

Topological Shadowing Methods in Arnold Diffusion: Weak Torsion and Multiple Time Scales

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Abstract

Consider a symplectic map which possesses a normally hyperbolic invariant manifold of any even dimension with transverse homoclinic channels. We develop a topological shadowing argument to prove the existence of Arnold diffusion along the invariant manifold, shadowing some iterations of the inner dynamics carried by the invariant manifold and the outer dynamics induced by the stable and unstable foliations. In doing so, we generalise an idea of Gidea and de la Llave in [26], based on the method of correctly aligned windows and a so-called transversality-torsion argument. Our proof permits that the dynamics on the invariant manifold satisfy only a non-uniform twist condition, and, most importantly for applications, that the splitting of separatrices be small in certain directions and thus the associated drift in actions very slow; diffusion occurs in the directions of the manifold having non-small splitting. Furthermore we provide estimates for the diffusion time.

1 Introduction

Arnold diffusion, as first exposed in [1] has become a major subject of study for nearly integrable Hamiltonian systems. This mechanism epitomises how an integrable (“stable”) system can become unstable, with actions varying slowly (as permitted by Nekhoroshev’s theorem) but substantially, as opposed to what would happen for an integrable perturbation [7, 13]. A key idea has been to focus on a normally hyperbolic invariant cylinder. The cylinder is indeed a more robust object than individual hyperbolic invariant tori that may lie inside the cylinder and which Arnold initially used. Another key addition to Arnold’s mechanism is the random iteration of the dynamics carried by the cylinder and the “outer” dynamics obtained at the limit by following unstable and stable

leaves of the cylinder (see figure 1). Both these ideas appear in Moeckel’s work [35]. These ideas were formalised by Delshams, de la Llave and Seara with the invention of the scattering map which encodes the outer dynamics. This allowed them to solve the large gap problem, when two (primary) invariant tori are too far away for their unstable and stable manifolds to meet [14]. In these last decades, this approach has been successfully implemented by a number of authors, using either a geometric description [14, 16, 24, 25, 38, 39] or a variational one [9, 10].

Another step consists in dealing with the a priori stable case, where the normally hyperbolic cylinder appears with the perturbation itself [2, 3, 8, 10, 34]. In this article we deal with the a priori unstable case only.

One of the main difficulties lies in proving the existence of orbits shadowing random iterations of the inner and outer maps. Shadowing results can be achieved using different tools: variational methods (see, for instance, [4, 5, 6]), modern versions of the Lambda lemma [23, 27, 37], or with topological techniques such as the the correctly aligned windows method.

The advantage of the topological methods is that they require relatively little information regarding the dynamics on the normally hyperbolic invariant cylinder. In particular, no knowledge of invariant quasiperiodic tori is expected. As far as the authors know, correctly aligned windows were imagined by Conley and Easton [12], and the first application of the correctly aligned windows method to Arnold diffusion problems is the paper [26]. In this paper, Gidea and de la Llave use the method of correctly aligned windows to prove the existence of diffusion orbits in a priori unstable Hamiltonian systems and to construct orbits with an unbounded growth of energy for the Mather problem (that is a geodesic flow with a time dependent potential). This second model is usually said to be *a priori chaotic*. In both settings the normally hyperbolic invariant cylinder is two-dimensional. Moreover, they assume that the induced dynamics on the cylinder satisfies a twist property, and the twist is uniform with respect to the perturbative parameter.

The purpose of this paper is to generalise [26] from several points of view:

- The normally hyperbolic cylinder may be of any (even) dimension.

More importantly:

- The twist property satisfied by the inner dynamics may be weak, i.e. the twist may vanish when the perturbative parameter goes to 0.
- We may split the actions between two groups, faster actions and slower actions, and ignore the latter. Indeed, proving “partial transversality” of the invariant manifolds of the normally hyperbolic cylinder along the fast ones is enough to achieve drift in these directions whereas the slow directions can be treated as a black box, with no control of their instability rate.

These improvements are crucial since such behaviour is exhibited naturally in physical models, particularly in models with multiple time scales as happens often, for instance, in Celestial Mechanics.

Many of the known shadowing mechanisms rely on rather strong assumptions, both on the inner dynamics in the cylinder and on the transversality of the associated stable and unstable invariant foliations. Often such hypotheses are difficult to verify when multiple time scales are present. The results presented in this paper are quite flexible and can be applied to rather general multiple time scale settings.

In particular, consider an analytic nearly integrable Hamiltonian system with multiple time scales such that some of the angles perform fast non-resonant rotation and therefore its conjugate actions are very slow whereas some other angles are not fast. It is well known that, in the fast directions, the transversality between the invariant manifolds is exponentially small and therefore very difficult to analyze. At the same time, thanks to averaging theory, one can make the dynamics in the conjugate actions much slower. The results in this paper allow us to obtain diffusing orbits along the actions conjugated to the slow angles even if one does not have “full transversality” of the invariant manifolds, i.e. no transversality in the fast directions (see Theorem 2.11 below).

All the improvements achieved in the present papers are needed to construct diffusing orbits in the 4 Body Problem along secular resonances. In the companion paper [11], the authors prove the existence of orbits of the 4 body problem in both the planetary regime (one massive body and three bodies with small mass) and hierarchical regime (bodies increasingly separated) such that some of the bodies have osculating eccentricities and some mutual inclinations drifting with no constraint. This is the first analytical result of unstable motions in an N body problem in the planetary regime.

Secular resonances are those given by the secular angles, that is the angles which are constant for the two body problem: the argument of the perihelion and the longitude of the ascending node of each of the bodies. Such angles are much slower than the mean anomalies of the bodies which perform fast rotation. Therefore, we are exactly in the multiple time-scale setting described above. In [11], we are able to construct diffusing orbits in the actions conjugated to the secular angles, i.e. the osculating eccentricities and mutual inclinations of the bodies, without having to control the dynamics on semimajor axis directions.

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2 Set-up, Assumptions, and Results

2.1 Definitions

Let M be a C^r manifold of dimension d where $r \geq 1$. Let $F \in \text{Diff}^1(M)$, and let Λ be a submanifold of M .

Definition 2.1. We call Λ a normally hyperbolic invariant manifold for F if Λ is F -invariant, and there are

$$0 < \lambda_- < \lambda_+ < \lambda_0 < 1 < \mu_0 < \mu_- < \mu_+ \quad (1)$$

and an invariant splitting of the tangent bundle

$$T_\Lambda M = TM \oplus E^s \oplus E^u$$

such that:

$$\begin{aligned} \lambda_-^n \|v\| &\leq \|DF^n(x)v\| \leq \lambda_+^n \|v\| \text{ for all } x \in \Lambda, v \in E_x^s, n \in \mathbb{N} \\ \mu_-^n \|v\| &\leq \|DF^n(x)v\| \leq \mu_+^n \|v\| \text{ for all } x \in \Lambda, v \in E_x^u, n \in \mathbb{N} \\ \lambda_0^{|n|} \|v\| &\leq \|DF^n(x)v\| \leq \mu_0^{|n|} \|v\| \text{ for all } x \in \Lambda, v \in T_x \Lambda, n \in \mathbb{Z}. \end{aligned} \quad (2)$$

The results presented in this paper also apply to the case where Λ is a manifold with boundary, and indeed the case where Λ is only *locally* invariant: this means that there is a neighbourhood V of Λ such that orbits of points in Λ stay in Λ until they leave V . In the latter case orbits may escape through the boundary.

This definition guarantees the existence of stable and unstable invariant manifolds $W^{s,u}(\Lambda) \subset M$ defined as follows. The local stable manifold $W_{\text{loc}}^s(\Lambda)$ is the set of points in a small neighbourhood of Λ whose forward iterates never leave the neighbourhood, and tend exponentially to Λ . The local unstable manifold $W_{\text{loc}}^u(\Lambda)$ is the set of points in the neighbourhood whose backward iterates stay in the neighbourhood and tend exponentially to Λ . We then define

$$W^s(\Lambda) = \bigcup_{i=0}^{\infty} F^{-i}(W_{\text{loc}}^s(\Lambda)), \quad W^u(\Lambda) = \bigcup_{i=0}^{\infty} F^i(W_{\text{loc}}^u(\Lambda)).$$

On the stable and unstable manifolds we have the strong stable and strong unstable foliations, the leaves of which we denote by $W^{s,u}(x)$ for $x \in \Lambda$. For each $x \in \Lambda$, the leaf $W^s(x)$ of the strong stable foliation is tangent at x to E_x^s , and the leaf $W^u(x)$ of the strong unstable foliation is tangent at x to E_x^u . Moreover the foliations are invariant in the sense that $F^i(W^s(x)) = W^s(F^i(x))$ and $F^i(W^u(x)) = W^u(F^i(x))$ for each $x \in \Lambda$ and $i \in \mathbb{Z}$. We thus define the *holonomy maps* $\pi^{s,u} : W^{s,u}(\Lambda) \rightarrow \Lambda$ to be projections along leaves of the strong stable and strong unstable foliations. That is to say, if $x \in W^s(\Lambda)$ then there is a unique $x_+ \in \Lambda$ such that $x \in W^s(x_+)$, and so $\pi^s(x) = x_+$. Similarly, if $x \in W^u(\Lambda)$ then there is a unique $x_- \in \Lambda$ such that $x \in W^u(x_-)$, in which case $\pi^u(x) = x_-$.

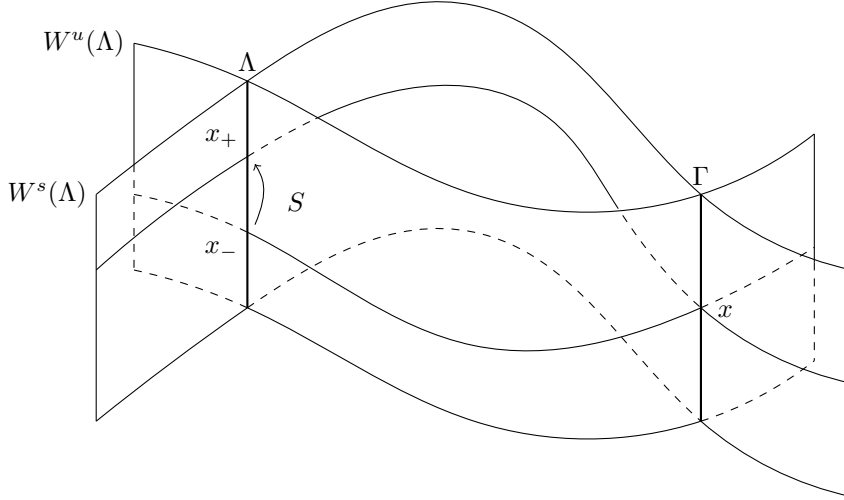


Figure 1: The scattering map S takes a point $x_- \in \Lambda$, follows the unique leaf of the strong unstable foliation passing through x_- to the point x in the homoclinic channel Γ , and from there follows the unique leaf of the strong stable foliation passing through x to the point x_+ on Λ .

Now, suppose that $x \in (W^s(\Lambda) \pitchfork W^u(\Lambda)) \setminus \Lambda$ is a transverse homoclinic point such that $x \in W^s(x_+) \cap W^u(x_-)$. We say that the homoclinic intersection at x is *strongly transverse* if

$$\begin{aligned} T_x W^s(x_+) \oplus T_x (W^s(\Lambda) \cap W^u(\Lambda)) &= T_x W^s(\Lambda), \\ T_x W^u(x_-) \oplus T_x (W^s(\Lambda) \cap W^u(\Lambda)) &= T_x W^u(\Lambda). \end{aligned} \quad (3)$$

In this case we can take a sufficiently small neighbourhood Γ of x in $W^s(\Lambda) \cap W^u(\Lambda)$ so that (3) holds at each point of Γ , and the restrictions to Γ of the holonomy maps are bijections onto their images. We call Γ a *homoclinic channel* (see Figure 1). We can then define the scattering map as follows [15].

Definition 2.2. Let $y_- \in \pi^u(\Gamma)$, let $y = (\pi^u|_\Gamma)^{-1}(y_-)$, and let $y_+ = \pi^s(y)$. The scattering map $S : \pi^u(\Gamma) \rightarrow \pi^s(\Gamma)$ is defined by

$$S = \pi^s \circ (\pi^u)^{-1} : y_- \mapsto y_+.$$

Suppose now that the smoothness r of M and F is at least 2, suppose the normally hyperbolic invariant manifold Λ is a C^r submanifold of M , and suppose instead of condition (1) we have the stronger condition

$$0 < \lambda_- < \lambda_+ < \lambda_0^r < 1 < \mu_0^r < \mu_- < \mu_+ \quad (4)$$

on the hyperbolicity parameters. This *large spectral gap condition* implies C^{r-1} regularity of the strong stable and strong unstable foliations [32], which in turn implies that the scattering map S is C^{r-1} [15].

Remark 2.3. *In general, the scattering map is not globally defined. The homoclinic intersection of invariant manifolds can give rise to very complicated domains of definition of the scattering map, and the general case is that there are many branches S_α of the map defined on sets $U_j \subset \Lambda$. The sets U_j may or may not overlap, and the scattering map may have singularities on ∂U_j . While the results of this paper apply to the general case, we simply write $S : U \rightarrow \Lambda$ to denote the scattering map to avoid awkward notation.*

Suppose the map F depends smoothly on a small parameter ϵ . We point out that all objects defined thus far may vary with ϵ , but, to simplify notation, we do not explicitly write ϵ as a subscript or argument. Indeed, as ϵ varies, the perturbed ($\epsilon > 0$) normally hyperbolic invariant manifold can be written as a graph over the unperturbed ($\epsilon = 0$) manifold as a result of Fenichel theory, and so we can continue to use the coordinates from the original unperturbed manifold [20, 21, 22].

Suppose the scattering map S is defined relative to a homoclinic channel Γ for all sufficiently small $\epsilon > 0$. We allow for the possibility that the angle between $W^{s,u}(\Lambda)$ along the homoclinic channel Γ goes to 0 as $\epsilon \rightarrow 0$. Denote by $\alpha(v_1, v_2)$ the angle between two vectors v_1, v_2 in the direction that yields the smallest result (i.e. $\alpha(v_1, v_2) \in [0, \pi]$). For $x \in \Gamma$, let

$$\alpha_\Gamma(x) = \inf \alpha(v_+, v_-)$$

where the infimum is over all $v_+ \in T_x W^s(\Lambda)^\perp$ and $v_- \in T_x W^u(\Lambda)^\perp$ such that $\|v_\pm\| = 1$.

Definition 2.4. *For $\sigma \geq 0$, we say that the angle of the splitting along Γ is of order ϵ^σ if there is a positive constant C (independent of ϵ) such that*

$$\alpha_\Gamma(x) \geq C\epsilon^\sigma$$

for all $x \in \Gamma$.

Suppose now that the normally hyperbolic invariant manifold Λ is diffeomorphic to $\mathbb{T}^n \times [0, 1]^n$, and denote by $(q, p) \in \mathbb{T}^n \times [0, 1]^n$ smooth coordinates on Λ . Suppose the maps F , and thus $f := F|_\Lambda$, depend on the small parameter ϵ .

Definition 2.5. *We say that $f : \Lambda \rightarrow \Lambda$ is a near-integrable twist map if there is some $k \in \mathbb{N}$ such that*

$$f : \begin{cases} \bar{q} = q + g(p) + O(\epsilon^k) \\ \bar{p} = p + O(\epsilon^k) \end{cases} \quad (5)$$

where

$$\det Dg(p) \neq 0$$

for all $p \in [0, 1]^n$, and where the higher order terms are uniformly bounded in the C^1 topology. If the higher order terms are 0 then f is an integrable twist map.

Remark 2.6. *The assumption that $\det Dg(p) \neq 0$ is sometimes referred to as a local twist property; see for example Section 4, Chapter 23 of [31]. Note that we do not require convexity.*

It follows from the definition that if $f : \Lambda \rightarrow \Lambda$ is a near-integrable twist map, then there exist twist parameters $T_+ > \tilde{T}_- > 0$ such that

$$\tilde{T}_- \|v\| \leq \|Dg(p)v\| \leq T_+ \|v\| \quad (6)$$

for all $p \in [0, 1]^n$ and all $v \in \mathbb{R}^n$. We can always choose T_+ to be independent of ϵ . Our formulation of the problem allows the parameter \tilde{T}_- to depend on ϵ : there is $\tau \in \mathbb{N}_0$ and a strictly positive constant T_- (independent of ϵ) such that $\tilde{T}_- = \epsilon^\tau T_-$.

Definition 2.7. *Suppose $f : \Lambda \rightarrow \Lambda$ is a near-integrable twist map. Denote by $T_+ > \tilde{T}_- = \epsilon^\tau T_- > 0$ the twist parameters. We say that f satisfies:*

- A uniform twist condition if $\tau = 0$;
- A non-uniform twist condition (of order ϵ^τ) if $\tau > 0$, and the order ϵ^k of the error terms in the definition of the near-integrable twist map f is such that $k > \tau$.

In the coordinates (q, p) , we may define a foliation of Λ , the leaves of which are given by

$$\mathcal{L}(p^*) = \{(q, p) \in \Lambda : p = p^*\}. \quad (7)$$

If $f : \Lambda \rightarrow \Lambda$ is a near-integrable twist map in the sense of Definition 2.5, then each leaf of the foliation is almost invariant under f , up to terms of order ϵ^k , where k is as in Definition 2.5.

Suppose we have a scattering map S defined on an open set U in Λ , and suppose the large spectral gap condition (4) holds, so S is C^1 .

Definition 2.8. *We say that the scattering map S is transverse to leaves along leaves (with respect to the leaves (7) of the foliation of Λ) if for all $p_0^* \in [0, 1]^n$ there is $c > 0$ (which may depend on ϵ) and $p_1^* \in [0, 1]^n$ such that for any $p^* \in [0, 1]^n$ satisfying*

$$\|p^* - p_1^*\| < c \quad (8)$$

we have

$$S(\mathcal{L}(p_0^*) \cap U) \cap \mathcal{L}(p^*) \neq \emptyset$$

and $S(\mathcal{L}(p_0^) \cap U)$ is transverse to $\mathcal{L}(p^*)$ at some point where they intersect. In this case we say that the angle of transversality is of order ϵ^v if there are positive constants C, c_* (independent of ϵ) such that*

$$\inf \alpha(v_0, v) \geq C\epsilon^v, \quad c = \epsilon^v c_*$$

where the infimum is taken over all $v_0 \in T_x S(\mathcal{L}(p_0^) \cap U)$ and $v \in T_x \mathcal{L}(p^*)$ such that $\|v_0\| = \|v\| = 1$, for some $x \in S(\mathcal{L}(p_0^*) \cap U) \cap \mathcal{L}(p^*)$, as well as for each $p_0^* \in [0, 1]^n$ and each $p^* \in [0, 1]^n$ satisfying (8).*

2.2 Statement of Theorem 2.9

Let M be a C^r manifold of dimension $2(m+n)$ where $r \geq 4$ and $m, n \in \mathbb{N}$. Suppose $F \in \text{Diff}^4(M)$ has a normally hyperbolic invariant manifold $\Lambda \subset M$ of dimension $2n$ satisfying the large spectral gap condition (4). Suppose $\dim W^s(\Lambda) = \dim W^u(\Lambda) = m+2n$. Suppose F depends on a small parameter ϵ . We make the following further assumptions.

- [A1] The stable and unstable manifolds $W^{s,u}(\Lambda)$ have a strongly transverse homoclinic intersection along a homoclinic channel Γ , and so we have an open set $U \subseteq \Lambda$ and a scattering map $S : U \rightarrow \Lambda$. The angle of the splitting along Γ is of order ϵ^σ .
- [A2] The inner map $f = F|_\Lambda$ is a near-integrable twist map with error terms of order ϵ^k satisfying a non-uniform (or uniform) twist condition of order ϵ^τ .
- [A3] The scattering map S is transverse to leaves along leaves (with respect to the leaves (7) of the foliation of Λ), and the angle of transversality is of order ϵ^ν .

Theorem 2.9. *Fix $\eta > 0$, let $\epsilon > 0$ be sufficiently small, and suppose*

$$k \geq 2(\rho + \tau) + 1 \quad (9)$$

where

$$\rho = \max\{2\sigma, 2\nu, \tau\}. \quad (10)$$

Choose $\{p_j^*\}_{j=1}^\infty \subset [0, 1]^n$ such that

$$S(\mathcal{L}_j \cap U) \cap \mathcal{L}_{j+1} \neq \emptyset,$$

and $S(\mathcal{L}_j \cap U)$ is transverse to \mathcal{L}_{j+1} , where $\mathcal{L}_j = \mathcal{L}(p_j^*)$. Suppose the distance between \mathcal{L}_j and \mathcal{L}_{j+1} is of order ϵ^ν for each j . Then there are $\{z_i\}_{i=1}^\infty \subset M$ and $n_i \in \mathbb{N}$ such that

$$z_{i+1} = F^{n_i}(z_i)$$

and

$$d(z_i, \mathcal{L}_i) < \eta.$$

Moreover, the time to move a distance of order 1 in the p -direction is bounded from below by a term of order

$$\epsilon^{-\rho-\tau-\nu}. \quad (11)$$

Remark 2.10. *Note that we do not assume symplecticity of the map F . We do, however, assume some properties displayed by symplectic maps; for example, the even dimension of the phase space and of the stable and unstable manifolds.*

2.3 Statement of Theorem 2.11

Let M be a C^r manifold of dimension $2(m+n)$ where $r \geq 4$ and $m, n \in \mathbb{N}$. Let $\Sigma = \mathbb{T}^{\ell_1} \times [0, 1]^{\ell_2}$ for some $\ell_1, \ell_2 \in \mathbb{N}_0$, and denote by $(\theta, \xi) \in \mathbb{T}^{\ell_1} \times [0, 1]^{\ell_2}$ coordinates on Σ . Write $\widetilde{M} = M \times \Sigma$. Suppose $\Psi \in \text{Diff}^4(\widetilde{M})$ such that

$$\Psi(z, \theta, \xi) = (G(z, \theta, \xi), \phi(z, \theta, \xi))$$

where $z \in M$, $G \in C^4(\widetilde{M}, M)$, and $\phi \in C^4(\widetilde{M}, \Sigma)$. Suppose Ψ depends on a small parameter ϵ . We make the following assumptions on Ψ .

[B1] There is some $L \in \mathbb{N}$ such that

$$G(z, \theta, \xi) = \widetilde{G}(z; \xi) + O(\epsilon^L)$$

where the higher order terms are uniformly bounded in the C^4 topology, and for each $\xi \in [0, 1]^{\ell_2}$ the map

$$\widetilde{G}(\cdot; \xi) : z \in M \mapsto \widetilde{G}(z; \xi) \in M$$

satisfies the assumptions [A1-3] of Theorem 2.9.

[B2] Moreover, the map ϕ has the form

$$\phi : \begin{cases} \bar{\theta} = \phi_1(z, \theta, \xi) \\ \bar{\xi} = \phi_2(z, \theta, \xi) = \xi + O(\epsilon^L) \end{cases}$$

where the higher order terms are uniformly bounded in the C^4 topology.

Results from [15] imply that Ψ has a normally hyperbolic invariant manifold $\widetilde{\Lambda}$ that is $O(\epsilon^L)$ close in the C^4 topology to $\Lambda \times \Sigma$ where $\Lambda \subset M$ is the normally hyperbolic invariant manifold of $\widetilde{G}(\cdot; \xi)$. Moreover there is an open set $\widetilde{U} \subset \widetilde{\Lambda}$ and a scattering map $\widetilde{S} : \widetilde{U} \rightarrow \widetilde{\Lambda}$ such that the z -component of $\widetilde{S}(z, \theta, \xi)$ is $O(\epsilon^L)$ close in the C^3 topology to $S(z; \xi)$ where $S(\cdot; \xi) : U \rightarrow \Lambda$ is the scattering map corresponding to $\widetilde{G}(\cdot; \xi)$.

We use the coordinates (q, p, θ, ξ) on $\widetilde{\Lambda}$ where (q, p) are the coordinates on Λ and (θ, ξ) are the coordinates on Σ . Notice that the sets

$$\widetilde{\mathcal{L}}(p^*, \xi^*) = \left\{ (q, p, \theta, \xi) \in \widetilde{\Lambda} : p = p^*, \xi = \xi^* \right\} = \mathcal{L}(p^*) \times \mathbb{T}^{\ell_1} \times \{\xi^*\}$$

for $p^* \in [0, 1]^n$ and $\xi^* \in [0, 1]^{\ell_2}$ define the leaves of a foliation of $\widetilde{\Lambda}$, where $\mathcal{L}(p^*)$ are the leaves of the foliation of Λ defined by (7).

Theorem 2.11. *Fix $\eta > 0$ and $K \in \mathbb{N}$ and let $\epsilon > 0$ be sufficiently small. Choose $N \in \mathbb{N}$ satisfying*

$$N \leq \frac{1}{\epsilon^K}$$

and $p_1^*, \dots, p_N^* \in [0, 1]^n$ as in Theorem 2.9, and choose $\xi_1^* \in \text{Int}([0, 1]^{\ell_2})$ such that

$$S(\mathcal{L}_j \cap U; \xi_1^*) \cap \mathcal{L}_{j+1} \neq \emptyset$$

and $S(\mathcal{L}_j \cap U; \xi_1^*)$ is transverse to \mathcal{L}_{j+1} , where $\mathcal{L}_j = \mathcal{L}(p_j^*)$. Suppose the distance between \mathcal{L}_j and \mathcal{L}_{j+1} is of order ϵ^v for each j , and $L > 0$ is sufficiently large. Then there are $\xi_2^*, \dots, \xi_N^* \in [0, 1]^{\ell_2}$ such that, with $\tilde{\mathcal{L}}_j = \tilde{\mathcal{L}}(p_j^*, \xi_j^*)$, there are $w_1, \dots, w_N \in \tilde{M}$ and $n_i \in \mathbb{N}$ such that the ξ component of w_1 is ξ_1^* ,

$$w_{i+1} = \Psi^{n_i}(w_i),$$

and

$$d(w_i, \tilde{\mathcal{L}}_i) < \eta$$

where ρ, σ, τ are as in the statement of Theorem 2.9. Moreover, the time to move a distance of order 1 in the p -direction is of order

$$\epsilon^{-\rho-\tau-v}. \quad (12)$$

Remark 2.12. Note that the transition chain obtained in Theorem 2.11 is only of finite length, while the one obtained in Theorem 2.9 may be infinite.

Remark 2.13. The parameter L in assumption [B1] can be thought of, in terms of applications, as the order of averaging. In celestial mechanics models, for example, there are often fast angles (denoted here by θ) that can be averaged out of the Hamiltonian function up to terms of order ϵ^L for any $L \geq 0$. As such, the conjugate momenta (denoted here by ξ) are constant up to terms of order ϵ^L . In this way, the parameter L can typically be chosen in applications to be as large as is required.

Theorem 2.11 provides a shadowing argument along the invariant manifolds of a normally hyperbolic invariant cylinder relying on long sequences of almost invariant leaves of a foliation. Often in Arnold diffusion results, one wants only to fix the initial and final points and not the whole sequence of leaves. Such statement in the following corollary which is a direct consequence of the assumed hypotheses and Theorem 2.11.

Corollary 2.14. Fix $\eta > 0$, constants ρ, τ and v satisfying (9) and $p_{\text{ini}}^*, p_{\text{fin}}^* \in [0, 1]^n$. Then, for any $L \geq 1$ large enough and $\epsilon > 0$ small enough there exists $\xi_1^* \in \text{Int}([0, 1]^{\ell_2})$, $N \in \mathbb{N}$ and $\{p_k^*\}_{k=1}^N$ such that

- $p_1^* = p_{\text{ini}}^*, p_N^* = p_{\text{fin}}^*$,
- $S(\mathcal{L}_j \cap U; \xi_1^*) \cap \mathcal{L}_{j+1} \neq \emptyset$ and $S(\mathcal{L}_j \cap U; \xi_1^*)$ is transverse to \mathcal{L}_{j+1} , where $\mathcal{L}_j = \mathcal{L}(p_j^*)$
- The distance between \mathcal{L}_j and \mathcal{L}_{j+1} is of order ϵ^v for each j .

Moreover, there are $\xi_2^*, \dots, \xi_N^* \in [0, 1]^{\ell_2}$, $w_1, \dots, w_N \in \widetilde{M}$ such that the ξ component of w_1 is ξ_1^* and natural numbers and $n_i \in \mathbb{N}$ satisfying and

$$n_1 \leq n_2 \leq \dots \leq n_N \lesssim \epsilon^{-\rho-\tau-v}.$$

such that

$$w_{i+1} = \Psi^{n_i}(w_i), \quad \text{and} \quad d(w_i, \widetilde{\mathcal{L}}_i) < \eta$$

where $\widetilde{\mathcal{L}}_j = \widetilde{\mathcal{L}}(p_j^*, \xi_j^*)$.

2.4 Heuristic Description of the Proof of Theorems 2.9 and 2.11

The key idea of the proof is the construction of a sequence of correctly aligned windows in a neighbourhood of the normally hyperbolic invariant cylinder and homoclinic channel. A window is a product of two rectangles, with each boundary component identified as belonging either to an *entry set* or to an *exit set*. Whether a boundary component belongs to the entry set or the exit set is a free choice; indeed, one can even choose the entry set to be empty, and the exit set to be the entire topological boundary of the window, or vice versa. Informally, two windows W_1 and W_2 are correctly aligned under a map f if $f(W_1)$ and W_2 fully overlap in such a way that the image of the exit set of W_1 under f does not intersect W_2 , and the entry set of W_2 does not intersect $f(W_1)$. The crux of this idea is that, if we have a (finite, infinite, or even doubly infinite) sequence of windows $\{W_n\}$ such that W_n is correctly aligned with W_{n+1} under f for each n , then there is a trajectory $\{x_n\}$ of f passing through this sequence of windows, in the sense that $f(x_n) = x_{n+1}$ and $x_n \in W_n$ for each n (see Section 3 for the formal definition and references).

The strategy of the proof, therefore, is to construct explicitly such a sequence of correctly aligned windows. This requires a suitable coordinate system in a neighbourhood of the normally hyperbolic manifold Λ (see Section 4.1) in which the map takes a particular form that allows us to see the twist condition. First we show how to construct a “short sequence” of correctly aligned windows (see Section 5.1), beginning in a neighbourhood of a point x_n in the homoclinic channel, approaching a point $y \in \Lambda$ along the stable manifold $W^s(\Lambda)$, moving around Λ for a (potentially large, depending on the order of the twist condition) number of iterates, and departing along the unstable manifold $W^u(\Lambda)$ towards another point x_{n+1} in the homoclinic channel. This part of the construction uses only normal hyperbolicity and the twist condition [A2].

The next step is to show that, given two such short sequences of correctly aligned windows, we can combine them at a homoclinic point x_n to obtain longer sequences (see Section 5.2). The difficulty here is that the windows at the homoclinic point x_n are expressed in different coordinates: one system of coordinates is obtained by iterating the coordinates near Λ forward along the unstable manifold $W^u(\Lambda)$, and the other by iterating the coordinates near Λ backward along the stable manifold $W^s(\Lambda)$. In order to guarantee that the

windows near x_n are correctly aligned, we need to obtain estimates on the coordinate transformation between these two systems. This step uses: [A1], the transversality of the stable and unstable manifolds at the homoclinic point x_n ; and [A3], the transversality of images under the scattering map of leaves of the foliation of Λ .

The final part of the proof of Theorem 2.9 consists in choosing the aspect ratios (i.e. the size of the constituent rectangles) of each window in the sequence in order to guarantee that the sequence can be continued indefinitely. This is done in Section 5.3.

In order to prove Theorem 2.11, we show that under conditions [B1] and [B2], the conditions [A1-3] of Theorem 2.9 are satisfied by a truncated version of the map. This allows us to consider the sequence of windows W_n in M constructed in the proof of Theorem 2.9; we then extend these windows to the extended phase space $\widehat{M} = M \times \Sigma$ by taking the product of W_n with two new rectangles Θ_n, Ξ_n where $\Theta_n \times \Xi_n \subset \Sigma$. In the Θ_n, Ξ_n directions, we choose the exit set to be empty and the entry set to be the entire topological boundary, and we just increase the size of the rectangles Θ_n, Ξ_n at each step. This guarantees the correct alignment of the windows $W_n \times \Theta_n \times \Xi_n$ at each step under iterates of the truncated map. Finally we show that the error terms of the full map do not spoil the correct alignment; this is true for finite sequences of windows, and it is not clear that it can be extended to infinite sequences.

The structure of the paper is as follows. In Section 3 we define windows, and what it means for them to be correctly aligned, and we state several necessary theorems regarding correct alignment. In Section 4 we establish a suitable system of coordinates in a neighbourhood of Λ , we show how it can be iterated along the stable and unstable manifolds to a homoclinic point, and we obtain estimates on the map in these coordinates. In Section 5 we prove Theorem 2.9, and finally Theorem 2.11 is proved in Section 6.

3 Correctly Aligned Windows

We follow the exposition in [26] (itself based on [28, 29, 30, 40]; see also the appendix of [27]), which elaborates on ideas introduced in [17, 18, 19].

3.1 Definitions and Main Ideas

Let M be a manifold of dimension m . A window is a subset of M that is a product of C^0 rectangles.

Definition 3.1. *Let $m_1, m_2 \in \mathbb{N}_0$ such that $m_1 + m_2 = m$.*

- *A set $W \subset M$ is an (m_1, m_2) window if there is an open neighbourhood V of W in M , an open neighbourhood \widehat{V} of $[0, 1]^{m_1} \times [0, 1]^{m_2}$ in \mathbb{R}^m , and a homeomorphism $\chi : \widehat{V} \rightarrow V$ such that*

$$W = \chi([0, 1]^{m_1} \times [0, 1]^{m_2}).$$

- Moreover there is a choice of entry set

$$W^+ = \chi([0, 1]^{m_1} \times \partial[0, 1]^{m_2}) \quad (13)$$

and exit set

$$W^- = \chi(\partial[0, 1]^{m_1} \times [0, 1]^{m_2}). \quad (14)$$

Remark 3.2. We say that there is a choice of entry and exit sets as we could equally have chosen (14) to be the entry set and (13) to be the exit set. In practice, every time we define a window we explicitly state our choice of its entry and exit sets, but we use the definitions (13) and (14) for the purposes of this exposition.

Definition 3.3. For $j = 1, 2$ let $W_j \subset M$ be an (m_1, m_2) window with parametrisation $\chi_j : \widehat{V}_j \rightarrow V_j$, and let $f \in C^0(M, M)$ such that $f(V_1) \subseteq V_2$. Let

$$\widehat{f} = \chi_2^{-1} \circ f \circ \chi_1 : \widehat{V}_1 \longrightarrow \widehat{V}_2.$$

Then W_1 is correctly aligned with W_2 under f if there is a homotopy

$$H : [0, 1] \times \widehat{V}_1 \longrightarrow \widehat{V}_2$$

such that:

1. We have

$$\begin{aligned} H(0, \cdot) &= \widehat{f}, \\ H([0, 1], \chi_1^{-1}(W_1^-)) \cap \chi_2^{-1}(W_2) &= \emptyset, \\ H([0, 1], \chi_1^{-1}(W_1)) \cap \chi_2^{-1}(W_2^+) &= \emptyset. \end{aligned}$$

2. If $m_1 = 0$ then $f(W_1) \subset \text{Int}(W_2)$. If $m_1 > 0$ then there is $y \in [0, 1]^{m_2}$ such that the map $A_y : [0, 1]^{m_1} \longrightarrow \mathbb{R}^{m_1}$ defined by

$$A_y(x) = \pi_1(H(1, (x, y)))$$

satisfies

$$A_y(\partial[0, 1]^{m_1}) \subset \mathbb{R}^{m_1} \setminus [0, 1]^{m_1}, \quad \deg(A_y, 0) \neq 0$$

where $\pi_1 : \mathbb{R}^{m_1} \times \mathbb{R}^{m_2} \rightarrow \mathbb{R}^{m_1}$ is the canonical projection onto the first component, and $\deg(A_y, 0)$ is the Brouwer degree of the map A_y at 0.

In part 2 of the definition, whenever $m_1 > 0$, if instead we have that there is a linear map $A : \mathbb{R}^{m_1} \rightarrow \mathbb{R}^{m_1}$ such that $H(1, (x, y)) = (Ax, 0)$ for all $x \in [0, 1]^{m_1}$, $y \in [0, 1]^{m_2}$, and $A(\partial[0, 1]^{m_1}) \subset \mathbb{R}^{m_1} \setminus [0, 1]^{m_1}$, then we say that W_1 is linearly correctly aligned with W_2 under f .

Remark 3.4. Observe that the property of two windows being linearly correctly aligned is stronger than simply being correctly aligned, in the sense that linearly correct alignment implies correct alignment (see Proposition 2 of [26]). This becomes useful when we consider products of windows (see Section 3.2).

The following result (Corollary 12 of [40]), whimsically characterised as the property that ‘one can see through a sequence of correctly aligned windows’, is the main point of this technique.

Theorem 3.5. *Let $\{W_i\}_{i \in \mathbb{Z}}$ be a collection of (m_1, m_2) windows in M , and $\{f_i\}_{i \in \mathbb{Z}} \subset C^0(M, M)$ a collection of continuous mappings such that W_i is correctly aligned with W_{i+1} under f_i for each $i \in \mathbb{Z}$. Then we can find $\{z_i\}_{i \in \mathbb{Z}} \subset M$ such that*

$$z_i \in W_i, \quad f_i(z_i) = z_{i+1}$$

for all $i \in \mathbb{Z}$.

The property of two windows being correctly aligned under a map is stable under perturbation in the sense of the next result (Theorem 13 of [40]).

Theorem 3.6. *Suppose $W_1, W_2 \subset M$ are (m_1, m_2) windows such that W_1 is correctly aligned with W_2 under a map $f \in C^0(M, M)$. Then there is an open neighbourhood \mathcal{U} of f in $C^0(M, M)$ with respect to the compact-open topology such that W_1 is correctly aligned with W_2 under \tilde{f} for all $\tilde{f} \in \mathcal{U}$.*

3.2 Products of Windows

Suppose M_j is a manifold of dimension k_j for $j = 1, 2$, and suppose the manifold

$$M = M_1 \times M_2 \tag{15}$$

is equipped with the product topology. For $j = 1, 2$ let $W_j \subset M_j$ be an (m_j, n_j) window with parametrisation $\chi_j : \tilde{V}_j \rightarrow V_j$ where $m_j, n_j \in \mathbb{N}_0$ such that $k_j = m_j + n_j$. Suppose moreover

$$W_j^+ = \chi_j([0, 1]^{m_j} \times \partial[0, 1]^{n_j}), \quad W_j^- = \chi_j(\partial[0, 1]^{m_j} \times [0, 1]^{n_j}).$$

Now, the sets $\widehat{V} = \widehat{V}_1 \times \widehat{V}_2 \subset \mathbb{R}^{k_1} \times \mathbb{R}^{k_2}$ and $V = V_1 \times V_2 \subset M$ are open in the product topology, and $W = W_1 \times W_2 \subset V$. Define $\chi : \widehat{V} \rightarrow V$ by

$$\chi(x, y) = (\chi_1(x), \chi_2(y))$$

where $x \in \widehat{V}_1, y \in \widehat{V}_2$. Define the entry and exit sets of W to be

$$\begin{cases} W^+ = (W_1^+ \times W_2) \cup (W_1 \times W_2^+) \\ W^- = (W_1^- \times W_2) \cup (W_1 \times W_2^-). \end{cases} \tag{16}$$

It can thus be seen that W is an $(m_1 + m_2, n_1 + n_2)$ window in M with parametrisation $\chi : \widehat{V} \rightarrow V$. In this case we say that W is the product of the windows W_1, W_2 , and we write $W = W_1 \times W_2$.

Suppose now that $f \in C^0(M, M)$, and for $(x, y) \in M_1 \times M_2$ we write

$$f(x, y) = (f_1(x, y), f_2(x, y))$$

where $f_j(x, y) \in M_j$. The following result was proved in [26].

Theorem 3.7. *Let $W = W_1 \times W_2$, $\widetilde{W} = \widetilde{W}_1 \times \widetilde{W}_2$ be $(m_1 + m_2, n_1 + n_2)$ windows in $M = M_1 \times M_2$. Suppose*

(i) W_1 is linearly correctly aligned with \widetilde{W}_1 under $f_1(\cdot, y)$ for each $y \in M_2$;
and

(ii) W_2 is linearly correctly aligned with \widetilde{W}_2 under $f_2(x, \cdot)$ for each $x \in M_1$.

Then W is correctly aligned with \widetilde{W} under f .

Remark 3.8. *In general, the manifolds M in our work have the structure (15) only locally. However this is enough to apply Theorem 3.7, since two windows being correctly aligned under a map is a local property.*

4 Coordinates and Estimates

4.1 A Suitable System of Coordinates

Suppose we are in the setting of Theorem 2.9, so M is a C^r manifold of dimension $2(m+n)$, and $F \in \text{Diff}^r(M)$ has a normally hyperbolic invariant manifold $\Lambda \simeq \mathbb{T}^n \times [0, 1]^n$ in M satisfying [A1-3]. In order to construct correctly aligned windows, we need estimates for the map in a neighbourhood of the invariant manifolds $W^{s,u}(\Lambda)$. This requires an appropriate system of coordinates in which the twist property of the inner map $F|_\Lambda$ is apparent; Fenichel coordinates (described below; see [33]) provide a starting point for the coordinate transformation. However, when we express the map F in Fenichel coordinates, there are error terms that complicate the estimates. We therefore seek a further coordinate transformation in which these error terms can be ignored. There is an analogue of the Hartman-Grobman Theorem for normally hyperbolic invariant manifolds, which says that there is a neighbourhood of Λ in which F is topologically conjugate to its so-called *normal map* [36]. These coordinates would be ideal for obtaining *local* estimates in a neighbourhood of the cylinder (as in Section 5.1), but later we need to analyse how the images of this coordinate chart fit together at a homoclinic point if we iterate it backwards (resp. forwards) along the stable (resp. unstable) manifold of Λ (see Section 5.2). This analysis requires a continuous second derivative; since the coordinates provided by [36] are purely topological, we instead follow Section 5.1 of [26] to construct a suitable system of coordinates that are as smooth as required.

Denote by (q, p) the coordinates on Λ . Assuming Λ is a C^r submanifold of M , and the spectral gap is of size r we can introduce, as in [33], C^{r-1} coordinates (s, u, q, p) in a neighbourhood U_0 of Λ , where s, u belong to a neighbourhood of the origin in \mathbb{R}^m such that:

- $W_{\text{loc}}^s(\Lambda) = \{u = 0\}$;
- $W_{\text{loc}}^u(\Lambda) = \{s = 0\}$;
- $W_{\text{loc}}^s(q_0, p_0) = \{u = 0, (q, p) = (q_0, p_0)\}$; and

- $W_{\text{loc}}^u(q_0, p_0) = \{s = 0, (q, p) = (q_0, p_0)\}$.

In these coordinates the map F takes the form

$$F : \begin{cases} (\bar{q}, \bar{p}) = f(q, p) + N_c(s, u, q, p) \\ \bar{s} = A_s(q, p) s + N_s(s, u, q, p) \\ \bar{u} = A_u(q, p) u + N_u(s, u, q, p) \end{cases}$$

where $f = F|_{\Lambda}$, $A_s(q, p) = DF(q, p)|_{E^s}$, $A_u(q, p) = DF(q, p)|_{E^u}$, and

$$N_c(0, u, q, p) = 0 = N_c(s, 0, q, p), \quad N_s(0, u, q, p) = 0 = N_u(s, 0, q, p).$$

In this paper we study the dynamics only in a small neighbourhood of the invariant manifolds $W^{s,u}(\Lambda)$. Therefore we can replace F by a map \tilde{F}_a defined as follows. Choose a function $\psi \in C^\infty(\mathbb{R})$ such that $\psi(x) = 1$ if $|x| \leq 1$ and $\psi(x) = 0$ if $|x| \geq 2$, and let $\psi_a(x) = \psi(ax)$. We then define the map

$$\tilde{F}_a : \begin{cases} (\bar{q}, \bar{p}) = f(q, p) + \psi_a(s^2 + u^2) N_c(s, u, q, p) \\ \bar{s} = A_s(q, p) s + \psi_a(s^2 + u^2) N_s(s, u, q, p) \\ \bar{u} = A_u(q, p) u + \psi_a(s^2 + u^2) N_u(s, u, q, p). \end{cases} \quad (17)$$

Note that the map \tilde{F}_a is only C^{r-1} , since the coordinates (s, u, q, p) are C^{r-1} . Moreover the normally hyperbolic invariant manifold Λ is a C^{r-1} submanifold, and the spectral gap remains of size r . The map \tilde{F}_a agrees with F whenever $s^2 + u^2 \leq a^{-1}$. Furthermore the error terms of \tilde{F}_a are uniformly small in C^0 , and can be made as small as necessary by increasing a . We thus guarantee, by choosing a large enough, that the map \tilde{F}_a is globally defined on $\mathbb{T}^n \times [0, 1]^n \times \mathbb{R}^{2m}$, since we have effectively set the error terms equal to 0 outside the domain of definition of N_c, N_s, N_u . Fix some sufficiently large a once and for all, and define $\Phi = \tilde{F}_a$.

Since Φ is defined globally, the local unstable manifold $W_{\text{loc}}^u(\Lambda)$ extends to a global unstable manifold $\tilde{W}^u(\Lambda) = \bigcup_{n=0}^{\infty} \Phi^n(W_{\text{loc}}^u(\Lambda))$. On $\tilde{W}^u(\Lambda)$ we have the strong unstable foliation, each leaf $\tilde{W}^u(q, p)$ of which is uniquely determined by a point $(q, p) \in \Lambda$ in the usual way (see Section 2.1 for definitions). Since $\tilde{W}^u(\Lambda) = \{s = 0\}$, the variables (u, q, p) define coordinates on $\tilde{W}^u(\Lambda)$. Now, it follows from (17) that the stable manifold $\mathcal{W}^s = W^s(\tilde{W}^u(\Lambda))$ of the unstable manifold $\tilde{W}^u(\Lambda)$ is the entire phase space $\Lambda \times \mathbb{R}^{2m}$. Moreover \mathcal{W}^s admits a C^{r-2} foliation by the leaves $\mathcal{W}^s(u, q, p) = W^s(0, u, q, p)$. Notice that the restriction of Φ to $\tilde{W}^u(\Lambda)$ is

$$\Phi|_{\tilde{W}^u(\Lambda)} : \begin{cases} (\bar{q}, \bar{p}) = f(q, p) \\ \bar{u} = A_u(q, p) u + \psi_a(u^2) N_u(0, u, q, p). \end{cases}$$

We can now take a coordinate s' on each $\mathcal{W}^s(u, q, p)$ so that (s', u, q, p) are C^{r-2}

coordinates in a neighbourhood of Λ , and

$$\Phi : \begin{cases} (\bar{q}, \bar{p}) = f(q, p) \\ \bar{s}' = A_s(q, p)s' + \tilde{N}_s(s', u, q, p) \\ \bar{u} = A_u(q, p)u + \psi_a(u^2) N_u(0, u, q, p) \end{cases} \quad (18)$$

where \tilde{N}_s is globally small, and $\tilde{N}_s(0, u, q, p) = 0$. Since F agrees with Φ in a neighbourhood of Λ , it also takes the form (18) in (s', u, q, p) coordinates. From now on we drop the prime notation and write (s, u, q, p) for this system of coordinates.

We have thus constructed C^{r-2} coordinates in a neighbourhood of Λ in which there are no error terms in the central directions, and the error terms in the hyperbolic directions are small: let $\delta_s, \delta_u > 0$ such that

$$\sup_{(s, u, q, p)} \left\| \tilde{N}_s(s, u, q, p) \right\| \leq \delta_s \|s\|, \quad \sup_{(u, q, p)} \left\| \psi_a(u^2) N_u(0, u, q, p) \right\| \leq \delta_u \|u\|$$

where we denote by $\|x\| = \max_{i=1, \dots, m} |x_i|$ the maximum norm. Since a is sufficiently large, δ_s, δ_u are sufficiently small. Due to the normal hyperbolicity estimates (2), we find that

$$\begin{aligned} \tilde{\lambda}_- \|s\| &\leq \left\| A_s(q, p)s + \tilde{N}_s(s, u, q, p) \right\| \leq \tilde{\lambda}_+ \|s\| \\ \tilde{\mu}_- \|u\| &\leq \left\| A_u(q, p)u + \psi_a(u^2) N(0, u, q, p) \right\| \leq \tilde{\mu}_+ \|u\| \end{aligned}$$

where $\tilde{\lambda}_\pm = \lambda_\pm \pm \delta_s$ and $\tilde{\mu}_\pm = \mu_\pm \pm \delta_u$. Since $\delta_{s, u}$ are sufficiently small, we still have

$$0 < \tilde{\lambda}_- < \tilde{\lambda}_+ < \lambda_0^r < 1 < \mu_0^r < \tilde{\mu}_- < \tilde{\mu}_+.$$

In this way we can use the linear estimates (2) for the nonlinear map F in the hyperbolic directions. We now drop the tilde notation, and simply write λ_\pm, μ_\pm for these adjusted hyperbolicity parameters. The following section deals with estimates in the central directions.

Remark 4.1. *Since the coordinates (s, u, q, p) are C^{r-2} , and since we require the coordinate transformation near a homoclinic point (see Section 5.2) to have two continuous derivatives, we take $r \geq 4$ in the statement of Theorem 2.9.*

We also require a system of coordinates near the homoclinic channel. Denote by \mathcal{U} the neighbourhood of Λ in which we have the (s, u, q, p) coordinates, by $\hat{\mathcal{U}}$ the neighbourhood of $\mathbb{T}^n \times [0, 1]^n \times \{0\}$ in $\mathbb{T}^n \times [0, 1]^n \times \mathbb{R}^{2m}$ to which the (s, u, q, p) variables belong, and by $h : \mathcal{U} \rightarrow \hat{\mathcal{U}}$ the C^2 coordinate transformation constructed above. Let $x \in \Gamma \subset (W^s(\Lambda) \pitchfork W^u(\Lambda)) \setminus \Lambda$ be a transverse homoclinic point. Then there are $N_\pm \in \mathbb{N}$ such that

$$F^{N_+}(x) \in W_{\text{loc}}^s(\Lambda) \cap \mathcal{U}, \quad F^{-N_-}(x) \in W_{\text{loc}}^u(\Lambda) \cap \mathcal{U}.$$

Choose a neighbourhood V of x in M such that $F^{N_+}(V) \subset \mathcal{U}$ and $F^{-N_-}(V) \subset \mathcal{U}$. We define two coordinate systems (s^+, u^+, q^+, p^+) and (s^-, u^-, q^-, p^-) on V via the diffeomorphisms $h^\pm : V \rightarrow \widehat{\mathcal{U}}$ defined by

$$h^+ = \Phi^{-N_+} \circ h \circ F^{N_+}, \quad h^- = \Phi^{N_-} \circ h \circ F^{-N_-}.$$

From (s^+, u^+, q^+, p^+) coordinates to (s, u, q, p) coordinates, the map F^{N_+} acts as Φ^{N_+} , in the sense that

$$h \circ F^{N_+} \circ (h^+)^{-1} = \Phi^{N_+}.$$

Similarly from (s^-, u^-, q^-, p^-) coordinates to (s, u, q, p) coordinates, the map F^{-N_-} acts as Φ^{-N_-} , in the sense that

$$h \circ F^{-N_-} \circ (h^-)^{-1} = \Phi^{-N_-}.$$

4.2 Estimates on the Shearing of a Window by the Twist Map

Recall the inner map $f = F|_\Lambda$ is given by (5), and we assume moreover the twist condition (6), where $\tilde{T}_- = \epsilon^\tau T_-$ for some $\tau \in \mathbb{N}_0$. We use the maximum norm:

$$\|x\| = \max_{i=1, \dots, n} |x_i|. \quad (19)$$

Fix $R > 0$ with the following property. For any $p^0, p \in [0, 1]^n$ we can write

$$g(p) = g(p^0) + Dg(p^0)(p - p^0) + R_*(p^0, p)$$

where $R_*(p^0, p)$ is a remainder term of order $O(\|p - p^0\|^2)$. Then the constant R is chosen so that

$$\|R_*(p^0, p)\| \leq R \|p - p^0\|^2 \quad (20)$$

for all $p^0, p \in [0, 1]^n$.

Define a window $W = [Q \times P] \subset \Lambda$ where

$$Q = [a, a + \gamma]^n \subset \mathbb{T}^n, \quad P = [b, b + \delta]^n \subset [0, 1]^n.$$

Choose the exit set W^- to be, say,

$$W^- = Q \times \partial P.$$

The following estimates apply equally if the exit set is chosen to be in the q -direction.

For each $j = 1, \dots, n$ define

$$B_j^0 = [b, b + \delta]^{j-1} \times \{b\} \times [b, b + \delta]^{n-j}, \quad B_j^1 = [b, b + \delta]^{j-1} \times \{b + \delta\} \times [b, b + \delta]^{n-j}$$

and let

$$E_j^0 = Q \times B_j^0, \quad E_j^1 = Q \times B_j^1$$

be the corresponding components of the exit set, in the sense that

$$W^- = \bigcup_{j=1}^n (E_j^0 \cup E_j^1).$$

Suppose now that $N \in \mathbb{Z}$, and suppose we take an iterate $f^N(W)$ of W . Define

$$\Delta_j^N = \min_{(q^0, p^0) \in E_j^0} \min_{(q^1, p^1) \in E_j^1} \|\pi \circ f^N(q^1, p^1) - \pi \circ f^N(q^0, p^0)\|$$

and

$$\Omega^N = \max_{(q^0, p^0), (q^1, p^1) \in W} \|\pi \circ f^N(q^1, p^1) - \pi \circ f^N(q^0, p^0)\|$$

where $\pi : \mathbb{T}^n \times [0, 1]^n \rightarrow \mathbb{T}^n$ is the canonical projection. The following lemma gives lower and upper bounds on the shearing of the window W under f^N .

Lemma 4.2. *There is a positive constant C such that*

$$\Delta_j^N \geq \epsilon^\tau |N| T_- \delta - |N| R \delta^2 - \gamma - C |N|^2 \epsilon^k \quad (21)$$

and

$$\Omega^N \leq \gamma + |N| T_+ \delta + C |N|^2 \epsilon^k \quad (22)$$

where ϵ^τ is the order of the twist condition, ϵ^k is the order of the error terms in the definition (5) of f , and the positive constant R is defined by (20).

Proof. We establish (21) for the case $j = 1$ as the remaining cases are analogous. Fix $q^0, q^1 \in Q$ and $p_*^0, p_*^1 \in [b, b + \delta]^{n-1}$. Define

$$p^0 = (b, p_*^0), \quad p^1 = (b + \delta, p_*^1)$$

Notice that

$$\|p^1 - p^0\| = \delta.$$

Write

$$(\bar{q}^j, \bar{p}^j) = f^N(q^j, p^j)$$

for $j = 0, 1$ so that

$$\bar{q}^j = q^j + Ng(p^j) + O(N^2 \epsilon^k), \quad \bar{p}^j = p^j + O(N \epsilon^k).$$

Let C be a uniform upper bound on the terms of order ϵ^k . Then we have

$$\begin{aligned} \|\bar{q}^1 - \bar{q}^0\| &\geq |N| \|g(p^1) - g(p^0)\| - \|q^1 - q^0\| - C |N|^2 \epsilon^k \\ &\geq |N| \|Dg(p^0)(p^1 - p^0) + R_*(p^0, p^1)\| - \gamma - C |N|^2 \epsilon^k \\ &\geq \epsilon^\tau |N| T_- \delta - |N| R \delta^2 - \gamma - C |N|^2 \epsilon^k. \end{aligned}$$

In order to prove (22), fix $(q^j, p^j) \in W$ and let $(\bar{q}^j, \bar{p}^j) = f^N(q^j, p^j)$ for $j = 0, 1$. Then

$$\begin{aligned} \|\bar{q}^1 - \bar{q}^0\| &\leq \|q^1 - q^0\| + |N| \|g(p^1) - g(p^0)\| + C |N|^2 \epsilon^k \\ &\leq \gamma + |N| T_+ \delta + C |N|^2 \epsilon^k. \end{aligned}$$

□

5 Proof of Theorem 2.9

Denote by $\mathcal{L}_1, \dots, \mathcal{L}_N$ a sequence of leaves of the foliation of Λ with transverse homoclinic connections, as in the statement of Theorem 2.9. Let $x_i \in \Gamma$ be a point in the homoclinic channel such that $\pi^s(x_i) \in \mathcal{L}_i$ and $\pi^u(x_i) \in \mathcal{L}_{i+1}$. The proof of Theorem 2.9 is divided into 3 parts. In the first part of the proof, we show how to construct a short chain of correctly aligned windows beginning at the homoclinic point x_i , moving near to \mathcal{L}_i , and ending at the homoclinic point x_{i+1} . In the second part of the proof we show that we can continue the sequence of windows across a transverse homoclinic intersection in order to glue together two short sequences. In the final part of the proof, it is shown that we can consistently choose the aspect ratios of the windows in a way that allows us to continue the sequence indefinitely.

5.1 Construction of a Short Sequence of Correctly Aligned Windows

Throughout the proof we use the coordinates (s, u, q, p) constructed in Section 4.1 near the normally hyperbolic invariant manifold Λ , as well as the coordinates (s^+, u^+, q^+, p^+) and (s^-, u^-, q^-, p^-) , defined in Section 4.1, near the homoclinic channel Γ .

This part of the proof is performed in 3 steps. In the first step, we construct a window W_i in (s^+, u^+, q^+, p^+) coordinates centred at $x_i \in \Gamma$, and a window \widetilde{W}_i in (s, u, q, p) coordinates centred at $f^{N_i}(y_i^s)$ where $y_i^s = \pi^s(x_i)$ such that W_i is correctly aligned with \widetilde{W}_i under F^{N_i} for some $N_i \in \mathbb{N}$. In the second step, we construct a window \widehat{W}_i in (s, u, q, p) coordinates centred at $f^{N_i+K_i}(y_i^s)$ such that \widetilde{W}_i is correctly aligned with \widehat{W}_i under F^{K_i} for some $K_i \in \mathbb{N}$. In the third step, we construct a window W'_i in (s^-, u^-, q^-, p^-) coordinates, centred at the transverse homoclinic point $x_{i+1} \in \Gamma$ such that \widehat{W}_i is correctly aligned with W'_i under F^{M_i} for some $M_i \in \mathbb{N}$.

In each step, the windows are defined as a product of two constituent windows. We ensure that the windows we construct are correctly aligned by ensuring that the constituent windows are *linearly* correctly aligned (see Definition 3.3 and Theorem 3.7). In the hyperbolic directions (i.e. the s, u variables) we simply use the normal hyperbolicity estimates (2). In the centre directions (i.e. the q, p variables) we use the shearing estimates provided by Lemma 4.2. Moreover at each step we state the entry and exit sets of the constituent windows, in which case the entry and exit sets of the product of these windows is given by (16).

Finally, let us define the *centre* of a window W in (s, u, q, p) coordinates. Suppose we define a window

$$W = [S \times U] \times [Q \times P]$$

where S, U, Q, P is a rectangle in s, u, q, p respectively, and suppose the rectan-

gles have size

$$\begin{aligned} |S| &= \max_{s_1, s_2 \in S} \|s_1 - s_2\| = \alpha, & |U| &= \max_{u_1, u_2 \in U} \|u_1 - u_2\| = \beta, \\ |Q| &= \max_{q_1, q_2 \in Q} \|q_1 - q_2\| = \gamma, & |P| &= \max_{p_1, p_2 \in P} \|p_1 - p_2\| = \delta, \end{aligned}$$

where the maximum norm $\|\cdot\|$ is defined by (19). The *centre* of W is the point

$$c^0 = (s^0, u^0, q^0, p^0)$$

such that

$$\begin{aligned} S &= \left[s_1^0 - \frac{\alpha}{2}, s_1^0 + \frac{\alpha}{2} \right] \times \cdots \times \left[s_m^0 - \frac{\alpha}{2}, s_m^0 + \frac{\alpha}{2} \right], \\ U &= \left[u_1^0 - \frac{\beta}{2}, u_1^0 + \frac{\beta}{2} \right] \times \cdots \times \left[u_m^0 - \frac{\beta}{2}, u_m^0 + \frac{\beta}{2} \right], \\ Q &= \left[q_1^0 - \frac{\gamma}{2}, q_1^0 + \frac{\gamma}{2} \right] \times \cdots \times \left[q_n^0 - \frac{\gamma}{2}, q_n^0 + \frac{\gamma}{2} \right], \\ P &= \left[p_1^0 - \frac{\delta}{2}, p_1^0 + \frac{\delta}{2} \right] \times \cdots \times \left[p_n^0 - \frac{\delta}{2}, p_n^0 + \frac{\delta}{2} \right]. \end{aligned}$$

5.1.1 Step 1

Define a window W_i , centred at the homoclinic point $x_i \in \Gamma$ and given in (s^+, u^+, q^+, p^+) coordinates by

$$W_i = [S_i \times U_i] \times [Q_i \times P_i]$$

where S_i, U_i, Q_i, P_i are rectangles in s^+, u^+, q^+, p^+ respectively, with sizes

$$|S_i| = \alpha_i, \quad |U_i| = \beta_i, \quad |Q_i| = \gamma_i, \quad |P_i| = \delta_i.$$

We choose the exit sets to be

$$[S_i \times U_i]^- = S_i \times \partial U_i, \quad [Q_i \times P_i]^- = Q_i \times \partial P_i.$$

Now choose $N_i \in \mathbb{N}$ such that $F^{N_i}(x_i) \in W_{\text{loc}}^s(\Lambda) \cap \mathcal{U}$ where \mathcal{U} is the neighbourhood of Λ in which (q, p, s, u) coordinates are defined. Let

$$\nu_i = d_{W^s(\Lambda)}(x_i, y_i^s),$$

where $d_{W^s(\Lambda)}$ is the distance measured along the stable manifold, and $y_i^s = \pi^s(x_i)$. Then

$$\nu_i \lambda_-^{N_i} \leq d_{W^s(\Lambda)}(F^{N_i}(x_i), F^{N_i}(y_i^s)) \leq \nu_i \lambda_+^{N_i}$$

and moreover by (2) and Lemma 4.2:

- In the s -direction, W_i is contracted by F^{N_i} to a size between $\alpha_i \lambda_-^{N_i}$ and $\alpha_i \lambda_+^{N_i}$;

- In the u -direction, W_i is expanded by F^{N_i} to a size between $\beta_i \mu_-^{N_i}$ and $\beta_i \mu_+^{N_i}$;
- In the q -direction, W_i is sheared by the twist map to a size at most $\gamma_i + N_i T_+ \delta_i + CN_i \epsilon^k$; and
- In the p -direction, the change is $\pm CN_i \epsilon^k$.

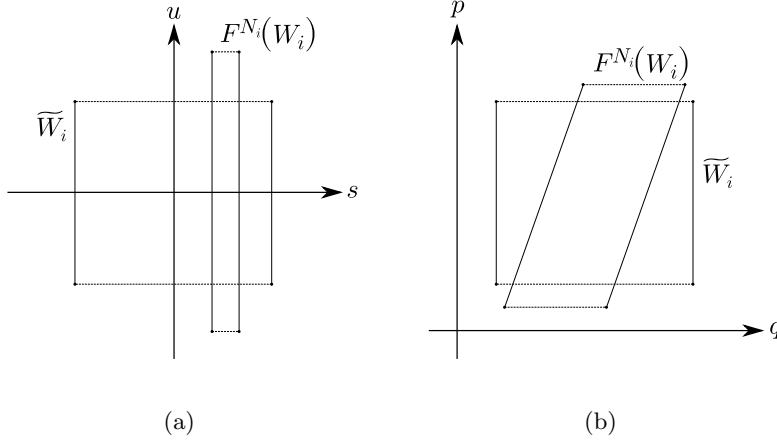


Figure 2: Step 1. The dashed edges represent exit sets.

Define a window \widetilde{W}_i centred at $f^{N_i}(y_i^s) \in \Lambda$ and given in (s, u, q, p) coordinates by

$$\widetilde{W}_i = [\widetilde{S}_i \times \widetilde{U}_i] \times [\widetilde{Q}_i \times \widetilde{P}_i]$$

where $\widetilde{S}_i, \widetilde{U}_i, \widetilde{Q}_i, \widetilde{P}_i$ are rectangles in s, u, q, p respectively, and where we choose the exit sets to be

$$[\widetilde{S}_i \times \widetilde{U}_i]^- = \widetilde{S}_i \times \partial \widetilde{U}_i, \quad [\widetilde{Q}_i \times \widetilde{P}_i]^- = \widetilde{Q}_i \times \partial \widetilde{P}_i.$$

If we choose the sizes of the constituent rectangles to be

$$|\widetilde{S}_i| = \tilde{\alpha}_i, \quad |\widetilde{U}_i| = \tilde{\beta}_i, \quad |\widetilde{Q}_i| = \tilde{\gamma}_i, \quad |\widetilde{P}_i| = \tilde{\delta}_i$$

where

$$\tilde{\alpha}_i > (\alpha_i + 2\nu_i) \lambda_+^{N_i}, \quad (23)$$

$$\tilde{\beta}_i < \beta_i \mu_-^{N_i}, \quad (24)$$

$$\tilde{\gamma}_i > \gamma_i + N_i T_+ \delta_i + CN_i^2 \epsilon^k, \quad (25)$$

$$\tilde{\delta}_i < \delta_i - CN_i \epsilon^k, \quad (26)$$

then W_i is correctly aligned with \widetilde{W}_i under F^{N_i} by Theorem 3.7. In Section 5.3 we explain how to choose $\alpha_i, \beta_i, \gamma_i, \delta_i$ and N_i so that inequalities (23), (24), (25), (26) are solvable.

Remark 5.1. *The reason ν_i carries a factor of 2 in (23) is the following. Since the hyperbolic rectangle of $F^{N_i}(W_i)$ lies to one side of $\{s=0\}$ (see Figure 2 (a)) we must choose $\tilde{\alpha}_i$ so that $\frac{\tilde{\alpha}_i}{2}$ is greater than the distance from $\{s=0\}$ to the outermost point on $F^{N_i}(W_i)$ in the s -direction, which is at most $\lambda_+^{N_i}(\frac{\alpha_i}{2} + \nu_i)$.*

5.1.2 Step 2

Take a forward iterate $F^{K_i}(\widetilde{W}_i)$ of \widetilde{W}_i that brings the centre $f^{N_i}(y_i^s)$ of \widetilde{W}_i to $f^{N_i+K_i}(y_i^s)$. Define a window \widehat{W}_i centred at $f^{N_i+K_i}(y_i^s)$, and given in (s, u, q, p) coordinates by

$$\widehat{W}_i = [\widehat{S}_i \times \widehat{U}_i] \times [\widehat{Q}_i \times \widehat{P}_i]$$

where $\widehat{S}_i, \widehat{U}_i, \widehat{Q}_i, \widehat{P}_i$ are rectangles in s, u, q, p respectively, and where we choose the exit sets to be

$$[\widehat{S}_i \times \widehat{U}_i]^- = \widehat{S}_i \times \partial\widehat{U}_i, \quad [\widehat{Q}_i \times \widehat{P}_i]^- = \partial\widehat{Q}_i \times \widehat{P}_i.$$

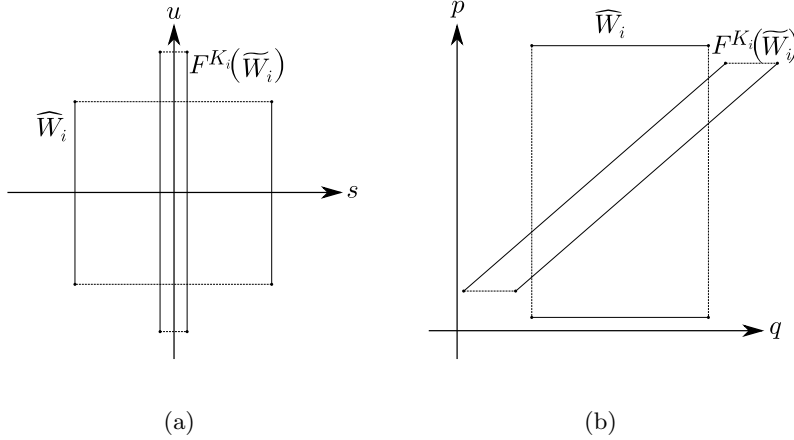


Figure 3: Step 2. The dashed edges represent exit sets.

Notice that the exit set of \widehat{W}_i is not in the same direction as the exit set of \widetilde{W}_i . Under the iteration F^{K_i} , the window \widetilde{W}_i is deformed, by (2) and Lemma 4.2, as follows:

- In the s -direction, \widetilde{W}_i is contracted by F^{K_i} to a size between $\tilde{\alpha}_i \lambda_-^{K_i}$ and $\tilde{\alpha}_i \lambda_+^{K_i}$;

- In the u -direction, \widetilde{W}_i is expanded by F^{K_i} to a size between $\widetilde{\beta}_i \mu_-^{K_i}$ and $\widetilde{\beta}_i \mu_+^{K_i}$;
- In the q -direction, \widetilde{W}_i is sheared by the twist map to a size at least $\epsilon^\tau K_i T_- \widetilde{\delta}_i - K_i R \widetilde{\delta}_i^2 - \widetilde{\gamma}_i - CK_i \epsilon^k$; and
- In the p -direction, the change is $\pm CK_i \epsilon^k$.

If we choose the sizes of the constituent rectangles of \widehat{W}_i to be

$$|\widehat{S}_i| = \widehat{\alpha}_i, \quad |\widehat{U}_i| = \widehat{\beta}_i, \quad |\widehat{Q}_i| = \widehat{\gamma}_i, \quad |\widehat{P}_i| = \widehat{\delta}_i$$

where

$$\widehat{\alpha}_i > \widetilde{\alpha}_i \lambda_+^{K_i}, \quad (27)$$

$$\widehat{\beta}_i < \widetilde{\beta}_i \mu_-^{K_i}, \quad (28)$$

$$\widehat{\gamma}_i < \epsilon^\tau K_i T_- \widetilde{\delta}_i - K_i R \widetilde{\delta}_i^2 - \widetilde{\gamma}_i - CK_i^2 \epsilon^k, \quad (29)$$

$$\widehat{\delta}_i > \widetilde{\delta}_i + CK_i \epsilon^k, \quad (30)$$

then \widetilde{W}_i is correctly aligned with \widehat{W}_i under F^{K_i} by Theorem 3.7. In Section 5.3 we explain how to choose $\widetilde{\alpha}_i, \widetilde{\beta}_i, \widetilde{\gamma}_i, \widetilde{\delta}_i$ and K_i so that inequalities (27), (28), (29), (30) are solvable.

5.1.3 Step 3

Let $x_{i+1} \in \Gamma$ be a point in the homoclinic channel for which there are $y_i^u \in \mathcal{L}_i$ and $y_{i+1}^s \in \mathcal{L}_{i+1}$ such that $\pi^u(x_{i+1}) = y_i^u$ and $\pi^s(x_{i+1}) = y_{i+1}^s$. Define a window W'_i centred at x_{i+1} and given in (s^-, u^-, q^-, p^-) coordinates by

$$W'_i = [S'_i \times U'_i] \times [Q'_i \times P'_i]$$

where we choose the exit sets to be

$$[S'_i \times U'_i]^- = S'_i \times \partial U'_i, \quad [Q'_i \times P'_i]^- = \partial Q'_i \times P'_i.$$

Suppose we choose the sizes of the constituent rectangles of W'_i to be

$$|S'_i| = \alpha'_i, \quad |U'_i| = \beta'_i, \quad |Q'_i| = \gamma'_i, \quad |P'_i| = \delta'_i.$$

We take an iterate $F^{M_i}(\widehat{W}_i)$ of \widehat{W}_i for some $M_i \in \mathbb{N}$, and require that $F^{M_i}(\widehat{W}_i)$ is correctly aligned with W'_i under F^{M_i} . Observe that this is equivalent to W'_i being correctly aligned with \widehat{W}_i under F^{-M_i} , but with the roles of entry and exit sets in the definition of correctly aligned windows reversed.

We can write $y_i^u = (q_i^u, p_i^u)$ and $f^{N_i+K_i}(y_i^s) = (q_i^s, p_i^s)$. Let

$$\omega'_i = \|q_i^s - q_i^u\|, \quad \nu'_i = d_{W^u(\Lambda)}(x_{i+1}, y_i^u) \quad (31)$$

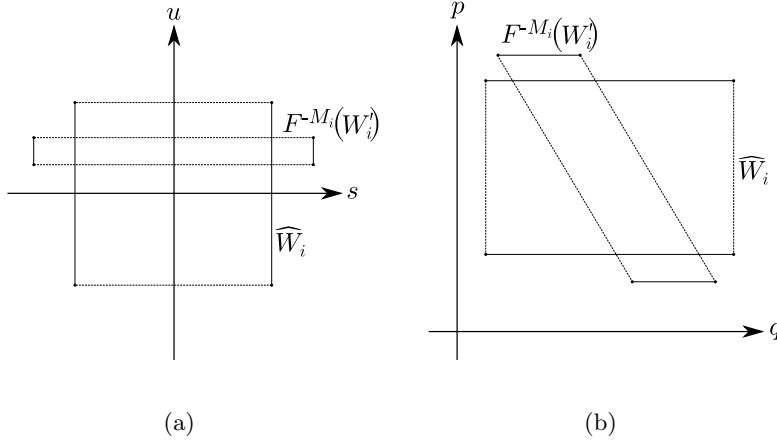


Figure 4: Step 3. The dashed edges represent exit sets.

where $d_{W^u(\Lambda)}$ is the distance measured along the unstable manifold. Then

$$\nu'_i \mu_+^{-M_i} \leq d_{W^u(\Lambda)}(F^{-M_i}(x_{i+1}), F^{-M_i}(y_i^u)) \leq \nu'_i \mu_-^{-M_i},$$

and moreover, by (2) and Lemma 4.2:

- In the s -direction, W_i' is expanded by F^{-M_i} to a size between $\alpha'_i \lambda_+^{-M_i}$ and $\alpha'_i \lambda_-^{-M_i}$;
- In the u -direction, W_i' is contracted by F^{-M_i} to a size between $\beta'_i \mu_+^{-M_i}$ and $\beta'_i \mu_-^{-M_i}$;
- In the q -direction, W_i' is sheared by the twist map to a size at most $\gamma'_i + M_i T_+ \delta'_i + C M_i \epsilon^k$; and
- In the p -direction, the change is $\pm C M_i \epsilon^k$.

If we choose $\alpha'_i, \beta'_i, \gamma'_i, \delta'_i$ such that

$$\alpha'_i \lambda_+^{-M_i} > \hat{\alpha}_i, \quad (32)$$

$$(\beta'_i + 2\nu'_i) \mu_-^{-M_i} < \hat{\beta}_i, \quad (33)$$

$$\gamma'_i + M_i T_+ \delta'_i + C M_i^2 \epsilon^k + 2\omega'_i < \hat{\gamma}_i, \quad (34)$$

$$\delta'_i - C M_i \epsilon^k > \hat{\delta}_i, \quad (35)$$

then \widehat{W}_i is correctly aligned with W_i' under F^{M_i} by Theorem 3.7. In Section 5.3 we explain how to choose $\hat{\alpha}_i, \hat{\beta}_i, \hat{\gamma}_i, \hat{\delta}_i$ and M_i so that inequalities (32), (33), (34), (35) are solvable.

Remark 5.2. *With regards to the factor of 2 carried both by ν'_i in (33) and by ω'_i in (34), compare Remark 5.1 with Figure 4.*

5.2 Continuing the Sequence Across a Transverse Homoclinic Intersection

So far, we have shown how to construct short sequences of windows $W_i, \widetilde{W}_i, \widehat{W}_i, W'_i$ beginning at the transverse homoclinic point x_i , passing near \mathcal{L}_i , and ending at another transverse homoclinic point x_{i+1} . We now want to show that this short sequence can be connected to the next short sequence $W_{i+1}, \widetilde{W}_{i+1}, \widehat{W}_{i+1}, W'_{i+1}$, which amounts to proving that W'_i is correctly aligned with W_{i+1} under the identity mapping. Both of these windows are centred at x_{i+1} , but the difficulty is that W'_i is defined in the (s^-, u^-, q^-, p^-) coordinates obtained by carrying the (s, u, q, p) coordinates from a neighbourhood of \mathcal{L}_i out along $W^u(\Lambda)$ to a neighbourhood of x_{i+1} , whereas W_{i+1} is defined in the (s^+, u^+, q^+, p^+) coordinates obtained by carrying the (s, u, q, p) coordinates from a neighbourhood of \mathcal{L}_{i+1} back along $W^s(\Lambda)$ to a neighbourhood of x_{i+1} .

Define the window W_{i+1} centred at x_{i+1} and given in (s^+, u^+, q^+, p^+) coordinates by

$$W_{i+1} = [S_{i+1} \times U_{i+1}] \times [Q_{i+1} \times P_{i+1}]$$

where the exit sets are

$$[S_{i+1} \times U_{i+1}]^- = S_{i+1} \times \partial U_{i+1}, \quad [Q_{i+1} \times P_{i+1}]^- = Q_{i+1} \times \partial P_{i+1}$$

and where the sizes of the constituent rectangles are

$$|S_{i+1}| = \alpha_{i+1}, \quad |U_{i+1}| = \beta_{i+1}, \quad |Q_{i+1}| = \gamma_{i+1}, \quad |P_{i+1}| = \delta_{i+1}.$$

The following lemma gives conditions under which W'_i is correctly aligned with W_{i+1} under the identity mapping.

Lemma 5.3. *There are nonnegative constants C_j for $j = 1, \dots, 8$ where $C_4, C_7 > 0$ and a constant $R' > 0$ such that, if*

$$\alpha_{i+1} > C_1 \epsilon^\sigma \alpha'_i + C_2 \beta'_i + R' \zeta_i^2, \quad (36)$$

$$\beta_{i+1} < -C_3 \alpha'_i + C_4 \epsilon^\sigma \beta'_i - R' \zeta_i^2, \quad (37)$$

$$\gamma_{i+1} > C_5 \gamma'_i + C_6 \epsilon^v \delta'_i + R' \zeta_i^2, \quad (38)$$

$$\delta_{i+1} < C_7 \epsilon^v \gamma'_i - C_8 \delta'_i - R' \zeta_i^2, \quad (39)$$

where

$$\zeta_i = \max \{ \alpha'_i, \beta'_i, \gamma'_i, \delta'_i \} \quad (40)$$

then W'_i is correctly aligned with W_{i+1} under the identity mapping.

Proof. Recall we assume that the stable and unstable manifolds $W^{s,u}(\Lambda)$ have equal dimension $m + 2n$, where we denote by m the dimension of the s, u variables, and by n the dimension of the q, p variables. By [A1], the angle between s^- and u^+ at x_{i+1} is of order ϵ^σ , as is the angle between u^- and s^+ . Moreover, by [A3], the angle between q^- and q^+ at x_{i+1} is of order ϵ^v , as is the angle

between p^- and p^+ . It follows that there is a well-defined coordinate transformation $\phi : (s^-, u^-, q^-, p^-) \mapsto (s^+, u^+, q^+, p^+)$ in a neighbourhood of x_{i+1} , and the linearisation of ϕ at x_{i+1} is of the form

$$D\phi(x_{i+1}) = \begin{pmatrix} A & 0 \\ 0 & B \end{pmatrix} \quad (41)$$

where

$$A = \begin{pmatrix} \epsilon^\sigma A_1 & A_2 \\ A_3 & \epsilon^\sigma A_4 \end{pmatrix}, \quad B = \begin{pmatrix} B_1 & \epsilon^\nu B_2 \\ \epsilon^\nu B_3 & B_4 \end{pmatrix}. \quad (42)$$

Here each matrix A_j is of dimension $m \times m$, each matrix B_j is of dimension $n \times n$, and A_1, A_4, B_2, B_3 are invertible.

Recall the maximum norm $\|\cdot\|$ is defined by (19). Fix $R' > 0$ with the following property. We can write

$$\phi(x) = \phi(x_{i+1}) + D\phi(x_{i+1})(x - x_{i+1}) + R'_*(x) \quad (43)$$

where $R'_*(x)$ is a remainder term of order $O(\|x - x_{i+1}\|^2)$. Then the constant R' is chosen so that

$$\|R'_*(x)\| \leq R' \|x - x_{i+1}\|^2 \quad (44)$$

for all x in some neighbourhood of x_{i+1} . Notice moreover that, since W'_i is centred at x_{i+1} , we have

$$\|x - x_{i+1}\| \leq \frac{\zeta_i}{2} \quad (45)$$

for all $x \in W'_i$, where ζ_i is defined by (40).

There are $a'_i, b'_i \in \mathbb{R}^m$ and $c'_i, d'_i \in \mathbb{R}^n$ such that

$$\begin{aligned} S'_i &= [a'_{i,1}, a'_{i,1} + \alpha'_i] \times \cdots \times [a'_{i,m}, a'_{i,m} + \alpha'_i], \\ U'_i &= [b'_{i,1}, b'_{i,1} + \beta'_i] \times \cdots \times [b'_{i,m}, b'_{i,m} + \beta'_i], \\ Q'_i &= [c'_{i,1}, c'_{i,1} + \gamma'_i] \times \cdots \times [c'_{i,n}, c'_{i,n} + \gamma'_i], \\ P'_i &= [d'_{i,1}, d'_{i,1} + \delta'_i] \times \cdots \times [d'_{i,n}, d'_{i,n} + \delta'_i]. \end{aligned}$$

For each $j = 1, \dots, m$ define

$$\begin{aligned} I_{j,0}^h &= [b'_{i,1}, b'_{i,1} + \beta'_i] \times \cdots \times \{b'_{i,j}\} \times \cdots \times [b'_{i,m}, b'_{i,m} + \beta'_i], \\ I_{j,1}^h &= [b'_{i,1}, b'_{i,1} + \beta'_i] \times \cdots \times \{b'_{i,j} + \beta'_i\} \times \cdots \times [b'_{i,m}, b'_{i,m} + \beta'_i], \end{aligned}$$

and for each $k = 1, \dots, n$ define

$$\begin{aligned} I_{k,0}^c &= [c'_{i,1}, c'_{i,1} + \gamma'_i] \times \cdots \times \{c'_{i,k}\} \times \cdots \times [c'_{i,n}, c'_{i,n} + \gamma'_i], \\ I_{k,1}^c &= [c'_{i,1}, c'_{i,1} + \gamma'_i] \times \cdots \times \{c'_{i,k} + \gamma'_i\} \times \cdots \times [c'_{i,n}, c'_{i,n} + \gamma'_i]. \end{aligned}$$

We then let

$$E_{j,0}^h = S'_i \times I_{j,0}^h \times Q'_i \times P'_i, \quad E_{j,1}^h = S'_i \times I_{j,1}^h \times Q'_i \times P'_i$$

be the corresponding components of the exit set of W'_i in the hyperbolic directions, and let

$$E_{k,0}^c = S'_i \times U'_i \times I_{k,0}^c \times P'_i, \quad E_{k,1}^c = S'_i \times U'_i \times I_{k,1}^c \times P'_i$$

be the corresponding components of the exit set of W'_i in the cylindrical directions. For each $\kappa \in \{s, u, q, p\}$ denote by $\Pi_\kappa : (s, u, q, p) \mapsto \kappa$ the projection onto the κ coordinate. Define

$$\Delta_j^h = \min_{z_0 \in E_{j,0}^h} \min_{z_1 \in E_{j,1}^h} \|\Pi_u \circ \phi(z_1) - \Pi_u \circ \phi(z_0)\|, \quad (46)$$

$$\Omega_j^h = \max_{z_0 \in E_{j,0}^h} \max_{z_1 \in E_{j,1}^h} \|\Pi_s \circ \phi(z_1) - \Pi_s \circ \phi(z_0)\|, \quad (47)$$

$$\Delta_k^c = \min_{z_0 \in E_{k,0}^c} \min_{z_1 \in E_{k,1}^c} \|\Pi_p \circ \phi(z_1) - \Pi_p \circ \phi(z_0)\|, \quad (48)$$

$$\Omega_k^c = \max_{z_0 \in E_{k,0}^c} \max_{z_1 \in E_{k,1}^c} \|\Pi_q \circ \phi(z_1) - \Pi_q \circ \phi(z_0)\|. \quad (49)$$

If we choose

$$\alpha_{i+1} > \Omega_j^h, \quad \beta_{i+1} < \Delta_j^h, \quad \gamma_{i+1} > \Omega_k^c, \quad \delta_{i+1} < \Delta_k^c \quad (50)$$

for each $j = 1, \dots, m$ and each $k = 1, \dots, n$ then W'_i is correctly aligned with W_{i+1} under the map ϕ (see Figure 5). Therefore, in what follows, we search for upper bounds for each Ω_j^h, Ω_k^c and lower bounds for each Δ_j^h, Δ_k^c . In fact, we compute bounds for Ω_1^h, Ω_1^c and Δ_1^h, Δ_1^c as the estimates for other values of j, k are analogous.

For $j = 0, 1$ let $s_j^h \in S'_i, q_j^h \in Q'_i, p_j^h \in P'_i$, and

$$u_j^* \in [b'_{i,2}, b'_{i,2} + \beta'_i] \times \dots \times [b'_{i,m}, b'_{i,m} + \beta'_i].$$

Define $u_0^h = (b'_{i,1}, u_0^*), u_1^h = (b'_{i,1} + \beta'_i, u_1^*)$, and $z_j^h = (s_j^h, u_j^h, q_j^h, p_j^h)$ for $j = 0, 1$. Observe that all points in $E_{1,j}^h$ are of the form z_j^h . Moreover, notice that

$$\|u_1^h - u_0^h\| = \beta'_i. \quad (51)$$

Due to (41) and (42), the hyperbolic part of $D\phi(x_{i+1})(z_1^h - z_0^h)$ is

$$\begin{pmatrix} \epsilon^\sigma A_1 & A_2 \\ A_3 & \epsilon^\sigma A_4 \end{pmatrix} \begin{pmatrix} s_1^h - s_0^h \\ u_1^h - u_0^h \end{pmatrix} = \begin{pmatrix} \epsilon^\sigma A_1 (s_1^h - s_0^h) + A_2 (u_1^h - u_0^h) \\ A_3 (s_1^h - s_0^h) + \epsilon^\sigma A_4 (u_1^h - u_0^h) \end{pmatrix}. \quad (52)$$

It follows from (43), (44), (45), (51), and (52) that

$$\begin{aligned} \|\Pi_s \circ \phi(z_1^h) - \Pi_s \circ \phi(z_0^h)\| &\leq \|\Pi_s(D\phi(x_{i+1})(z_1^h - z_0^h))\| \\ &\quad + R' \left(\|z_1^h - x_{i+1}\|^2 + \|z_0^h - x_{i+1}\|^2 \right) \\ &\leq \|\epsilon^\sigma A_1 (s_1^h - s_0^h) + A_2 (u_1^h - u_0^h)\| + \frac{1}{2} R' \zeta_i^2 \\ &\leq C_1 \epsilon^\sigma \alpha'_i + C_2 \beta'_i + R' \zeta_i^2 \end{aligned}$$

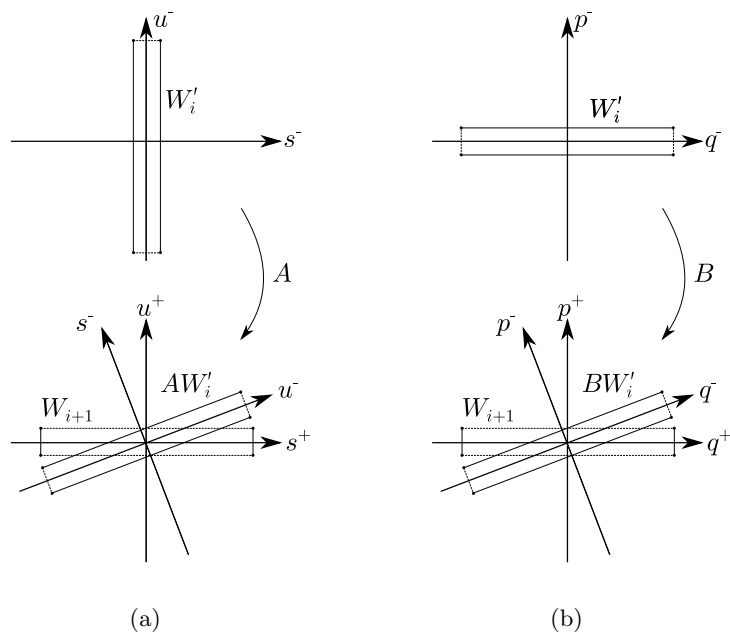


Figure 5: The figure illustrates how the window W'_i is transformed under the linearisation $D\phi(x_{i+1})$ of the change of coordinates $\phi : (s^-, u^-, q^-, p^-) \mapsto (s^+, u^+, q^+, p^+)$ at the homoclinic point x_{i+1} . In (a), we see how the hyperbolic rectangle is transformed assuming [A1], i.e. that there is an angle of order ϵ^σ between the stable and unstable manifolds at the homoclinic point x_{i+1} . In (b), we see how the inner (cylindrical) rectangle is transformed assuming [A3], i.e. that the angle between the image of a leaf \mathcal{L}_i of the foliation under the scattering map and another leaf \mathcal{L}_{i+1} is of order ϵ^v . Note that the rectangles are not centred at the origin. The dashed edges of the rectangles represent the exit sets.

and

$$\begin{aligned}
\|\Pi_u \circ \phi(z_1^h) - \Pi_u \circ \phi(z_0^h)\| &\geq \|\Pi_u(D\phi(x_{i+1})(z_1^h - z_0^h))\| \\
&\quad - R' \left(\|z_1^h - x_{i+1}\|^2 + \|z_0^h - x_{i+1}\|^2 \right) \\
&\geq \|A_3(s_1^h - s_0^h) + \epsilon^\sigma A_4(u_1^h - u_0^h)\| - \frac{1}{2}R'\zeta_i^2 \\
&\geq -\|A_3(s_1^h - s_0^h)\| + \epsilon^\sigma \|A_4(u_1^h - u_0^h)\| - R'\zeta_i^2 \\
&\geq -C_3\alpha'_i + C_4\epsilon^\sigma\beta'_i - R'\zeta_i^2
\end{aligned}$$

where ζ_i is defined by (40), where C_j is the matrix norm of A_j for $j = 1, 2, 3$, and where $C_4 > 0$ is such that $\|A_4 v\| \geq C_4 \|v\|$ for all $v \in \mathbb{R}^m$. That we can choose such a strictly positive C_4 follows from the invertibility of A_4 . Combining the previous inequalities with (46), (47), and (50) yields (36) and (37).

Now, for $j = 0, 1$, let $s_j^c \in S'_i$, $u_j^c \in U'_i$, $p_j^c \in P'_i$, and

$$q_j^* \in [c'_{i,2}, c'_{i,2} + \gamma'_i] \times \cdots \times [c'_{i,n}, c'_{i,n} + \gamma'_i].$$

Define $q_0^c = (c'_{i,1}, q_0^*)$, $q_1^c = (c'_{i,1} + \gamma'_i, q_1^*)$, and $z_j^c = (s_j^c, u_j^c, q_j^c, p_j^c)$ for $j = 0, 1$. Then all points in $E_{1,j}^c$ are of the form z_j^c , and furthermore

$$\|q_1^c - q_0^c\| = \gamma'_i \quad (53)$$

Due to (41) and (42), the cylindrical part of $D\phi(x_{i+1})(z_1^c - z_0^c)$ is

$$\begin{pmatrix} B_1 & \epsilon^v B_2 \\ \epsilon^v B_3 & B_4 \end{pmatrix} \begin{pmatrix} q_1^c - q_0^c \\ p_1^c - p_0^c \end{pmatrix} = \begin{pmatrix} B_1(q_1^c - q_0^c) + \epsilon^v B_2(p_1^c - p_0^c) \\ \epsilon^v B_3(q_1^c - q_0^c) + B_4(p_1^c - p_0^c) \end{pmatrix}. \quad (54)$$

It follows from (43), (44), (45), (53), and (54) that

$$\begin{aligned}
\|\Pi_q \circ \phi(z_1^c) - \Pi_q \circ \phi(z_0^c)\| &\leq \|\Pi_q(D\phi(x_{i+1})(z_1^c - z_0^c))\| \\
&\quad + R' \left(\|z_1^c - x_{i+1}\|^2 + \|z_0^c - x_{i+1}\|^2 \right) \\
&\leq \|B_1(q_1^c - q_0^c) + \epsilon^v B_2(p_1^c - p_0^c)\| + \frac{1}{2}R'\zeta_i^2 \\
&\leq C_5\gamma'_i + C_6\epsilon^v\delta'_i + R'\zeta_i^2
\end{aligned}$$

and

$$\begin{aligned}
\|\Pi_p \circ \phi(z_1^c) - \Pi_p \circ \phi(z_0^c)\| &\geq \|\Pi_p(D\phi(x_{i+1})(z_1^c - z_0^c))\| \\
&\quad - R' \left(\|z_1^c - x_{i+1}\|^2 + \|z_0^c - x_{i+1}\|^2 \right) \\
&\geq \|\epsilon^v B_3(q_1^c - q_0^c) + B_4(p_1^c - p_0^c)\| - \frac{1}{2}R'\zeta_i^2 \\
&\geq \epsilon^v \|B_3(q_1^c - q_0^c)\| - \|B_4(p_1^c - p_0^c)\| - R'\zeta_i^2 \\
&\geq C_7\epsilon^v\gamma'_i - C_8\delta'_i - R'\zeta_i^2
\end{aligned}$$

where ζ_i is defined by (40), where C_j is the matrix norm of B_j for $j = 5, 6, 8$, and where $C_7 > 0$ is such that $\|B_3 v\| \geq C_7 \|v\|$ for all $v \in \mathbb{R}^n$. As above, the reason that we can choose C_7 to be strictly positive is due to the invertibility of B_3 . The previous inequalities combined with (48), (49), and (50) give (38) and (39). \square

5.3 Construction of Long Sequences of Correctly Aligned Windows

In this section we show how to choose the aspect ratios of the windows at each step (see Table 1 for a summary) so that the process can be continued indefinitely, completing the proof of Theorem 2.9.

Fix $\eta > 0$ as in the statement of Theorem 2.9. Let us first show how to choose $\alpha'_i, \beta'_i, \gamma'_i, \delta'_i$ so that the inequalities of Lemma 5.3 are solvable; in particular, we require the right-hand side of inequalities (37) and (39) to be positive. Define $\kappa = \max\{\sigma, \nu\}$ and suppose we choose

$$\alpha'_i = \epsilon^{2\kappa} \alpha_*, \quad \beta'_i = \gamma'_i = \epsilon^\kappa \zeta_*, \quad \delta'_i = \epsilon^{2\kappa} \delta_*$$

where

$$\begin{cases} 0 < \zeta_* < \frac{1}{R'} \min\{C_4, C_7\}, \\ 0 < \alpha_* < \min\left\{\zeta_*, \frac{C_4}{C_3} \zeta_* \left(1 - \frac{R'}{C_4} \zeta_*\right)\right\}, \\ 0 < \delta_* < \min\left\{\zeta_*, \frac{C_7}{C_8} \zeta_* \left(1 - \frac{R'}{C_7} \zeta_*\right)\right\}. \end{cases}$$

It follows that $\beta'_i = \gamma'_i = \zeta_i$, where ζ_i is defined by (40), and so the right-hand side of (37) is

$$-C_3 \alpha'_i + C_4 \epsilon^\sigma \beta'_i - R' (\beta'_i)^2 \geq \epsilon^{2\kappa} (C_4 \beta_* - C_3 \alpha_* - R' \beta_*^2) > 0.$$

Similarly, the right-hand side of (39) is

$$C_7 \epsilon^\nu \gamma'_i - C_8 \delta'_i - R' (\gamma'_i)^2 \geq \epsilon^{2\kappa} (C_7 \gamma_* - C_8 \delta_* - R' \gamma_*^2) > 0.$$

Therefore we can choose $\beta_{i+1}, \delta_{i+1} > 0$ of order $O(\epsilon^{2\kappa})$, and $\alpha_{i+1}, \gamma_{i+1} > 0$ of order $O(\epsilon^\kappa)$ (or larger, but not smaller) so that inequalities (36), (37), (38), (39) are satisfied. Moreover, we choose

$$\alpha_{i+1}, \beta_{i+1}, \delta_{i+1} < \eta.$$

Notice that $\beta_{i+1}, \delta_{i+1}$ can be chosen as small as we like.

The next major constraint when choosing the aspect ratios is that $\hat{\gamma}_i$ must be chosen in (29) to be of order $O(1)$, and we must be able to choose it as close to 1 as we like, for the following reason. By assumption, there is a point $z \in \mathcal{L}_i$ such that $S(z) \in \mathcal{L}_{i+1}$, and $S(U \cap \mathcal{L}_i)$ is transverse to \mathcal{L}_{i+1} at $S(z)$, where

	“ ”	~	^	'
α	1	1	$\epsilon^{2\kappa}$	$\epsilon^{2\kappa}$
β	$\epsilon^{2\kappa}$	$\epsilon^{2\kappa}$	1	ϵ^κ
γ	ϵ^κ	ϵ^κ	1	ϵ^κ
δ	$\epsilon^{2\kappa}$	ϵ^ρ	$\epsilon^{2\kappa}$	$\epsilon^{2\kappa}$

Table 1: The order of the size of each rectangle at each step in the proof. In addition, we choose the iterates $N_i, M_i = O(1)$ and $K_i = O(\epsilon^{-\rho-\tau})$ to be sufficiently large.

$S : U \subseteq \Lambda \rightarrow \Lambda$ is a branch of the scattering map. It is essential that $z \in \widehat{W}_i$, and so we must be able to choose \widehat{W}_i as wide as necessary in the q -direction. This requirement can be seen mathematically in (34): if $\hat{\gamma}_i$ does not dominate $M_i T_+ \delta'_i + C M_i^2 \epsilon^k + 2\omega'_i$, then γ'_i cannot be chosen to be positive.

Suppose we have $\tilde{\delta}_i = \epsilon^\rho \delta_*$ where ρ, δ_* are to be determined. It follows from (26) that $\tilde{\delta}_i = O(\epsilon^{2\kappa})$, and so

$$\rho \geq 2\kappa. \quad (55)$$

Moreover, the first two terms on the right-hand side of (29) are

$$K_i \epsilon^{\rho+\tau} \delta_* (T_- - R \delta_* \epsilon^{\rho-\tau}). \quad (56)$$

Since we require that this is positive, we must have

$$\rho \geq \tau \quad (57)$$

and $\delta_* < \frac{T_-}{R}$. It follows from (55) and (57) that the value $\rho = \max\{2\sigma, 2\nu, \tau\}$ defined in (10) suffices.

Now, for the right-hand side of (29) to be positive we require that (56) dominates $\tilde{\gamma}_i + C K_i^2 \epsilon^k$, which is true whenever $K_i = O(\epsilon^{-\rho-\tau})$ is sufficiently large, where we have used (9). Due to (9), the inequality (30) is also solvable for $\hat{\delta}_i = O(\epsilon^{2\kappa})$. Therefore we can again choose $\delta'_i > 0$ in (35) to be of order $O(\epsilon^{2\kappa})$. Since we can choose $\tilde{\delta}_i$ to be as small as we like (due to (26)), we can subsequently choose each $\hat{\delta}_i, \delta'_i$ in (30), (35) respectively to be as small as required.

If we choose $M_i = O(1)$, then (35) is solvable. Let us explain how to find $\hat{\gamma}_i$ satisfying (34). If we define $\mathbb{T} = \mathbb{R}/\mathbb{Z}$, then $\omega'_i \leq \frac{1}{2}$, where ω'_i is defined in (31). We may assume, by slightly shifting y_i^u if necessary (to some point \tilde{y}_i^u such that $S(\tilde{y}_i^u)$ lies in a sufficiently small $O(\eta)$ -neighbourhood of \mathcal{L}_{i+1}), that $\omega'_i < \frac{1}{2}$. Notice that, by increasing K_i and shrinking $\tilde{\gamma}_i$ if necessary, we can choose $\hat{\gamma}_i$ to be as close to 1 as required. In this way we can find some positive $\gamma'_i = O(1)$ so that the left-hand side of (34) is less than 1. Therefore (34) is solvable, and we can choose $\gamma'_i = O(\epsilon^\kappa)$.

Suppose we choose $N_i = O(1)$. Then, by (23), $\tilde{\alpha}_i$ is of order 1. Since $K_i = O(\epsilon^{-\rho-\tau})$ and since $\lambda_+ \in (0, 1)$, we can choose $\hat{\alpha}_i$ to be of order $\epsilon^{2\kappa}$ in (27). Since $M_i = O(1)$, we can in turn choose α'_i to be of order $\epsilon^{2\kappa}$ due to (32). Moreover, shrinking $\hat{\alpha}_i$ in (27) allows us to shrink α'_i in (32) if necessary.

Now, since β_i is of order $O(\epsilon^{2\kappa})$, so too is $\hat{\beta}_i$ due to (24). Due to (28), since $K_i = O(\epsilon^{-\rho-\tau})$ and since $\mu_- > 1$, we can choose $\hat{\beta}_i$ of order 1. Finally, choosing $M_i = O(1)$ large enough, we can ensure that $\hat{\beta}_i - 2\nu'_i \mu_-^{-M_i} > 0$ so that (33) is solvable by a positive choice of $\beta'_i = O(\epsilon^\kappa)$.

We have thus shown that these choices can be made consistently. Combining this construction with Theorem 3.5 implies the existence of a trajectory $\{z_i\}$ as in Theorem 2.9. Since we choose M_i, N_i to be of order $O(1)$, K_i of order $O(\epsilon^{-\rho-\tau})$, the time taken to move from a neighbourhood of x_i to a neighbourhood of x_{i+1} is

$$N_i + K_i + M_i = O(\epsilon^{-\rho-\tau}).$$

In order to move a distance of order 1 in the p -direction, we must choose N consecutive leaves of the foliation connected by the scattering map where $N = O(\epsilon^{-\nu})$, and so the time is of order $\epsilon^{-\rho-\tau-\nu}$. This concludes the proof of Theorem 2.9.

6 Proof of Theorem 2.11

Let the notation be as in the statement of Theorem 2.11, and fix $\eta > 0$. Denote by $\tilde{\Sigma} = \mathbb{R}^{\ell_1} \times [0, 1]^{\ell_2}$ the universal cover of Σ . Suppose we have lifted the dynamics to the covering space $\tilde{M} = M \times \tilde{\Sigma}$ of \tilde{M} . For convenience, we do not change the notation of the lifted mappings. Fix $p_1^*, \dots, p_N^* \in [0, 1]^n$ and $\xi_1^* \in \text{Int}([0, 1]^{\ell_2})$ as in the statement of Theorem 2.11. Let $F = \tilde{G}(\cdot; \xi_1^*) \in \text{Diff}^4(M)$. By [B1], F satisfies the assumptions of Theorem 2.9. Therefore, by the proof of Theorem 2.9, there are windows $\tilde{W}_1, \dots, \tilde{W}_N \subset M$ and $n_j \in \mathbb{N}$ such that

$$d(z, \mathcal{L}_j) < \frac{\eta}{2}$$

for all $z \in \tilde{W}_j$, and \tilde{W}_j is correctly aligned with \tilde{W}_{j+1} under F^{n_j} . For convenience we drop the tilde notation and write simply W_j . By Theorem 3.6, there is a neighbourhood \mathcal{V} of F in $\text{Diff}^4(M)$ such that W_j is correctly aligned with W_{j+1} under \tilde{F}^{n_j} for each $j = 1, \dots, N-1$ and each $\tilde{F} \in \mathcal{V}$. Moreover, our assumptions on F imply that there is a $K > 0$ such that the neighbourhood \mathcal{V} is of order ϵ^K , in the sense that there is some $R > 0$ independent of ϵ such that the ball of radius $R\epsilon^K$ centred at F in $\text{Diff}^4(M)$ is contained in \mathcal{V} . Therefore there is some $a^* > 0$ (that may depend on ϵ) such that if we define

$$\Xi^* = [\xi_{1,1}^* - a^*, \xi_{1,1}^* + a^*] \times \dots \times [\xi_{1,\ell_2}^* - a^*, \xi_{1,\ell_2}^* + a^*] \subset \text{Int}([0, 1]^{\ell_2}),$$

where $\xi_1^* = (\xi_{1,1}^*, \dots, \xi_{1,\ell_2}^*)$, then $\tilde{G}(\cdot; \xi) \in \mathcal{V}$ for all $\xi \in \Xi^*$. Moreover, if we choose $L \in \mathbb{N}$ large enough, we may assume that $G(\cdot, \theta, \xi) = \tilde{G}(\cdot; \xi) + O(\epsilon^L) \in \mathcal{V}$ for all $\theta \in \mathbb{R}^{\ell_1}, \xi \in \Xi^*$.

For each $j = 2, \dots, N$ let

$$\Omega_j = \sum_{k=1}^{j-1} n_k.$$

Since $L \in \mathbb{N}$ is sufficiently large, we can find positive constants C_j (independent of ϵ) such that, if we let $\Xi_1 = \{\xi_1^*\}$, and

$$\Xi_j = [\xi_{1,1}^* - C_j \Omega_j \epsilon^L, \xi_{1,1}^* + C_j \Omega_j \epsilon^L] \times \cdots \times [\xi_{1,\ell_2}^* - C_j \Omega_j \epsilon^L, \xi_{1,\ell_2}^* + C_j \Omega_j \epsilon^L]$$

for $j = 2, \dots, N$, then

$$\Xi_j \subset \Xi^*$$

for each $j = 1, \dots, N$, and the ξ component of $\Psi^{n_j}(z, \theta, \xi)$ lies in $\text{Int}(\Xi_{j+1})$ for each $z \in W_j, \theta \in \mathbb{R}^{\ell_1}, \xi \in \Xi_j$.

Choose any $\zeta_1^\pm \in \mathbb{R}$ such that $\zeta_1^- < \zeta_1^+$, and consider the rectangle

$$\Theta_1 = [\zeta_1^-, \zeta_1^+]^{\ell_1} \subset \mathbb{R}^{\ell_1}.$$

For each $j = 2, \dots, N$ we can find $\zeta_j^\pm \in \mathbb{R}$ such that $\zeta_j^- < \zeta_j^+$, and such that if

$$\Theta_j = [\zeta_j^-, \zeta_j^+]^{\ell_1},$$

then the θ component of $\Psi^{n_j}(z, \theta, \xi)$ lies in $\text{Int}(\Theta_{j+1})$ for each $z \in W_j, \theta \in \Theta_j, \xi \in [0, 1]^{\ell_2}$.

Now, for each $j = 1, \dots, N$ choose the entry and exit sets of the rectangles Θ_j, Ξ_j to be

$$\Theta_j^+ = \partial\Theta_j, \quad \Theta_j^- = \emptyset, \quad \Xi_j^+ = \partial\Xi_j, \quad \Xi_j^- = \emptyset.$$

Define

$$\mathcal{W}_j = W_j \times \Theta_j \times \Xi_j \subset \widehat{M}$$

where the entry and exit sets \mathcal{W}_j^\pm are given by the product formula (16). Clearly the sets \mathcal{W}_j are windows. By construction, \mathcal{W}_j is correctly aligned with \mathcal{W}_{j+1} under Ψ^{n_j} for each $j = 1, \dots, N$ since the error terms are small (by assumption [B1]), and

$$d\left((z, \theta, \xi), \tilde{\mathcal{L}}_j\right) < \eta \tag{58}$$

for each $(z, \theta, \xi) \in \mathcal{W}_j$ if L is sufficiently large, where $\tilde{\mathcal{L}}_j = \tilde{\mathcal{L}}(p_j^*, \xi_j^*)$ with p_1^*, \dots, p_N^* and ξ_1^* as chosen earlier, and for some $\xi_2^*, \dots, \xi_N^* \in [0, 1]^{\ell_2}$. Therefore, combining (58) and Theorem 3.5, we see that there are $w_1, \dots, w_N \in \widehat{M}$ such that

$$w_{j+1} = \Psi^{n_j}(w_j)$$

and

$$d\left(w_j, \tilde{\mathcal{L}}_j\right) < \eta.$$

Moreover the time estimate (12) follows from the time estimate (11).

References

- [1] V. I. Arnold. Instability of dynamical systems with many degrees of freedom. *Dokl. Akad. Nauk SSSR*, 156:9–12, 1964.
- [2] P. Bernard. The dynamics of pseudographs in convex Hamiltonian systems. *J. Amer. Math. Soc.*, 21(3):615–669, 2008.
- [3] P. Bernard, V. Kaloshin, and K. Zhang. Arnold diffusion in arbitrary degrees of freedom and normally hyperbolic invariant cylinders. *Acta Math.*, 217(1):1–79, 2016.
- [4] M. Berti, L. Biasco, and P. Bolle. Drift in phase space: a new variational mechanism with optimal diffusion time. *J. Math. Pures Appl. (9)*, 82(6):613–664, 2003.
- [5] M. Berti and P. Bolle. A functional analysis approach to Arnold diffusion. *Ann. Inst. H. Poincaré Anal. Non Linéaire*, 19(4):395–450, 2002.
- [6] U. Bessi. An approach to Arnol’d’s diffusion through the calculus of variations. *Nonlinear Anal.*, 26(6):1115–1135, 1996.
- [7] S. Bolotin and D. Treschev. Unbounded growth of energy in nonautonomous Hamiltonian systems. *Nonlinearity*, 12(2):365–388, 1999.
- [8] C. Cheng. Dynamics around the double resonance. *Camb. J. Math.*, 5(2):153–228, 2017.
- [9] C. Cheng and J. Yan. Existence of diffusion orbits in a priori unstable Hamiltonian systems. *J. Differential Geom.*, 67(3):457–517, 2004.
- [10] C. Cheng and J. Yan. Arnold diffusion in Hamiltonian systems: a priori unstable case. *J. Differential Geom.*, 82(2):229–277, 2009.
- [11] A. Clarke, J. Féjóz, and M. Guardia. Drift in eccentricity and inclination in the 4 body problem (in preparation).
- [12] C. Conley and R. Easton. Isolated invariant sets and isolating blocks. *Trans. Amer. Math. Soc.*, 158:35–61, 1971.
- [13] A. Delshams, R. de la Llave, and T. Seara. A geometric approach to the existence of orbits with unbounded energy in generic periodic perturbations by a potential of generic geodesic flows of \mathbb{T}^2 . *Comm. Math. Phys.*, 209(2):353–392, 2000.
- [14] A. Delshams, R. de la Llave, and T. Seara. A geometric mechanism for diffusion in hamiltonian systems overcoming the large gap problem: heuristics and rigorous verification on a model. *Mem. Amer. Math. Soc.*, 2006.
- [15] A. Delshams, R. de la Llave, and T. M. Seara. Geometric properties of the scattering map of a normally hyperbolic invariant manifold. *Advances in Mathematics*, 217(3):1096–1153, 2008.

- [16] A. Delshams, R. de la Llave, and T. M. Seara. Instability of high dimensional Hamiltonian systems: multiple resonances do not impede diffusion. *Adv. Math.*, 294:689–755, 2016.
- [17] R. Easton. Homoclinic phenomena in hamiltonian systems with several degrees of freedom. *Journal of Differential Equations*, 29(2):241–252, 1978.
- [18] R. Easton. Orbit structure near trajectories biasymptotic to invariant tori. *Classical mechanics and dynamical systems*, 70, 1981.
- [19] R. Easton and R. McGehee. Homoclinic phenomena for orbits doubly asymptotic to an invariant three-sphere. *Indiana University Mathematics Journal*, 28(2):211–240, 1979.
- [20] N. Fenichel. Persistence and smoothness of invariant manifolds for flows. *Indiana University Mathematics Journal*, 21(3):193–226, 1971.
- [21] N. Fenichel. Asymptotic stability with rate conditions. *Indiana University Mathematics Journal*, 23(12):1109–1137, 1974.
- [22] N. Fenichel. Asymptotic stability with rate conditions, ii. *Indiana University Mathematics Journal*, 26(1):81–93, 1977.
- [23] E. Fontich and P. Martín. Differentiable invariant manifolds for partially hyperbolic tori and a lambda lemma. *Nonlinearity*, 13(5):1561–1593, 2000.
- [24] V. Gelfreich and D. Turaev. Unbounded energy growth in Hamiltonian systems with a slowly varying parameter. *Comm. Math. Phys.*, 283(3):769–794, 2008.
- [25] V. Gelfreich and D. Turaev. Arnold diffusion in a priori chaotic symplectic maps. *Comm. Math. Phys.*, 353(2):507–547, 2017.
- [26] M. Gidea and R. de la Llave. Topological methods in the instability problem of hamiltonian systems. *Discrete & Continuous Dynamical Systems*, 14(2):295, 2006.
- [27] M. Gidea, R. de la Llave, and T. M. Seara. A general mechanism of diffusion in hamiltonian systems: qualitative results. *Communications on Pure and Applied Mathematics*, 73(1):150–209, 2020.
- [28] M. Gidea and C. Robinson. Topologically crossing heteroclinic connections to invariant tori. *Journal of Differential Equations*, 193(1):49–74, 2003.
- [29] M. Gidea and C. Robinson. Symbolic dynamics for transition tori-II. In *New advances in celestial mechanics and Hamiltonian systems*, pages 95–108. Springer, 2004.
- [30] M. Gidea and P. Zgliczyński. Covering relations for multidimensional dynamical systems—II. *Journal of Differential Equations*, 202(1):59–80, 2004.

- [31] C. Golé. *Symplectic twist maps: global variational techniques*, volume 18. World Scientific, 2001.
- [32] M. W. Hirsch, C. C. Pugh, and M. Shub. Invariant manifolds. *Bulletin of the American Mathematical Society*, 76(5):1015–1019, 1970.
- [33] C. K. Jones and S.-K. Tin. Generalized exchange lemmas and orbits heteroclinic to invariant manifolds. *Discrete & Continuous Dynamical Systems-S*, 2(4):967, 2009.
- [34] V. Kaloshin and K. Zhang. *Arnold diffusion for smooth systems of two and a half degrees of freedom*, volume 208 of *Annals of Mathematics Studies*. Princeton University Press, Princeton, NJ, 2020.
- [35] R. Moeckel. Generic drift on Cantor sets of annuli. In *Celestial mechanics (Evanston, IL, 1999)*, volume 292 of *Contemp. Math.*, pages 163–171. Amer. Math. Soc., Providence, RI, 2002.
- [36] C. Pugh and M. Shub. Linearization of normally hyperbolic diffeomorphisms and flows. *Inventiones mathematicae*, 10(3):187–198, 1970.
- [37] L. Sabbagh. An inclination lemma for normally hyperbolic manifolds with an application to diffusion. *Ergodic Theory Dynam. Systems*, 35(7):2269–2291, 2015.
- [38] D. Treschev. Evolution of slow variables in a priori unstable Hamiltonian systems. *Nonlinearity*, 17(5):1803–1841, 2004.
- [39] D. Treschev. Arnold diffusion far from strong resonances in multidimensional a priori unstable Hamiltonian systems. *Nonlinearity*, 25(9):2717–2757, 2012.
- [40] P. Zgliczyński and M. Gidea. Covering relations for multidimensional dynamical systems. *Journal of Differential Equations*, 202(1):32–58, 2004.