tau_l)^{p_{l+1}} (fK_l)(\tau_l) \int_{-\infty}^{\tau_l} (\tau_l-\tau_{l-1})^{p_l} (\bar{K}_l M)(\tau_{l-1}) \cdots K_1(\tau_1) \int_{-\infty}^{\tau_1} (\tau_1-\tau_0)^{p_1} (\bar{K}_1 M)(\tau_0) \tau_0^{p_0} K_0(\tau_0),$$

where $p_{l+1} = 0$. If $p_0 = 1$ (instead of 0), we use $\tau_0 = (\tau_0 - \tau_1) + \cdots + (\tau_{l-1} - \tau_l) - (-t - \tau_l) - t$, getting l+2 terms of the original form except that $p_0=0$ and either there is a factor t or precisely one change $p_i \mapsto p_i + 1$ occurs. Due to $M = \mathcal{O}(\epsilon g^2)$, (4.11) and (4.12), each $K_i \bar{K}_i M$ produces a factor $CA|\epsilon|g^{p_i+1}$ to the upper bound, whereas analyticity yields $e^{-\sigma|q_i|}$ (which we prefer although at each order we reduce to trigonometric polynomials); $\prod_{i=0}^l e^{-\sigma|q_i|} \leq e^{-\sigma|q|}$. Last, the integral involves a total of $\sum_{i=1}^{l+1} (1+p_i) \leq 2l+3$ denominators of the form $\tilde{k}_i \gamma + i\omega \cdot Q_i$, where $\tilde{k}_i \in \mathbb{Z}$ and $Q_i := \sum_{j=0}^i q_j$ with $Q_l = q \neq 0$. We bound these by

$$\left(\frac{1}{\min\{\{\gamma\} \cup \{|\omega \cdot Q_i| \mid Q_i \neq 0\}\}}\right)^{\sum_{i=1}^{l+1}(1+p_i)} \leq \left(\max(\gamma^{-1}, a^{-1}(N\ell)^{\nu})\right)^{\sum_{i=1}^{l+1}(1+p_i)}.$$

Altogether, the bounds above yield $C_N^{\ell} |\epsilon|^l g^{-4} e^{-\sigma|q|} (\ell!)^{2\nu}$.

Skipping further details, this is the upper bound on the integral. It is smaller than what was derived for the integral $\int_{-\infty}^{\infty} F$ —and thus for \mathcal{E}_{ij}^p in (3.12)—above. In conclusion, among the contributions $(c_{ij}^p)^{\ell_1} (\mathcal{E}_{ij}^p)^{\ell_2} (\ell_1 + \ell_2 = \ell - 1)$ to Υ^{ℓ} , $(\mathcal{E}_{ij}^p)^{\ell - 1}$ is the most dangerous one.

Remark 4.11. Above, the sums over q_i , with $\sum_i q_i = q$ and $\sum_i |q_i| \leq \ell N$ were dealt with as follows. Since the analyticity domain with respect to θ is the compact $\{|\Im \mathfrak{m}\,\theta|\leq\sigma\}$, it can be substituted by some $\{|\Im \mathfrak{m} \theta| \leq \sigma'\}$ with $\sigma < \sigma'$. Then, for every i we actually have the factor $e^{-\sigma'|q_i|}$ in the estimates above. While $e^{-\sigma\sum_i|q_i|} \leq e^{-\sigma|q|}$, we get rid of the sums over q_i :

$$\sum_{q_i} e^{-(\sigma' - \sigma)|q_i|} \le C_{\sigma' - \sigma}.$$

A similar remark applies to the Fourier indices in Subsection 4.5.

5. Discussion

Each term in the asymptotic expansion (3.12) of the splitting matrix Υ is proportional to an integral of the form $\int_{-\infty}^{\infty} \partial_{\theta} F(X^u; \tau, z, \theta) d\tau$. These we showed to be exponentially small in the limit $g \to 0$ at all orders ℓ , by extending the integrands analytically into a wedge $\{|\arg z| \le \vartheta\}$ on the complex plane and then shifting the contour of integration. A key point is that we had to do this *order by order*, because the series $\sum_{\ell=0}^{\infty} \epsilon^{\ell} X^{u,\ell}(z,\theta)$ is not expected to converge for large values of |z| if ϵ is fixed.

The large powers of the factorial $\ell!$, associated with the regularized integrals, are produced by accumulation of poles at the origin in the R plane. To some extent the factorials are artifact, as is shown by the following simple example. In order to study, say, the integral

$$\int_0^\infty u(\theta + \omega t) \int_{-\infty}^t v(\theta + \omega \tau) z e^{\gamma \tau} d\tau dt,$$

we have to show that the sums

$$\sum_{p} \hat{u}(p)\hat{v}(q-p) \int_{0}^{\infty} e^{(ip\cdot\omega - R)t} \int_{-\infty}^{t} e^{(i(q-p)\cdot\omega + \gamma)\tau} d\tau dt$$

extend analytically to a (punctured) neighbourhood of the origin, R = 0. We integrate by parts, just as in (4.20) when extending analytically the tree integral (4.10), and get

$$\sum_{p} \hat{u}(p)\hat{v}(q-p)\frac{1}{R-ip\cdot\omega} \left\{ \frac{1}{(i(q-p)\cdot\omega+\gamma)} + \frac{1}{R-(\gamma+iq\cdot\omega)} \right\}.$$

Here the pole at $ip \cdot \omega$ gets arbitrarily close to the origin, unless we restrict p somehow—for instance, by considering trigonometric polynomials. Either by simplification, or by computing the same expression directly by starting from the inner integral, we obtain

$$\sum_{p} \hat{u}(p)\hat{v}(q-p) \frac{1}{i(q-p)\cdot\omega + \gamma} \cdot \frac{1}{R - (\gamma + iq \cdot \omega)}.$$

In the latter form there is no problem; the pole has cancelled. Of course, this is a naive example and in general it is hard to see whether a given pole popping out of the integration-by-parts procedure should really be there.

Also the coefficients c_{ij}^p appearing in Proposition 3.12 produced large powers of ℓ !. Even though integration by parts was not exploited, the source of the factorials was again the accumulation of poles at the origin in the R plane. In both this case and the previous, the type of "divergence" is very similar to what is encountered in KAM theory. There repeated resonances, or arbitrarily many occurrences of the operator \mathcal{D}^{-1} in convolutions, ruin absolute convergence of the Fourier–Taylor expansion of a solution by producing high powers of the factorial ℓ !. On the other hand, the state of affairs can be cured by well-known resummations, as in [Gal94]. Such resummations still escape us in the context of homoclinic splitting.

APPENDIX A. SOME COMPUTATIONS AND PROOFS

Proof of Theorem 3. It is known that the whiskers are Lagrangian manifolds [Eli94, DG00, LMS03]; they are graphs over the angles and, moreover, the actions are gradients of potentials. In other words, there exist generating functions $S^{u,s}$ such that

$$I^{u,s} \equiv \partial_{\phi} S^{u,s}(\phi, \psi)$$
 and $A^{u,s} \equiv \partial_{\psi} S^{u,s}(\phi, \psi)$.

We consider the latter row vectors. Differentiating both sides of

$$\mathcal{H}^{u,s} := \mathcal{H}(\phi, \psi, \partial_{\phi} S^{u,s}, \partial_{\psi} S^{u,s}) \equiv E$$

with respect to the *i*th component of $(\phi, \psi) =: \varphi = (\varphi_0, \varphi_1, \dots, \varphi_d)$ yields

$$0 = \mathcal{H}_i^{u,s} + \sum_{j=0}^d \mathcal{H}_{d+1+j}^{u,s} \, \partial_{\varphi_i} \partial_{\varphi_j} S^{u,s} = \mathcal{H}_i^{u,s} + \sum_{j=0}^d (\partial_{\varphi}^2 S^{u,s})_{ij} \partial_{J_j} \mathcal{H}^{u,s}. \tag{A.1}$$

Here J := (I, A) and $\partial_{J_j} \mathcal{H}^{u,s}$ means $\mathcal{H}^{u,s}_{d+1+j} = \partial_{J_j} \mathcal{H}(\phi, \psi, \partial_{\phi} S^{u,s}, \partial_{\psi} S^{u,s})$.

From (A.1) it follows that, on the homoclinic trajectory,

$$\partial_J \mathcal{H}^{u,s} \,\partial_{\varphi}^2 (S^u - S^s) = 0 \tag{A.2}$$

because one has $\partial_{\varphi}S^{u}=\partial_{\varphi}S^{s}$ and therefore also $\mathcal{H}_{i}^{u}=\mathcal{H}_{i}^{s}$ for $i=0,\ldots,2d+1$. For our particular Hamiltonian, $\partial_{J}\mathcal{H}^{u,s}=(I,A)^{u,s}$, where we now drop the superscripts u,s and evaluate everything at the homoclinic point

$$\varphi = (\phi, \psi) = (0, \omega t) + X^u(e^{\gamma t}, \omega t) = (0, \omega t) + X^s(e^{\gamma t}, \omega t)$$
(A.3)

below. The ψ component of (A.2) reads

$$I \partial_{u} \partial_{\phi} (S^u - S^s) + A \partial_{u}^2 (S^u - S^s) = 0,$$

or, recalling $I = I^0 + \mathcal{O}(\epsilon) = 2g/\cosh gt + \mathcal{O}(\epsilon) \neq 0$ (taking $\tilde{\epsilon} := g^{-1}e^{g|t|}\epsilon$ small),

$$\partial_{\psi}\partial_{\phi}(S^u - S^s) = -I^{-1}A\,\partial_{\psi}^2(S^u - S^s). \tag{A.4}$$

Here A is the row vector (A_1, \ldots, A_d) , such that both sides of the equality are row vectors.

Next, let us perform the coordinate transformations $(\phi, \psi) = F^{u,s}(z, \theta) \mapsto (z, \theta)$. To this end, we observe that the splitting vector satisfies

$$\Delta^{T}(\phi, \psi) = A^{u} - A^{s} = \partial_{\psi}(S^{u} - S^{s})(\phi, \psi).$$

Then, the θ derivative of the column vector Δ is the square matrix

$$\partial_{\theta} \Delta = \partial_{\psi} \Delta \partial_{\theta} \psi + \partial_{\phi} \Delta \partial_{\theta} \phi$$

by the chain rule. In particular, by (A.4),

$$\partial_{\phi}\Delta = \left[\partial_{\psi}\partial_{\phi}(S^{u} - S^{s})\right]^{T} = -\partial_{\psi}^{2}(S^{u} - S^{s})I^{-1}A^{T} = -\partial_{\psi}\Delta I^{-1}A^{T}$$
(A.5)

holds at a homoclinic point (see (A.3)), such that

$$\partial_{\theta} \Delta = \partial_{\psi} \Delta \left(\partial_{\theta} \psi - I^{-1} A^{T} \partial_{\theta} \phi \right).$$

The matrix $M := \partial_{\theta} \psi - I^{-1} A^{T} \partial_{\theta} \phi$ has the asymptotic expression

$$M = \left(\mathbb{1} + \mathcal{O}(\epsilon)\right) - \left(\frac{2g}{\cosh at} + \mathcal{O}(\epsilon)\right)^{-1} \cdot \left(\omega + \mathcal{O}(\epsilon)\right)^{T} \cdot \mathcal{O}(\epsilon)$$

with $\epsilon = ge^{-g|t|}\tilde{\epsilon}$, as before, and $\tilde{\epsilon}$ small.

Proof of Lemma 3.1. (Adapted from [Sau01]). The Fourier transform of (3.2) yields

$$\hat{F}(s,q) = \hat{F}(0,q) e^{-i(g^{-1}\omega \cdot q)s},$$
(A.6)

which is entire in s. But $|\hat{F}(0,q)| e^{(g^{-1}\omega \cdot q) \Im \mathfrak{m} s} = |\hat{F}(s,q)| \le B(g)e^{-\eta|q|}$ for $s \in [-i\vartheta, i\vartheta]$, and

$$|\hat{F}(0,q)| \le B(g)e^{-\vartheta g^{-1}|\omega \cdot q| - \eta|q|}$$
.

Finally, plugging this into (A.6), we get

$$|\hat{F}(s,q)| \le B(g)e^{-(\vartheta - |\Im \mathfrak{m} s|)g^{-1}|\omega \cdot q| - \eta|q|}$$
 $(s \in \mathbb{C}).$

For $|\Im \mathfrak{m} s| \leq \vartheta$ and $|\Im \mathfrak{m} \theta| \leq \eta' < \eta$, the series $\sum_{q \in \mathbb{Z}^d} \hat{F}(s,q)e^{iq\cdot\theta}$ is uniformly convergent and, as such, provides the analytic extension.

Since $\alpha x^{-\nu} + \beta x \ge \alpha(\nu+1) \left(\frac{\alpha\nu}{\beta}\right)^{-\nu/(\nu+1)}$ for positive α , β , ν , and x, we get from the Diophantine condition (1.4) that

$$|\hat{F}(s,q)| \le B(g)e^{-\vartheta g^{-1}a|q|^{-\nu} - \eta|q|} \le B(g)e^{-\delta|q|}e^{-w(\vartheta,\eta-\delta)g^{-1/(\nu+1)}}$$

holds if $s \in \mathbb{R}$, $q \in \mathbb{Z}^d \setminus \{0\}$, $0 < \delta < \eta$, and $w(\vartheta, \eta - \delta) := (\vartheta a)^{1/(\nu+1)} (\eta - \delta)^{\nu/(\nu+1)} (\nu + 1)\nu^{-\nu/(\nu+1)}$. Moreover, $\sum_{q \in \mathbb{Z}^d \setminus \{0\}} e^{-\delta|q|} \leq C\delta^{-d}$, where C only depends on the dimension d.

By (A.6),
$$\widetilde{F} := \langle F(s, \cdot) \rangle = \widehat{F}(s, 0) = \widehat{F}(0, 0)$$
, such that \widetilde{F} does not depend on (s, θ) .

Proof of Lemma 3.13. By shifting the contour of integration from \mathbb{R} to the complex plane by $it_q := i \operatorname{sgn}(\omega \cdot q) \vartheta g^{-1}$ units, we compute

$$\int_{-\infty}^{\infty} t^p \hat{h}(e^{gt}, q) e^{iq \cdot \omega t} dt = \underset{R=0}{\text{res}} \frac{1}{R} \int_{-\infty}^{\infty} e^{-R|t|} t^p \hat{h}(e^{gt}, q) e^{iq \cdot \omega t} dt$$

$$= \underset{R=0}{\text{res}} \frac{1}{R} \left\{ e^{-iRt_q} H_q(R) + e^{iRt_q} I_q(R) \right\} e^{-\vartheta g^{-1} |\omega \cdot q|},$$

where H_q is defined in (3.17) and $I_q(R) := \int_{-\infty}^0 e^{(iq \cdot \omega + R)t} (t + it_q)^p \hat{h}(e^{g(t+it_q)}, q) dt$. There are, a priori, two additional line integrals $\int_0^{it_q}$, but they cancel due to the residue at R = 0, as is easily checked. Because $\hat{h}(\cdot, 0) = 0$, $I_q(R)$ does not have a pole at R = 0. Hence,

$$\mathop{\rm res}_{R=0} \frac{e^{iRt_q} I_q(R)}{R} = \int_{-\infty}^{0} e^{iq \cdot \omega t} (t + it_q)^p \hat{h}(e^{g(t + it_q)}, q) \, dt.$$

If $H_q(R)$ has a pole of order k at R=0, then

$$\left| \underset{R=0}{\operatorname{res}} \frac{e^{-iRt_q} H_q(R)}{R} \right| = \left| \sum_{j=0}^k \frac{(-it_q)^j H_{q,-j}}{j!} \right| \le \sum_{j=0}^k \frac{1}{j!} \left(\frac{\rho \vartheta}{g} \right)^j \sup_{|R|=\rho} |H_q(R)|,$$

because the Laurent coefficients $H_{q,-j} := \frac{1}{2\pi} \oint_{|R|=\rho} \frac{H_q(R)}{R^{-j+1}} dR$ can be bounded from above by $|H_{q,-j}| \le \rho^j \sup_{|R|=\rho} |H_q(R)|$, whenever the circle $|R|=\rho$ is inside the domain of H_q .

Under the assumptions of the lemma, for any $0 < \delta < \sigma$ and $q \in \mathbb{Z}^d \setminus \{0\}$,

$$\left| \int_{-\infty}^{\infty} t^p \hat{h}(e^{gt}, q) e^{iq \cdot \omega t} dt \right| \leq \left[A(g) + B(g) \sum_{j=0}^{k} \frac{1}{j!} \left(\frac{\rho \vartheta}{g} \right)^j \right] e^{-\delta|q|} e^{-w(\vartheta, \sigma - \delta)g^{-1/(\nu + 1)}},$$

by mimicking the proof of Lemma 3.1. Here $w(\vartheta, \sigma - \delta) := (\vartheta a)^{1/(\nu+1)} \left(\frac{\sigma - \delta}{\nu}\right)^{\nu/(\nu+1)} (\nu + 1)$. Summation over q produces a factor $C\delta^{-d}$.

Proof of Lemma 4.9. Here we usually omit the subindices i, v, and v_0 . According to Lemma 4.6, u_v is analytic on $\{|z| \geq \tau^{-1}\} \times \{|\Im \mathfrak{m} \theta| \leq \sigma\}$ with a pole of order r at $z = \infty$. A Cauchy estimate then reads $|u_{v,k}(\theta)| \leq \tau^k \sup_{|z|=\tau^{-1}} |u_v(z,\theta)| \leq C\tau^{k-r}$. Now the first bound follows from $|u_{v,k}(q)| \leq e^{-\sigma|q|} \sup_{|\Im \mathfrak{m} \theta| < \sigma} |u_{v,k}(\theta)|$.

The second inequality is a trivial consequence of the first one for $|z| > 2\tau^{-1}$. For the other values of z, one uses (4.19) and (4.26) as well as $|z|^{-r}|u_v(z,q)| \leq Ce^{-\sigma|q|}$ to bound

$$|u_{v,<-\tilde{k}-s}(z,q)| \le |z|^{-\tilde{k}-s-1} \Big\{ |z|^{\tilde{k}+s+1} |u_v(z,q)| + Ce^{-\sigma|q|} \tau^{-\tilde{k}-s-1} \sum_{l=0}^{r+k+s} |z\tau|^{l+1} \Big\}$$

$$\le Ce^{-\sigma|q|} |z|^{-\tilde{k}-s-1} \Big\{ (2\tau^{-1})^{\tilde{k}+s+r+1} + \tau^{-\tilde{k}-s-r-1} 2^{r+\tilde{k}+s} \Big\}.$$

The last inequality in the lemma is obvious and is stated for completeness.

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