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Pierre Berger, Sebastien Biebler

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# EMERGENCE OF WANDERING STABLE COMPONENTS 

PIERRE BERGER AND SÉBASTIEN BIEBLER*

To Mikhail Lyubich on his 60th birthday.


#### Abstract

We prove the existence of a locally dense set of real polynomial automorphisms of $\mathbb{C}^{2}$ displaying a wandering Fatou component; in particular this solves the problem of their existence, reported by Bedford and Smillie in 1991. These Fatou components have non-empty real trace and their statistical behavior is historical with high emergence. The proof is based on a geometric model for parameter families of surface real mappings. At a dense set of parameters, we show that the dynamics of the model displays a historical, high emergent, stable domain. We show that this model can be embedded into families of Hénon maps of explicit degree and also in an open and dense set of 5-parameter $C^{r}$-families of surface diffeomorphisms in the Newhouse domain, for every $2 \leq r \leq \infty$ and $r=\omega$. This implies a complement of the work of Kiriki and Soma (2017), a proof of the last Taken's problem in the $C^{\infty}$ and $C^{\omega}$-case. The main difficulty is that here perturbations are done only along finite-dimensional parameter families. The proof is based on the multi-renormalization introduced in [Ber18].


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Introduction: State of the art and main results
The aim of this paper is to reveal two new phenomena in the dynamics of analytic surface diffeomorphisms. In Sections 0.1 and 0.2 , we will describe these two phenomena in the iconic Hénon family of diffeomorphisms of the (real and complex) plane. Then in Section 0.3, we will present the general set up of our study. Finally an outline of the proof and of the manuscript is given in Section 0.4.
0.1. Wandering Fatou components. Given a holomorphic endomorphism $f$ of a complex manifold $X$, the Fatou set consists of the set of points $x \in X$ which have a neighborhood $U$ such that $\left(f^{n} \mid U\right)_{n}$ is normal. By definition, it is open and invariant by the dynamics. In particular the connected components of the Fatou set, called Fatou components, are mapped to each other under the dynamics. A main question is:

Question 0.1. Is the dynamics restricted to the Fatou set always "simple"?
To answer this "philosophical" question it is fundamental to know whether a dynamics may admit a wandering Fatou component, that is a component which is not sent by an iterate of the dynamics to a periodic one.

When $X$ is the Riemann sphere $\mathbb{P}^{1}(\mathbb{C})$, in a seminal work [Sul85], Sullivan showed that a rational mapping does not have any wandering Fatou component. This result, together with the classification of Fatou components due to Fatou himself, Siegel and Herman, allows for a complete description of the dynamics restricted to the Fatou set: the orbit of every point in the Fatou set eventually lands in an attracting basin, a parabolic basin or a rotation domain (that is, a Siegel disk or a Herman ring). Hence, in this context the answer to Question 0.1 is positive. This was later generalized by Eremenko and Lyubich [EL87] and Goldberg and Keen [GK86] who showed that entire mappings with finitely many singular values have no wandering Fatou component. The absence of wandering Fatou component for another class of entire maps has been studied in [MBRG13].

On the other hand, when $X=\mathbb{C}$, Baker [Bak76] gave (prior to Sullivan's result) the first example of an entire map displaying a wandering Fatou component, see also[EL87], [Sul85], [Her84]. In all these examples, the singular set is unbounded. Bishop introduced in [Bis15] another example displaying this time a bounded singular set. The rich history of wandering Fatou components for entire maps includes the more recent examples [MPS18, FJL17, FGJ14, KS08, Ber07].

In higher dimension, the fundamental problem of the existence of wandering Fatou components was first studied in 1991 in the work of Bedford and Smillie [Bed91, Th. 5.6], in the context of polynomial automorphisms of $\mathbb{C}^{2}$ (see also [FsS99, p. 275]). These are the polynomial maps of two complex variables which are invertible with a polynomial inverse. In this setting, a Fatou component is either the open set of points with unbounded forward orbit (which is invariant by the dynamics) or a connected component of the interior of the set $K^{+}$of points with bounded forward orbit. Later the question of the existence of wandering Fatou components was generalized by Fornaess and Sibony in [FS01, Ques. 2.2] to a holomorphic endomorphism of $\mathbb{P}^{k}(\mathbb{C})$, for $k \geq 2$. This question was completely open until the recent breakthrough of Astorg, Buff, Dujardin, Peters and Raissy $\left[\mathrm{ABD}^{+} 16\right]$ who proved, following an idea of Lyubich, the existence of a wandering Fatou component for holomorphic endomorphisms of $\mathbb{P}^{2}(\mathbb{C})$. Their example is a skew product of real polynomial maps and possesses a parabolic point (two eigenvalues equal to 1 ). Their proof uses parabolic implosion techniques. Developments of this construction allowed recently Hahn and Peters, and then Astorg, Boc Thaler and Peters to build examples of polynomial automorphisms of $\mathbb{C}^{4}$ displaying a wandering Fatou component in [HP18, ABTP19b] or to replace the parabolic fixed point by an elliptic point in [ABFP17]. However it does not seem possible to adapt these techniques to the long standing problem of the existence of such components for polynomial automorphisms of $\mathbb{C}^{2}$, for two reasons. First by [FM89], there is no polynomial automorphism of $\mathbb{C}^{2}$ which is a skew product and has a wandering Fatou component. Secondly, the Jacobian of an automorphism is constant and so the existence of a parabolic fixed point implies that the Jacobian is 1, which is incompatible with their construction (a real disk is shrunk along its orbit). In the same topic, let us mention the example of a transcendental biholomorphic map in $\mathbb{C}^{2}$ with a wandering Fatou component oscillating to infinity by Fornæss and Sibony in [FS98] and other examples of wandering domains for transcendental mappings in higher dimension in [ABTP19a]. About the non-existence of wandering Fatou components in higher dimension, a few cases study have been done proved [Lil04, Ji18] in some particular cases (skew products with a super-attracting invariant fiber).

The problem of the existence of a wandering Fatou component for polynomial automorphisms of $\mathbb{C}^{2}$ increased its interest with the recent conjecture of Bedford [Bed15] stating the existence of wandering Fatou components for polynomial automorphisms of $\mathbb{C}^{2}$ and the spectacular development [LP14] giving a complete classification of the periodic Fatou components for some polynomial automorphisms ${ }^{1}$, see also [Wei03].

The (first) main result of the present work is an answer to the original problem raised by Bedford and Smillie [BS91]. We recall that a subset is locally dense if its closure has nonempty interior.

1. called moderately dissipative, but in a different meaning as we will use.

Theorem A. For any sufficiently small real number $b \neq 0$, there exists a locally dense set of parameters $p=\left(p_{i}\right)_{0 \leq i \leq 4} \in \mathbb{R}^{5}$ for which the following polynomial automorphism displays a wandering Fatou component $\mathcal{C}$ in $\mathbb{C}^{2}$ :

$$
f_{p}:(z, w) \in \mathbb{C}^{2} \mapsto\left(z^{6}+\sum_{i=0}^{4} p_{i} \cdot z^{i}-w, b \cdot z\right) \in \mathbb{C}^{2}
$$

Moreover the Fatou component $\mathcal{C}$ satisfies the following properties:
(1) the real trace $\mathcal{C} \cap \mathbb{R}^{2}$ of $\mathcal{C}$ is non-empty,
(2) for any compact set $\mathcal{C}^{\prime} \Subset \mathcal{C}$, it holds $\lim _{n \rightarrow \infty} \operatorname{diam} f_{p}^{n}\left(\mathcal{C}^{\prime}\right)=0$,
(3) for every $(z, w) \in \mathcal{C}$, the limit set of the orbit of $(z, w)$ contains a (real) horseshoe $K$,
(4) for every $(z, w) \in \mathcal{C}$, the sequence $\left(\mathrm{e}_{n}\right)_{n}$ of empirical measures $\mathrm{e}_{n}:=\frac{1}{n} \sum_{i=0}^{n-1} \delta_{f_{p}^{i}(z, w)}$ diverges. Furthermore, there exists $\mu$ in the set of invariant probability measures $\mathcal{M}_{p}(K)$ of $f_{p} \mid K$ such that the limit set of $\left(\mathrm{e}_{n}\right)_{n}$ contains $t \cdot \mu+(1-t) \cdot \mathcal{M}_{p}(K)$ for some $0<t<1$.

Remark 0.2. Our techniques give also new examples of endomorphisms of $\mathbb{P}^{2}(\mathbb{C})$. For instance, the conclusion of Theorem A holds true also for the family of maps:

$$
f_{p}:[z, w, t] \in \mathbb{P}^{2}(\mathbb{C}) \mapsto\left[z^{6}+\sum_{i=0}^{4} p_{i} \cdot z^{i} \cdot t^{6-i}-w^{6}, b \cdot z^{6}, t^{6}\right]
$$

Remark 0.3. The set of parameters $\left(\left(p_{i}\right)_{0 \leq i \leq 4}, b\right) \in \mathbb{R}^{5} \times \mathbb{R}$ for which there is a wandering Fatou component is locally dense. Indeed, by assumption the set $P_{b}:=P \cap \mathbb{R}^{5} \times\{b\}$ has non-empty interior for every $b \neq 0$ small. And since $P$ is closed, the map $b \mapsto P_{b}$ is semi-continuous, and so continuous at generic parameter $b$.

The last item in Theorem A indicates that the statistical complexity of the dynamics is high. More specifically this implies that its emergence is high, and even that its order is positive and at least the order of the unstable dimension of $K$. We explain this in the sequel.
0.2. Statistical complexity of the dynamics: Emergence. Given a (general) differentiable dynamical system $f$ of a compact manifold $M$, since the seminal works of Boltzmann and Birkhoff, one is interested in describing the statistical behavior of most (Lebesgue) of the points . The statistical behavior of a point $x$ is described by the sequence $\left(\mathrm{e}_{n}(x)\right)_{n}$ of empirical measures

$$
\mathrm{e}_{n}(x):=\frac{1}{n} \sum_{i=0}^{n-1} \delta_{f^{i}(x)}
$$

In particular, one can wonder whether it is possible to approximate the behavior of most of the points with a finite number of probability measures [Pal00]. A natural tool to quantify this approximation is the Wasserstein distance. Recall that by the Kantorovich-Rubinstein's theorem, the $\left(1^{s t}\right)$ Wasserstein distance between two probability measures $\mu, \nu$ on $M$ is equal to:

$$
d_{W}(\mu, \nu)=\max _{\phi \in \operatorname{Lip}^{1}} \int \phi d(\mu-\nu)
$$

where Lip ${ }^{1}$ stands for the space of 1-Lipschitz real functions. Let us also recall that on the set of probability measures on any compact metric space, the Wasserstein distance induces the weak

* topology. To quantify the complexity of the statistical behavior of typical orbits for general dynamical systems, the notion of emergence has been introduced in [Ber17].

Definition. The emergence $\mathcal{E}_{f}(\epsilon)$ of $f$ at scale $\epsilon$ is the minimal number $N$ of probability measures $\left(\mu_{i}\right)_{1 \leq i \leq N}$ satisfying:

$$
\limsup _{n \rightarrow \infty} \int \min _{1 \leq i \leq N} d_{W}\left(\mathrm{e}_{n}(x), \mu_{i}\right) d \text { Leb }<\epsilon
$$

We are interested in understanding the asymptotic behavior of $\mathcal{E}_{f}(\epsilon)$ when $\epsilon \rightarrow 0$. For example, the emergence of a uniformly hyperbolic system is finite ( since there are finitely many physical probability measures whose basins cover a full measure subset of $M$ ). On the other hand, the emergence of the geodesic flow on the unit tangent bundle of the flat torus $\mathbb{R}^{n} / \mathbb{Z}^{n}$ is of the order of $\epsilon^{-n}$, hence polynomial. When the emergence is not polynomial:

$$
\limsup _{\epsilon \rightarrow 0} \frac{\log \mathcal{E}_{f}(\epsilon)}{-\log \epsilon}=\infty
$$

then the global statistical behavior of the system is deemed very complex. Let us recall:
Conjecture ([Ber17]). Super polynomial emergence is typical in many senses and in many categories of dynamical systems.

In the recent work [BB19], it has been shown that the order of the emergence $\mathcal{O E}_{f}$ :

$$
\mathcal{O} \mathcal{E}_{f}:=\limsup _{\epsilon \rightarrow 0} \frac{\log \log \mathcal{E}_{f}(\epsilon)}{-\log \epsilon}
$$

of a system $f$ on a compact manifold of dimension $d$ is at most $d$. Furthermore, it has been shown that this upper-bound is generically attained among conservative, surface $C^{\infty}$-diffeomorphisms displaying an elliptic periodic point. Note that if a dynamic has emergence of positive order, then its emergence is not polynomial. Thus this confirmed the main Conjecture of [Ber17] in the category of surface conservative diffeomorphisms with local genericity as a version of typicality.

In the situation of Theorem A, we are able to estimate the order of emergence using Theorem A.(4) and the following consequence of [BB19, Thm B]. Recall that $\mathcal{M}_{f}(K)$ is the set of invariant probability measures of $f \mid K$.

Theorem 0.4. If $K$ is a horseshoe of a $C^{1+\alpha}$-surface diffeomorphism $f$, the covering number $\mathcal{N}(\epsilon)$ at scale $\epsilon$ of $\mathcal{M}_{f}(K)$ for the Wasserstein distance has order at least the unstable dimension ${ }^{2}$ $d_{u}>0$ of $K$ :

$$
\liminf _{\epsilon \rightarrow 0} \frac{\log \log \mathcal{N}(\epsilon)}{-\log \epsilon} \geq d_{u}
$$

This together with Theorem A and Remark 0.3 imply:
Corollary A. There exists a locally dense set in the space of real polynomial automorphisms of $\mathbb{C}^{2}$ whose emergence has positive order $\mathcal{O} \mathcal{E}_{f}$ (for the canonical normalized, volume form on $\mathbb{P}^{2}(\mathbb{C})$ ):
2. This is the Hausdorff dimension of $K \cap W_{l o c}^{u}(x)$ for any $x \in K$.

To this extent, Corollary A gives a negative answer to Question 0.1 in the case of polynomial automorphisms of $\mathbb{C}^{2}$ and confirms the main conjecture of [Ber17], in the category of real polynomial automorphisms with local density as a (weak) version of typicality.

Note that in contrast to the nowhere dense example of $\left[A B D^{+} 16\right]$ whose emergence of the wandering domain is linear (as for the well-known Bowen eye), Theorem A and Corollary A exhibit a locally dense set of polynomial automorphisms of $\mathbb{R}^{2}$ whose emergence has positive order.

Let us communicate the following question raised by Ledrappier:
Question 0.5. Are there [typical] examples of transcendental maps of $\mathbb{C}$ with super-polynomial emergence?

In the same spirit, a positive answer to this question would lead to a negative answer to Question 0.1 in the category of transcendental maps of $\mathbb{C}$.
0.3 . Emergence of wandering stable components in the Newhouse domain. We recall that a point $x$ has a historic behavior if its sequence of empirical measures $\left(\mathrm{e}_{n}(x)\right)_{n}$ does not converge. The method used to prove Theorem A basically comes from real, smooth dynamical systems and is related to the following well-known:

Takens' Last Problem ([Tak08]). Are there persistent classes of smooth dynamical systems for which the set of initial states which give rise to orbits with historic behavior has positive Lebesgue measure?

A strategy to solve this problem is deduced from a scenario introduced by Colli and Vargas [CV01], and takes place in the dissipative Newhouse domain. The dissipative $C^{r}$-Newhouse domain $\mathcal{N}^{r}$ is the open set of surface diffeomorphisms leaving invariant a hyperbolic basic set displaying a $C^{r}$-robust quadratic homoclinic tangency and an area contracting periodic point (see Definition 2.13 for more details). This strategy was recently implemented by Kiriki and Soma [KS17] to solve this problem in the case of finitely differentiable diffeomorphisms of surface. For every $2 \leq r<\infty$, they proved the existence of a dense set in $\mathcal{N}^{r}$ of diffeomorphisms $f$ displaying a wandering open set $U$ whose points are asymptotic, but do not lie in the basin of a periodic orbit. Let us formalize this by introducing the following terminology:

Definition 0.6. A point $x$ is asymptotically stable if it has a neighborhood $U$ formed by asymptotic points: for every $y \in U$, it holds $d\left(f^{n}(x), f^{n}(y)\right) \rightarrow 0$ as $n \rightarrow+\infty$. A stable domain is a connected open subset formed by stable points. A stable component is a stable domain which is maximal. In other words, a stable component is a component of the set of asymptotically stable points. A stable component is wandering if it does not intersect its iterates.

Note that the first example of a wandering stable component goes back to Denjoy [Den32].
Question 0.7. Is a stable component either wandering or the basin of a periodic point?
It can be shown that the construction of Colli, Vargas, Kiriki and Soma yields a wandering stable component ${ }^{3}$. Moreover they showed that the sequence of empirical measures of points in

[^0]their stable domain does not converge. Very recently, pushing forward the construction of [KS17], Kiriki, Nakano and Soma [KNS19] showed that a dense subset of this dense set of diffeomorphisms has a super-polynomial emergence.

We introduce a geometric model for parameter families of dynamics, and we show that the scenario of [CV01, KS17] occurs at a dense set of parameters of such families, using a new perturbation method. The proof is based on a technique of composed Hénon-like renormalizations [Ber18, Thm D and Rem. 3.5] originally devoted to such pathological behaviors ${ }^{4}$. Also, we prove that this geometric model appears both in the setting of Theorem A and densely inside the Newhouse domain. This implies the second main result of this paper, which gives a positive answer to Takens' last problem in the smooth and even the real analytic category:

Theorem B. For every $r \in[2, \infty] \cup\{\omega\}$, every $f \in \mathcal{N}^{r}$ can be $C^{r}$-approximated by $\tilde{f} \in \mathcal{N}^{r}$ such that for each $C^{r}$-family $\left(f_{p}\right)_{p \in \mathbb{R}^{5}}$ in an open dense set of $\tilde{f}$ 's unfolding (i.e. $f_{0}=\tilde{f}$ ), there exists a subset of parameters $p \in \mathbb{R}^{5}$ which is dense in a neighborhood of 0 for which $f_{p}$ displays a wandering stable component $\mathcal{C}$ satisfying the following:
(1) for every $x \in \mathcal{C}$, the limit set of the orbit of $x$ contains a horseshoe $K$,
(2) for every $x \in \mathcal{C}$, the sequence $\left(\mathrm{e}_{n}\right)_{n}$ of empirical measures $\mathrm{e}_{n}:=\frac{1}{n} \sum_{i=0}^{n-1} \delta_{f_{p}^{i}(x)}$ diverges. Furthermore, there exists $\mu$ in the set of invariant probability measures $\mathcal{M}_{p}(K)$ of $f_{p} \mid K$ such that the limit set of $\left(\mathrm{e}_{n}\right)_{n}$ contains $t \cdot \mu+(1-t) \cdot \mathcal{M}_{p}(K)$ for some $0<t<1$.

Hence by Theorem 0.4 we obtain the following contribution to the main conjecture of [Ber17]:
Corollary B. In the setting of Theorem B, the emergence of $f_{p} \mid \bigcup_{n \geq 0} f_{p}^{n}(\mathcal{C})$ has order $\geq d_{u}(K)$.
This theorem extends the works [KS17, KNS19] in many respects. First it solves the $C^{\infty}$-case and the real analytic case (among analytic mappings displaying a complex extension on a uniform complex strip) among surface diffeomorphisms. Secondly, it does not only show that the set of dynamics displaying a wandering domain with historical behavior is dense in $\mathcal{N}^{r}$, but also that this phenomenon is of codimension $\leq 5$. Thirdly, it gives a stronger description of the statistical behavior of the wandering domain, proving that the emergence of this wandering component $\mathcal{C}$ has positive order (whereas they showed that the emergence is super polynomial). Finally we show that the set $\mathcal{N}^{r}$ contains families of polynomial automorphisms.

In this sense, Theorem B and its corollary confirm the above conjecture on high emergence in the category of $C^{r}$-surface diffeomorphisms, $2 \leq r \leq \infty, r=\omega$ and even in the category of real polynomial automorphisms of $\mathbb{R}^{2}$, with density of codimension $\leq 5$ as a weak version of typicality.

Theorems A, Remark 0.2 and Theorem B are actually consequences of the following result:
Theorem C. Let $\left(f_{p}\right)_{p \in \mathbb{R}^{5}}$ be a $C^{2}$-family of surface diffeomorphisms. Assume that $f_{0}$ has an area contracting periodic saddle point displaying 5 different quadratic homoclinic tangencies.
4. The renormalizations we will deal with have unbounded combinatorics, in the case of bounded combinatorics, Ou [Ou19] showed recently the nonexistence of wandering stable components.

If these quadratic tangencies unfold ${ }^{5}$ non-degenerately with $\left(f_{p}\right)_{p}$, then there is a parameter subset $\mathcal{D} \subset \mathbb{R}^{5}$ such that $\operatorname{cl}(\mathcal{D}) \ni 0, \operatorname{int}(\operatorname{cl}(\mathcal{D})) \neq \emptyset$ and for every $p \in \mathcal{D}$ the map $f_{p}$ displays a wandering stable component $\mathcal{C}$ satisfying:
(1) for every $x \in \mathcal{C}$, the limit set of the orbit of $x$ contains a horseshoe $K$,
(2) every $x \in \mathcal{C}$ has its sequence $\left(\mathrm{e}_{n}(x)\right)_{n \geq 0}$ of empirical measures which diverges. Furthermore, the limit set of $\left(\mathrm{e}_{n}\right)_{n}$ contains $t \cdot \mu+(1-t) \cdot \mathcal{M}_{p}(K)$ for some $0<t<1$ and $\mu \in \mathcal{M}_{p}(K)$.
(3) If moreover $f_{p}$ is real analytic, then $\mathcal{C}$ is the real trace of a wandering (complex) Fatou component which is a stable component of the complex extension of $f_{p}$.

All the holomorphic mappings we found with a wandering Fatou component are actually real (even if the Fatou component is an open subset of $\mathbb{C}^{2}$ or $\mathbb{P}^{2}(\mathbb{C})$ ). Nevertheless, it seems to us possible to develop the techniques presented in this work to show the following:

Conjecture 0.8. There is a locally dense set in the space of polynomial automorphisms Aut $\left(\mathbb{C}^{2}\right)$ of $\mathbb{C}^{2}$ formed by dynamics displaying a wandering Fatou component with high emergence.

Even if Theorem A is about Hénon maps of degree 6, we know that there exists a family of polynomial automorphisms of degree 5 satisfying the assumptions of Theorem C, hence possessing a wandering Fatou component. We do not know if it is possible for lower degree:

Question 0.9. Does there exist a polynomial automorphism of $\mathbb{C}^{2}$ with degree $\leq 4$ which displays a wandering Fatou component?

To answer to this question, one might follow from a mix of the parabolic renormalization techniques of $\left[\mathrm{ABD}^{+} 16\right]$ together with those of the present paper. A reasonable problem to develop these techniques is:

Problem 0.10. Give an explicit lower bound on the Hausdorff dimension of the set of parameters having a wandering Fatou component in the space of polynomial automorphisms of $\mathbb{C}^{2}$ of degree 5.
0.4. Outline of the proof and organization of the manuscript. Theorem C follows from a general result (Theorem 1.45) on a geometric model. We are going to briefly describe the mechanism behind this result. The proofs that Theorem 1.45 implies Theorem C which implies Theorems A and B, are done in Section 2.

The starting point of the proof is to work with a wild hyperbolic basic set $\Lambda$. This is a hyperbolic set such that for any perturbation of the dynamics, there is a local unstable manifold which displays a quadratic tangency with a stable manifold. For such dynamics there are sequences of points $\left(P_{i}\right)_{i \geq 0}$ in $\Lambda^{\mathbb{N}}$ which define chains of heteroclinic quadratic tangencies: the local unstable manifold $W_{l o c}^{u}\left(P_{i}\right)$ of $P_{i}$ is sent by $f$ tangent to the local stable manifold $W^{s}\left(P_{i+1}\right)$ and the tangency is quadratic. For the sake of simplicity, let us begin by assuming that each $P_{i}$ is periodic of some period $q_{i}$ and the dynamics $f^{q_{i}}$ near $P_{i}$ is affine with diagonal linear part:

$$
L_{i}: P_{i}+(x, y) \mapsto P_{i}+\left(\frac{x}{\sigma_{i}}, \lambda_{i} \cdot y\right) \quad \text { with } 0<\sigma_{i}, \lambda_{i}<1
$$

5. See Definition 2.9.

In particular the local stable manifold of $P_{i}$ is vertical and the local unstable manifold of $P_{i}$ is horizontal. Now we consider the tangencies slightly unfolded such that $f$ folds the local unstable manifold $W_{l o c}^{u}\left(P_{i}\right)$ close to a point $T_{i}$ with $L_{i+1} \circ f\left(T_{i}\right)$ close to $T_{i+1}$ (see Fig. 1). A toy model of the dynamics nearby a homoclinic tangency is given by a Hénon map:

$$
f: T_{i}+(x, y) \mapsto L_{i+1}^{-1}\left(T_{i+1}\right)+\left(x^{2}-y / 2+a_{i}, b x+c_{i}\right), \text { with } \quad\left(a_{i}, c_{i}\right)=f\left(T_{i}\right)-L_{i+1}^{-1}\left(T_{i+1}\right) .
$$



Figure 1. Toy model for the heteroclinic tangencies chain leading to a wandering domain.

Then we have:

$$
L_{i+1} \circ f: T_{i}+(x, y) \mapsto T_{i+1}+\left(\frac{x^{2}-y / 2+a_{i}}{\sigma_{i+1}}, \lambda_{i+1}\left(b x+c_{i}\right)\right)
$$

In this toy model, if there exists a sequence of neighborhoods $\left(B_{i}\right)_{i}$ of $\left(T_{i}\right)_{i}$ such that diam $B_{i} \rightarrow 0$ and $L_{i+1} \circ f\left(B_{i}\right) \subset B_{i+1}$, then the point $T_{0}$ belongs to a stable domain. To find the necessary conditions on $a_{i}, c_{i}, \sigma_{i}, \lambda_{i}$ to achieve this, we use a variation of the rescaling formula of [Ber18, Theorem D$]$. We define the rescaling factors and maps by:

$$
\gamma_{i}=\prod_{j=1}^{\infty} \sigma_{i+j}^{2^{-j}} \quad \text { and } \quad \phi_{i}:(X, Y) \mapsto T_{i}+\left(\gamma_{i} X, \gamma_{i}^{2} Y\right)
$$

We assume that $\gamma_{i}>0$. Observe that $\gamma_{i}^{2}=\sigma_{i+1} \gamma_{i+1}$. We now consider the renormalized map $\mathcal{R}_{i} f(X, Y):=\phi_{i+1}^{-1} \circ L_{i+1} \circ f \circ \phi_{i}(X, Y)$. A simple computation gives:

$$
\mathcal{R}_{i} f(X, Y)=\left(X^{2}-\frac{Y}{2}+\frac{a_{i}}{\gamma_{i}^{2}}, \frac{\lambda_{i+1}}{\gamma_{i+1}^{2}}\left(b \gamma_{i} X+c_{i}\right)\right) .
$$

This is a quadratic Hénon map. Thus the idea is to state conditions implying that its coefficients are small, so that the image by $\phi_{i}$ of the ball of radius $1 / 3$ centered at 0 is contracted and sent into the image by $\phi_{i+1}$ of the same ball by these iterations (since $\left.1 / 3^{2}+1 /(2 \times 3)<1 / 3\right)$. Namely our conditions express as:
$\left(C_{1}\right) \prod_{j \geq 0} \sigma_{j}^{2^{-j}}>0$ or equivalently $\sum_{j \geq 0} 2^{-j} \log \sigma_{j}>-\infty$.
$\left(C_{2}\right) a_{i}=o\left(\gamma_{i}^{2}\right)$ when $i$ is large; this implies that the first coordinate of $\mathcal{R}_{i} f$ is close to $X^{2}-Y / 2$.
$\left(C_{3}\right) \gamma_{i} \rightarrow 0$ and $\lambda_{i+1}=o\left(\gamma_{i+1}^{2}\right)$; this ensures that the second coordinate of $\mathcal{R}_{i} f$ is small.
To go from this toy model to actual surface diffeomorphisms (including polynomial automorphisms), we first work with a more general model where (in particular) the mapping is not necessarily linearizable at its periodic points. This model of geometrical and combinatorial type is a variation of the one used in [Ber18]. The precise definition is given in Section 1.1 under the name of systems of type $(\mathscr{A}, \mathscr{C})$. A basic and simple example of such system is depicted Fig. 4 Page 20. Then an analogue of the above conditions $\left(C_{1}\right)-\left(C_{2}\right)-\left(C_{3}\right)$ is shown to imply the conclusions of Theorem C in Theorem 1.32 (real case) and Theorem 1.34 (complex case). Their proofs occupy the whole Section 3. In Section 3.1 we recall the notion of the implicit representations for hyperbolic transformations of Palis-Yoccoz [PY01] and introduce a real-analytic generalization. These representations allow us to obtain sharp uniform distortion bounds, without assuming any linearization (i.e. non-resonance) conditions, which is very crucial for our problem. In Section 3.2, we show that the assumptions of Theorems 1.32 and 1.34 imply a normal form for the folding maps which is 'uniformly' close to be Hénon-like (and more precisely close to the assumptions of [Ber18, Theorem D and Remark 3.1]). In Section 3.3, we prove Theorems 1.32 and 1.34. The main point is to show that the domain $B_{j}$ is well included in the boxes where the implicit representations are defined and where the folding maps are defined. To handle the complex analytic case which is presently not included in [Ber18, Theorem D], we use the Cauchy inequality as a shortcut to obtain the complex distortion bounds. Details are given in Appendix B.

At this point, it is not straightforward at all to see that the assumptions of Theorems 1.32 and 1.34 are satisfied for infinitely smooth surface diffeomorphisms, let alone for polynomial automorphisms of $\mathbb{R}^{2}$. This context requires to work with a fixed $d$-parameter family of diffeomorphisms and consider perturbations only along such a family (contrarily to [KS17] where bump functions are used to perturb the dynamics). Thus in Section 1.2, we consider families $\left(F_{p}\right)_{p \in \mathbb{R}^{d}}$ of systems of type $(\mathscr{A}, \mathscr{C})$. We ask that at every parameter $p$, there are $d$ pairs of points $\left(Q_{i}^{u}, Q_{i+1}^{s}\right)_{i}$ in $\Lambda$ such that $F_{p}$ maps the local unstable manifold $W_{l o c}^{u}\left(Q_{i}^{u}\right)$ of $Q_{i}^{u}$ tangent to the local stable manifold $W_{l o c}^{s}\left(Q_{i+1}^{s}\right)$ and the tangency is quadratic. Then to apply Theorems 1.32 and 1.34 as above, we shadow each $\left(Q_{i}^{s}, Q_{i}^{u}\right)_{i}$ by a periodic orbit $P_{i}$ such that if the eigenvalues of $P_{i}$ are $\left(1 / \sigma_{i}, \lambda_{i}\right)$, the curves $W_{l o c}^{s}\left(Q_{i}^{s}\right)$ and $W_{l o c}^{s}\left(P_{i}\right)$ are $O\left(\sigma_{i}\right)$-close whereas $W_{l o c}^{u}\left(Q_{i}^{u}\right)$ and $W_{l o c}^{u}\left(P_{i}\right)$ are $O\left(\lambda_{i}\right)$-close. To satisfy $\left(C_{1}\right)$, we take:

$$
\sigma_{i} \asymp \delta^{\beta^{i}}, \quad \text { with } \beta \in(1,2) \text { and } \delta \text { small. }
$$

This implies that

$$
\gamma_{i} \asymp \prod_{j \geq 1} \sigma_{i}^{(\beta / 2)^{j}}=\sigma_{i}^{\beta /(2-\beta)}
$$

So if we do not perturb the parameter, a priori $a_{i}$ is of the order of $\sigma_{i+1}+\lambda_{i}$, which is often large compared to $\gamma_{i}^{2}=\sigma_{i}^{2 \beta /(2-\beta)}$ (since $\left.\gamma_{i} \rightarrow 0\right)$. So we have to perturb the parameter $p$ to obtain a wandering stable component (otherwise we would have found an open set of mappings with a wandering Fatou component!). But then a new difficulty appears: when we perturb the parameter $p$ so that the $i^{t h}$ - homoclinic tangency is neatly unfolded, the others are possibly badly unfolded
by a length $\sigma_{i+1}+\lambda_{i}>\sigma_{i+1}$. We recall that all $a_{j}$ must be small compared to $\gamma_{j}^{2}$. So if the points $\left(P_{j}\right)_{j}$ are inductively constructed, for $j<i$ we must have at step $i$, for every $j<i$ :

$$
\sigma_{i+1}=o\left(\gamma_{j}^{2}\right)
$$

Note that $\gamma_{j}^{2}=\gamma_{j+1} \cdot \sigma_{j+1}$. So this is equivalent to ask for:

$$
\sigma_{i+1}=o\left(\gamma_{j+1} \cdot \sigma_{j+1}\right)=o\left(\sigma_{j+1}^{2 /(2-\beta)}\right)
$$

And with $k=i-j$, we have $\sigma_{i+1}=\sigma_{j+1}^{\beta^{k}}$. So this is equivalent to ask:

$$
\beta^{k}>2 /(2-\beta)
$$

This inequality does not have any solution $\beta \in(1,2)$ for $k=0,1,2,3$, but $\beta=3 / 2$ is a solution of this inequality for every $k \geq 4$. One main new trick of this work is to work with a several parameters family: while we unfold the $i^{\text {th }}$ tangency to plug it at a neat position using one parameter, we use 3 other parameters to maintain the $(i-1)^{t h},(i-2)^{t h},(i-3)^{t h}$ folds in their neat position (see Fig. 2). But there is an extra difficulty. Once we have perturbed the parameter in this way, $W_{l o c}^{u}\left(Q_{i+1}\right)$ might be not anymore tangent to the stable local manifold $W_{l o c}^{s}\left(Q_{i+2}\right)$. Even worst, $W_{l o c}^{u}\left(Q_{i+1}\right)$ might be $\sigma_{i+1}$-far from being tangent to any stable manifold. So we need to work with a 5 -parameter family. The fifth parameters is used to restore this heteroclinic tangency. This explains why all our main results on wandering stable components are stated for 5 -parameters families. To handle such an argument we will need five groups of persistent homoclinic tangencies which unfold non-degenerately at every parameter. These are assumptions $\left(\mathbf{H}_{\mathbf{1}}\right)$ and $\left(\mathbf{H}_{\mathbf{2}}\right)$ of Definition 1.43 of unfolding of wild type $(\mathscr{A}, \mathscr{C})$, with Card $\mathscr{C}=5$.

To satisfy $\left(C_{3}\right)$ we must also have that $\lambda_{j}=o\left(\gamma_{j}^{2}\right)$. Using that $\gamma_{j} \asymp \sigma_{j}^{\beta /(2-\beta)}=\sigma_{j}^{3}$, this condition is equivalent to $\left(\lambda_{j} / \sigma_{j}\right) \cdot \sigma_{j}^{-5}=o(1)$. Actually $\lambda_{j} / \sigma_{j}$ is the Jacobian determinant along the orbit of $P_{j}$ while $\sigma_{j}^{-5}$ is approximately the fifth power of the differential of the dynamics along the orbit of $P_{j}$. Hence in order to satisfy this equation, we will require the following moderate dissipativeness hypothesis in Section 1.2:

$$
\|D f\| \cdot\|\operatorname{det} D f\|^{\epsilon}<1 \quad \text { with } \epsilon=\frac{1}{10}
$$

Let us now explain a tricky point of the parameter selection. The number $a_{j}$ is a function of $T_{j+1}$ which is itself a function of the manifold $W^{s}\left(P_{j+2}\right)$. However this manifold is not yet defined by the induction, and depends on the parameter which is moving at the step in progress. So we will need to define the point $T_{j+1}$ independently of $W^{s}\left(P_{j+2}\right)$ in a very sharp way. This will be done using an adapted family $\left(\pi_{p}\right)_{p}$ of projections, which is basically a $C^{1}$-extension of the stable lamination which is equivariant by the dynamics nearby the set $\Lambda$. Such a family of projections is actually part of the structure of unfolding of wild type $(\mathscr{A}, \mathscr{C})$. In Appendix C, we show the existence of such a family of projections given any hyperbolic set for a $C^{2}$-surface diffeomorphism, following classical techniques.

The main technical theorem in the paper is Theorem 1.45. It asserts that given a moderately dissipative, unfolding $\left(F_{p}, \pi_{p}\right)$ of wild type $(\mathscr{A}, \mathscr{C})$ with $\operatorname{Card} \mathscr{C}=5$, there is a dense set of parameters at which there is a stable component $\mathcal{C}$ satisfying the conclusions of Theorems 1.32 and 1.34


Figure 2. This figure depicts the parameter selection; the motion blur reflects how the dynamics varies when the parameter is changed at each inductive step. We assume by induction the existence of a chain of periodic points $\left(P_{j}\right)_{j<i}$ such that $W_{l o c}^{u}\left(P_{j}\right)$ is in nice position w.r.t. $W_{l o c}^{s}\left(P_{j+1}\right)$ (i.e. satisfying $\left.\left(C_{2}\right)\right), W_{l o c}^{u}\left(P_{i}\right)$ is tangent to $W_{l o c}^{s}\left(Q_{i+1}^{s}\right)$ and $W_{l o c}^{u}\left(Q_{i+1}^{u}\right)$ is tangent to the stable lamination of $\Lambda$. We consider a periodic point $P_{i+1}$ such that $W_{l o c}^{s}\left(P_{i+1}\right)$ and $W_{l o c}^{u}\left(P_{i+1}\right)$ are $\sigma_{i+1^{-}}$ close to $W_{l o c}^{s}\left(Q_{i+1}^{s}\right)$ and $W_{l o c}^{u}\left(Q_{i+1}^{u}\right)$. We handle a $\sigma_{i+1}$-small perturbation of two parameters such that $W_{l o c}^{u}\left(P_{i}\right)$ is in nice position w.r.t. $W_{l o c}^{s}\left(P_{i+1}\right)$ (i.e. satisfying $\left.\left(C_{2}\right)\right)$ and $W_{l o c}^{u}\left(P_{i+1}\right)$ is tangent to the stable lamination of $\Lambda$. Then we use three others parameters in order that $W_{l o c}^{u}\left(P_{j}\right)$ is in nice position w.r.t. $W_{l o c}^{s}\left(P_{j+1}\right)$ for every $i-3 \leq j \leq i-1$. We need a fifth parameter to make $W_{l o c}^{u}\left(P_{i+1}\right)$ tangent to a leaf $W^{s}\left(Q_{i+2}^{s}\right)$ of the stable lamination of $\Lambda$. By construction, the curve $W_{l o c}^{u}\left(P_{j}\right)$ is in nice position w.r.t. $W_{l o c}^{s}\left(P_{j+1}\right)$ for every $j \leq i-4$ since $\sigma_{i+1}=o\left(\gamma_{j}^{2}\right)$.
(wandering combinatorics) and moreover the points in $\mathcal{C}$ satisfy conclusions (1) and (2) of Theorem C ( $\omega$-limit set of $x$ containing $\Lambda$ and large emergence). Its proof occupies Section 4 . We start in Section 4.1 by showing some uniform bounds on the parameter family for the iterations near a hyperbolic set. Then in Section 4.2, we give uniform bounds on the parameter family for the normal form nearby the folding map. In Section 4.3, we proceed to the parameter selection using the argument explained above. In Section 4.4, we prove the first part of Theorem 1.45 by showing the existence of a wandering stable component by applying Theorems 1.32 and 1.34 . Finally in Section 4.4, we will use some room in our estimates to make the stable component having a large emergence and accumulate onto all the hyperbolic set $\Lambda$.

At this stage what remains to do is to show that the families of Theorems A, B and C can be put in the framework of Theorem 1.45. We say that a dynamics satisfies the geometric model if an iterate of it leaves invariant a moderately dissipative wild unfolding with Card $\mathscr{C}=5$. Corollary 1.50 of Theorem 1.45 then asserts that such a system displays a wandering stable component which accumulates onto the embedding of $\Lambda$ and has stretched exponential emergence. Moreover if the embedding and the dynamics are real analytic this wandering domain is the real trace of a wandering Fatou component.

In Section 2, we show that the families of Theorems A, B and C satisfy the geometric model. We start in Section 2.1 with a proof that the basic and simple example depicted in Fig. 4 defines a moderately dissipative wild unfolding. We remark that this simple example already enables
to find a family of polynomial automorphisms which display a wandering Fatou component, of some controlled degree. In Section 2.2 we show that Theorem C implies Theorems A and B. In order to do so we recall some celebrated results of Newhouse [New79], and its extension by [Kra92, HKY93]. These results imply that in the unfolding of a homoclinic tangency there is a wild hyperbolic set with uncountably many tangencies. In Section 2.3 we show that under the assumptions of Theorem C the geometric model occurs. The techniques of the proof are classical but rely on rather sophisticated techniques from real bifurcations theory and uniformly hyperbolic dynamics.

Let us finally emphasize that Section 2 and Appendix C using classical tools from real uniform hyperbolicity theory are completely independent from Sections 3 and 4 and Appendices A and B using tools from both real and complex analysis. All these are independent of Appendix E which is devoted to the proof of Theorem 0.4.

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## 1. The geometric model

The geometric model is devoted to offer a nice framework for showing the main theorem on the existence of a wandering stable component and which applies to the following situation: a wild horseshoe which is strongly dissipative and displays five robust homoclinic tangencies which unfold non-degenerately for a 5-parameter family, as we will see in Remark 2.11. Actually, we will show that the geometric model implies the density of the parameters at which there is a wandering stable component, in Corollary 1.50.

Let us define the functional spaces involved in this model. Let $K \subset \mathbb{R}^{n}$ and $K^{\prime} \subset \mathbb{R}^{m}$ be two compact subsets, with $n, m \geq 1$. For $r \in[1, \infty] \cup\{\omega\}$, a map $f: K \rightarrow K^{\prime}$ is of class $C^{r}$ if it can be extended to a $C^{r}$-map from an open neighborhood of $K$ to $\mathbb{R}^{m}$. It is a $C^{r}$-diffeomorphism if $f$ is bijective and $f^{-1}$ is of class $C^{r}$. Then $n=m$. We recall that $C^{\omega}$ is the class of real analytic maps. Likewise, given two compact subsets $K \subset \mathbb{C}^{n}$ and $K^{\prime} \subset \mathbb{C}^{m}$, a map $f: K \rightarrow K^{\prime}$ is holomorphic if it can be extended to a holomorphic map from an open neighborhood of $K$ to $\mathbb{C}^{m}$. It is a biholomorphism if $f$ is bijective and $f^{-1}$ is holomorphic (and again $n=m$ ).
1.1. System of type $(\mathscr{A}, \mathscr{C})$. In this subsection, we introduce the notion of hyperbolic map of type $\mathscr{A}$ through the definition of hyperbolic transformations. One should think about a hyperbolic map of type $\mathscr{A}$ as a horseshoe $\Lambda$ and about hyperbolic transformations as the restriction of the dynamics to a neighborhood of a rectangle of a Markov partition (see Propositions 1.13 and 2.25).

Let $I:=[-1,1]$ with topological boundary $\partial I:=\{-1,1\}$. We fix for this subsection $\rho>0$ and we put $\tilde{I}:=I+i \cdot[-\rho, \rho]$ with topological boundary $\partial \tilde{I} \subset \mathbb{C}$. Put $Y^{e}:=I^{2}, \tilde{Y}^{e}:=\tilde{I}^{2} \subset \mathbb{C}^{2}$ and:

$$
\partial^{s} Y^{e}:=(\partial I) \times I \quad, \quad \partial^{u} Y^{e}:=I \times \partial I \quad, \quad \partial^{s} \tilde{Y}^{e}:=(\partial \tilde{I}) \times \tilde{I} \quad, \quad \partial^{u} \tilde{Y}^{e}:=\tilde{I} \times \partial \tilde{I}
$$

The sets $Y^{e}$ and $\tilde{Y}^{e}$ will be associated to the transformations $F^{e}$ equal to the identity. The definition of hyperbolic transformations depends on fixed $\theta>0$ and $\lambda>1$. For convenience, fix:

$$
\theta=\frac{1}{2} \quad \text { and } \quad \lambda=2
$$

Definition 1.1 (Box). A box is a subset $Y$ of $\mathbb{R}^{2}$ which is diffeomorphic to $Y^{e}$ and of the form:

$$
Y=\left\{(x, y) \in Y^{e}: \phi^{-}(y) \leq x \leq \phi^{+}(y)\right\}
$$

where $\phi^{-}$and $\phi^{+}$are $C^{1}$-functions on I such that $\phi^{-}<\phi^{+}$with $\left\|\phi^{ \pm}\right\|_{C^{0}} \leq 1$ and $\left\|D \phi^{ \pm}\right\|_{C^{0}}<\theta$, for $\pm \in\{-,+\}$. We put $\partial^{u} Y=\partial^{u} Y^{e} \cap Y$ and $\partial^{s} Y=\operatorname{cl}\left(\partial Y \backslash \partial^{u} Y\right)$.

We notice that the set $Y^{e}$ is a box. Also in general, a box $Y$ satisfies:

$$
\partial^{u} Y=\left\{(x, y) \in I \times \partial I: \phi^{-}(y) \leq x \leq \phi^{+}(y)\right\} \quad \text { and } \quad \partial^{s} Y=\left\{(x, y) \in I^{2}: x \in\left\{\phi^{-}(y), \phi^{+}(y)\right\}\right\}
$$

For analytic tranformations we will consider the following complex extensions of boxes.
Definition $1.2\left(C_{\rho}^{\omega}\right.$-box). A $C_{\rho}^{\omega}$-box $\tilde{Y}$ is a subset of $\tilde{Y}^{e}$ such that $\tilde{Y} \cap Y^{e}$ is a box and $\tilde{Y}$ is the range of a biholomorphism of the form $\zeta:(z, w) \in \tilde{I}^{2} \mapsto(\mathcal{Z}(z, w), w) \in \tilde{Y}$ satisfying $\left|\partial_{w} \mathcal{Z}\right|<\theta$. Let $\partial^{u} \tilde{Y}=\partial^{u} \tilde{Y}^{e} \cap \tilde{Y}$ and $\partial^{s} \tilde{Y}=\operatorname{cl}\left(\tilde{Y} \backslash \partial^{u} \tilde{Y}\right)$.

Similarly we notice that $\tilde{Y}^{e}$ is a $C_{\rho}^{\omega}$-box. Also in general, a $C_{\rho}^{\omega}$-box $\tilde{Y}$ satisfies:

$$
\partial^{u} \tilde{Y}:=\zeta(\tilde{I} \times \partial \tilde{I}) \quad \text { and } \quad \partial^{s} \tilde{Y}:=\zeta((\partial \tilde{I}) \times \tilde{I})
$$

The definition of hyperbolic transformation will involve cones.
Definition 1.3 (Cones). We define the two following cones $\chi_{h}$ and $\chi_{v}$ :

$$
\chi_{h}:=\left\{\left(u_{x}, u_{y}\right) \in \mathbb{R}^{2}:\left|u_{y}\right|<\theta \cdot\left|u_{x}\right|\right\} \quad \text { and } \quad \chi_{v}:=\left\{\left(u_{x}, u_{y}\right) \in \mathbb{R}^{2}:\left|u_{x}\right|<\theta \cdot\left|u_{y}\right|\right\} .
$$

The real cones $\chi_{h}$ and $\chi_{v}$ admit a canonical extension to $\mathbb{C}^{2}$ :

$$
\tilde{\chi}_{h}=\left\{\left(u_{z}, u_{w}\right) \in \mathbb{C}^{2}:\left|u_{w}\right|<\theta \cdot\left|u_{z}\right|\right\} \quad \text { and } \quad \tilde{\chi}_{v}=\left\{\left(u_{z}, u_{w}\right) \in \mathbb{C}^{2}:\left|u_{z}\right|<\theta \cdot\left|u_{w}\right|\right\} .
$$

Definition 1.4. A hyperbolic transformation $(Y, F)$ is the data of a box $Y$ and a $C^{2}$-diffeomorphism from $Y$ onto its image in $Y^{e}$ such that either $Y=Y^{e}$ and $F$ is the identity or:
(1) $F(Y) \subset Y^{e}$ and $F\left(\partial^{s} Y\right) \subset \partial^{s} Y^{e}$ whereas $Y \cap \partial^{s} Y^{e}=\emptyset \quad$ and $\quad F(Y) \cap \partial^{u} Y^{e}=\emptyset$.
(2) For every $z \in Y$, every non-zero vector in the complement of $\chi_{v}$ is sent into $\chi_{h}$ by $D_{z} F$ and has its first coordinate which is more than $\lambda$-expanded by $D_{z} F$.
(3) For every $z \in F(Y)$, every non-zero vector in the complement of $\chi_{h}$ is sent into $\chi_{v}$ by $D_{z} F^{-1}$ and has its second coordinate which is more than $\lambda$-expanded by $D_{z} F^{-1}$.

Here is the complex counterpart of the latter notion.
Definition 1.5. A $C_{\rho}^{\omega}$-hyperbolic transformation $(\tilde{Y}, F)$ is the data of a $C_{\rho}^{\omega}$-box $\tilde{Y}$ and a biholomorphism $F$ from $\tilde{Y}$ onto its image in $\tilde{Y}^{e}$ such that either $(\tilde{Y}, F)=\left(\tilde{Y}^{e}, i d\right)$ or:
(1) $F(\tilde{Y}) \subset \tilde{Y}^{e}$ and $F\left(\partial^{s} \tilde{Y}\right) \subset \partial^{s} \tilde{Y}^{e}$ whereas $\tilde{Y} \cap \partial^{s} \tilde{Y}^{e}=\emptyset$ and $F(\tilde{Y}) \cap \partial^{u} \tilde{Y}^{e}=\emptyset$.
(2) For every $z \in \tilde{Y}$, every non-zero vector in the complement of $\tilde{\chi}_{v}$ is sent into $\tilde{\chi}_{h}$ by $D_{z} F$ and has its first coordinate's modulus which is more than $\lambda$-expanded by $D_{z} F$.
(3) For every $z \in F(\tilde{Y})$, every non-zero vector in the complement of $\tilde{\chi}_{h}$ is sent into $\tilde{\chi}_{v}$ by $D_{z} F^{-1}$ and has its second coordinate modulus which is more than $\lambda$-expanded by $D_{z} F^{-1}$.
(4) The restriction of $F$ to the box $Y:=\tilde{Y} \cap \mathbb{R}^{2}$ is a real analytic hyperbolic transformation.


Figure 3. A hyperbolic transformation.
We denote also by $F$ its restriction $F \mid Y$. The pair $(\tilde{Y}, F)$ is called the complex extension or $C_{\rho}^{\omega}$-extension of $(Y, F)$.

We define the following operation on hyperbolic transformations.
Definition 1.6 (*-product). Given two hyperbolic transformations $(Y, F)$ and $\left(Y^{\prime}, F^{\prime}\right)$, we define:

$$
(Y, F) \star\left(Y^{\prime}, F^{\prime}\right):=\left(Y \cap F^{-1}\left(Y^{\prime}\right), F^{\prime} \circ F\right)
$$

Given two $C_{\rho}^{\omega}$-hyperbolic transformations $(\tilde{Y}, F)$ and $\left(\tilde{Y}^{\prime}, F^{\prime}\right)$, we define:

$$
(\tilde{Y}, F) \star\left(\tilde{Y}^{\prime}, F^{\prime}\right):=\left(\tilde{Y} \cap F^{-1}\left(\tilde{Y}^{\prime}\right), F^{\prime} \circ F\right)
$$

The $\star$-product is a binary operation on the set of hyperbolic transformations:
Proposition 1.7. The $\star$-product $(Y, F) \star\left(Y^{\prime}, F^{\prime}\right)$ of any pair of hyperbolic transformations is a hyperbolic transformation. The $\star$-product $(\tilde{Y}, F) \star\left(\tilde{Y}^{\prime}, F^{\prime}\right)$ of any pair of $C_{\rho}^{\omega}$-hyperbolic transformations is a $C_{\rho}^{\omega}$-hyperbolic transformation, moreover $(\tilde{Y}, F) \star\left(\tilde{Y}^{\prime}, F^{\prime}\right)$ is a complex extension of $(Y, F) \star\left(Y^{\prime}, F^{\prime}\right)$ if $(\tilde{Y}, F)$ and $\left(\tilde{Y}^{\prime}, F^{\prime}\right)$ are respective complex extensions of $(Y, F)$ and $\left(Y^{\prime}, F^{\prime}\right)$.

Proof. In [PY01, $\S 3.2 .1]$ it is noticed that the proof is immediate in the real case. The proof for the complex case is postponed to Appendix A.

The *-product is associative. In particular, it will not be necessary to use brackets while taking $\star$-products of more than two factors. Also for every hyperbolic transformation $(Y, F)$ it holds:

$$
(Y, F) \star\left(Y^{e}, F^{e}\right)=\left(Y^{e}, F^{e}\right) \star(Y, F)=(Y, F),
$$

where $F^{e}$ is the identity of $Y^{e}$. The same occurs for $C_{\rho}^{\omega}$-hyperbolic transformation. This implies:

Proposition 1.8. The set of (resp. $C_{\rho}^{\omega}-$ ) hyperbolic transformations endowed with the operation $\star$ is a monoid with neutral element $\left(Y^{e}, F^{e}\right)\left(\right.$ resp. $\left.\left(\tilde{Y}^{e}, F^{e}\right)\right)$.

Let $(\mathscr{V}, \mathscr{A})$ be an oriented graph where $\mathscr{V}$ is a finite set of vertexes and $\mathscr{A}$ a finite set of arrows. The maps $a: \mathscr{A} \rightarrow \mathscr{V}$ and $t: \mathscr{A} \rightarrow \mathscr{V}$ associate to each arrow $a \in \mathscr{A}$, its origin vertex $a(a)$ and its target vertex $t(a) \in \mathscr{V}$. Let $\mathscr{B} \subset \mathscr{A}^{2}$ be the set of pairs $\left(a, a^{\prime}\right)$ of arrows such that $t(a)=a\left(a^{\prime}\right)$. Then we will always assume that the subshift of finite type defined by $(\mathscr{A}, \mathscr{B})$ is transitive and has positive entropy. This is equivalent to say that the graph is connected and there are two different paths with the same length and the same origin and target vertexes.

The following is inspired from the notion of regular Cantor set [MY10] and of the Markovian structure of [PY09].

Definition 1.9. A hyperbolic map of type $\mathscr{A}$ is a map of the form:

$$
F^{\mathscr{A}}:(z, a(a)) \in \bigsqcup_{a \in \mathscr{A}} Y^{a} \times\{a(a)\} \mapsto\left(F^{a}(z), t(a)\right) \in Y^{e} \times \mathscr{V}
$$

where $\left(Y^{a}, F^{a}\right)_{a \in \mathscr{A}}$ is a finite family of hyperbolic transformations $\left(Y^{a}, F^{a}\right)$ satisfying that $Y^{a}$ is disjoint from $Y^{a^{\prime}}$ for every $a \neq a^{\prime} \in \mathscr{A}$ such that $a(a)=a\left(a^{\prime}\right)$.

The map $F^{\mathscr{A}}$ is $C_{\rho}^{\omega}$-hyperbolic if there exists a family $\left(\tilde{Y}^{a} \times\{a(a)\}\right)_{a \in \mathscr{A}}$ of complex extensions $\left(\tilde{Y}^{a}, F^{a}\right)$ of $\left(Y^{a}, F^{a}\right)$ satisfying that $\tilde{Y}^{a}$ is disjoint from $\tilde{Y}^{a^{\prime}}$ for every $a \neq a^{\prime} \in \mathscr{A}$ such that $a(a)=a\left(a^{\prime}\right)$. We put:

$$
D(\mathscr{A}):=\bigsqcup_{a \in \mathscr{A}} Y^{a} \times\{a(a)\} \quad \text { and } \quad \text { in the } C_{\rho}^{\omega}-\operatorname{case} \tilde{D}(\mathscr{A}):=\bigsqcup_{a \in \mathscr{A}} \tilde{Y}^{a} \times\{a(a)\}
$$

We notice that a hyperbolic map of type $\mathscr{A}$ is a local $C^{2}$-diffeomorphism.
Example 1.10. Let $(\mathscr{V}, \mathscr{A})$ be the graph with a single vertex $\mathscr{V}=\{\diamond\}$ and $N \geq 1$ distinct arrows $\mathscr{A}=\left\{a_{1}, \cdots, a_{N}\right\}$ with both origin and target vertexes equal to $\diamond$. Let $\delta>0$ be small compared to $1 / N$. For every $1 \leq j \leq N$, put:

$$
I_{j}:=\left[\frac{2}{N}(j-1)-1+\delta^{2}, \frac{2}{N} j-1-\delta^{2}\right]
$$

For $1 \leq j \leq N$, let $C_{j}$ be the affine, orientation preserving map from $I$ onto $I_{j}$, let:

$$
Y^{a_{j}}:=I_{j} \times I \quad \text { and } \quad F^{a_{j}}:=(x, y) \in Y^{a_{j}} \mapsto\left(C_{j}^{-1}(x), \sqrt{\delta} \cdot C_{j}(y)\right)
$$

We notice that $\left(Y^{a_{j}}, F^{a_{j}}\right)$ is a hyperbolic transformation for every $1 \leq j \leq N$. The domains $\left(Y^{a_{j}}\right)_{1 \leq j \leq N}$ are disjoint since the intervals $\left(I_{j}\right)_{j}$ are disjoint. Also the ranges of $\left(F^{a_{j}}\right)_{j}$ are $\left(I \times \sqrt{\delta} \cdot I_{j}\right)_{j}$ which are disjoint from the one other. Thus the following diffeomorphism is hyperbolic of type $\mathscr{A}$ :

$$
F^{s l}:(z, \diamond) \in \bigsqcup_{a \in \mathscr{A}} Y^{a} \times\{\diamond\} \mapsto\left(F^{a}(z), \diamond\right) \text { if } z \in Y^{a}
$$

This hyperbolic map with $N=6$ is depicted on Fig. 4.
We are going to associate geometric and dynamical objects to the following symbolic objects:

Definition 1.11. Let $\mathscr{A}^{*} \subset \mathscr{A}^{(\mathbb{N})}$ be the set of admissible finite words $c:=a_{1} \cdots a_{k} \in \mathscr{A}^{k}:$ for each $j<k$ we have $\left(a_{j}, a_{j+1}\right) \in \mathscr{B}$. We denote by e the empty word in $\mathscr{A}^{*}$. The origin and target vertexes of $c$ are defined by the maps:

$$
a: a_{1} \cdots a_{k} \in \mathscr{A}^{*} \backslash\{e\} \mapsto a\left(a_{1}\right) \in \mathscr{V} \quad \text { and } \quad t: a_{1} \cdots a_{k} \in \mathscr{A}^{*} \backslash\{e\} \mapsto t\left(a_{k}\right) \in \mathscr{V}
$$

If the pair $\left(c, c^{\prime}\right) \in\left(\mathscr{A}^{*}\right)^{2}$ is admissible $\left(t(c)=a\left(c^{\prime}\right)\right)$, we denote by $c \cdot c^{\prime} \in \mathscr{A}^{*}$ the concatenation of these two words. The length $|c|$ of $c \in \mathscr{A}$ is its number of letter.

There are three ways of taking limits of these admissible words: by both sides, the right side and the left side. This defines three sets.
Definition 1.12. Let $\overleftrightarrow{A}$ be the set of sequences $\left(a_{i}\right)_{i \in \mathbb{Z}} \in \mathscr{A}^{\mathbb{Z}}$ which are admissible: $\left(a_{j}, a_{j+1}\right) \in$ $\mathscr{B}$ for every $j$; it is the space of the orbits of the aforementioned transitive subshift of finite type.

Let $\overrightarrow{\mathscr{A}}$ be the set of sequences $s:=\left(a_{i}\right)_{i \geq 0} \in \mathscr{A}^{\mathbb{N}}$ such that $\left(a_{j}, a_{j+1}\right) \in \mathscr{B}$ for every $j \geq 0$.
Let $\overleftarrow{\mathscr{A}}$ be the set of sequences $u:=\left(a_{i}\right)_{i<0} \in \mathscr{A}^{\mathbb{Z}}$. such that $\left(a_{-j-1}, a_{-j}\right) \in \mathscr{B}$ for every $j>0$
We endow the sets $\overleftarrow{\mathscr{A}}, \overrightarrow{\mathscr{A}}$ and $\overleftrightarrow{\mathscr{A}}$ with the topologies induced by the product topologies of $\mathscr{A}^{\mathbb{Z}_{-}^{*}}$, $\mathscr{A}^{\mathbb{N}}$ and $\mathscr{A}^{\mathbb{Z}}$ which are compact by Tychonoff's theorem.

The map $a$ can be canonically extended from $\mathscr{A}^{*}$ to $\overrightarrow{\mathscr{A}}$ by $a\left(\left(\omega_{i}\right)_{i \geq 0}\right):=a\left(\omega_{0}\right)$. We extend similarly $t$ to $\overleftarrow{\mathscr{A}}$ by $t\left(\left(a_{i}\right)_{i<0}\right):=t\left(a_{-1}\right)$. We extend canonically the concatenation operation to the sets $\mathscr{B}^{*}, \overrightarrow{\mathscr{B}}$ and $\overleftarrow{\mathscr{B}}$ of pairs $(c, d)$ in respectively $\left(\mathscr{A}^{*}\right)^{2}, \mathscr{A}^{*} \times \overrightarrow{\mathscr{A}}$ and $\overleftarrow{\mathscr{A}} \times \mathscr{A}^{*}$ such that $t(c)=a(d)$. For $b \in \mathscr{A}^{*}$, we denote:

$$
\begin{align*}
& \mathscr{A}^{*} \cdot b:=\left\{c \cdot b \in \mathscr{A}^{*}:(c, b) \in \mathscr{B}^{*}\right\} \quad \text { and } \quad b \cdot \overrightarrow{\mathscr{A}}:=\{b \cdot s\in \overrightarrow{\mathscr{A}}:(b, s) \in \overrightarrow{\mathscr{B}}\}  \tag{1.1}\\
& \text { and } \quad \overleftarrow{\mathscr{A}} \cdot b:=\{u \cdot b \in \overleftarrow{\mathscr{A}}:(u, b) \in \overleftarrow{\mathscr{B}}\}
\end{align*}
$$

Given a hyperbolic map of type $\mathscr{A}$, by associativity of the $\star$-product, we denote for every $c:=a_{1} \cdots a_{k} \in \mathscr{A}^{*}:$

$$
\begin{equation*}
\left(Y^{e}, F^{e}\right):=\left(Y^{e}, i d\right) \quad \text { and } \quad\left(Y^{c}, F^{c}\right):=\left(Y^{a_{1}}, F^{a_{1}}\right) \star \cdots \star\left(Y^{a_{k}}, F^{a_{k}}\right) \tag{1.2}
\end{equation*}
$$

which is a hyperbolic transformation by Proposition 1.7. If $F$ is $C_{\rho}^{\omega}$-hyperbolic of type $\mathscr{A}$, we denote similarly $\left(\tilde{Y}^{e}, F^{e}\right):=\left(\tilde{Y}^{e}, i d\right)$ and $\left(\tilde{Y}^{c}, F^{c}\right)$ the $\star$-product of the $C_{\rho}^{\omega}$-hyperbolic transformations associated to the letters of $c$.

We recall that a hyperbolic basic set is a hyperbolic compact set which is transitive and locally maximal.

Proposition 1.13. The maximal invariant $\Lambda=\bigcap_{n \in \mathbb{Z}}\left(F^{\mathscr{A}}\right)^{n}\left(Y^{e} \times \mathscr{V}\right)$ is a hyperbolic basic set.
Proof. First note that $\Lambda$ is compact since it is an intersection of compact sets. Furthermore, it is locally maximal since:

$$
\Lambda \subset F^{\mathscr{A}}\left(Y^{e} \times \mathscr{V}\right) \cap\left(F^{\mathscr{A}}\right)^{-1}\left(Y^{e} \times \mathscr{V}\right)=\bigcup_{\left(a, a^{\prime}\right) \in \mathscr{A}^{2}: a(a)=t\left(a^{\prime}\right)} Y^{a} \cap F^{a^{\prime}}\left(Y^{a^{\prime}}\right) \times\{a(a)\}
$$

which is strictly included in $Y^{e} \times \mathscr{V}$ by property (1) of hyperbolic transformation's Definition 1.4. Furthermore it is hyperbolic by the cone properties (2) and (3) of Definition 1.4. The transitivity of $\Lambda$ comes from the one of the shift $\overleftrightarrow{\Lambda}$ defined by $(\mathscr{A}, \mathscr{B})$ and the next remark.

Remark 1.14. If $F^{\mathscr{A}}$ is a diffeomorphism, the dynamics of $\Lambda$ is conjugate to the shift on $\overleftrightarrow{\mathscr{A}}$. In general, the dynamics of the inverse limit $\overleftrightarrow{\Lambda}:=\left\{\left(z_{i}, v_{i}\right)_{i \in \mathbb{Z}} \in \Lambda^{\mathbb{Z}}:\left(z_{i+1}, v_{i+1}\right)=F^{\mathscr{A}}\left(z_{i}, v_{i}\right)\right\}$ of $\Lambda$ is conjugate to the shift on $\overleftrightarrow{\mathscr{A}}$ via the map

$$
h:\left(z_{i}, v_{i}\right)_{i \in \mathbb{Z}} \in \overleftrightarrow{\Lambda} \mapsto\left(a_{j}\right)_{j \in \mathbb{Z}} \in \overleftrightarrow{\mathscr{A}} \quad \text { with } a_{j} \in \mathscr{A} \text { such that } z_{j} \in Y^{a_{j}} \text { and } v_{j}=\sigma\left(a_{j}\right)
$$

Conversely, we will see in Proposition 2.25 Page 38 that modulo iterations, any $C^{2}$-horseshoe is conjugate to a hyperbolic map $F^{\mathscr{A}}$ of a certain type $\mathscr{A}$.

Definition 1.15. Given a hyperbolic map of type $\mathscr{A}$, for $s=\left(a_{i}\right)_{i \geq 0} \in \overrightarrow{\mathscr{A}}$ and $u=\left(a_{i}\right)_{i<0} \in \overleftarrow{\mathscr{A}}$, we denote:

$$
W^{j}:=\bigcap_{i \geq 0} Y^{a_{0} \cdots a_{i}} \quad \text { and } \quad W^{u}:=\bigcap_{i>0} F^{a_{-i} \cdots a_{-1}}\left(Y^{a_{-i} \cdots a_{-1}}\right) .
$$

The following is immediate:
Proposition 1.16. For every $s \in \vec{A}$, the set $W^{s} \times\{a(\jmath)\}$ is a local stable manifold of the hyperbolic set $\Lambda$ with tangent spaces in $\chi_{v}$ and endpoints in the two different components of $\partial^{u} Y^{e}$.

For every $u \in \overleftarrow{\mathscr{A}}$, the set $W^{u} \times\{t(u)\}$ is a local unstable manifold of $\Lambda$ with tangent spaces in $\chi_{h}$ and with endpoints in the two different components of $\partial^{s} Y^{e}$.

By Proposition 1.16, for every $u \in \overleftarrow{\mathscr{A}}$ and $s \in \overrightarrow{\mathscr{A}}$, there are $\theta$-Lipschitz functions $w^{u}$ and $w^{s}$ on $I$ whose graphs satisfy Graph $w^{u}=W^{u}$ and ${ }^{t}$ Graph $w^{s}=W^{s}$.

Proposition 1.17. The maps $x \in I \mapsto w^{u}(x) \in \mathbb{R}$ and $y \in I \mapsto w^{s}(y) \in \mathbb{R}$ are of class $C^{2}$ and depend continuously on $u$ and $s$ in the $C^{2}$-topology.

Proof. This is the continuity of the local stable and unstable manifolds of the hyperbolic set $\Lambda$.
The following enables to consider quadratic homoclinic tangencies of the basic set $\Lambda$. It is defined using a finite set $\mathscr{C}$ of (new) arrows $\square$ between vertexes of $\mathscr{V}$.

Definition 1.18. Given $\square \in \mathscr{C}$, a folding transformation $\left(Y^{\square}, F^{\square}\right)$ is the data of a filled square $Y^{\square}=I^{\square} \times J^{\square}$ included in $Y^{e} \backslash \partial^{s} Y^{e}$ and a $C^{2}$-diffeomorphism $F^{\square}$ from $Y^{\square}$ onto its image in $Y^{e} \backslash \partial^{u} Y^{e}$, which satisfies:
(1) The following subset is nonempty and clopen:
$\overleftarrow{\mathscr{A}}_{\square}:=\left\{u \in \overleftarrow{\mathscr{A}}: t(u)=\sigma(\square), W^{u} \cap Y^{\square} \neq \emptyset \& W^{u} \cap \partial^{u} Y^{\square}=\emptyset\right\} \quad$ with $\partial^{u} Y^{\square}:=I^{\square} \times \partial J^{\square}$
(2) With $\overrightarrow{\mathscr{A}}_{\square}:=\{s \in \overrightarrow{\mathscr{A}}: \sigma(s)=t(\square)\}$, for any $u \in \overleftarrow{\mathscr{A}}_{\square}$ and $s \in \overrightarrow{\mathscr{A}}_{\square}$, if the curves $W^{u}$ and $\left(F^{\unlhd}\right)^{-1}\left(W^{s}\right)$ are tangent, then the tangency is unique, quadratic and lies in int $Y_{\square}$.
The folding transformation is of class $C_{\rho}^{\omega}$ if $F^{\square}$ extends to a biholomorphism from $\tilde{Y}^{\boxminus}:=Y^{\square}+$ $[-i \cdot \rho, i \cdot \rho]^{2}$ onto its image in $\mathbb{C}^{2}$.

Here is the main dynamical object of this subsection.
Definition 1.19. A system of type $(\mathscr{A}, \mathscr{C})$ is a $C^{2}$-local diffeomorphism $F$ of the form:

$$
F:(x, v) \in D(\mathscr{A}) \sqcup D(\mathscr{C})= \begin{cases}F^{\mathscr{A}}(x, v) & \text { if }(x, v) \in D(\mathscr{A}) \\ F^{\mathscr{C}}(x, v) & \text { if }(x, v) \in D(\mathscr{C})\end{cases}
$$

for a hyperbolic map $F^{\mathscr{A}}$ of type $\mathscr{A}$ and the map $F^{\mathscr{C}}$ is the disjoint union of a finite family of folding transformations $\left(Y^{\boxminus}, F^{\boxminus}\right)_{\square \in \mathscr{C}}$ for $F^{\mathscr{Q}}$ by

$$
F^{\mathscr{C}}:(z, \sigma(\square)) \in D(\mathscr{C}):=\bigsqcup_{\square \in \mathscr{C}} Y^{\bullet} \times\{\sigma(\square)\} \mapsto\left(F^{\bullet}(z), t(\square)\right) \in Y^{e} \times \mathscr{V}
$$

and such that for every $\square \neq \square^{\prime} \in \mathscr{C}$, the sets $\overleftarrow{\mathscr{A}}_{\square}$ and $\overleftarrow{\mathscr{A}}_{\square}$, are disjoint
Example 1.20. Let $F^{s}$ be the example of hyperbolic map of type $\mathscr{A}$ of Example 1.10 with $N$ even and $\delta>0$ small. We recall that $(\mathscr{V}, \mathscr{A})$ is a graph with a unique vertex $\diamond$ and $N$ arrows $a_{j}$. Let $\mathscr{C}=\left\{\square_{j}: 1 \leq j \leq N-1\right\}$ be a set of $N-1$ other arrows satisfying $a\left(\square_{j}\right)=t\left(\square_{j}\right)=\diamond$. For $p$ in:

$$
\mathcal{P}:=\left\{\left(p_{i}\right)_{i} \in \mathbb{R}^{N-1}: \forall 1 \leq i \leq N-1, \quad p_{i} \in I_{i+1} \text { and } p_{i} \text { is } \sqrt[3]{\delta} \text {-distant from } \partial I_{i+1}\right\}
$$

we define a folding transformation $F_{p}^{\square_{j}}$ for $F^{\mathscr{A}}$ by:

$$
\begin{aligned}
& F_{p}^{®_{j}}: z=(x, y) \in Y^{\square_{j}} \mapsto\left(x^{2}+y+p_{j},-\sqrt{\delta} \cdot x\right) \\
& \text { with } \quad Y^{\unrhd_{j}}=I^{\boxminus} \times J^{®_{j}}, \quad I^{\bowtie}=\left[-\delta^{2} / 4, \delta^{2} / 4\right] \quad \text { and } \quad J^{®_{j}}:=\sqrt{\delta} \cdot I_{j} .
\end{aligned}
$$

Note that $F_{p}^{\square_{j}}$ depends on a parameter $p \in \mathcal{P}$ but its definition domain $Y^{\square_{j}}$ is independent of $p$, see Fig. 4.
Fact 1.21. The map $F_{p}^{\square_{j}}$ is a folding transformation at $\diamond$ adapted to $F^{\mathscr{A}}$ with $\overleftarrow{\mathscr{A}}_{\square_{j}}:=\overleftarrow{\mathscr{A}} \cdot a_{j}$
Proof. First note that the segment $I^{\square}$ is disjoint from $\bigsqcup_{k} I_{k}=\bigsqcup_{k}\left[\frac{2}{N}(k-1)-1+\delta^{2}, \frac{2}{N} k-1-\delta^{2}\right]$. Thus $Y^{\unlhd_{j}}$ is disjoint from $\bigcup_{k=1}^{N} Y^{a_{k}}$. Also the image of $F_{p}^{\sqcup_{j}}$ is included in $(-1,1) \times\left[-\delta^{2} / 4, \delta^{2} / 4\right]$, which is in the interior of $Y^{e}$. We notice that for every $u \in \overleftarrow{\mathscr{A}}$, there exists $y^{u}$ such that $W^{u}=I \times\left\{y^{u}\right\}$. Note also that $y^{u}$ belongs to the interior of $\bigsqcup_{k} \sqrt{\delta} \cdot I_{k}$ and so $W^{u}$ is disjoint from $I \times \bigsqcup_{k} \sqrt{\delta} \cdot \partial I_{k} \supset \partial^{u} Y^{\complement_{j}}$. Also $W^{u}$ intersects $Y^{\unlhd_{j}}$ iff $u \in \overleftarrow{\mathscr{A}}_{\unrhd_{j}}:=\overleftarrow{\mathscr{A}} \cdot a_{j}$ which is indeed a clopen set. Also, for $u \in \overleftarrow{\mathscr{A}}_{\square_{j}}$, the image by $F_{p}^{\square_{j}}$ is:

$$
\begin{equation*}
F_{p}^{\unlhd_{j}}\left(W^{u} \cap Y^{\bowtie_{j}}\right)=\left\{\left(x^{2}+y^{u}+p_{j},-\sqrt{\delta} \cdot x\right):-\frac{\delta^{2}}{4} \leq x \leq \frac{\delta^{2}}{4}\right\} \tag{1.3}
\end{equation*}
$$

As every stable manifold $W^{s}$, for $s \in \overrightarrow{\mathscr{A}}$ is of the form $\left\{x_{s}\right\} \times I$, if $F_{p}^{\triangleleft_{j}}\left(W^{u} \cap Y^{\square_{j}}\right)$ is tangent to $W^{s}$ then the tangency is quadratic.

As the boxes $\left(Y^{\square_{j}}\right)_{1 \leq j \leq N-1}$ are disjoint, we can define:

$$
F_{p}:(x, \diamond) \in \bigsqcup_{j=1}^{N} Y^{a_{j}} \sqcup \bigsqcup_{j=1}^{N-1} Y^{®_{j}} \times\{\diamond\}= \begin{cases}\left(F^{a}(x), \diamond\right) & \text { if } x \in Y^{a}, a \in \mathscr{A} \\ \left(F_{p}^{\bullet}(x), \diamond\right) & \text { if } x \in Y^{\bullet}, \boxtimes \in \mathscr{C} .\end{cases}
$$

Fact 1.22. The map $F_{p}$ of is a system of type $(\mathscr{A}, \mathscr{C})$ which is a diffeomorphism for every $p \in \mathcal{P}$.
Proof. We already saw that $F^{\mathscr{A}}$ is a diffeomorphism. Looking at the definition of $\mathcal{P}$ and $F_{p}^{\mathscr{C}}$, we see that the folding transformations $F_{p}^{\bullet}$ have disjoint images from the one other. These images are included in $(-1,1) \times\left[-\delta^{2} / 4, \delta^{2} / 4\right]$ which is disjoint from the range of $F^{s l}$.

In order to do the parameter selection in Section 4, we will need the following object enabling to follow the tangency point of the homoclinic tangencies in a neat way:


Figure 4. An example of a system of type $(\mathscr{A}, \mathscr{C})$.
Definition 1.23. An adapted projection to a map $F$ of type $(\mathscr{A}, \mathscr{C})$ is a $C^{1}$-submersion:

$$
\pi:(z, v) \in \mathbb{R} \times I \times \mathscr{V} \mapsto\left(\pi^{v}(z), v\right) \in \mathbb{R} \times \mathscr{V}
$$

such that with $\mathbb{P}\left(\chi_{v}\right) \subset \mathbb{P}(\mathbb{R})$ the set of lines spanned by vectors in $\chi_{v}$ :
(0) For every $(z, d) \in \bigcup_{d \in \mathscr{Q} \cup \boldsymbol{A}} Y^{d} \times\{\sigma(d)\}$, the $\pi^{o(d)}$-fiber of $z$ is included in $Y^{e}$.
(1) The map $(z, v) \in Y^{e} \times \mathscr{V} \mapsto \operatorname{ker} D_{z} \pi^{v}$ is of class $C^{1}$ and takes its values in $\mathbb{P}\left(\chi_{v}\right)$.
(2) For any $v \in \mathscr{V}$, for every $x \in \mathbb{R}$, it holds $\pi^{v}(x, 0)=x$.
(3) For all $a \in \mathscr{A}$ and $z, z^{\prime} \in Y^{a}$, if $\pi^{\circ(a)}(z)=\pi^{\circ(a)}\left(z^{\prime}\right)$ then $\pi^{t(a)} \circ F^{a}(z)=\pi^{t(a)} \circ F^{a}\left(z^{\prime}\right)$.
(4) For every $u \in \overleftarrow{\mathscr{A}}_{\square}$ with $\square \in \mathscr{C}$, the curve $W^{u}$ is tangent to a unique fiber of $\pi^{t(\square)} \circ F^{\square}$ at a unique point $\zeta^{u} \in \operatorname{int} Y^{\square}$, and the tangency is quadratic.

Let us comment this definition. Item (0) says that for any $a \in \mathscr{C} \cup \mathscr{A}$, any $\pi^{a(a)}$-fiber intersecting $Y^{a}$ is fully included in $Y^{e}$. Item (1) implies that the fibers of $\pi^{v}$ are $C^{2}$-curves with tangent space in $\chi_{v}$ and with two endpoints, one in each of the components of $\mathbb{R} \times \partial I$. Item (2) implies that $\pi$ is a retraction. Item (3) implies, as we will prove in Proposition 1.25 that each curve $W^{3}$ is a fiber of $\pi^{o(\jmath)}$ for $s$. Thus $\pi$ defines an extension of the stable lamination which is invariant by the dynamics $F^{\mathscr{L}}$. The last item is equivalent to say that for every $u \in \overleftarrow{\mathscr{A}}_{\square}$, the map $x \mapsto \pi^{t(\square)} \circ F^{\unlhd} \circ\left(x, w^{u}(x)\right)$ has a unique critical point which is non-degenerate. It enables to describe how far an unstable curve $W^{u}, u \in \overleftarrow{\mathscr{A}}$ is to be folded tangent to a stable curve. The definition of $\zeta^{u}$ is depicted in Fig. 5.

Example 1.24. The projection $\pi:(x, y, \diamond) \mapsto(x, \diamond)$ is adapted to the map $F_{p}$ of Example 1.20.


Figure 5. Definition of $\zeta^{u}$.

Proof. The three first items of Definition 1.23 are obviously satisfied. Item (3) as well since $F^{\mathscr{A}}$ leaves invariant the vertical foliation. Finally we observe that by Eq. (1.3), for every $1 \leq j \leq N-1$, for every $u \in \overleftarrow{\mathscr{A}}_{\square_{j}}$ the curve $F_{p}^{\square_{j}}\left(W^{u} \cap Y^{®_{j}}\right)$ has a unique quadratic tangency with any vertical line. This implies the last item.

More generally, we will see in Proposition 2.25 that any horseshoe of a surface $C^{2}$-diffeomorphism can be modeled by a hyperbolic map of a certain type endowed with a map $\pi$ satifying items ( 0 -$1-2-3-4$ ) of Definition 1.23. As promised let us show:
Proposition 1.25. If $\pi$ is adapted to $F$ then for every $s \in \overrightarrow{\mathscr{A}}$, the curve $W^{s}$ is a fiber of $\pi^{0(s)}$.
Proof. Let $s=a_{0} \cdots a_{i} \cdots \in \vec{A}, v=a(s), z \in W^{s}$. As $W^{s}=\bigcap_{n \geq 0} Y^{a_{0} \cdots a_{n}}$, it comes that $\left(F^{\mathscr{A}}\right)^{n}(z, v)$ is in $Y^{a_{n}} \times\left\{a\left(a_{n}\right)\right\}$. Thus by Definition 1.23-(0), the $\pi$-fiber of $\left(F^{d}\right)^{n}(z, v)$ is included in $Y^{e} \times \mathscr{V}$ for every $n$. By Definition 1.23 (3), the $\pi$-fiber of $(z, v)$ is sent into the $\pi$-fiber of $\left(F^{d}\right)^{n}(z, v)$ by $\left(F^{d}\right)^{n}$. The tangent spaces of the fibers of $\pi$ are in $\chi_{v}$ by Definition 1.23-(1), and every vector in $\chi_{v}$ has its second coordinate which is more than $\lambda$-expanded by $D_{z}\left(F^{a}\right)^{-1}$ for any $a \in \mathscr{A}$ and $z \in F^{a}\left(Y^{a}\right)$ by Definition 1.4 (3). Thus the $\pi$-fiber of $(z, v)$ is a local stable manifold of $(z, v)$. By Proposition 1.16 and uniqueness of the local stable manifold of the hyperbolic set $\Lambda$, it comes that $W^{s} \times\{v\}$ is equal to the $\pi$-fiber of $(z, v)$.

We are going to give sufficient conditions for a system $F$ of type $(\mathscr{A}, \mathscr{C})$ to display a wandering stable component. We will use the following:

Definition 1.26. For every $c \in \mathscr{A}^{*}$, let $H^{c}:=F^{c}\left(I \times\{0\} \cap Y^{c}\right)$.
Definition 1.27. For any $c \in \mathscr{A}^{*}$, let $\left|Y^{c}\right|$ be the maximum of the lengths of $I \times\{y\} \cap Y^{c}$ among $y \in I$ and $\left|F^{c}\left(Y^{c}\right)\right|$ be the maximum of the lengths of $\{x\} \times I \cap F^{c}\left(Y^{c}\right)$ among $x \in I$.

The following enables to extend the tangency point $\zeta^{u}$ for some finite words $c \in \mathscr{A}^{*}$.

Proposition 1.28 (definition of $\mathscr{A}_{\square}^{*}$ and $\zeta^{c}$ ). For every $\boxtimes \in \mathscr{C}$, there exists $n_{\square}>0$ minimal such that for any c in:

$$
\mathscr{A}_{\square}^{*}:=\left\{c \in \mathscr{A}^{*}: c \text { coincides with the }|c| \geq n_{\square} \text { last letters of a certain } u \in \overleftarrow{\mathscr{A}}_{\square}\right\}
$$

the curve $H^{c}$ intersects $Y^{\square}$ but not $\partial^{u} Y^{\square}$, the curve $H^{c}$ is tangent to a unique fiber of $\pi^{t(\boxminus)} \circ F^{\square}$ at a unique point $\zeta^{c} \in \operatorname{int} Y^{\square}$, and the tangency is quadratic.

Proof. It is sufficient to prove that for every sequence $\left(c_{i}\right)_{\geq 0}$ of words $c_{i} \in \mathscr{A}^{*}$ such that $\left|c_{i}\right| \rightarrow \infty$ and there exists $u_{i} \in \overleftarrow{\mathscr{A}}_{\square}$ which finishes by $c_{i}$, then for $i$ large enough the curve $H^{c_{i}}$ intersects $Y^{\square}$ but not $\partial^{u} Y^{\square}$, the curve $H^{c_{i}}$ is tangent to a unique fiber of $\pi^{t(\square)} \circ F^{\square}$ at a unique point in int $Y^{\square}$, and the tangency is quadratic. By compactness of $\overleftarrow{\mathscr{A}}_{\square}$, we can assume that $\left(u_{i}\right)_{i}$ converges to $u \in \overleftarrow{\mathscr{A}}_{\square}$. By continuity of $u \mapsto W^{u}$, the curve $W^{u_{i}}$ is $C^{2}$-close to $W^{u}$ when $i$ is large

We notice that each of the $H^{c_{i}}$ is a component of the image by $\left(F^{\mathscr{A}}\right)^{\left|c_{i}\right|}$ of the curve $I \times\{0\} \times \mathscr{V}$ which is uniformly traverse to the stable lamination of $\Lambda$. When $i$ is large, by the Inclination Lemma, the curve $H^{c_{i}}$ is $C^{2}$-close to $W^{u_{i}}$ and so to $W^{u}$.

By definition of $\overleftarrow{\mathscr{A}}_{\square}$, the curve $W^{u}$ intersects $Y^{\square}$ but not $\partial^{u} Y^{\unrhd}$. This equivalent to say that $W^{u}$ intersects the interior of $Y^{\square}$ but not $\partial^{u} Y^{\square}$ which is an open condition on curves. Thus for $i$ sufficiently large $H^{c_{i}}$ intersects the interior of $Y^{\square}$ but not $\partial^{u} Y^{\square}$.

Also $W^{u}$ has a tangency with the foliation defined by $\pi^{t(\square)} \circ F^{\square}$, which is unique, in int $Y^{\square}$ and quadratic. This is equivalent to say that the curve $T W^{u} \subset Y^{e} \times \mathbb{P}(\mathbb{R})$ is transverse to the surface $\left\{\left(z, \operatorname{ker} D \pi^{t(\boxminus)} \circ D_{z} F^{\boxminus}\right): z \in Y^{\unlhd_{i}}\right\} \subset Y^{e} \times \mathbb{P}(\mathbb{R})$ and intersects it at a unique point in int $Y^{\square} \times \mathbb{P}(\mathbb{R})$. This is an open condition on $C^{2}$-curves. Thus $H^{c_{i}}$ satisfies the same property when $i$ is large.
Remark 1.29. We have showed that the set of $C^{2}$-curves $\left\{H^{c}: c \in \mathscr{A}_{\square}^{*}\right\} \cup\left\{W^{u}: u \in \overleftarrow{\mathscr{A}}_{\square}\right\}$ is compact. Also the subset $\left\{\zeta^{c}: c \in \mathscr{A}_{\square}^{*} \cup \overleftarrow{\mathscr{A}}_{\square}\right\}$ of $\mathbb{R}^{2}$ is compact.

When $F$ is real analytic, we define $\rho>0$ and the sets $\left(\tilde{Y}^{a}\right)_{a \in \mathscr{A}}$ and $\tilde{Y}^{\square}$ by the following:
Proposition 1.30 (and definition of $\rho>0$ ). If $F$ is real analytic, then there exists $\rho>0$ such that for every $c \in \mathscr{A}^{*}$, the hyperbolic transformation $\left(Y^{c}, F^{c}\right)$ extends to a $C_{\rho}^{\omega}$-hyperbolic transformation $\left(\tilde{Y}^{c}, F^{c}\right)$ and the map $F^{\square}$ is of class $C_{\rho}^{\omega}$ on $\tilde{Y}^{\boxminus}$ for every $\boxtimes \in \mathscr{C}$.

Proof. The first condition on $F^{s l}$ is proved in Proposition A. 3 and the second condition on $F^{\square}$ is obvious by Definition 1.18.

Definition 1.31. If $F$ is analytic, for any $c \in \mathscr{A}^{*}$, let $\left|F^{c}\left(\tilde{Y}^{c}\right)\right|$ be the maximum of the diameters of $\{z\} \times \tilde{I} \cap F^{c}\left(\tilde{Y}^{c}\right)$ among $z \in \tilde{I}$.

Here is the proposition which gives sufficient conditions to ensure the existence of a wandering stable component. It is a variation of Theorem D of [Ber18] and its Remark 3.5.

Theorem 1.32. Let $F$ be a system of type $(\mathscr{A}, \mathscr{C})$ endowed with an adapted projection $\pi$. Let $\left(\square_{j}\right)_{j \geq 1} \in \mathscr{C}^{\mathbb{N}}$ and let $\left(c_{j}\right)_{j \geq 1}$ be a sequence of words $c_{j} \in \mathscr{A}_{\square_{j}}^{*}$ such that $a\left(c_{j}\right)=t\left(\square_{j-1}\right)$ and satisfying the following assumptions when $j$ is large:
(i) the limit exists and is non-zero $\breve{\gamma}_{j}:=\left|Y^{c_{j+1}}\right|^{1 / 2} \cdot\left|Y^{c_{j+2}}\right|^{1 / 2^{2}} \cdots\left|Y^{c_{j+k}}\right|^{1 / 2^{k}} \cdots$,
(ii) the distance between $\pi^{t\left(c_{j+1}\right)}$-fibers of $F^{c_{j+1}} \circ F^{\unlhd_{j}}\left(\zeta^{c_{j}}\right)$ and $\zeta^{c_{j+1}}$ is small compared to $\breve{\gamma}_{j+1}$, (iii) the width $\left|Y^{c_{j}}\right|$ is small and the height $\left|F^{c_{j}}\left(Y^{c_{j}}\right)\right|$ small compared to $\breve{\gamma}_{j}^{2}$.

Then there exist $J \geq 0$ and a sequence of nonempty open subsets $\left(B_{j}\right)_{j \geq J}$ of $Y^{e}$ such that:
$B_{j} \subset Y^{\unlhd_{j}}, \quad F^{\unlhd_{j}}\left(B_{j}\right) \subset Y^{c_{j+1}}, \quad F^{\iota_{j+1}} \circ F^{\unlhd_{j}}\left(B_{j}\right) \subset B_{j+1} \quad$ and $\quad \lim _{k \rightarrow \infty} \operatorname{diam} F^{k}\left(B_{j} \times\left\{a\left(\square_{j}\right)\right\}\right)=0$.
Remark 1.33. We can show that (i) is equivalent to ask that $\log \left|Y^{c_{j}}\right| \asymp\left|c_{j}\right|$ satisfies that $\sum_{j} 2^{-j}\left|c_{j}\right|<\infty$. We will see, using distortion bound, that assumption (ii) is equivalent to say that the distance between the $\pi^{0\left(c_{j+1}\right)}$-fibers of the points $F^{®_{j}}\left(\zeta^{c_{j}}\right)$ and $\left(F^{c_{j+1}}\right)^{-1}\left(\zeta^{c_{j+1}}\right)$ is small compared to $\breve{\gamma}_{j}^{2}$. The first part of assumption (iii) implies that $\left|c_{j}\right|$ is large when $j$ is large.

Theorem 1.34. Under the assumptions of Theorem 1.32, if moreover $F$ is analytic and if $(\widetilde{i i i}) \breve{\gamma}_{j}$ is small compared to $\left|Y^{c_{j+1}}\right|$ and $\left|F^{c_{j}}\left(\tilde{Y}^{c_{j}}\right)\right|$ is small compared to $\breve{\gamma}_{j}^{2}$ when $j$ is large, then each $B_{j}$ is the real trace of an open subset $\tilde{B}_{j} \subset \mathbb{C}^{2}$ such that:
$\tilde{B}_{j} \subset \tilde{Y}^{\unlhd_{j}}, \quad F^{\complement_{j}}\left(\tilde{B}_{j}\right) \subset \tilde{Y}^{c_{j+1}}, \quad F^{c_{j+1}} \circ F^{\complement_{j}}\left(\tilde{B}_{j}\right) \subset \tilde{B}_{j+1} \quad$ and $\quad \lim _{k \rightarrow \infty} \operatorname{diam} F^{k}\left(\tilde{B}_{j} \times\left\{a\left(\square_{j}\right)\right\}\right)=0$, where $\tilde{Y}^{\unlhd_{j}}$ and $\tilde{Y}^{c_{j+1}}$ are defined by Proposition 1.30.

We notice that $B_{j} \times\left\{a\left(\square_{j}\right)\right\}$ and $\tilde{B}_{j} \times\left\{a\left(\square_{j}\right)\right\}$ are stable domains for the dynamics $F$ which are in wandering stable components (for their combinatorics are wandering).

Remark 1.35. In the analytic case, assumption ( $(\widetilde{i i i})$ is actually not necessary. It could be replaced by the complex counterpart of the computational distortion bound of [PY09, (MP6), p. 21] (see Theorem 3.9), whose proof lasts a couple of pages and is rather computational.

The proof of Theorems 1.32 and 1.34 will occupy all Section 3.3 and will need a development of the techniques of implicit represention of Section 3.1 and a normal form defined in Section 3.2.
1.2. Unfolding of wild type. In this section we are going to state natural conditions on families of systems of type $(\mathscr{A}, \mathscr{C})$ so that at a dense set of parameters the map will satisfy the assumptions of Theorems 1.32 and 1.34.

We recall that a closed subset is regular if it is nonempty and equal to the closure of its interior. Let $\mathcal{P}$ be a non-empty, regular compact subset of $\mathbb{R}^{d}$. For the proofs of our main theorems, $\mathcal{P}$ will be a hypercube of $\mathbb{R}^{5}$. Let $n, m \geq 1$ and let $\left(K_{p}\right)_{p}$ be a family of compact sets of $\mathbb{R}^{n}$. For $r \geq 0$, a family $\left(f_{p}\right)_{p \in \mathcal{P}}$ of maps $f_{p}: K_{p} \rightarrow \mathbb{R}^{m}$ is of class $C^{r}$ if $p \in \mathcal{P} \mapsto K_{p}$ is continuous for the Hausdorff topology and the map $(p, z) \in \hat{K}:=\bigcup_{p \in \mathcal{P}}\{p\} \times K_{p} \mapsto f_{p}(z) \in \mathbb{R}^{m}$ is of class $C^{r}$ in the sense we defined Page 13.

We are now going to consider an unfolding of the stable and unstable laminations of a $C^{2}$-family of maps $\left(F_{p}\right)_{p \in \mathcal{P}}$ of type $(\mathscr{A}, \mathscr{C})$.

Definition 1.36. A family $\left(F_{p}\right)_{p \in \mathcal{P}}$ of systems of type $(\mathscr{A}, \mathscr{C})$ is regular if each $F_{p}$ is a system of type $(\mathscr{A}, \mathscr{C})$ with a same set $\overleftarrow{\mathscr{A}}_{\square}$ for every $p \in \mathcal{P}$ at each fixed $\square \in \mathscr{C}$ and so that the family $\left(F_{p}\right)_{p \in \mathcal{P}}$ is of class $C^{2}$. Note that for every $p \in \mathcal{P}$, the map $F_{p}$ is of the form:
$F_{p}:(x, v) \in \bigsqcup_{a \in \mathscr{A}} Y_{p}^{a} \times\{a(a)\} \sqcup \bigsqcup_{\square \in \mathscr{C}} Y_{p}^{\boxminus} \times\{a(\square)\} \mapsto \begin{cases}F_{p}^{\mathscr{A}}(x, v) & \text { if }(x, v) \in \bigsqcup_{a \in \mathscr{A}} Y_{p}^{a} \times\{a(a)\}, \\ F_{p}^{\mathscr{C}}(x, v) & \text { if }(x, v) \in \bigsqcup_{\square \in \mathscr{C}} Y_{p}^{\square} \times\{a(\square)\} .\end{cases}$

A family of projections $\left(\pi_{p}\right)_{p}$ is adapted to $\left(F_{p}\right)_{p}$ if each $\pi_{p}$ is a projection adapted to $F_{p}$, the family $\left(\pi_{p}\right)_{p}$ is of class $C^{1}$ and the family $\left(\operatorname{ker} D_{z} \pi_{p}\right)_{p \in \mathcal{P}}$ is of class $C^{1}$.

Example 1.37. For instance the family $\left(F_{p}\right)_{p}$ of systems of type $(\mathscr{A}, \mathscr{C})$ of Example 1.20 is regular and the trivial family of projections $(\pi)_{p \in \mathcal{P}}$ of Example 1.24 is adapted to it.

Let us fix a regular family $\left(F_{p}\right)_{p}$ of type $(\mathscr{A}, \mathscr{C})$ endowed with an adapted family $\left(\pi_{p}\right)_{p}$ of projections. As in Definition 1.36, for every $p \in \mathcal{P}$, we add an index $p$ to all the geometric objects defined in the last subsection. For instance, we denote also by $\Lambda_{p}$ the maximal invariant set of $F_{p}^{\mathscr{A}}$ defined in Proposition 1.13. For every $c=a_{1} \cdots a_{k} \in \mathscr{A}^{*}$, we denote $\left(Y_{p}^{c}, F_{p}^{c}\right):=$ $\left(Y_{p}^{a_{1}}, F_{p}^{a_{1}}\right) \star \cdots\left(Y_{p}^{a_{k}}, F_{p}^{a_{k}}\right)$. For $s \in \overrightarrow{\mathscr{A}}$ and $u \in \overleftarrow{\mathscr{A}}$, we denote by $W_{p}^{s}$ and $W_{p}^{u}$ the corresponding local stable and unstable manifolds of $F_{p}^{d}$.

The following is immediate:
Proposition 1.38. For every $c \in \mathscr{A}^{*}$, the family $\left(\left(Y_{p}^{c}, F_{p}^{c}\right)\right)_{p \in \mathcal{P}}$ is of class $C^{2}$.
By Proposition 1.16, for every $u \in \overleftarrow{\mathscr{A}}$ and $s \in \overrightarrow{\mathscr{A}}$, for every $p \in \mathcal{P}$, there are $\theta$-Lipschitz functions $w_{p}^{u}$ and $w_{p}^{s}$ on $I$ whose graphs satisfy:

$$
\begin{equation*}
\operatorname{Graph} w_{p}^{u}=W_{p}^{u} \quad \text { and } \quad{ }^{t} \operatorname{Graph} w_{p}^{s}=W_{p}^{s} \tag{1.4}
\end{equation*}
$$

The following states that the stable and unstable laminations depend $C^{2}$ on the parameter:
Proposition 1.39. The maps $(p, x) \in \mathcal{P} \times I \mapsto w_{p}^{u}(x) \in \mathbb{R}$ and $(p, y) \in \mathcal{P} \times I \mapsto w_{p}^{s}(y) \in \mathbb{R}$ are of class $C^{2}$. Moreover, the two following maps are continuous:

$$
u \in \overleftarrow{\mathscr{A}} \mapsto\left(w_{p}^{u}\right)_{p \in \mathcal{P}} \in C^{2}(\mathcal{P} \times I, I) \text { and } s \in \overrightarrow{\mathscr{A}} \mapsto\left(w_{p}^{s}\right)_{p \in \mathcal{P}} \in C^{2}(\mathcal{P} \times I, I)
$$

Proof. This is an immediate consequence of Theorem C. 5 of the Appendix applied to the family of hyperbolic sets $\left(\Lambda_{p}\right)_{p}$.

Definition 1.40. For every $\square \in \mathscr{C}$, we define:

- The map $p \in \mathcal{P} \mapsto a_{p}^{u}:=\pi_{p}^{t(\square)} \circ F_{p}^{\bullet}\left(\zeta_{p}^{u}\right)$, for every $u \in \overleftarrow{\mathscr{A}}_{\square}$.
- The map $p \in \mathcal{P} \mapsto b_{p}^{s}:=w_{p}^{s}(0)$, for every $s \in \overrightarrow{\mathcal{A}}$.
- The map $p \in \mathcal{P} \mapsto \mathcal{V}_{p}^{u, s}=a_{p}^{u}-b_{p}^{s}$, for all $u \in \overleftarrow{\mathscr{A}}_{\square}$ and $s \in \overrightarrow{\mathscr{A}}_{\square}$.

As $\pi_{p}^{v}(x, 0)=x$ for every $(x, v) \in I \times \mathscr{V}$, it holds $\pi_{p}^{\circ(s)}\left(W_{p}^{s}\right)=\left\{b_{p}^{s}\right\}$. We recall that by Definition 1.18, $\zeta_{p}^{u}$ is the unique critical point of $\pi_{p}^{t(\square)} \circ F_{p}^{\bullet}$. Thus $a_{p}^{u}$ is the critical value of $\pi_{p}^{t(\square)} \circ F_{p}^{\unlhd} \circ w_{p}^{u}$ and $\mathcal{V}_{p}^{u, s}$ quantifies how far $F_{p}^{\square}\left(W_{p}^{u}\right)$ is to be tangent to $W_{p}^{s}=\left(\pi^{t(\square)}\right)^{-1}\left(\left\{\hat{e}_{p}^{s}\right\}\right)$. In particular we have:

Remark 1.41. The curve $F_{p}^{\square}\left(W_{p}^{u} \cap Y_{p}^{\bowtie}\right)$ is tangent to $W_{p}^{s}$ if and only if $\mathcal{V}_{p}^{u, s}=0$.

Proposition 1.42. For every $\square \in \mathscr{C}$, the following are $C^{1}$-functions depending continuously on $u \in \overleftarrow{\mathscr{A}}_{\square}, s \in \overrightarrow{\mathscr{A}}_{\square}$ :

$$
p \mapsto \zeta_{p}^{u} \in \operatorname{int} Y_{p}^{\boxminus}, \quad p \in \mathcal{P} \mapsto a_{p}^{u} \in I, \quad p \in \mathcal{P} \mapsto b_{p}^{s} \in I \quad \text { and } \quad p \mapsto \mathcal{V}_{p}^{u, s} \in \mathbb{R}
$$



Figure 6. Definition of $\mathcal{V}^{u, s}$.

Proof. We recall that $(p, x) \mapsto w_{p}^{u}(x)$ and $(p, y) \mapsto w_{p}^{s}(y)$ are $C^{2}$-functions depending continuously on $u \in \overleftarrow{\mathscr{A}}$ and $s \in \overrightarrow{\mathscr{A}}$ by Proposition 1.39. Thus the map $p \mapsto b_{p}^{s}=w_{p}^{s}(0)$ is of class $C^{1}$ and depends continuously on s. By Definition 1.23 (4), the curve $W_{p}^{u} \cap Y_{p}^{\square}$ displays a unique tangency with a fiber of $\pi_{p}^{t(\square)} \circ F_{p}^{\square}$ and this tangency is quadratic. This means that the curve $T W_{p}^{u} \cap Y_{p}^{\square} \times \mathbb{P}(\mathbb{R})$ of $Y^{e} \times \mathbb{P}(\mathbb{R})$ intersects the surface $\left\{\left(z, \operatorname{ker} D \pi^{t(\square)} \circ D_{z} F^{\square}\right): z \in Y^{\square_{i}}\right\}$ of $Y^{e} \times \mathbb{P}(\mathbb{R})$ at a unique point $\left(\zeta_{p}^{u}, T_{\zeta_{p}^{u}} W^{u}\right)$ and this intersection is transverse. As this curve and this surface vary $C^{1}$ with $p$ and this continuously with $u$, by transversality we obtain the sought regularity of $p \in \mathcal{P} \mapsto \zeta_{p}^{u}$. Then the sought regularities of $p \in \mathcal{P} \mapsto a_{p}^{u}=\pi_{p}^{t(\boxminus)} \circ F_{p}^{\boxminus}\left(\zeta_{p}^{u}\right)$ and $p \mapsto \mathcal{V}_{p}^{u, s}=a_{p}^{u}-b_{p}^{s}$ follow.

The next definition regards a regular family $\left(F_{p}\right)_{p \in \mathcal{P}}$ of systems of type $(\mathscr{A}, \mathscr{C})$ endowed with an adapted family $\left(\pi_{p}\right)_{p}$ of projections.

Definition 1.43. The family $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ is an unfolding of wild type $(\mathscr{A}, \mathscr{C})$ if:
$\left(\mathbf{H}_{1}\right)$ For every $p \in \mathcal{P}$, for every $\square \in \mathscr{C}$, there exist $u \in \overleftarrow{\mathscr{A}}_{\square}$ and $s \in \overrightarrow{\mathscr{A}}_{\square}$ satisfying that $F_{p}^{\boxminus}\left(W_{p}^{u} \cap Y_{p}^{\square}\right)$ has a quadratic tangency with $W_{p}^{\jmath}$.
$\left(\mathbf{H}_{2}\right)$ For all $\left(u_{\square}, s_{\square}\right)_{\square} \in \prod_{\square \in \mathscr{C}} \overleftarrow{\mathscr{A}}_{\square} \times \overrightarrow{\mathscr{A}}_{\square}$, the following is a $C^{1}$-diffeomorphism onto its image:

$$
p \in \operatorname{int} \mathcal{P} \mapsto\left(\mathcal{V}_{p}^{u_{\square}, \mathfrak{s}_{\square}}\right)_{\square \in \mathscr{C}} \in \mathbb{R}^{\operatorname{Card} \mathscr{C}}
$$

Let us comment this definition. Property $\left(\mathbf{H}_{1}\right)$ asks for Card $\mathscr{C}$ simultaneous homoclinic tangencies, each of which being associated to a different $\square \in \mathscr{C}$. By Proposition 1.42, we already now that the map involved in $\left(\mathbf{H}_{2}\right)$ is of class $C^{1}$. What is required is its invertibility. It means roughly speaking that the unfoldings of the tangencies are independent and non-degenerated along the parameter space $\operatorname{int} \mathcal{P}$ whose dimension is equal to Card $\mathscr{C}$.

We introduced the notion of unfolding of wild type $(\mathscr{A}, \mathscr{C})$ to show the existence of a dense set of parameter $p$ for which the the system $F_{p}$ has a wandering stable domain. To this end, a form of dissipation is necessary:

Definition 1.44. A hyperbolic map $F^{\mathscr{A}}$ of type $\mathscr{A}$ is moderately dissipative if it holds:

$$
\left\|D F^{\mathscr{A}}\right\| \cdot\left\|\operatorname{det} D F^{\mathscr{A}}\right\|^{\epsilon}<1 \quad \text { with } \epsilon=\frac{1}{10}
$$

An unfolding $\left(F_{p}, \pi_{p}\right)_{p}$ of wild type is moderately dissipative if $F_{p}^{\&}$ is moderately dissipative for every $p \in \mathcal{P}$.

Here is the main (abstract) result of this manuscript:
Theorem 1.45 (Main). Let $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ be a moderately dissipative unfolding of wild type ( $\left.\mathscr{A}, \mathscr{C}\right)$ with Card $\mathscr{C}=5$. Then there exists a dense set of parameters $p \in \mathcal{P}$ for which $\left(F_{p}, \pi_{p}\right)$ satisfies the conclusions of Theorem 1.32 (and of Theorem 1.34 if $F_{p}$ is analytic) with subsets $\left(B_{j} \times\left\{\sigma\left(\square_{j}\right)\right\}\right)_{j \geq J}$ such that every point $(z, v) \in B_{J} \times\left\{\sigma\left(\square_{J}\right)\right\}$ satisfies:
$(\downarrow)$ There exist $0<t<1$ and $\mu$ in the set of invariant probability measures $\mathcal{M}_{p}\left(\Lambda_{p}\right)$ of $F_{p} \mid \Lambda_{p}$, such that the limit set of the sequence of empirical measures $\mathrm{e}_{n}:=\frac{1}{n} \sum_{i=0}^{n-1} \delta_{F_{p}^{i}(z, v)}$ contains $t \cdot \mu+(1-t) \cdot \mathcal{M}_{p}\left(\Lambda_{p}\right)$.

The proof of this theorem will occupy the whole Section 4. In this proof we will first show that the assumptions of Theorem 1.32 are satisfied at a dense set of parameter. Finally in Section 4.4, we will show conclusion $(\boldsymbol{\wedge})$ by selecting the combinatorics of sequences of words $\left(c_{j}\right)_{j}$.

Theorem 1.45 being admitted for now, let us show how it implies the main theorems. To this end, we define the following geometric model which applies to surface maps families.

Definition 1.46 ( $C^{r}$-Geometric model). A $C^{r}$-family $\left(f_{p}\right)_{p \in \mathcal{P}}$ of surface diffeomorphisms $f_{p} \in$ Diff ${ }^{2}(M)$ displays the geometric model if there is a moderately dissipative unfolding $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of wild type $(\mathscr{A}, \mathscr{C})$ with $\operatorname{Card} \mathscr{C}=5$ so that $\left(F_{p}\right)_{p}$ is embedded into $M$ via a $C^{2}$-family $\left(H_{p}\right)_{p}$ of $C^{r}$-embeddings $H_{p}: Y^{e} \times \mathscr{V} \hookrightarrow M$ :

$$
f_{p} \circ H_{p} \mid \bigsqcup_{d \in \mathscr{A} \sqcup \mathscr{C}} Y_{p}^{d} \times\{a(d)\}=H_{p} \circ F_{p}, \quad \forall p \in \mathcal{P}
$$

In Section 2.1 we will prove that the families of maps defined in Example 1.37 define a moderately dissipative unfolding of wild type $(\mathscr{A}, \mathscr{C})$ for any $\delta>0$ small enough depending on an even $N \geq 2$, and so the conclusions of Theorem 1.45 when moreover $N=6$.

In Section 2.2, we will state a general Theorem 2.10 implying that the geometric model is displayed in any non-degenerate unfolding of five homoclinic tangencies of a same dissipative saddle periodic point. This will imply Theorem C. Then we will show that this implies the two following theorems. The first one states that the geometric model appears locally densely among generalized Hénon maps:

Theorem 1.47. For every $b \neq 0$ small, there is a regular compact set $\mathcal{P} \subset \mathbb{R}^{5}$ such that with:

$$
f_{p}:(x, y) \in \mathbb{R}^{2} \mapsto\left(x^{6}+\sum_{i=0}^{4} p_{i} \cdot x^{i}-y, b \cdot x\right) \text { for every } p=\left(p_{i}\right)_{0 \leq i \leq 4} \in \mathcal{P}
$$

the family $\left(f_{p}\right)_{p \in \mathcal{P}}$ displays the $C^{\omega}$-geometric model 1.46 .

Remark 1.48. The same holds true for the family of maps:

$$
f_{p}:[z, w, t] \in \mathbb{P}^{2}(\mathbb{C}) \mapsto\left[z^{6}+\sum_{i=0}^{4} p_{i} \cdot z^{i} \cdot t^{6-i}-w^{6}, b \cdot z^{6}, t^{6}\right] .
$$

The second states that the geometric model appears openly and densely in the space of unfolding of a dissipative homoclinic tangency. We recall that a saddle periodic point displays a homoclinic tangency if its unstable manifold is tangent to its stable manifold. The homoclinic tangency is dissipative if the determinant of the differential at the periodic cycle has modulus less than 1.

Theorem 1.49. For every $r \in[2, \infty] \cup\{\omega\}$, if $f \in \operatorname{Diff}^{r}(M)$ displays a dissipative, quadratic homoclinic tangency, then there exist a $C^{r}$-perturbation $\tilde{f}$ of $f, a C^{r}$-open-dense set of families $\left(f_{p}\right)_{p \in \mathbb{R}^{5}}$ satisfying that $f_{0}=\tilde{f}$ such that for $n \geq 1$ and a regular compact subset $\mathcal{P}$ close to $0 \in \mathbb{R}^{5}$, the family $\left(f_{p}^{n}\right)_{p \in \mathcal{P}}$ satisfies the $C^{r}$-geometric model 1.46.

These theorems together with Corollary 1.50 imply the first proof of the existence of a wandering Fatou component of a polynomial automorphism of $\mathbb{C}^{2}$ and a proof of the last Takens' problem in any regularity:

Proof of Theorems A and B. These theorem are direct consequences of respectively Theorems 1.47 and 1.49 and the following consequence of main Theorem 1.45.

Corollary 1.50 (Main). For every $r \in[2, \infty] \cup\{\omega\}$, if a family $\left(f_{p}\right)_{p \in \mathcal{P}}$ satisfies the $C^{r}$-geometric model 1.46 with a family of embeddings $\left(H_{p}\right)_{p \in \mathcal{P}}$, then there exists a dense subset of $\mathcal{P}$ of parameters $p$ at which the map $f_{p}$ has the following properties:
(1) There exists a wandering stable component $\mathcal{C}$. If $r=\omega$, then the holomorphic extension of $f_{p}$ has a wandering Fatou component containing $\mathcal{C}$.
(2) For every $x \in \mathcal{C}$, the limit set of the orbit of $x$ contains the horseshoe $K_{p}:=H_{p}\left(\Lambda_{p}\right)$.
(3) For every $x \in \mathcal{C}$, the sequence $\left(\mathrm{e}_{n}^{f}\right)_{n}$ of empirical measures $\mathrm{e}_{n}^{f}:=\frac{1}{n} \sum_{i=0}^{n-1} \delta_{f_{p}^{i}(x)}$ diverges. Furthermore, there exist $0<t<1$ and $\mu^{\prime}$ in the set of invariant probability measures $\mathcal{M}_{p}\left(K_{p}\right)$ of $f_{p} \mid K_{p}$ such that the limit set of the sequence $\left(e_{n}^{f}\right)_{n}$ contains $t \cdot \mu^{\prime}+(1-t) \cdot \mathcal{M}_{p}\left(K_{p}\right)$.
$\operatorname{Proof}$. Let $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ be the moderately dissipative unfolding of a certain wild type $(\mathscr{A}, \mathscr{C})$ with $\operatorname{Card} \mathscr{C}=5$ which is embedded into $M$ via $\left(H_{p}\right)_{p}$ and such that $f_{p} \circ H_{p}=H_{p} \circ F_{p}$. Let $p \in \mathcal{P}$ be in the dense set of parameters given by Theorem 1.45. Let $\left(B_{j} \times\left\{a\left(\square_{j}\right)\right\}\right)_{j \geq J}$ and $\left(c_{j}\right) \in\left(\mathscr{A}^{*}\right)^{\mathbb{N}}$ be the associated sequences of stable domains and words for the fixed parameter $p$ satisfying both the conclusions of Theorem 1.32 (resp. Theorem 1.34 if $F_{p}$ is analytic) and condition ( $\boldsymbol{*}$ ). Put $D_{j}:=H_{p}\left(B_{j} \times\left\{a\left(\square_{j}\right)\right\}\right)$ which is the real trace of $\tilde{D}_{j}:=H_{p}\left(\tilde{B}_{j} \times\left\{\sigma\left(\square_{j}\right)\right\}\right)$ in the analytic case. By the conclusions of Theorem 1.32, it holds for every $j \geq J$ :

$$
D_{j} \subset H_{p}\left(Y_{p}^{\square_{j}} \times a\left(\square_{j}\right)\right), \quad f_{p}\left(D_{j}\right) \subset H_{p}\left(Y^{c_{j+1}} \times\left\{a\left(c_{j+1}\right)\right\}\right) \quad \text { and } \quad f_{p}^{\left|c_{j+1}\right|+1}\left(D_{j}\right) \subset D_{j+1}
$$

and respectively in the analytic case by the conclusions of Theorem 1.34, it holds for every $j \geq J$ :

$$
\tilde{D}_{j} \subset H_{p}\left(\tilde{Y}_{p}^{\square_{j}} \times o\left(\oplus_{j}\right)\right), \quad f_{p}\left(\tilde{D}_{j}\right) \subset H_{p}\left(\tilde{Y}^{c_{j+1}} \times\left\{\theta\left(c_{j+1}\right)\right\}\right) \quad \text { and } \quad f_{p}^{\left|c_{j+1}\right|+1}\left(\tilde{D}_{j}\right) \subset \tilde{D}_{j+1} .
$$

The last conclusion of Theorem 1.32 (resp. 1.34) implies that $D_{j}$ (resp. $\tilde{D}_{j}$ ) is a stable domain:

$$
\lim _{k \rightarrow \infty} \operatorname{diam} f_{p}^{k}\left(D_{j}\right)=0 \quad \text { and } \quad \text { resp. } \quad \lim _{k \rightarrow \infty} \operatorname{diam} f_{p}^{k}\left(\tilde{D}_{j}\right)=0
$$

Thus $D_{j}$ (resp. $\tilde{D}_{j}$ ) is included in a stable component $\mathcal{C}$ (resp. $\tilde{\mathcal{C}}$ ). We now prove (1)-(2)-(3) for this fixed parameter $p$.

Proof of (1). Assume for the sake of contradiction that $\mathcal{C}$ is preperiodic. By taking $j$ larger, we can assume it $q$ periodic for some $q \geq 1: f_{p}^{q}(\mathcal{C})=\mathcal{C}$. Let $D_{j}^{\prime} \Subset D_{j}$ be with non-empty interior. We have $D_{j}^{\prime} \Subset \mathcal{C}$ and so

$$
A:=D_{j}^{\prime} \cup f_{p}^{q}\left(D_{j}^{\prime}\right) \Subset \mathcal{C}
$$

Since $\mathcal{C}$ is a stable component, it holds $\lim _{n \rightarrow \infty} \operatorname{diam} f_{p}^{n}(A)=0$.
We recall that $D_{p}(\mathscr{A}):=\bigsqcup_{a \in \mathscr{A}} Y_{p}^{a} \times\{a(\mathfrak{a})\}$ and $D_{p}(\mathscr{C}):=\bigsqcup_{\square \in \mathscr{C}} Y_{p}^{\boxminus} \times\{a(\square)\}$. We observe that for $k \geq j+1$, with $N_{k}:=1+\left|c_{j+1}\right|+\cdots+1+\left|c_{k}\right|$, it holds that:

$$
\left\{\begin{array}{cl}
f_{p}^{l}\left(D_{j}\right) \subset D_{k} \subset H_{p}\left(D_{p}(\mathscr{C})\right) & \text { if } l=N_{k} \\
f_{p}^{l}\left(D_{j}\right) \subset H_{p}\left(D_{p}(\mathscr{A})\right) & \text { if } l \in\left\{N_{k}+1, \ldots, N_{k+1}-1\right\}
\end{array}\right.
$$

The sets $H_{p}\left(D_{p}(\mathscr{A})\right)$ and $H_{p}\left(D_{p}(\mathscr{C})\right)$ are disjoint and compact. As $N_{k} \rightarrow \infty$ and $N_{k}-N_{k-1} \rightarrow \infty$, for $k$ large enough, it holds:
(i) the diameter of $f_{p}^{N_{k}-q}(A)$ is smaller than the distance between $H_{p}\left(D_{p}(\mathscr{A})\right)$ and $H_{p}\left(D_{p}(\mathscr{C})\right)$.
(ii) the set $f_{p}^{N_{k}-q}(A)$ contains $f_{p}^{N_{k}-q}\left(f_{p}^{q}\left(D_{j}^{\prime}\right)\right)=f_{p}^{N_{k}}\left(D_{j}^{\prime}\right)$ and so intersects $H_{p}\left(D_{p}(\mathscr{C})\right)$.
(iii) the set $f_{p}^{N_{k}-q}(A)$ contains $f_{p}^{N_{k}-q}\left(D_{j}^{\prime}\right)$ and so intersects $H_{p}\left(D_{p}(\mathcal{A})\right)$ since $N_{k}-q \in\left\{N_{k-1}+\right.$ $\left.1, \ldots, N_{k}-1\right\}$ for $k$ sufficiently large.
Conclusions (i) and (ii) contradict (iii). Thus $\mathcal{C}$ is a wandering stable component.
Literally the same argument (by replacing $\mathcal{C}$ by $\tilde{\mathcal{C}}$ ) proves that the open set $\tilde{\mathcal{C}} \subset \tilde{M}$ is wandering. We observe that $\tilde{\mathcal{C}}$ is included in a Fatou component $\mathcal{F}$. Let $\mathcal{F}^{\prime} \Subset \mathcal{F}$ be a precompact open subset of $\mathcal{F}$ which intersects the interior of $\tilde{B}_{j}$. Then any cluster value for $\left(f_{p}^{n} \mid \mathcal{F}^{\prime}\right)_{n}$ is a holomorphic map constant on $\tilde{B}_{j} \cap \mathcal{F}^{\prime}$ and so is constant. Thus any cluster value of ( $\left.\operatorname{diam} f_{p}^{n}\left(\mathcal{F}^{\prime}\right)\right)_{n}$ is zero and so $\mathcal{F}^{\prime}$ is included in $\tilde{\mathcal{C}}$. Since $\mathcal{F}$ can be written as an union of such open sets $\mathcal{F}^{\prime}$ it comes $\mathcal{F}=\tilde{\mathcal{C}}$ and so $\mathcal{F}$ is a wandering Fatou component.

Proof of (3). Let $j=J$. By (1) the stable component $\mathcal{C}$ containing $D_{J}$ is wandering. Let $(z, v)$ be a point of $B_{J} \times\left\{a\left(\square_{J}\right)\right\}$ and $x:=H_{p}(z, v) \in D_{J} \subset \mathcal{C}$. By Theorem $1.45(\boldsymbol{)})$, there exist $0<t<1$ and $\mu \in \mathcal{M}_{p}\left(\Lambda_{p}\right)$ such that the limit set of the sequence of empirical measures $\mathrm{e}_{n}:=\frac{1}{n} \sum_{i=0}^{n-1} \delta_{F_{p}^{i}(z, v)}$ contains $t \cdot \mu+(1-t) \cdot \mathcal{M}_{p}\left(\Lambda_{p}\right)$. We notice that the pushforward of $\mathrm{e}_{n}$ by $H_{p}$ is equal to $\mathrm{e}_{n}^{f}:=\frac{1}{n} \sum_{i=0}^{n-1} \delta_{f_{p}^{i}(x)}$ since $f_{p} \circ H_{p}=H_{p} \circ F_{p}$. Also the pushforward by $H_{p}$ of $\mathcal{M}_{p}\left(\Lambda_{p}\right)$ is equal to $\mathcal{M}_{p}\left(K_{p}\right)$ since $K_{p}$ is the image of $\Lambda_{p}$ by the embedding $H_{p}$. As the pushforward operation on probability measures by an embedding is an embedding in the spaces of probability measures, the limit set of $\left(\mathrm{e}_{n}^{f}\right)_{n}$ contains $t \cdot \mu^{\prime}+(1-t) \cdot \mathcal{M}_{p}\left(K_{p}\right)$ with $\mu^{\prime}:=H_{p *} \mu \in \mathcal{M}_{p}\left(K_{p}\right)$. In particular $\left(e_{n}^{f}\right)_{n}$ diverges. Finally, for every $x^{\prime}$ in $\mathcal{C}$, the limit set of $\left(\frac{1}{n} \sum_{i=0}^{n-1} \delta_{f_{p}^{i}\left(x^{\prime}\right)}\right)_{n}$ is equal to the limit set of $\left(\mathrm{e}_{n}^{f}\right)_{n \geq 0}$ and so contains $t \cdot \mu^{\prime}+(1-t) \cdot \mathcal{M}_{p}\left(K_{p}\right)$, since $\mathcal{C}$ is a stable component.

Proof of (2). We now recall that any open set $U$ of $K_{p}$ has positive mass for a measure $\mu_{U} \in \mathcal{M}_{p}\left(K_{p}\right)$, and so for the measure $t \cdot \mu^{\prime}+(1-t) \cdot \mu_{U}$. Hence by (3) the limit set of any point in $\mathcal{C}$ intersects $U$ and so is dense in $K_{p}$. As the limit set is compact, it must contain $K_{p}$.

## 2. EXAMPLES OF FAMILIES DISPLAYING THE GEOMETRIC MODEL

In this Section, we give examples satisfying the geometric model (Definition 1.46).
First we show that the family $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of Example 1.37 is an unfolding of wild type $(\mathscr{A}, \mathscr{C})$ and we give a sketch of proof of how to deduce simply from this an example of wandering Fatou component for a real polynomial automorphism of $\mathbb{C}^{2}$.

The main result of this Section is Theorem 2.10 from which we will deduce Theorem 1.47 (and thus examples of wandering Fatou components for real polynomial automorphisms of $\mathbb{C}^{2}$ ) and Theorem 1.49. Namely, Theorem 2.10 states that an unfolding of wild type is embedded into some iterate of any $d$-parameter family which displays a non-degenerate unfolding of $d$-quadratic homoclinic tangencies. When $d=5$ and up to a dissipativeness assumption, the geometric model is satisfied. To prove both Theorems 1.49 and 2.10 we will need to recall classical results about Newhouse phenomenon (Theorems 2.15 and 2.16). Let us point out that the proof of Theorem 2.10 is technical and could be omitted in a first reading.
2.1. A simple example of unfolding of wild type $(\mathscr{A}, \mathscr{C})$. We are going to prove that for any even $N \geq 2$ and $\delta>0$ small enough, the family $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of Example 1.37 is a moderately dissipative unfolding of wild type ( $\mathscr{A}, \mathscr{C})$. To this end we must show that it satisfies (in particular) Property $\left(\mathbf{H}_{\mathbf{1}}\right)$ of Definition 1.43 , which by Remark 1.41 is equivalent to say that the Cantor sets $\left\{a_{p}^{u}: u \in \overleftarrow{\mathscr{A}}_{\square_{j}}\right\}$ and $\left\{b_{p}^{s}: s \in \overrightarrow{\mathscr{A}}_{\varpi_{j}}\right\}$ intersect for all $p \in \mathcal{P}$ and $1 \leq j \leq N-1$. This robust intersection will be obtained using the following well-known notion:

Definition 2.1 (Thickness of a Cantor set). Given a Cantor set $K \subset \mathbb{R}$, a gap of $K$ is a connected component of $\mathbb{R} \backslash K$. Given a bounded gap $G$ of $K$ and $x$ in the boundary of $G$, the bridge $B$ of $K$ at $x$ is the connected component of $x$ in the complement of the union of the gaps larger or equal than $|G|$. The thickness of $K$ at $x$ is $\tau(K, u)=|B| /|G|$. The thickness of $K$, denoted by $\tau(K)$ is the infimum of these $\tau(K, u)$ among all boundary points $x$ of bounded gaps.

The following is the celebrated Newhouse gap Lemma.
Lemma 2.2 ([New79, Lemma 4]). Let $K_{1}, K_{2} \subset \mathbb{R}$ be Cantor sets with thicknesses $\tau_{1}$ and $\tau_{2}$. If $\tau_{1} \cdot \tau_{2}>1$, then one of the three following possibilities occurs: $K_{1}$ is contained in a component of $\mathbb{R} \backslash K_{2}, K_{2}$ is contained in a component of $\mathbb{R} \backslash K_{1}$ or $K_{1} \cap K_{2} \neq \emptyset$.

We are now ready to prove:
Proposition 2.3. For every even integer $N \geq 2$ and $\delta>0$ smaller than a positive function of $N$, the family $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of Example 1.37 is an unfolding of wild type $(\mathscr{A}, \mathscr{C})$ which is moderately dissipative.

Proof. We recall that $\mathscr{A}=\left\{a_{j}: 1 \leq j \leq N\right\}$ and the hyperbolic transformations are

$$
Y^{a_{j}}:=I_{j} \times I \quad \text { and } \quad F^{a_{j}}:=(x, y) \in Y^{a_{j}} \mapsto\left(C_{j}^{-1}(x), \sqrt{\delta} \cdot C_{j}(y)\right)
$$

with $I_{j}:=\left[\frac{2}{N}(j-1)-1+\delta^{2}, \frac{2}{N} j-1-\delta^{2}\right]$ and $C_{j}$ the affine, orientation preserving map from $I$ onto $I_{j}$. We notice that with $\epsilon=0.1$, if $\delta$ is sufficiently small, then for every $a \in \mathscr{A}$,

$$
\max _{p \in \mathcal{P}, z \in Y^{a}}\left\|D_{z} F^{a}\right\| \cdot\left|\operatorname{det} D_{z} F^{a}\right|^{\epsilon}=\left(\frac{1}{N}-\delta^{2}\right)^{-1} \delta^{\epsilon / 2}<1 .
$$

Thus $F^{\mathscr{A}}$ is moderately dissipative. Also the family of systems $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of type $(\mathscr{A}, \mathscr{C})$ satisfies property $\left(\mathbf{H}_{2}\right)$. Indeed, for $\left(u_{j}, s_{j}\right)_{j} \in \prod_{i=1}^{N-1} \overleftarrow{\mathscr{A}}_{\square_{j}} \times \overrightarrow{\mathscr{A}}_{\unlhd_{j}}$, by Eq. (1.3) of Example 1.20, we have for every $1 \leq j \leq N-1$ :

$$
\begin{equation*}
a^{u_{j}}=y^{u_{j}}+p_{j}, \quad b_{p}^{s_{j}}=x^{s_{j}} \quad \text { and } \quad \mathcal{V}_{p}^{u_{j},,_{j+1}}=y^{u_{j}}-x^{s_{j+1}}+p_{j} \tag{2.1}
\end{equation*}
$$

Thus the map $p \in \mathcal{P} \mapsto\left(\mathcal{V}^{u_{1}, s_{2}}, \cdots, \mathcal{V}_{p}^{u_{N-1}, s_{N}}\right) \in \mathbb{R}^{N-1}$ is an affine map of linear part Id, and so a diffeomorphism.

Let us show that the regular family $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of systems satisfies $\left(\mathbf{H}_{1}\right)$. To this end, we shall study the following Cantor sets:

$$
K_{s}:=\left\{x^{s}: s \in \overrightarrow{\mathscr{A}}\right\} \quad \text { and } \quad K_{u}:=\left\{y^{u}: u \in \overleftarrow{\mathscr{A}}\right\}
$$

The Cantor set $K_{s}$ is the limit set of the IFS given by the contractions $\left(C_{j}\right)_{j=1}^{N}$. The Cantor set $K_{u}$ is the limit set of the IFS given by the contractions $\left(\sqrt{\delta} \cdot C_{j}\right)_{j=1}^{N}$. It suffices to show that for every $p=\left(p_{j}\right)_{j} \in \mathcal{P}$ and $1 \leq k \leq N-1$, the Cantor set $\sqrt{\delta} \cdot C_{k}\left(K_{u}\right)+p_{k}$ intersects $K_{s}$. To show this we compute the thicknesses of these Cantor sets and their diameters. To this end, it is worth to recall:

Fact 2.4. The intervals $\left(I_{j}\right)_{j}$ are disjoint and display the same length $2 / N-2 \delta^{2}$. The complement of their union in $\left[-1+\delta^{2}, 1-\delta^{2}\right]$ is made by intervals of length $2 \delta^{2}$. The left endpoint of $I_{1}$ is $-1+\delta^{2}$ and the right endpoint of $I_{N}$ is $1-\delta^{2}$. The component of 0 in $I \backslash \bigcup_{j=1}^{N} I_{j}$ is $V^{\square}:=\left(-\delta^{2}, \delta^{2}\right)$.

We are ready to show:
Lemma 2.5. The Cantor set $K_{s}$ is at distance $\frac{N \delta^{2}}{N-1+N \delta^{2}} \sim \frac{\delta^{2} N}{N-1}$ from the boundary of $I$ and its thickness is $\tau\left(K_{s}\right):=\frac{N-1+N \delta^{2}}{\delta^{2} N^{2}}-1 \sim \frac{N-1}{\delta^{2} N^{2}}$ when $\delta \rightarrow 0$.
Proof. Indeed, the Cantor set $K_{s}$ is the limit set of the IFS given by the contractions $\left(C_{j}\right)_{j=1}^{N}$ and each $C_{j}$ is affine, so we only need to compute the ratio $|B| /|G|$ for the $(N-1)$ gaps between $C_{j}(I)=I_{j}$ and $C_{j+1}(I)=I_{j+1}$ for $1 \leq j \leq N-1$. By Fact 2.4, each of these gaps has length:

$$
\begin{equation*}
|G|:=2 \delta^{2} \sum_{j \geq 0}\left(\frac{1}{N}-\delta^{2}\right)^{j}=\frac{2 \delta^{2}}{1-\frac{1}{N}+\delta^{2}}=\frac{2 \delta^{2} N}{N-1+N \delta^{2}}=2 d\left(K_{s}, I^{c}\right) \tag{2.2}
\end{equation*}
$$

The bridges between these gaps have length $|B|:=\frac{2}{N}-|G|$. Then the proposition follows from the quotient of these two equalities:

$$
\tau\left(K_{s}\right):=\frac{2}{N|G|}-1=\frac{N-1+N \delta^{2}}{\delta^{2} N^{2}}-1
$$

Remark 2.6. By Eq. (2.2), for every $j$, the lengths of gaps of $K_{s}^{j}:=K_{s} \cap I_{j}$ are $\leq \frac{2 \delta^{2} N}{N-1} \cdot\left(\frac{1}{N}-\delta^{2}\right) \sim$ $\frac{2 \delta^{2}}{N-1}$. Also d $\left(I_{j}^{c}, K_{s}^{j}\right) \sim\left(\frac{1}{N}-\delta^{2}\right) \cdot \frac{\delta^{2} N}{N-1} \sim \frac{\delta^{2}}{N-1}$.

Lemma 2.7. The Cantor set $K_{u}$ has diameter equivalent to $2 \sqrt{\delta} \frac{N-1}{N}$ and its thickness is equivalent to $\sqrt{\delta} \frac{N-1}{N}$ when $\delta \rightarrow 0$.

Proof. The Cantor set $K_{u}$ is the limit set of the IFS given by the contractions $\left(\sqrt{\delta} \cdot C_{j}\right)_{j=1}^{N}$ and each $\sqrt{\delta} \cdot C_{j}$ is affine. Let $[-a, a]$ be the convex hull of $K_{u}$. We observe that $a$ is the unique fixed point of $\sqrt{\delta} \cdot C_{N}$, whereas $-a$ is the fixed point of $\sqrt{\delta} \cdot C_{1}=x \mapsto-\sqrt{\delta} \cdot C_{N}(-x)$. Since by definition $C_{N}$ is the orientation preserving affine map sending $I$ to $I_{N}$, this yields that $\sqrt{\delta} \cdot C_{N}(x)=\sqrt{\delta}\left(\left(\frac{1}{N}-\delta^{2}\right)(x-1)+1-\delta^{2}\right)$ for $x \in I$. Thus the unique fixed point of $\sqrt{\delta} \cdot C_{N}$ is:

$$
a:=\sqrt{\delta}\left(1-\frac{1}{N}\right) \cdot\left(1-\frac{\sqrt{\delta}}{N}+\delta^{5 / 2}\right)^{-1} \sim \sqrt{\delta} \frac{N-1}{N} \quad \text { when } \delta \rightarrow 0
$$

This implies the estimate on the diameter of $K_{u}$. To compute the thickness, we only need to compute the ratio $|B| /|G|$ for the $(N-1)$ gaps between $\sqrt{\delta} \cdot C_{j}([-a, a])$ and $\sqrt{\delta} \cdot C_{j+1}([-a, a])$ for $1 \leq j \leq N-1$. We denote $\tilde{I}_{j}=\sqrt{\delta} \cdot C_{j}([-a, a]) \subset[-a, a]$. We observe that each $\tilde{I}_{j}$ is of length $b:=2 a \cdot \sqrt{\delta}\left(\frac{1}{N}-\delta^{2}\right)$. There are $N-1$ gaps between the segments $\tilde{I}_{j}$ which are equal. Thus their size is $\frac{2 a-N b}{N-1}$. Thus the thickness is:

$$
b \cdot\left(\frac{2 a-N b}{N-1}\right)^{-1}=(N-1) \frac{2 a \cdot \sqrt{\delta}\left(\frac{1}{N}-\delta^{2}\right)}{2 a-N \cdot 2 a \cdot \sqrt{\delta}\left(\frac{1}{N}-\delta^{2}\right)} \sim \sqrt{\delta} \frac{N-1}{N}
$$

For $1 \leq j \leq N$, let us denote:

$$
\begin{equation*}
K_{s}^{j}:=K_{s} \cap I_{j}=C_{j}\left(K_{s}\right) \quad \text { and } \quad K_{u}^{j}:=\sqrt{\delta} \cdot C_{j}\left(K_{u}\right) \tag{2.3}
\end{equation*}
$$

The sets $K_{s}^{j}$ and $K_{u}^{j}$ are Cantor sets of the same thickness as respectively $K_{s}$ and $K_{u}$. Hence by Lemmas 2.5 and 2.7, the product of the thicknesses of $K_{s}^{j}$ and $K_{u}^{j}$ is greater than 1. By Remark 2.6, all the gaps of $K_{s}^{j}$ are smaller than $\sim \frac{2 \delta^{2}}{N-1}$ and $\sim \frac{\delta^{2}}{N-1}$ distant to the boundary of $I_{j}$. By Lemma 2.7, this is small compared to the diameter of $K_{u}^{j}$ which is equivalent to $\sim 2 \frac{\sqrt{\delta}(N-1)}{N} \cdot \sqrt{\delta}\left(\frac{1}{N}-\delta^{2}\right)$. As $K_{u}^{j}+p_{j}$ is included in $I_{j+1}$, by the gap Lemma 2.2, it must intersect $K_{s}^{j+1}$ for every $p \in \mathcal{P}$. This proves Property $\left(\mathbf{H}_{1}\right)$.

Remark 2.8. It is already possible to deduce an example of wandering Fatou component for a real polynomial automorphism of $\mathbb{C}^{2}$.

Here we just give a sketch of proof of this affirmation since in the next subsection we will prove the same conclusion for a Hénon family of degree 6 (to prove Theorem A).

Sketch of proof. The map $F_{p}$ has constant Jacobian determinant for every $p$. Thus by DacarognaMoser theorem [DM90], it is possible to extend the family $\left(F_{p}\right)_{p}$ to a family $\left(f_{p}\right)_{p}$ of diffeomorphisms $f_{p}$ of $\mathbb{R}^{2}$ such that $\left(f_{p}\right)_{p}$ satisfies the geometric model when $N=6$. Thus the map $f_{p}$ has a wandering stable component at a dense set of parameter $p$. Moreover, it is possible to show the existence of a neighborhood $\mathcal{U}$ of $\left(f_{p}\right)_{p \in \mathcal{P}}$ formed by families displaying the geometric model (the proof is not complicated but rather boring since many items must be checked) via an embedding $H_{p}$ equal to the canonical inclusion: $Y^{e} \times\{\diamond\} \hookrightarrow \mathbb{R}^{2}$. Then by noting that $\left(\delta^{-1 / 4} \cdot f_{p}\right)_{p \in \mathcal{P}}$ has constant Jacobian equal to 1, by Turaev's Theorem [Tur02, Thm. 2 and rem. 1], this family can be approximated by a smooth family of compositions of generalized Hénon-maps and thus by a smooth family of compositions of polynomial generalized Hénon-maps $\left(Q_{p}\right)_{p \in \mathcal{P}}$. These are polynomial automorphisms of $\mathbb{R}^{2}[X, Y]$. Remark that $\tilde{f}_{p}:=\delta^{1 / 4} \cdot Q_{p}$ is a polynomial automorphism
whose restriction to $[-1,1]^{2}$ is $C^{2}$-close to $f_{p}$. Moreover $\left(\tilde{f}_{p} \mid[-1,1]^{2}\right)_{p}$ is $C^{2}$-close to $\left(f_{p} \mid[-1,1]^{2}\right)_{p}$ and so it displays the geometric model via a smooth family of analytic embeddings. Thus by Corollary 1.50 , there exists a dense set of parameters $p \in \mathcal{P}$ such that the biholomorphic extension of $f_{p}$ to $\mathbb{C}^{2}$ has a wandering Fatou component.
2.2. Natural examples satisfying the geometric model. In this subsection, we are going to state Theorem 2.10 from which we will deduce Theorems 1.47 and 1.49. Theorem 2.10 states that an unfolding of a wild type is embedded into some iterate of any $d$-parameter family which displays a non-degenerate unfolding of $d$-quadratic homoclinic tangencies:

Definition 2.9 (Non-degenerate unfolding). Let $d \geq 1, \hat{\mathcal{P}}$ an open set of $\mathbb{R}^{d}$, and $\left(f_{p}\right)_{p \in \hat{\mathcal{P}}}$ be a $C^{2}$-family of surface diffeomorphisms leaving invariant a continuation ${ }^{6}\left(K_{p}\right)_{p \in \hat{\mathcal{P}}}$ of hyperbolic basic sets. Given $p_{0} \in \hat{\mathcal{P}}$ and d points of quadratic tangencies between the stable and unstable manifolds of $K_{p_{0}}$, their unfolding along $\left(f_{p}\right)_{p}$ is non-degenerated if the matrix with the following coefficients $\left(a_{i, j}\right)_{(i, j) \in\{1, \ldots, d\}^{2}}$ is invertible: $a_{i, j}$ is the derivative following the $j^{\text {th }}$-coordinate of $\mathbb{R}^{d}=T_{p_{0}} \mathcal{P}$ of the relative position of the local stable and unstable manifolds associated to the $i^{\text {th }}$-quadratic tangency point.

We recall that a periodic point $P$ of period $n$ is non-conservative if $\left|\operatorname{det} D_{P} f^{n}\right| \neq 1$. Here is the general result enabling the aforementioned applications:

Theorem 2.10. Let $r \in[2, \infty] \cup\{\omega\}$, let $d \geq 1$, let $\hat{\mathcal{P}}$ be a non empty open subset of $\mathbb{R}^{d}$ and let $\left(f_{p}\right)_{p \in \hat{\mathcal{P}}}$ be a $C^{r}$-family of surface diffeomorphisms. Given $p_{0} \in \hat{\mathcal{P}}$, assume that $f_{p_{0}}$ has a non-conservative ${ }^{7}$ periodic saddle point $P$ displaying d different quadratic homoclinic tangencies.

If these quadratic tangencies unfold non-degenerately with $\left(f_{p}\right)_{p}$, then there exist $k \geq 1$, a regular compact subset $\mathcal{P} \subset \hat{\mathcal{P}}$ arbitrarily close to $p_{0}$, an unfolding $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of a certain wild type $(\mathscr{A}, \mathscr{C})$ with $\operatorname{Card} \mathscr{C}=d,\left(F_{p}\right)_{p}$ being of class $C^{r}$ and a $C^{r}$-family of embeddings $\left(H_{p}\right)_{p}$ such that

$$
f_{p}^{k} \circ H_{p}=H_{p} \circ F_{p} \quad \forall p \in \mathcal{P} .
$$

Moreover $F_{p}^{s t}$ is moderately dissipative for every $p \in \mathcal{P}$ if $P$ is dissipative.
Remark 2.11. Hence if $\operatorname{Card} \mathscr{C}=5$, the family $\left(f_{p}\right)_{p \in \mathcal{P}}$ satisfies the $C^{r}$-geometric model 1.46 .
The proof of this theorem is given in the next subsection 2.3. Let us deduce now rather quickly Theorems 1.47 and 1.49 and Theorem C.

Proof of Theorem C. This is an immediate consequence of Corollary 1.50 and Remark 2.11.
Proof of Theorem 1.47 and Remark 1.48. Let us consider the Tchebychev polynomial of degree 6:

$$
P_{0}(X)=32 \cdot X^{6}-48 \cdot X^{4}+18 \cdot X^{2}-1
$$

The critical points and values of this polynomial are respectively:

$$
\left(\zeta_{i 0}\right)_{1 \leq i \leq 5}=(-\sqrt{3} / 2,-1 / 2,0,1 / 2, \sqrt{3} / 2) \quad \text { and } \quad\left(a_{i 0}\right)_{1 \leq i \leq 5}=(-1,1,-1,1,-1)
$$

[^1]They are all quadratic. All critical values are sent to $\beta_{0}=1$ which is a repelling fixed point of $P_{0}$. We notice that $\left[-\beta_{0}, \beta_{0}\right]$ is the complement of the basin of $\infty$. For $p=\left(p_{0}, \ldots, p_{4}\right) \in \mathbb{R}^{5}$, we consider:

$$
P_{p}(X)=P_{0}(X)+\sum_{j=0}^{4} p_{j} \cdot X^{j}
$$

Let $\beta_{p},\left(\zeta_{i p}\right)_{1 \leq i \leq 5}$ and $\left(a_{i p}\right)_{1 \leq i \leq 5}$ be the respective continuations of the fixed point $\beta_{0}$, of the critical points and of the critical values for $P_{p}$. We have the following:

Lemma 2.12. The unfolding of the images $P_{p}\left(a_{i p}\right)$ of the critical values $a_{i p}$ w.r.t. the fixed point $\beta_{p}$ is non-degenerated:

$$
\operatorname{det}\left[\partial_{p_{j}}\left(P_{p}\left(a_{i p}\right)\right)-\partial_{p_{j}} \beta_{p}\right]_{i, j} \neq 0 \quad \text { at } p=0
$$

Proof. One way to prove this lemma is to use Epstein's transversality (see [BE09], Thm 1 and the subsequent remark). Let us give a direct proof. We have $D P_{0}\left(\beta_{0}\right)=36$ and $\partial_{p} P_{p}(X)=$ $\sum_{j=0}^{4} X^{j} \cdot d p_{j}$. Thus we have $\left(D P_{0}\left(\beta_{0}\right) \cdot \partial_{p} \beta_{p}+\partial_{p} P_{p}\left(\beta_{0}\right)\right)_{\mid p=0}=\left(\partial_{p} \beta_{p}\right)_{\mid p=0}$ and so:

$$
\left(\partial_{p} \beta_{p}\right)_{\mid p=0}=-\frac{1}{35} \sum_{j=0}^{4} d p_{j}
$$

Also $\left(\partial_{p} a_{i p}\right)_{\mid p=0}=\left(D P_{0}\left(\zeta_{i 0}\right) \cdot \partial_{p} \zeta_{i p}+\partial_{p} P_{p}\left(\zeta_{i 0}\right)\right)_{\mid p=0}=\left(\partial_{p} P_{p}\left(\zeta_{i 0}\right)\right)_{\mid p=0}$ implies:

$$
\left(\partial_{p} a_{i p}\right)_{\mid p=0}=\sum_{j=0}^{4}\left(\zeta_{i 0}\right)^{j} \cdot d p_{j} \quad \text { and } \quad\left(\partial_{p} P_{p}\left(a_{i p}\right)\right)_{\mid p=0}=D P_{0}\left(a_{i 0}\right) \cdot\left(\partial_{p} a_{i p}\right)_{\mid p=0}+\left(\partial_{p} P_{p}\left(a_{i 0}\right)\right)_{\mid p=0}
$$

and so:

$$
\left(\partial_{p} P_{p}\left(a_{i p}\right)\right)_{\mid p=0}=(-1)^{i} \cdot 36 \cdot \sum_{j=0}^{4}\left(\zeta_{i 0}\right)^{j} \cdot d p_{j}+\sum_{j=0}^{4}(-1)^{i \cdot j} \cdot d p_{j}
$$

Then to show the lemma, it suffices to see the invertibility of the $5 \times 5$-matrix whose $(i, j)$ entry is $\left[\partial_{p_{j}}\left(P_{p}\left(a_{i p}\right)\right)-\partial_{p_{j}} \beta_{p}\right]=\left[(-1)^{i} \cdot 36\left(\zeta_{i 0}\right)^{j}+(-1)^{i \cdot j}+1 / 35\right]$.

For $b \geq 0$ small and $p \in \mathbb{R}^{5}$ small, we consider the Hénon map:

$$
f_{p b}:(x, y) \mapsto\left(P_{p}(x)-y, b \cdot x\right)
$$

For $b=0$, the family $\left(f_{p 0}\right)_{p \in \mathbb{R}^{5}}$ is semi-conjugate to $\left(P_{p}\right)_{p \in \mathbb{R}^{5}}$ via the maps $(x, y) \mapsto P_{p}(x)-y$, while for $b>0$, the maps $f_{p b}$ are polynomial automorphisms. For $b=0$, the point $\beta_{00}:=\left(\beta_{0}, 0\right)$ is fixed and hyperbolic for $f_{00}$. Thus it persists as a hyperbolic fixed point $\beta_{p b}$ of $f_{p b}$ for $b>0$. We notice that $W_{l o c}^{u}\left(\beta_{00} ; f_{00}\right):=(-0.9,1.1) \times\{0\}$ is a local unstable manifold of $\beta_{00}$ : there exists an arbitrarily small neighborhood of $\beta_{00}$ in $\mathbb{R} \times\{0\}$ which is sent diffeomorphically to $W_{l o c}^{u}\left(\beta_{00} ; f_{00}\right)$ by an iterate of $f_{00}$. For $p$ and $b$ small, $W_{l o c}^{u}\left(\beta_{00} ; f_{00}\right)$ persists as a curve $W_{l o c}^{u}\left(\beta_{p b} ; f_{p b}\right) C^{2}$-close to $W_{l o c}^{u}\left(\beta_{00} ; f_{00}\right)$. On the other hand, the set

$$
W_{l o c}^{s}\left(\beta_{00} ; f_{00}\right):=\left\{(x, y) \in[-2,2]^{2}: P_{0}(x)-y=\beta_{0}\right\}
$$

is a local stable manifold of $f_{000}$. For $p$ and $b$ small, it persists as a curve $W_{l o c}^{s}\left(\beta_{p b} ; f_{p b}\right) C^{2}$ close to $W_{l o c}^{s}\left(\beta_{00} ; f_{00}\right)$. As the critical points $\left(\zeta_{i 0}, 0\right)_{1 \leq i \leq 5}$ belong to $W_{l o c}^{u}\left(\beta_{00} ; f_{00}\right)$ and are sent by $f_{00}^{2}$ to quadratic tangency points between $f_{00}^{2}\left(W_{l o c}^{u}\left(\beta_{00} ; f_{00}\right)\right.$ and $W_{l o c}^{s}\left(\beta_{00} ; f_{00}\right)$ which unfold
non-degerately by Lemma 2.12, for $b$ small there is a proper compact subset $\hat{\mathcal{P}}$ close to 0 and $p_{b} \in \hat{\mathcal{P}}$ close to 0 such that $f_{p_{b} b}^{2}\left(W_{l o c}^{u}\left(\beta_{p_{b} b} ; f_{p_{b} b}\right)\right.$ and $W_{l o c}^{s}\left(\beta_{p_{b} b} ; f_{p_{b} b}\right)$ have five points of quadratic tangency $\left(f_{p_{b} b}^{2}\left(\tilde{\zeta}_{i b}\right)\right)_{1 \leq i \leq 5}$ which unfold non-degenerately with $\left(f_{p b}\right)_{p \in \hat{\mathcal{P}}}$. Moreover $\left(\tilde{\zeta}_{i b}\right)_{1 \leq i \leq 5}$ are close to $\left(\left(\zeta_{i 0}, 0\right)\right)_{1 \leq i \leq 5}$. Hence we can apply Remark 2.11 of Theorem 2.10 with $P=\beta_{p b}$ to deduce Theorem 1.47. The same argument proves Remark 1.48.

We recall that that a point $x$ of a hyperbolic compact set $K$ displays a homoclinic tangency if its stable manifold is tangent to its unstable manifold. Theorem 2.10 enables also to deduce rather quickly Theorem 1.49 on the Newhouse domain:

Definition 2.13. Let $r \in[2, \infty] \cup\{\omega\}$. A basic set $K$ for a surface diffeomorphism $f$ is $C^{r}$ wild if for every $C^{r}$-perturbation $\tilde{f}$ of $f$, there is a point in the hyperbolic continuation of $K$ for $\tilde{f}$ displaying a quadratic homoclinic tangency. The dissipative $C^{r}$-Newhouse domain $\mathcal{N}^{r}$ is the open set of surface diffeomorphisms leaving invariant a hyperbolic basic set displaying a $C^{r}$-robust quadratic homoclinic tangency and an area contracting periodic point ${ }^{8}$.

The existence of wild hyperbolic basic sets was proved by Newhouse [New74] by extending the concept of thickness to hyperbolic basic sets:

Definition 2.14 (Stable and unstable thicknesses). The stable thickness $\tau^{s}(K)$ of a horseshoe $K$ for a dynamics $f$ is the infimum of the thicknesses of $K \cap W_{\text {loc }}^{u}(z)$ among $z \in K$. The unstable thickness $\tau^{u}(K)$ of $K$ is is the infimum of the thicknesses of $K \cap W_{\text {loc }}^{s}(z)$ among $z \in K$.

A point $z \in K$ does not bound an unstable gap if $z$ does not bound a gap of the Cantor set $W_{\mathrm{loc}}^{s}(z) \cap K$ of $W_{\mathrm{loc}}^{s}(z)$. A point $z \in K$ does not bound a stable gap if $z$ does not bound a gap of the Cantor set $W_{\mathrm{loc}}^{u}(z) \cap K$ of $W_{\mathrm{loc}}^{u}(z)$.

Here is Newhouse's theorem implying the existence of wild basic sets:
Theorem 2.15 (Newhouse [New79]). Let $f$ be a surface $C^{2}$-dynamics having a hyperbolic basic set $\Lambda$ containing two basic sets $\Lambda_{1}$ and $\Lambda_{2}$ satisfying the following properties:
(1) $\tau^{u}\left(\Lambda_{1}\right) \cdot \tau^{s}\left(\Lambda_{2}\right)>1$.
(2) There is a point $P_{1} \in \Lambda_{1}$ which does not bound an unstable gap of $\Lambda_{1}$ and a point $P_{2} \in \Lambda_{2}$ which does not bound a stable gap of $\Lambda_{2}$, such that a local unstable manifold $W_{l o c}^{u}\left(P_{1}\right)$ has a quadratic tangency with a local stable manifold $W_{\text {loc }}^{s}\left(P_{2}\right)$ at a point $T$,
then for any $C^{2}$-perturbation $\tilde{f}$ of $f$, the continuation of $\Lambda$ has local unstable and stable manifolds $W_{l o c}^{u}\left(\tilde{P}_{1}\right)$ and $W_{l o c}^{s}\left(\tilde{P}_{2}\right)$ close to respectively $W_{l o c}^{u}\left(P_{1}\right)$ and $W_{l o c}^{s}\left(P_{2}\right)$ and displaying a quadratic tangency at a point $\tilde{T}$ close to $T$.

Proof. Let us briefly recall Newhouse's argument and see how we obtain that $\tilde{T}$ is close to $T$ (this is not explicitly stated in [New79] but it is an easy consequence of his construction). Indeed, there exists a foliation $\mathcal{F}$ whose tangent space is of class $C^{1}$ and which extends the stable lamination of $\Lambda$ (see Theorem C.7) on a neighborhood $L$ of $\Lambda$. Up to looking an iterate, we can assume $T \in L$. We use a local unstable manifold $W$ which is transverse to this foliation and use it to define the

[^2]holonomy hol from a neighborhood of $T$ in $L$ into $W$ along $\mathcal{F}$. The image of the stable lamination of $\Lambda_{2}$ is a Cantor set $K_{2}$ with thickness at least $\tau^{s}\left(\Lambda_{2}\right)$. Nearby $T$, the tangency points between the leaves of $\mathcal{F}$ and the unstable lamination of $\Lambda_{1}$ is sent by hol onto a Cantor set $K_{1}$ of thickness locally at least close to $\tau^{u}\left(\Lambda_{1}\right)$. By assumption $K_{1}$ intersects $K_{2}$ at hol $(T)$ which does not bound a gap of $K_{1}$ nor of $K_{2}$. Thus any perturbation of $K_{1}$ and $K_{2}$ cannot be included in the gap of the one other. When we perturb $f$ the tangent space to the foliation $\mathcal{F}$ persists to one which is $C^{1}$-close. Thus the compact sets $K_{1}$ and $K_{2}$ persist. By [New79, prop 6] their thicknesses vary also continuously and so by the Gap Lemma 2.2 we get that $\tilde{T}$ is close to $T$.

Newhouse showed also that wild hyperbolic sets appear nearby homoclinic tangencies:
Theorem 2.16 (Newhouse [New79, Lem. 7 and Thm 3]). Let $\left(f_{t}\right)_{t \in(-1,1)}$ be a $C^{r}$-family of surface diffeomorphisms such that $f_{0}$ has a homoclinic quadratic tangency associated to a $n_{0}$-periodic point $P_{0}$, which unfolds non-degenerately and such that $\left|\operatorname{det} D f^{n_{0}}\left(P_{0}\right)\right|<1$. Then arbitrarily close to 0 , there exists a nonempty open interval I and there is a continuation $\left(K_{t}\right)_{t \in I}$ of hyperbolic basic sets such that for every $t \in I$, the basic set $K_{t}$ has a homoclinic tangency which is non-degenerately unfolded and $K_{t}$ contains the continuation $P_{t}$ of $P_{0}$.

To show this, Newhouse constructed a first horseshoe $\Lambda_{1}(t) \ni P_{t}$ with unstable thickness $\tau^{u}>0$ (see [New79, Prop 6]) which persists for $t$ in an open interval $\hat{I}$ nearby 0 . Then for any $T>0$ he found for a certain $t \in \hat{I}$ a second horseshoe $\Lambda_{2}(t)$ such that:

- the stable thickness of $\Lambda_{2}(t)$ is $>T$ (see [New79, $\left.\ell 16 \mathrm{P} .124\right]$ ).
- A local unstable manifold of $\Lambda_{1}(t)$ has a quadratic homoclinic tangency with a local stable manifold of $\Lambda_{2}(t)$ which is non-degenerately unfolded by $\left(f_{t}\right)_{t}$ (see [New79, Lem. 7]).
- the basic set $\Lambda_{1}(t)$ and $\Lambda_{2}(t)$ are included in a basic set $K_{t}$ (see [New79, Lem. 7-8]).

Then Theorem 2.16 is a consequence of this construction with $T>1 / \tau_{u}$ and of Theorem 2.15. Taking $T$ greater than $\max \left(\left(1+2 \tau_{u}\right)^{2} / \tau_{u}^{3},\left(\tau_{u}^{2}+3 \tau_{u}+1\right) / \tau_{u}^{2}\right)$, the main theorem of [Kra92, HKY93] implies furthermore:

Corollary 2.17. Under the assumptions of Theorem 2.16, arbitrarily close to 0 , there exists a nonempty open interval I and there is a continuation of a hyperbolic basic set $\left(K_{t}\right)_{t \in I}$ containing two basic sets $\Lambda_{1}(t) \ni P_{t}$ and $\Lambda_{2}(t)$ such that
(i) $\tau^{u}\left(\Lambda_{1}(t)\right) \cdot \tau^{s}\left(\Lambda_{2}(t)\right)>1$, moreover $K_{t}$ is a Cantor set,
(ii) there is a Cantor set of quadratic homoclinic tangencies between the leaves of an unstable lamination $W^{u}\left(\Lambda_{1}(t)\right)$ and a stable lamination $W^{s}\left(\Lambda_{2}(t)\right)$ which unfold non-degenerately.

Remark 2.18. Using modern techniques (see [GST08, Lemma 1], [Ber18, example 2.3], and historically [TLY86, Eqs. (3.2) and (3.3)]), the set $\Lambda_{2}(t)$ constructed by Newhouse is given by a renormalization of a horseshoe of a perturbation of $(x, y) \in[-4,4]^{2} \mapsto\left(x^{2}-2,0\right)$. Thus we can assume that $f_{t}$ is asymptotically moderately dissipative at $\Lambda_{2}(t)$ :

$$
\limsup _{n \rightarrow \infty}\left\|\operatorname{det} D f^{n}\left|\Lambda_{2}(t)\left\|^{1 / \epsilon} \cdot\right\| D f^{n}\right| \Lambda_{2}(t)\right\|<1
$$

We are now ready to show that the geometric model appears typically in the Newhouse domain:

Proof of Theorem 1.49. Let $f$ be a $C^{r}$-surface diffeomorphism having a dissipative saddle periodic point $P$ displaying a quadratic homoclinic tangency at a point $T$. We first apply Proposition 2.19 given below with $d=5$ to perturb $f$ into $\tilde{f}$ so that the map $\tilde{f}$ has 5 different quadratic tangencies between the stable and unstable manifolds of a dissipative periodic saddle point $z$, which unfold non-degenerately with $\left(f_{p}\right)_{p \in \mathbb{R}^{5}}$ for any $\left(f_{p}\right)_{p \in \mathbb{R}^{5}}$ in a $C^{r}$-open and dense set of $C^{r}$-families containing $\tilde{f}: f_{0}=\tilde{f}$. To conclude, we apply Remark 2.11 of Theorem 2.10 to this open and dense set of 5 -dimensional unfoldings of $\tilde{f}$.

Proposition 2.19. Let $f$ be a $C^{r}$-surface diffeomorphism having a dissipative saddle periodic point $P$ displaying a quadratic homoclinic tangency at a point $T$. Then for any $d>1$, there exists $\tilde{f} C^{r}$-close to $f$ such that $\tilde{f}$ has d different quadratic homoclinic tangencies associated to some dissipative periodic saddle point $z$. These tangencies unfold non-degenerately with $\left(f_{p}\right)_{p \in \mathbb{R}^{d}}$ for any $\left(f_{p}\right)_{p \in \mathbb{R}^{d}}$ in a $C^{r}$-open and dense set of $C^{r}$-families containing $\tilde{f}: f_{0}=\tilde{f}$.

Proof of Proposition 2.19. We begin by a lemma which allows to unfold the initial tangency:
Lemma 2.20. There exists a $C^{r}$-flow $\left(\phi^{t}\right)_{t}$ such that the quadratic homoclinic tangency of $P$ is non-degenerately unfolded by $\left(\phi^{t} \circ f\right)_{t}$.

Proof. For the sake of simplicity we assume that $P$ is a fixed point; the case where $P$ is $n \geq 2$ periodic is left to the reader. First we choose a neighborhood $N$ of $T$ sufficiently small to intersect $W_{l o c}^{s}(P, f) \cup W_{l o c}^{u}(P, f)$ only at the union of two small neighborhoods of $T$ in $W_{l o c}^{u}(P, f)$ and $W_{l o c}^{s}(P, f)$. Now we consider a $C^{\infty}$-vector field $\chi$ supported by $N$ such that $\chi(T)$ is not in the tangent space of $W_{l o c}^{u}(P, f)$ (or equivalently of $W_{l o c}^{s}(P, f)$ ) at $T$. Let $\left(\breve{\phi}^{t}\right)_{t}$ be the flow of $\chi$. We notice that the family $\left(\breve{\phi}^{t} \circ f\right)_{t}$ unfolds non-degenerately the homoclinic tangency at $T$. This gives the fact when $r \neq \omega$. Otherwise, let $\tilde{\chi}$ be an analytic vector field which is $C^{\infty}$-close to $\chi$. Let $\left(\phi^{t}\right)_{t}$ be its flow. Then the family $\left(\phi^{t} \circ f\right)_{t}$ unfolds non-degenerately the homoclinic tangency at $T$.

Now we use Corollary 2.17 which asserts the existence of $t_{0} \in \mathbb{R}$ small such that $\tilde{f}:=\phi^{t_{0}} \circ f$ has a hyperbolic basic set $K_{0}$ containing two basic sets $\Lambda_{1}$ and $\Lambda_{2}$ such that
(i) $\tau^{u}\left(\Lambda_{1}\right) \cdot \tau^{s}\left(\Lambda_{2}\right)>1$,
(ii) there is a Cantor set of quadratic tangencies between the leaves of $W^{u}\left(\Lambda_{1}\right)$ and $W^{s}\left(\Lambda_{2}\right)$

Note that the map $\tilde{f}$ is $C^{r}$-close to $f$ since $t_{0}$ is small. As a Cantor set is not countable, there exist $d$ points $\left(P_{i}\right)_{1 \leq i \leq d}$ with local unstable and stable manifolds $W^{u}\left(P_{i}, \tilde{f}\right)$ and $W^{s}\left(P_{i}, \tilde{f}\right)$ displaying a quadratic tangency at a point $T_{i}$, such that the orbits of $T_{i}$ and $T_{j}$ are disjoint when $j \neq i$.

Fact 2.21. The set of $C^{r}$-families $\left(f_{p}\right)_{p \in \mathbb{R}^{d}}$ such that $f_{0}=\tilde{f}$ and for which the quadratic tangencies at $\left(T_{i}\right)_{1 \leq i \leq d}$ unfold non-degenerately is $C^{r}$-open and dense.

Proof. Clearly the set of non-degenerate unfoldings is $C^{r}$-open. Let us show that it is $C^{r}$-dense. Let $\left(f_{p}\right)_{p \in \mathbb{R}^{d}}$ be a $C^{r}$-family containing $\tilde{f}$. Let us show that it can be $C^{r}$-approximated by one for which the $d$ homoclinic tangencies unfold non-degenerately. A proof like the one of Lemma 2.20 shows the existence of $d$ flows $\left(\phi_{i}^{t}\right)_{1 \leq i \leq d}$ with pairwise disjoint supports centered at $\left(T_{i}\right)_{1 \leq i \leq d}$ such that with:

$$
t=\left(t_{j}\right)_{j} \in \mathbb{R}^{d} \mapsto \phi^{t}:=\phi_{d}^{t_{d}} \circ \cdots \circ \phi_{1}^{t_{1}}
$$

the family $\left(\phi^{t} \circ \tilde{f}\right)_{p \in \mathbb{R}^{d}}$ unfolds non-degenerately the $d$ homoclinic tangencies at $\left(T_{i}\right)_{1 \leq i \leq d}$. Recall that given a matrix $A \in \mathrm{M}_{d}(\mathbb{R})$ and a matrix $B \in \mathrm{Gl}_{d}(\mathbb{R})$ there is an open and dense set of $\eta \in \mathbb{R}$ such that $A+\eta B$ is invertible. Thus for an open and dense set of small $\eta \in \mathbb{R}$, the $C^{r}$-perturbation $\left(\phi^{\eta \cdot p} \circ f_{p}\right)_{p \in \mathbb{R}^{d}}$ of $\left(f_{p}\right)_{p \in \mathbb{R}^{d}}$ unfolds non-degenerately the homoclinic tangencies at $\left(T_{i}\right)_{1 \leq i \leq d}$.

Using the continuity of the parametrized stable and unstable laminations Theorem C.5 Page 72, and the density of stable and unstable manifolds of a given periodic point $z \in K_{0}$ in them, up to a small parameter translation, we can assume that $P_{i}=z$ for every $i$. Again by Theorem C.5, the unfolding is still non degenerate.
2.3. Proof of Theorem 2.10. Let $P$ be a non-conservative periodic point of $f_{p_{0}}$ such that $W^{u}\left(P ; f_{p_{0}}\right)$ and $W^{s}\left(P ; f_{p_{0}}\right)$ have $d$ quadratic tangencies at points $T_{1}, \ldots, T_{d}$, which unfold nondegenerately with $\left(f_{p}\right)_{p}$. Up to inversing the dynamics, we can assume that $P$ is dissipative.

The aim of this subsection is to show that there exist $k \geq 1$, a regular compact subset $\mathcal{P} \subset \hat{\mathcal{P}}$ arbitrarily close to $p_{0}$, a moderately dissipative unfolding $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of a certain wild type $(\mathscr{A}, \mathscr{C})$ with Card $\mathscr{C}=d$, with $\left(F_{p}\right)_{p}$ being of class $C^{r}$, and a $C^{r}$-family of embeddings $\left(H_{p}\right)_{p}$ such that:

$$
f_{p}^{k} \circ H_{p}=H_{p} \circ F_{p} .
$$

In order to do so, we first construct a wild basic set $K$ with $d$ homoclinic tangencies which unfold non-degenerately, then we construct a family of hyperbolic maps $\left(F_{p}^{s l}\right)_{p}$ and finally we will construct a family of folding maps $\left(F_{p}^{\mathscr{C}}\right)_{p}$.

Step 1: Construction of the wild basic set $K$. The following uses Newhouse's Corollary 2.17.
Lemma 2.22. There is $p_{1} \in \hat{\mathcal{P}}$ close to $p_{0}$ such that:
(1) The map $f_{p_{1}}$ has a hyperbolic basic set $K$ containing the continuation of $P$ and two basic sets $\Lambda_{1}$ and $\Lambda_{2}$ satisfying $\tau^{u}\left(\Lambda_{1}\right) \cdot \tau^{s}\left(\Lambda_{2}\right)>1$. Moreover $K$ is a Cantor set.
(2) There are $d$ points of quadratic tangencies between local unstable and stable manifolds of points of $\Lambda_{1}$ and $\Lambda_{2}$ which do not bound unstable and stable gaps and which unfold nondegenerately with $\left(f_{p}\right)_{p}$.
(3) $f_{p_{1}}$ is asymptotically moderately dissiptative at $K$ :

$$
\limsup _{n \rightarrow \infty}\left\|\operatorname{det} D f_{p_{1}}^{n}\left|K\left\|^{1 / \epsilon} \cdot\right\| D f_{p_{1}}^{n}\right| K\right\|<1
$$

Proof. Up to a coordinate change of the parameter set $\hat{\mathcal{P}}$, we can assume that $p_{0}=0$ and that along the one-parameter family $\left(f_{t}\right)_{t \in J}$, where $J:=\mathbb{R} \times\{0\} \cap \hat{\mathcal{P}}$, the first homoclinic tangency of $P$ unfolds non degenerately while the $d$-1-other homoclinic tangencies remain tangent. We notice that with $\left[a_{i, j}\right]_{1 \leq i, j \leq d}$ the matrix of the unfolding (see Definition 2.9), we have $a_{1,1} \neq 0$ and $a_{i, 1}=0$ for every $2 \leq i \leq d$. As the matrix $\left[a_{i, j}\right]_{1 \leq i, j \leq d}$ is invertible, it holds that $\left[a_{i, j}\right]_{2 \leq i, j \leq d}$ is invertible. Now by Corollary 2.17, there exists $t_{0} \in J$ small and a hyperbolic basic set $K$ containing two basic sets $\Lambda_{1}$ and $\Lambda_{2}$ such that:
(1) the set $\Lambda_{1}$ contains the hyperbolic continuation of $P$,
(2) $\tau^{u}\left(\Lambda_{1}\right) \cdot \tau^{s}\left(\Lambda_{2}\right)>1$.
(3) the local unstable and stable manifolds of $\Lambda_{1}$ and $\Lambda_{2}$ have uncountably many points $T_{i}^{\prime}$ of quadratic tangency which unfold non-degenerately with $\left(f_{t}\right)_{t \in J}$.

As there are only countably many gaps in a Cantor set, we can assume that $T_{1}^{\prime}$ satisfies the no-ending gap condition (2) of Theorem 2.15.

At $\left(t_{0}, 0\right) \in \hat{\mathcal{P}}$, the matrix of the unfolding of $\left(T_{1}^{\prime}, T_{2}, T_{3}, \ldots, T_{d}\right)$ is of the form $\left[a_{i, j}^{\prime}\right]_{1 \leq i, j \leq d}$ with $\left[a_{i, j}^{\prime}\right]_{2 \leq i, j \leq d} \approx\left[a_{i, j}\right]_{2 \leq i, j \leq d}$ and $a_{i, 1}^{\prime} \neq 0$ iff $i=1$. Thus $\left[a_{i, j}^{\prime}\right]_{1 \leq i, j \leq d}$ is invertible and the unfolding of these $d$-homoclinic tangencies is non-degenerated. By density of the local unstable and stable manifolds of points in $\Lambda_{1}$ and $\Lambda_{2}$ which do not bound gaps, there is a parameter $p_{1} \in \hat{\mathcal{P}}$ nearby $\left(t_{0}, 0\right)$ which displays $d$ homoclinic tangencies at points $\left(T_{1}^{\prime}, T_{2}^{\prime}, T_{3}^{\prime}, \ldots, T_{d}^{\prime}\right)$ between local unstable and stable manifolds of points of $\Lambda_{1}$ and $\Lambda_{2}$ which satisfy the no-ending gap condition (2) of Theorem 2.15. Moreover by continuation, their unfolding is non-degenerate. By [New79, Prop 6] the unstable and stable thicknesses of these basic sets vary continuously. Thus at the parameter $p_{1}$, we still have $\tau^{u}\left(\Lambda_{1}\right) \cdot \tau^{s}\left(\Lambda_{2}\right)>1$. Then the following achieves the proof:

Fact 2.23. We can assume that $f_{p_{1}}$ is asymptotically moderately dissiptative at $K$.
Proof. Otherwise, as the unstable and stable manifolds of a given periodic point $\tilde{P} \in \Lambda_{2}$ is dense in $W^{u}(K)$ and $W^{s}(K)$, using the non-degeneracy of the unfolding, there is a parameter close to $p_{0}$ such that $\tilde{P}$ has $d$ homoclinic tangencies which unfold non-degenerately. By Remark 2.18 , the point $\tilde{P}$ is moderately dissipative. So we can redo the construction starting with this point $\tilde{P}$ instead of $P$, which is moderately dissipative and so is its dynamics nearby.

Remark 2.24. By the proof, we can assume furthermore that the unstable manifolds containing the tangency points have disjoint orbits.

Step 2: Construction of families of hyperbolic maps $\left(F_{p}^{\mathcal{d}}\right)_{p}$ and conjugacy maps $\left(H_{p}\right)_{p}$. The following is a general (but somehow classical) Proposition 2.25 which might be useful in other contexts than the proof of Theorem 2.10. Hence it will be proved in Appendix D.

Let $r \in[2, \infty] \cup\{\omega\}$, let $f_{0}$ be a $C^{r}$-diffeomorphism of a surface $M$ leaving invariant a horseshoe ${ }^{9} K_{0}$. Let $d \geq 0, \hat{\mathcal{P}}$ be an open neighborhood of $0 \in \mathbb{R}^{d}$ and let $\left(f_{p}\right)_{p \in \hat{\mathcal{P}}}$ be a $C^{r}$-family of diffeomorphisms of $M$ containing $f_{0}$. Let $\left(K_{p}\right)_{p}$ be the hyperbolic continuation of $K_{0}$.

Proposition 2.25. There exist $k \geq 1$ arbitrarily large, a regular compact neighborhood $\mathcal{P}$ of $0 \in \hat{\mathcal{P}}$, an oriented graph $(\mathscr{V}, \mathscr{A})$ and a $C^{1}$-family $\left(\breve{H}_{p}\right)_{p \in \mathcal{P}}$ of embeddings $\breve{H}_{p}: Y^{e} \times \mathscr{V} \hookrightarrow M$ such that
(0) There exists a $C^{1}$-family $\left(\breve{F}_{p}^{\mathcal{A}}\right)_{p}$ of maps $\breve{F}_{p}^{\mathscr{A}}: \breve{D}_{p}(\mathscr{A}) \rightarrow Y^{e} \times \mathscr{V}$ which leaves invariant the horizontal and vertical directions and such that $f_{p}^{k} \circ \breve{H}_{p} \mid \breve{D}_{p}(\mathscr{A})=\breve{H}_{p} \circ \breve{F}_{p}^{s d}$. and for any $C^{r}$-family $\left(H_{p}\right)_{p \in \mathcal{P}}$ which is $C^{1}$-close to $\left(\breve{H}_{p}\right)_{p \in \mathcal{P}}$ it holds:
(1) There exists a $C^{r}-$ family $\left(F_{p}^{\mathscr{A}}\right)_{p}$ of hyperbolic maps $F_{p}^{\mathscr{A}}: D_{p}(\mathscr{A}) \rightarrow Y^{e} \times \mathscr{V}$ of type $\mathscr{A}$ such that $H_{p} \circ F_{p}^{\mathscr{A}}:=f_{p}^{k} \circ H_{p} \mid D_{p}(\mathscr{A})$, and the set $H_{p}\left(D_{p}(\mathscr{A})\right)$ contains $K_{p}$ and is included in a small neighborhood of $\breve{H}_{p}\left(\breve{D}_{p}(\mathscr{A})\right)$,
(2) if $\lim \sup _{n \rightarrow \infty}\left\|\operatorname{det} D f_{0}^{n}\left|K_{0}\left\|^{1 / \epsilon} \cdot\right\| D f_{0}^{n}\right| K_{0}\right\|<1$, then $\left(F_{p}^{s A}\right)_{p}$ is moderately dissipative,
(3) there is a $C^{1}$-family $\left(\pi_{p}\right)_{p \in \mathcal{P}}$ of maps $\pi_{p}: \mathbb{R} \times I \times \mathscr{V} \rightarrow \mathbb{R} \times \mathscr{V}$ which satisfies items (0)-(1)-(2)-(3) of Definition 1.23 and such that $\left(\operatorname{ker} D \pi_{p}\right)_{p}$ is of class $C^{1}$.

A simple consequence of the proof will be:
Corollary 2.26. If under the assumptions of Proposition 2.25, the family $\left(f_{p}\right)_{p \in \hat{\mathcal{P}}}$ is of class $C^{2+}$, then the same conclusion holds true and moreover $\left(\pi_{p}\right)_{p}$ is of class $C^{1+}$.

Step 3: Construction of the folding transformations. Let us come back to the proof of Theorem 2.10. Let $p_{1} \in \hat{\mathcal{P}}$ and let $K \supset \Lambda_{1} \sqcup \Lambda_{2}$ be the basic set of $f_{p_{1}}$ given by Lemma 2.22. Let $\mathcal{P}_{1}$ be a neighborhood of $p_{1}$ for which this hyperbolic set persists as $\left(K_{p}\right)_{p \in \mathcal{P}_{1}}$. Let $\left(T_{i}\right)_{1 \leq i \leq d}$ be the $d$ tangencies points between the local unstable and stable manifolds of $K$ satisfying the no-ending gap condition (2) of Theorem 2.15. Also it holds $\tau^{u}\left(\Lambda_{1}\right) \cdot \tau^{s}\left(\Lambda_{2}\right)>1$. Thus by Theorem 2.15, we have:

Fact 2.27. For every $p$ close to $p_{1}$, there are $d$ tangencies points between the local unstable and stable manifolds of $K_{p}$ which are close to $\left(T_{i}\right)_{1 \leq i \leq d}$.

Let $\left(F_{p}^{\mathscr{A}}\right)_{p \in \mathcal{P}_{1}}$ be the family of hyperbolic maps of type $\mathscr{A}$ given by Proposition 2.25 for the continuation $\left(K_{p}\right)_{p \in \mathcal{P}_{1}}$. Let $k \geq 1$ and let $\left(H_{p}\right)_{p}$ be the family of $C^{r}$-embeddings such that:

$$
f_{p}^{k} \circ H_{p}=H_{p} \circ F_{p}^{\mathscr{A}}
$$

Let $\Lambda_{p}$ be the maximal invariant set of $F_{p}^{d}$. Let $\left(\pi_{p}\right)_{p}$ be the family of projections given by Proposition 2.25.(3). We put also:

$$
W^{s}\left(\Lambda_{p}\right):=\bigcup_{s \in \vec{A}} W_{p}^{s} \times\{a(s)\} \quad \text { and } \quad W^{u}\left(\Lambda_{p}\right):=\bigcup_{u \in \overleftarrow{A}} W_{p}^{u} \times\{t(u)\}
$$

Let us construct the folding transformations. As each $T_{i}$ is in the unstable lamination of $K_{p_{1}}$, up to changing $T_{i}$ by $f_{p_{1}}^{-m}\left(T_{i}\right)$ for every $i$, we can assume that each $T_{i}$ belongs to $H_{p_{1}}\left(W_{p_{1}}^{u}\left(\Lambda_{p_{1}}\right) \cap\right.$ int $\left.Y^{e} \times \mathscr{V}\right)$. Let $u_{i} \in \mathscr{A}$ be such that $H_{p_{1}}^{-1}\left(T_{i}\right)$ belongs to $W_{p_{1}}^{u_{i}} \times\left\{t\left(u_{i}\right)\right\}$. By Remark 2.24, we have $u_{i} \neq u_{j}$ for $i \neq j$. Put $W_{l o c}^{u}\left(T_{i}, f_{p}\right):=H_{p}\left(W_{p}^{u_{i}} \times\left\{t\left(u_{i}\right)\right\}\right)$.

Also each $T_{i}$ is in the stable lamination of $K_{p_{1}}$. Thus for every $n \geq 1$ large enough, each point $f_{p_{1}}^{n}\left(T_{i}\right)$ is in $H_{p_{1}}\left(W_{p_{1}}^{s}\left(\Lambda_{p_{1}}\right) \cap \operatorname{int} Y^{e} \times \mathscr{V}\right)$. Moreover there exists $s_{i} \in \overrightarrow{\mathscr{A}}$ such that $f_{p_{1}}^{n}\left(W_{l o c}^{u}\left(T_{i}, f_{p_{1}}\right)\right)$ has a quadratic tangency with $H_{p_{1}}\left(W_{p_{1}}^{\jmath_{i}} \times \vartheta\left(s_{i}\right)\right)$ at $f_{p_{1}}^{n}\left(T_{i}\right)$. Also the unfolding of these $d$ tangencies induced by $\left(f_{p}\right)_{p}$ is non degenerated by Lemma 2.22 (2). By regularity of $\pi$, there is a small neighborhood of $p_{1}$ formed by parameters $p$ for which the curve $W_{l o c}^{u}\left(T_{i}, f_{p}\right)$ has a quadratic tangency with a fiber of $\pi_{p}^{o\left(s_{i}\right)} \circ\left(H_{p}^{o\left(s_{i}\right)}\right)^{-1} \circ f_{p}^{n}$ at a point $T_{i}(p)$ nearby $T_{i}$. Let $\zeta_{p}^{u_{i}} \in W_{p}^{u_{i}}$ be the preimage by $H_{p}^{t\left(u_{i}\right)}$ of $T_{i}(p)$. As the unfolding is non degenerate, the following map has invertible differential at $p=p_{1}$ :

$$
p \mapsto\left(\pi_{p}^{o\left(s_{i}\right)} \circ\left(H_{p}^{\sigma\left(s_{i}\right)}\right)^{-1} \circ f_{p}^{n}\left(T_{i}(p)\right)-b_{p}^{s_{i}}\right)_{1 \leq i \leq d} \quad \text { with }\left\{b_{p}^{s_{i}}\right\}:=\pi_{p}^{o\left(s_{i}\right)}\left(W_{p}^{s_{i}}\right)
$$

By regularity of the stable laminations, there is a clopen neighborhood $W_{i} \subset \overrightarrow{\mathscr{A}}$ of each $s_{i}$, such that for any $\left(\tilde{\jmath}_{i}\right)_{1 \leq i \leq d} \in \prod_{1 \leq i \leq d} W_{i}$, the following map has invertible differential at $p=p_{1}$ :

$$
p \mapsto\left(\pi_{p}^{\theta\left(s_{i}\right)} \circ\left(H_{p}^{\theta\left(s_{i}\right)}\right)^{-1} \circ f_{p}^{n}\left(T_{i}(p)\right)-b_{p}^{\tilde{s}_{i}}\right)_{1 \leq i \leq d} \quad \text { with }\left\{b_{p}^{\tilde{\tilde{s}}_{i}}\right\}:=\pi_{p}^{\theta \cdot\left(\tilde{j}_{i}\right)}\left(W_{p}^{\tilde{s}_{i}}\right)
$$

Up to reducing each $W_{i}$, we can assume it of the form $W_{i}=c_{i} \cdot \overrightarrow{\mathscr{A}}$ for admissible words $c_{i} \in \mathscr{A}^{*}$ of the same length $|c|$. Thus up to replacing $n$ by $n+|c|$, we can assume that $c_{i}=e$ and obtain:

Fact 2.28. The following maps have invertible differential at $p=p_{1}$ :

$$
p \mapsto\left(\pi_{p}^{\circ\left(s_{i}\right)} \circ\left(H_{p}^{a\left(s_{i}\right)}\right)^{-1} \circ f_{p}^{n}\left(T_{i}(p)\right)-b_{p}^{\tilde{s}_{i}}\right)_{1 \leq i \leq d}, \quad \forall\left(\tilde{\jmath}_{i}\right)_{1 \leq i \leq d} \in \prod_{i=1}^{d}\left\{\jmath \in \overrightarrow{\mathscr{A}}: 0(\jmath)=\sigma\left(s_{i}\right)\right\}
$$

We recall that $k$ is a large integer such that $f_{p}^{k} \circ H_{p}=H_{p} \circ F_{p}^{\mathscr{A}}$ by Proposition 2.25. As $n$ is any large number, we can assume that $n=m \cdot k$ for a certain $m \geq 1$. By replacing $\mathscr{A}$ by $\mathscr{A}^{*} \cap \mathscr{A}^{m}$ and $F_{p}^{s l}$ by $\left(F_{p}^{s l}\right)^{m}$, we can assume that $k=n$.

For every $1 \leq i \leq d$, let $\square_{i}$ be a new arrow from $t\left(u_{i}\right)$ to $\sigma\left(s_{i}\right)$. Let $\mathscr{C}:=\left\{\square_{i}: 1 \leq i \leq d\right\}$ and:

$$
\overrightarrow{\mathscr{A}}_{Ð_{i}}:=\left\{s \in \overrightarrow{\mathscr{A}}: a(\jmath)=t\left(\square_{i}\right)\right\} \quad \text { and } \quad F_{p}^{\oplus_{i}}:=\left(H_{p}^{t\left(Ð_{i}\right)}\right)^{-1} \circ f_{p}^{k} \circ H_{p}^{o\left(Ð_{i}\right)}
$$

The latter map is defined on a neighborhood of $\zeta_{p_{1}}^{u_{i}}$. The following is a rephrasing of Fact 2.28:
Fact 2.29. For any $\left(\tilde{\jmath}_{i}\right)_{1 \leq i \leq d} \in \prod_{i=1}^{d} \overrightarrow{\mathscr{A}}_{\square_{i}}$, the following map has invertible differential at $p=p_{1}$ :

$$
p \mapsto\left(\pi_{p}^{t\left(Ð_{i}\right)} \circ F_{p}^{\oplus_{i}}\left(\zeta_{p}^{u_{i}}\right)-b_{p}^{\tilde{s}_{i}}\right)_{1 \leq i \leq d}
$$

As $\zeta_{p}^{u_{i}}$ is a point of quadratic tangency between $W_{p}^{u_{i}}$ and a fiber of $\pi_{p}^{t\left(\square_{i}\right)} \circ F_{p}^{\oplus_{i}}$, there is a small clopen neighborhood $\overleftarrow{\mathscr{A}}_{\square_{i}}$ of $u_{i}$ such that for $u \in \overleftarrow{\mathscr{A}}_{\square_{i}}$ and $p$ close to $p_{1}$, there is a quadratic tangency between $W_{p}^{u}$ and a fiber of $\pi_{p}^{t\left(Ð_{i}\right)} \circ F_{p}^{Ð_{i}}$ at a point $\zeta_{p}^{u}$ close to $\zeta_{p}^{u_{i}}$. By regularity of the unstable lamination and compactness of $\prod_{i=1}^{d} \overrightarrow{\mathscr{A}}_{\square_{i}}$, by reducing each $\overleftarrow{\mathscr{A}}_{\square_{i}}$, we obtain:
Fact 2.30. For every $\left(\tilde{u}_{i}, \tilde{\jmath}_{i}\right)_{1 \leq i \leq d} \in \prod_{i=1}^{d} \overleftarrow{\mathscr{A}}_{\square_{i}} \times \overrightarrow{\mathscr{A}}_{Ð_{i}}$, the following map has invertible differential for every $p$ close to $p_{1}$ :

$$
p \mapsto\left(\pi_{p}^{t\left(\oplus_{i}\right)} \circ F_{p}^{\oplus_{i}}\left(\zeta_{p}^{\tilde{u}_{i}}\right)-b_{p}^{\tilde{s}_{i}}\right)_{1 \leq i \leq d} .
$$

Up to reducing each $\overleftarrow{\mathscr{A}}_{\varpi_{i}}$, there is a segment $I^{\boxtimes_{i}}$ (independent of $p$ close to $p_{1}$ ) such that for every $\tilde{u}_{i} \in \overleftarrow{\mathscr{A}}_{\unlhd_{i}}$ and $p$ close to $p_{1}$, the point $\zeta_{p}^{\tilde{u}_{i}}$ is the unique critical point of $\pi_{p}^{t\left(Ð_{i}\right)} \circ F_{p}^{\square_{i}} \mid I^{\square_{i}} \times I \cap W_{p}^{\tilde{u}_{i}}$ and belongs to $\operatorname{int} I^{®_{i}} \times I$. Moreover we can assume that the convex hull of $I^{\square_{i}} \times I \cap \bigcup_{\overleftarrow{\mathscr{A}}_{\unlhd_{i}}} W_{p}^{\tilde{u}_{i}}$ is sent into int $Y^{e}$ by $F_{p}^{®_{i}}$.

By Proposition 2.25.(0), we can assume that the unstable manifolds $\left(W_{p}^{u}\right)_{p}$ are arbitrarily close to be horizontal, without changing any structure. Thus for every $i$, we can find a segment $J^{\square_{i}} \Subset I$ (independent of $p$ close to $p_{1}$ ) such that, with $Y_{p}^{\square_{i}}=I^{\square_{i}} \times J^{\square_{i}}$ and $\partial^{u} Y_{p}^{\square_{i}}=I^{\square_{i}} \times \partial J^{\square_{i}}$, it holds:

$$
\overleftarrow{\mathscr{A}}_{\square_{i}}=\left\{u \in \overleftarrow{\mathscr{A}}: t(u)=\sigma\left(\square_{i}\right), W_{p}^{u} \cap Y_{p}^{\square_{i}} \neq \emptyset \& W_{p}^{u} \cap \partial^{u} Y_{p}^{\square_{i}}=\emptyset\right\}
$$

for any $p$ close to $p_{1}$. Up to reducing each $\overleftarrow{\mathscr{A}}_{\unlhd_{i}}$ a new time and consequently $J^{®_{i}}$ and then reducing $I^{\unlhd_{i}}$, we can suppose that $F_{p}^{@_{i}}$ sends $Y_{p}^{\rrbracket_{i}}$ into the interior of $Y^{e}$.

By Fact 2.27 , for every $p$ close to $p_{1}$, there are $\left(u_{i}^{\prime}, s_{i}^{\prime}\right)_{1 \leq i \leq d} \in \prod_{i=1}^{d} \overleftarrow{\mathscr{A}}_{\oplus_{i}} \times \overrightarrow{\mathscr{A}}_{\oplus_{i}}$ close to $\left(u_{i}, s_{i}\right)_{1 \leq i \leq d}$ such that each $W_{p}^{u_{i}^{\prime}}$ has a quadratic tangency with the preimage of $W_{p}^{s_{i}^{\prime}}$ by $F_{p}^{@_{i}}$ at a point close to $\zeta_{p}^{u_{i}}$. By uniqueness such a point must be $\zeta_{p}^{u_{i}^{\prime}}$. This proves $\left(\mathbf{H}_{\mathbf{1}}\right)$.

Property $\left(\mathbf{H}_{\mathbf{2}}\right)$ is given by Fact 2.30 and the local inversion theorem together with the compactness of $\prod_{i=1}^{d} \overleftarrow{\mathscr{A}}_{\varpi_{i}} \times \overrightarrow{\mathscr{A}}_{\square_{i}}$ to define the regular compact neighborhood $\mathcal{P}$ of $p_{1}$. Clearly, the

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families $\left(F^{\mathscr{A}}\right)_{p \in \mathcal{P}}$ and $\left(F^{\mathscr{C}}\right)_{p \in \mathcal{P}}$ define a regular family of systems of type $(\mathscr{A}, \mathscr{C})$. The proof of Theorem 2.10 is complete.

## 3. Sufficient conditions for a wandering stable domain

In this section, we consider a system of type ( $\mathscr{A}, \mathscr{C}$ ) endowed with an adapted projection $\pi$ which satisfies the assumptions of Theorems 1.32 and 1.34 and we prove the conclusions of these Theorems. Before performing these proofs, we recall and develop the formalism of implicit representations that will be used throughout this paper.
3.1. Implicit representations and initial bounds. In the following, we will use the formalism of implicit representations. It was introduced by Palis and Yoccoz in [PY01] and [PY09] in order to get bounds on the iterations of hyperbolic transformations ${ }^{10}$. In fact, this formalism goes back to the generating functions and the Shilnikov cross coordinates, appeared in the work of Shilnikov [Shi67] and his school [GST08].

Proposition 3.1. Let $(Y, F)$ be a hyperbolic transformation. Then, there are functions $\mathcal{X}_{0}$ and $\mathcal{Y}_{1}$ in $C^{2}\left(I^{2}, I\right)$ satisfying for every $\left(x_{0}, y_{0}\right) \in Y$ and $\left(x_{1}, y_{1}\right) \in F(Y)$ :

$$
F\left(x_{0}, y_{0}\right)=\left(x_{1}, y_{1}\right) \Leftrightarrow\left\{\begin{array}{l}
x_{0}=\mathcal{X}_{0}\left(x_{1}, y_{0}\right)  \tag{A}\\
y_{1}=\mathcal{Y}_{1}\left(x_{1}, y_{0}\right)
\end{array}\right.
$$

Proof. Let $\left(x_{0}, y_{0}\right) \in Y$ and $\left(x_{1}, y_{1}\right) \in F(Y)$. By properties (1) and (2) of Definition 1.4, the image by $F$ of $\left\{y=y_{0}\right\} \cap Y$ is the graph of a function $x \in I \mapsto \mathcal{Y}_{1}\left(x, y_{0}\right) \in I$. Likewise, by properties (1) and (3) of Definition 1.4, the image by $F^{-1}$ of $\left\{x=x_{1}\right\} \cap F(Y)$ is the transpose of the graph of a function $y \in I \mapsto \mathcal{X}_{0}\left(x_{1}, y\right) \in I$. We notice that the functions $\left(x_{1}, y_{0}\right) \mapsto \mathcal{X}_{0}\left(x_{1}, y_{0}\right)$ and $\left(x_{1}, y_{0}\right) \mapsto \mathcal{Y}_{1}\left(x_{1}, y_{0}\right)$ are of class $C^{2}$. Observe also that by construction, $F\left(x_{0}, y_{0}\right)=\left(x_{1}, y_{1}\right)$ iff the intersection point $\left\{\left(x_{1}, \mathcal{Y}_{1}\left(x_{1}, y_{0}\right)\right)\right\}=\operatorname{Graph} \mathcal{Y}_{1}\left(\cdot, y_{0}\right) \cap\left\{x=x_{1}\right\}$ is the image by $F$ of $\left\{\left(\mathcal{X}_{0}\left(x_{1}, y_{0}\right), y_{0}\right)\right\}={ }^{t} \operatorname{Graph} \mathcal{X}_{0}\left(x_{1}, \cdot\right) \cap\left\{y=y_{0}\right\}$.

Definition 3.2. The pair $\left(\mathcal{X}_{0}, \mathcal{Y}_{1}\right)$ is the implicit representation of $(Y, F)$.
Let us introduce the natural counterpart of implicit representation in the complex setting:
Proposition 3.3. Let $(\tilde{Y}, F)$ be a $C_{\rho}^{\omega}$-hyperbolic transformation. Then there are holomorphic functions $\mathcal{Z}_{0}$ and $\mathcal{W}_{1}$ defined on $\tilde{I}^{2}$ satisfying for every $\left(z_{0}, w_{0}\right) \in \tilde{Y}$ and $\left(z_{1}, w_{1}\right) \in F(\tilde{Y})$ :

$$
F\left(z_{0}, w_{0}\right)=\left(z_{1}, w_{1}\right) \Leftrightarrow\left\{\begin{array}{rr}
z_{0}= & \mathcal{Z}_{0}\left(z_{1}, w_{0}\right) \\
w_{1}= & \mathcal{W}_{1}\left(z_{1}, w_{0}\right)
\end{array}\right.
$$

Proof. Let $\left(z_{0}, w_{0}\right) \in \tilde{Y}$ and $\left(z_{1}, w_{1}\right) \in F(\tilde{Y})$. By properties (1) and (2) of Definition 1.5, the image by $F$ of $\left\{w=w_{0}\right\} \cap \tilde{Y}$ is the graph of a function $z \in \tilde{I} \mapsto \mathcal{W}_{1}\left(z, w_{0}\right) \in \tilde{I}$. Likewise, by properties (1) and (3) of Definition 1.5, the image by $F^{-1}$ of $\left\{z=z_{1}\right\} \cap F(\tilde{Y})$ is the transpose of the graph of a function $w \in \tilde{I} \mapsto \mathcal{Z}_{0}\left(z_{1}, w\right)$. By transversality, the functions $\left(z_{1}, w_{0}\right) \mapsto \mathcal{Z}_{0}\left(z_{1}, w_{0}\right) \in \tilde{I}$ and $\left(z_{1}, w_{0}\right) \mapsto \mathcal{W}_{1}\left(z_{1}, w_{0}\right)$ are holomorphic. Observe also that by construction, $F\left(z_{0}, w_{0}\right)=\left(z_{1}, w_{1}\right)$ iff the intersection point $\left\{\left(z_{1}, \mathcal{W}_{1}\left(z_{1}, w_{0}\right)\right)\right\}=\operatorname{Graph} \mathcal{W}_{1}\left(\cdot, w_{0}\right) \cap\left\{z=z_{1}\right\}$ is the image by $F$ of $\left\{\left(\mathcal{Z}_{0}\left(z_{1}, w_{0}\right), w_{0}\right)\right\}={ }^{t} \operatorname{Graph} \mathcal{Z}_{0}\left(z_{1}, \cdot\right) \cap\left\{w=w_{0}\right\}$.
10. They call such maps affine-like maps.

Definition 3.4. The pair $\left(\mathcal{Z}_{0}, \mathcal{W}_{1}\right)$ is the implicit representation of $(\tilde{Y}, F)$.
Remark 3.5. If $(\tilde{Y}, F)$ is a complex extension of a hyperbolic transformation $(Y, F)$, with implicit representations $\left(\mathcal{Z}_{0}, \mathcal{W}_{1}\right)$ and $\left(\mathcal{X}_{0}, \mathcal{Y}_{1}\right)$ respectively, it holds: $\left(\mathcal{Z}_{0}, \mathcal{W}_{1}\right) \mid I^{2}=\left(\mathcal{X}_{0}, \mathcal{Y}_{1}\right)$.

Proposition 3.6. If $\left(\mathcal{X}_{0}, \mathcal{Y}_{1}\right)$ is the implicit representation of a hyperbolic transformation $(Y, F)$, for every $\left(x_{0}, y_{0}\right) \in Y$ sent to $\left(x_{1}, y_{1}\right)$ by $F$, it holds:

$$
\operatorname{det} D F\left(x_{0}, y_{0}\right) \cdot \partial_{x_{1}} \mathcal{X}_{0}\left(x_{1}, y_{0}\right)=\partial_{y_{0}} \mathcal{Y}_{1}\left(x_{1}, y_{0}\right)
$$

Proof. By $(\mathcal{A})$, we have $F\left(\mathcal{X}_{0}\left(x_{1}, y_{0}\right), y_{0}\right)=\left(x_{1}, \mathcal{Y}_{1}\left(x_{1}, y_{0}\right)\right)$ and so

$$
\operatorname{det} D_{x_{0}, y_{0}} F \cdot \operatorname{det} D_{x_{1}, y_{0}}\left(\mathcal{X}_{0}\left(x_{1}, y_{0}\right), y_{0}\right)=\operatorname{det} D_{x_{1}, y_{0}}\left(x_{1}, \mathcal{Y}_{1}\left(x_{1}, y_{0}\right)\right)
$$

Then the proposition follows from the following equalities:

$$
\operatorname{det} D_{x_{1}, y_{0}}\left(\mathcal{X}_{0}\left(x_{1}, y_{0}\right), y_{0}\right)=\partial_{x_{1}} \mathcal{X}_{0}\left(x_{1}, y_{0}\right) \quad \text { and } \quad \operatorname{det} D_{x_{1}, y_{0}}\left(x_{1}, \mathcal{Y}_{1}\left(x_{1}, y_{0}\right)\right)=\partial_{y_{0}} \mathcal{Y}_{1}\left(x_{1}, y_{0}\right)
$$

An immediate consequence of the latter proof is:
Corollary 3.7. If $(\tilde{Y}, F)$ is a $C_{\rho}^{\omega}$-hyperbolic transformation with implicit representation $\left(\mathcal{Z}_{0}, \mathcal{W}_{1}\right)$, then for every $\left(z_{0}, w_{0}\right) \in \tilde{Y}$ sent to $\left(z_{1}, w_{1}\right)$ by $F$, it holds:

$$
\operatorname{det} D F\left(z_{0}, w_{0}\right) \cdot \partial_{z_{1}} \mathcal{Z}_{0}\left(z_{1}, w_{0}\right)=\partial_{w_{0}} \mathcal{W}_{1}\left(z_{1}, w_{0}\right)
$$

The following distortion was introduced in [PY09, p. 19].
Definition 3.8. The distortion $\mathbf{D}(Y, F)$ of a hyperbolic transformation $(Y, F)$ with implicit representation $(\mathcal{X}, \mathcal{Y})$ is the maximum absolute value attained on $I^{2}$ by the six following functions:

$$
\partial_{x} \log \left|\partial_{x} \mathcal{X}\right|, \quad \partial_{y} \log \left|\partial_{x} \mathcal{X}\right|, \quad \partial_{x} \log \left|\partial_{y} \mathcal{Y}\right|, \quad \partial_{y} \log \left|\partial_{y} \mathcal{Y}\right|, \quad \partial_{y^{2}}^{2} \mathcal{X}, \quad \partial_{x^{2}}^{2} \mathcal{Y} .
$$

Let us recall:
Theorem 3.9 ([PY09, (MP6), p. 21]). There exists $C_{1}>0$ which depends only on $\theta$ and $\lambda$ such that for every $N \geq 1$, if $\left(\left(Y_{i}, F_{i}\right)\right)_{1 \leq i \leq N}$ are hyperbolic transformations, then the distortion of $(Y, F)=\left(Y_{1}, F_{1}\right) \star \cdots \star\left(Y_{N}, F_{N}\right)$ satisfies:

$$
\mathbf{D}(Y, f) \leq C_{1} \cdot \max _{1 \leq i \leq N} \mathbf{D}\left(Y_{i}, F_{i}\right)
$$

Corollary 3.10. If $F^{\mathscr{A}}$ is a hyperbolic map of type $\mathscr{A}$, there exists $\mathbf{D}>\mathbf{0}$ and $\mathbf{B}:=2 \exp (2 \sqrt{2} \mathbf{D})$ such that for any $c \in \mathscr{A}^{*}$ and $(x, y) \in I^{2}$, it holds both $\mathbf{D}\left(Y^{c}, F^{c}\right) \leq \mathbf{D}$ and the following:

$$
\left|Y^{c}\right| \cdot \mathbf{B}^{-1} \leq\left|\partial_{x} \mathcal{X}^{c}(x, y)\right| \leq\left|Y^{c}\right| \cdot \mathbf{B} \quad \text { and } \quad\left|F^{c}\left(Y^{c}\right)\right| \cdot B^{-1} \leq\left|\partial_{y} \mathcal{Y}^{c}(x, y)\right| \leq\left|F^{c}\left(Y^{c}\right)\right| \cdot \mathbf{B}
$$

Proof. We fix $c \in \mathscr{A}^{*}$ and $(x, y) \in I^{2}$. A consequence of Theorem 3.9 and of the finiteness of $\mathscr{A}$ is:

$$
\begin{align*}
& \left|\partial_{x} \mathcal{X}^{c}(x, y)\right| \cdot(\mathbf{B} / 2)^{-1} \leq\left|\partial_{x} \mathcal{X}^{c}\left(x^{\prime}, y^{\prime}\right)\right| \leq\left|\partial_{x} \mathcal{X}^{c}(x, y)\right| \cdot(\mathbf{B} / 2) \quad \forall\left(x^{\prime}, y^{\prime}\right) \in I^{2}  \tag{3.1}\\
& \text { with } \mathbf{D}:=C_{1} \cdot \max _{a \in \mathscr{A}} \mathbf{D}\left(Y^{a}, F^{a}\right) \quad \text { and } \quad \mathbf{B}:=2 \exp \left(\operatorname{diam}\left(I^{2}\right) \mathbf{D}\right)=2 \exp (2 \sqrt{2} \mathbf{D})
\end{align*}
$$

By Definition 1.27, the number $\left|Y^{c}\right|$ is the maximum of the lengths of $I \times\left\{y^{\prime}\right\} \cap Y^{c}$ among $y^{\prime} \in I$. Thus there exists $y^{\prime} \in I$ such that $\left|Y^{c}\right|=\left|\int_{I} \partial_{x} \mathcal{X}^{c}\left(x^{\prime}, y^{\prime}\right) d x^{\prime}\right|$. Then Eq. (3.1) together with $\operatorname{diam}(I)=2$ imply the first bound. The proof of the second one is the same by symmetry.
3.2. Normal form and definitions of the sets $B_{j}$ and $\tilde{B}_{j}$. In this subsection, we consider a system $F$ of type $(\mathscr{A}, \mathscr{C})$ endowed with an adapted projection $\pi$. Our aim is to prove a normal form useful for the proofs of Theorems 1.32 and 1.34. We start with general initial bounds.

We recall that for $c \in \mathscr{A}^{*}$, the curve $H^{c}$ was defined in Definition 1.26 , and for $\square \in \mathscr{C}$, the set $\mathscr{A}_{\square}^{*}$ and the critical points $\left(\zeta^{c}\right)_{c \in \mathscr{A}_{\square}^{*}}$ were defined in Proposition 1.28.
Proposition 3.11 (definition of $\epsilon_{\mathscr{C}}$ ). There exists $\epsilon_{\mathscr{C}}>0$ such that for every $\square \in \mathscr{C}$ :
(1) $F^{\unlhd}\left(Y^{\square}\right) \subset Y^{e}$ is $2 \epsilon_{\square}$-distant from $\partial^{u} Y^{e}$,
(2) $\zeta^{c} \in Y^{\square}$ is $2 \epsilon_{\square}$-distant to $\partial Y^{\square}$, for every $c \in \mathscr{A}_{\square}^{*}$,
(3) $H^{c}$ is $2 \epsilon_{\mathscr{G}}$-distant to $\partial^{u} Y^{\square}$, for every $c \in \mathscr{A}_{\square}^{*}$.

Proof. By finiteness of $\mathscr{C}$, it suffices to show that there exists such a constant for every $\square \in \mathscr{C}$.
By Definition 1.18, we have $F^{\square}\left(Y^{\square}\right) \subset Y^{e} \backslash \partial^{u} Y^{e}$ and so the first statement is easy.
By Remark 1.29, the set $\left\{\zeta^{c}: c \in \mathscr{A}_{\square}^{*} \sqcup \overleftarrow{A}_{\square}\right\}$ is compact. It is inclucled in the interior of $Y^{\boxminus}$ by Definition 1.23 (4) and Proposition 1.28. This proves the second item.

By Remark 1.29, the set $\bigcup_{u \in \overleftarrow{\mathscr{S}}_{\square}} W^{u} \cup \bigcup_{c \in \mathbb{P}_{\square}^{*}} H^{c}$ is compact. By Definition 1.18 (1) and Proposition 1.28, it does not intersect $\partial^{u} Y^{\square}$. Thus the last item of the proposition follows.

From now on, we work under the assumptions of Theorem 1.32 (resp. Theorem 1.34), which defined the sequences $\left(c_{j}\right)$ (and so $\left.\left(\zeta^{\complement_{j}}\right)_{j}\right)$ and $\left(\square_{j}\right)_{j}$. The idea is to define the domains $B_{j}$ and $\tilde{B}_{j}$ using a normal form for the folding of the curve $H^{c_{j}}$ with respect to the $\pi$-fibers. However it is rather complicated to work with $\pi^{0\left(e_{j+1}\right)}$-fibers, for they are a priori not analytic. Hence we will work with:

Definition 3.12. Let $\breve{V}^{j+1}$ be the pullback of $\left\{x_{j+1}\right\} \times I$ by $F^{c_{j+1}}$, with $x_{j+1}$ the first coordinate of $\zeta^{c_{j+1}}$.

We note that:

$$
\begin{equation*}
\breve{V}^{j+1}=\left\{\left(\breve{u}_{j+1}(y), y\right): y \in I\right\} \quad \text { with } \quad \breve{u}_{j+1}(\cdot):=\mathcal{X}^{c_{j+1}}\left(x_{j+1}, \cdot\right), \quad \forall j \geq 1 . \tag{3.2}
\end{equation*}
$$

We will also use the following notation $h_{j}$ whose expression is similar to the one of $\breve{u}_{j+1}$ :

$$
\begin{equation*}
H^{c_{j}}=\left\{\left(x, \hbar_{j}(x)\right): x \in I\right\} \quad \text { with } \quad h_{j}(\cdot):=\mathcal{Y}^{c_{j}}(\cdot, 0), \quad \forall j \geq 1 \tag{3.3}
\end{equation*}
$$

Let $F_{1}^{\square_{j}}$ and $F_{2}^{\square_{j}}$ be the $x$-coordinate and the $y$-coordinate of $F^{\unlhd_{j}}$. The $x$-projection $I^{\square_{j}} \subset I$ of $Y^{\square_{j}}$ is the definition domain of $x \mapsto F^{\unlhd_{j}} \circ\left(x, \hbar_{j}(x)\right)$. For every $r>0$, let $\mathbb{D}(r):=\{z \in \mathbb{C}:|z|<r\}$.

Here is the aforementioned normal form:
Proposition 3.13. There exists $K>0$ such that for every $j$ large enough, there exist $\breve{x}_{j} \in I^{\square_{j}}$ which is $\epsilon_{\mathscr{G}}$-distant to the boundary of $I^{\square_{j}}$, real numbers $b_{j}, q_{j}$ and a $C^{2}$-function $r_{j}$ such that for every $(x, y) \in\left[-\epsilon_{\mathscr{G}}, \epsilon_{\mathscr{G}}\right]^{2}$ the following is well-defined and holds true:

$$
\begin{equation*}
\left(F_{1}^{\square_{j}}-\breve{u}_{j+1} \circ F_{2}^{\square_{j}}\right) \circ \varphi_{j}(x, y)=q_{j} \cdot x^{2}+b_{j} \cdot y+r_{j}(x, y) \tag{3.4}
\end{equation*}
$$

with $\varphi_{j}(x, y)=\left(\breve{x}_{j}+x, \hbar_{j}\left(\breve{x}_{j}+x\right)+y\right), 1 / K \leq\left|q_{j}\right|,\left|b_{j}\right| \leq K$ and $r_{j}$ satisfying:

$$
\begin{equation*}
0=\partial_{x} r_{j}(0)=\partial_{y} r_{j}(0)=\partial_{x}^{2} r_{j}(0) \text { and } r_{j}(0)=o\left(\breve{\gamma}_{j}^{2}\right) . \tag{3.5}
\end{equation*}
$$

Also if $F$ is analytic, there exists $\tilde{\epsilon}_{\mathscr{C}}>0$ independent of $j$, such that $\varphi_{j}$ and $r_{j}$ extend holomorphicaly to $\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}$.

We prove this proposition below. This normal form enables us to define rescaling coordinates and the domains $B_{j}$ and $\tilde{B}_{j}$. We need just to define the rescaling factors:

Definition 3.14. Let $y_{j}$ be the $y$-coordinate of $F^{\unlhd_{j}}\left(\zeta^{c_{j}}\right)$. We define:

$$
\sigma_{j}:=\partial_{x} \mathcal{X}^{c_{j}}\left(x_{j}, y_{j}\right) \quad \text { and } \quad \gamma_{j}:=\left|\frac{\sigma_{j+1}}{q_{j}}\right|^{1 / 2} \cdot\left|\frac{\sigma_{j+2}}{q_{j+1}}\right|^{1 / 2^{2}} \cdots\left|\frac{\sigma_{j+k}}{q_{j+k-1}}\right|^{1 / 2^{k}} \cdots
$$

By the latter proposition, $1 / K \leq\left|q_{j}\right| \leq K$ and so $\prod_{k}\left|1 / q_{k}\right|^{2^{k}}$ is well-defined and in $[1 / K, K]$. Also by Corollary 3.10, it holds $\left|\sigma_{j}\right| / \mathbf{B} \leq\left|Y^{c_{j}}\right| \leq \mathbf{B}\left|\sigma_{j}\right|$. Thus the limit $\gamma_{j}$ exists and satisfies:

$$
\begin{equation*}
\frac{\breve{\gamma}_{j}}{\mathbf{B} K} \leq \gamma_{j} \leq \mathbf{B} K \breve{\gamma}_{j} \tag{3.6}
\end{equation*}
$$

Therefore $\gamma_{j}$ is well-defined, small and nonzero because $\breve{\gamma}_{j}$ is so by assumptions (i) and (iii) of Theorem 1.32. In particular, since the domain of $\varphi_{j}$ contains $\left[-\epsilon_{\mathscr{C}}, \epsilon_{\mathscr{C}}\right]^{2}$, the following is welldefined:

Definition 3.15. The rescaling coordinates $\Phi_{j}$ are defined for $j$ large by:
$\Phi_{j}:(X, Y) \in(-0.3,0.3)^{2} \mapsto \varphi_{j}\left(\gamma_{j} \cdot X, \frac{q_{j} \cdot \gamma_{j}^{2}}{2 b_{j}} \cdot Y\right)=\left(\breve{x}_{j}+\gamma_{j} \cdot X, h_{j}\left(\breve{x}_{j}+\gamma_{j} \cdot X\right)+\frac{q_{j} \cdot \gamma_{j}^{2}}{2 b_{j}} \cdot Y\right)$.
Similarly, if $F$ is analytic, the definition domain of $\varphi_{j}$ contains $\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}$. So it holds:
Fact 3.16. If $F$ is analytic, the rescaling coordinates $\Phi_{j}$ extend holomorphically to $\mathbb{D}(0.3)^{2}$.
Definition 3.17. The real image of $\Phi_{j}$ is denoted by $B_{j}:=\Phi_{j}\left((-0.3,0.3)^{2}\right)$. If $F$ is analytic, the complex image of $\Phi_{j}$ is denoted by $\tilde{B}_{j}:=\Phi_{j}\left(\mathbb{D}(0.3)^{2}\right)$.

Remark 3.18. Actually for $j$ large enough, the map $\varphi_{j}$ is also well defined on $[-3,3]^{2}$. Using [Ber18, Thm C and Rem. 3.5.], one can show that when $j$ is large, the following map:

$$
F_{j}:=(X, Y) \in(-0.3,0.3)^{2} \mapsto\left(\Phi_{j+1}\right)^{-1} \circ F^{c_{j+1}} \circ F^{\unlhd_{j}} \circ \Phi_{j}(X, Y),
$$

is $C^{2}$-close to $(X, Y) \in(-0.3,0.3)^{2} \mapsto \operatorname{sgn}\left(\sigma_{j+1} q_{j}\right)\left(X^{2}+\frac{1}{2} Y, 0\right)$, and obtain Theorem 1.32.(1). However we will give a variation of this proof which works in our complex case.

Proof of Proposition 3.13. To understand the position of the curve $\breve{V}^{j+1}=\left(F^{c_{j+1}}\right)^{-1}\left(\left\{x_{j+1}\right\} \times I\right)$ w.r.t. $F^{\unlhd_{j}}\left(H^{c_{j}} \cap Y^{\rrbracket_{j}}\right.$ ), we are going to use the two following curves:

- The $\pi^{0\left(c_{j+1}\right)}$-fiber $\hat{V}^{j+1}$ of $\left(F^{c_{j+1}}\right)^{-1}\left(\zeta^{c_{j+1}}\right)$.
- The $\pi^{t\left(Ð_{j}\right)}$-fiber $V^{j+1}$ of $F^{\unlhd_{j}}\left(\zeta^{\iota_{j}}\right)$.

The two next lemmas states that these curves are close to $\breve{V}^{j+1}$.
Lemma 3.19. The two curves $V^{j+1}$ and $\hat{V}^{j+1}$ are $C^{2}$-close and $o\left(\breve{\gamma}_{j}^{2}\right)-C^{0}$-close when $j$ is large.
Proof. By Theorem 1.32 (ii), the distance between the $\pi^{t\left(c_{j+1}\right)}$-fibers of the points $F^{c_{j+1}} \circ F^{\complement_{j}}\left(\zeta^{c_{j}}\right)$ and $\zeta^{c_{j+1}}$ is small compared to $\breve{\gamma}_{j+1}$. As the horizontal direction is at least $\asymp\left|Y^{c_{j+1}}\right|^{-1}$ expanded by $F^{c_{j+1}}$ by Corollary 3.10, it comes that the pull back of these fibers by $F^{c_{j+1}}$ are at $C^{0}$-distance


Figure 7. Notations for the proof of Proposition 3.13
small compared to $\breve{\gamma}_{j+1} \cdot\left|Y^{c_{j+1}}\right|=\breve{\gamma}_{j}^{2}$. By the equivariance of the fibers given by Definition 1.23.(3), these two pullbacks are equal to respectively $V^{j+1}$ and $\hat{V}^{j+1}$ respectively. Since the $\pi$-fibers depend continuously on the based points in the $C^{2}$-topology, these two curves are also $C^{2}$-close.
Lemma 3.20. The two curves $V^{j+1}$ and $\breve{V}^{j+1}$ are $C^{2}$-close and o $\left(\breve{\gamma}_{j}^{2}\right)-C^{0}$-close when $j$ is large.
Proof. By Lemma 3.19, it suffices to show that the two curves $\hat{V}^{j+1}$ and $\breve{V}^{j+1}$ are $C^{2}$-close and $o\left(\breve{\gamma}_{j}^{2}\right)$ - $C^{0}$-close when $j$ is large. The $\pi^{t\left(c_{j+1}\right)}$-fiber going through $\zeta^{c_{j+1}}$ is a curve included in $Y^{e}$ with endpoints in $\partial^{u} Y^{e}$ and tangent spaces in $\chi_{v}$ by Definition 1.23 (0) and (1). The same occurs obviously for the curve $\left\{x_{j+1}\right\} \times I$. Hence the curves $\hat{V}^{j+1}$ and $\breve{V}^{j+1}$ are the pullbacks by $F^{c_{j+1}}$ of two curves which are uniformly transverse to the unstable lamination of the maximal invariant set $\Lambda$ associated to the hyperbolic map $F^{s t}$. Hence by the Inclination Lemma, the curves $\hat{V}^{j+1}$ and $\breve{V}^{j+1}$ are $C^{2}$-close when $\left|c_{j+1}\right|$ is large and so when $j$ is large.

We now show that the Hausdorff distance between $\hat{V}^{j+1}$ and $\breve{V}^{j+1}$ is small compared to $\breve{\gamma}_{j}^{2}$. Note that the images by $F^{c_{j+1}}$ of $\hat{V}^{j+1}$ and $\breve{V}^{j+1}$ intersects at $\zeta^{c_{j+1}}$. These curves have tangent spaces in $\chi_{v}$ and are included in $F^{c_{j+1}}\left(Y^{c_{j+1}}\right)$, whose vertical height is smaller than $\left|F^{c_{j+1}}\left(Y^{c_{j+1}}\right)\right|$ and thus small compared to $\breve{\gamma}_{j+1}^{2}$ by assumption (iii) of Theorem 1.32. Thus the $C^{0}$-distance between the image by $F^{c_{j+1}}$ of $\hat{V}^{j+1}$ and $\breve{V}^{j+1}$ is small compared to $\breve{\gamma}_{j+1}^{2} \leq \breve{\gamma}_{j+1}$. As the horizontal direction is at least $\asymp\left|Y^{c_{j+1}}\right|^{-1}$ expanded by $F^{c_{j+1}}$ by Corollary 3.10, it comes that the $C^{0}$-distance between $\hat{V}^{j+1}$ and $\breve{V}^{j+1}$ is small compared to $\breve{\gamma}_{j+1} \cdot\left|Y^{c_{j+1}}\right|=\breve{\gamma}_{j}^{2}$.

The following lemma is the second main step of the proof of Proposition 3.13, it defines the new critical point $\breve{x}_{j}$ defined by $\breve{V}^{j+1}$, which is close to the first coordinate $x_{j}$ of $\zeta^{c_{j}}$.

Lemma 3.21. For every $j$ large enough, the following map:

$$
\breve{\Delta}_{j}: x \in I^{\unlhd_{j}} \mapsto\left(F_{1}^{\unlhd_{j}}-\breve{v}_{j+1} \circ F_{2}^{\unlhd_{j}}\right)\left(x, h_{j}(x)\right)
$$

has a unique critical point $\breve{x}_{j}$ in the $\breve{\gamma}_{j}$-neighborhood of $x_{j}$ in $I^{\unlhd_{j}}$ which is $\epsilon_{\mathscr{G}}$-distant to $\partial I^{\unlhd_{j}}$. Moreover $\left|\breve{x}_{j}-x_{j}\right|$ is small compared to $\breve{\gamma}_{j}$ and $\left|\breve{\Delta}_{j}\left(\breve{x}_{j}\right)\right|$ is small compared to $\breve{\gamma}_{j}^{2}$.

Proof. By Definition 1.23 (4), for every $\boxtimes \in \mathscr{C}$ and $u \in \overleftarrow{\mathscr{A}}_{\square}$, the curve $W^{u}$ is sent by $F^{\square}$ tangent to a unique fiber of $\pi^{t(『)}$ at a unique point $F^{\boxminus}\left(\zeta^{u}\right)$, with $\zeta^{u} \in \operatorname{int} Y^{\square}$, and the tangency is quadratic. By Definition 1.23 (1) there exists a $C^{2}$-function $v_{u}: I \rightarrow \mathbb{R}$ such that this fiber is equal to $\left\{\left(v_{u}(y), y\right): y \in I\right\}$. Thus there exists $r \in\left(0, \epsilon_{\mathscr{C}} / 3\right)$ such that $\left[x_{u}-3 r, x_{u}+3 r\right] \Subset I^{\square}$ and the following map

$$
\Delta_{u}: x \in I^{\complement} \mapsto F_{1}^{\bowtie}\left(x, w^{u}(x)\right)-v_{u} \circ F_{2}^{\bowtie}\left(x, w^{u}(x)\right)
$$

has second derivative of absolute value $>3 r$ on $\left[x_{u}-3 r, x_{u}+3 r\right]$ and a unique critical point in $\left(x_{u}-r / 3, x_{u}+r / 3\right)$ equal to $x_{u}$, with $x_{u}$ the first coordinate of $\zeta^{u}$. Note that $\Delta_{u}$ is indeed well-defined on $I^{\square}$ because $u \in \overleftarrow{\mathscr{A}}_{\square}$. Since $w^{u}$ depends continuously on $u$ in the $C^{2}$-topology and $\bigcup_{\boxminus \in \mathscr{C}} \overleftarrow{\mathscr{A}}_{\square}$ is compact, we can take a uniform value of $r$ among $\boxtimes \in \mathscr{C}$ and $u \in \overleftarrow{\mathscr{A}}_{\square}$

On the other hand, by Proposition 1.28, the curve $H^{c_{j}}$ is tangent to a unique fiber $V^{j+1}$ of $\pi^{t\left(\oplus_{j}\right)} \circ F^{\unlhd_{j}}$ at a unique point $\zeta^{c_{j}} \in \operatorname{int} Y^{\unlhd_{j}}$, and the tangency is quadratic. Still by Definition 1.23 (1) there exists a $C^{2}$-function $v_{j+1}: I \rightarrow \mathbb{R}$ such that $V^{j+1}=\left\{\left(v_{j+1}(y), y\right): y \in I\right\}$.

When $j$ is large we recall that $\left|c_{j}\right|$ is large. Still by Proposition 1.28, the word $c_{j}$ is made by the last letters of some $u_{j} \in \overleftarrow{\mathscr{A}}_{\mathrm{Đ}_{j}}$. By the inclination lemma (as in the proof of Proposition 1.28), the map $h_{j}$, defined on $I^{\square_{j}}$, is $C^{2}$-close to $w^{u_{j}}$ (also defined on $I^{\square_{j}}$ ) and $\zeta^{c_{j}}$ is close to $\zeta^{u_{j}}$ when $j$ is large. Thus $x_{j}$ is close to $x_{u_{j}}$ (in particular $\left[x_{j}-2 r, x_{j}+2 r\right] \Subset\left[x_{u_{j}}-3 r, x_{u_{j}}+3 r\right]$ ), and the maps $v_{u_{j}}$ and $v_{j+1}$ are $C^{2}$-close. Then the following map is defined on $I^{\square_{j}}$ and $C^{2}$-close to $\Delta_{u_{j}}$ on $\left[x_{j}-2 r, x_{j}+2 r\right]$ :

$$
\Delta_{j}: x \in I^{®_{j}} \mapsto F_{1}^{\bullet}\left(x, h_{j}(x)\right)-v_{j+1} \circ F_{2}^{\bullet}\left(x, h_{j}(x)\right) .
$$

Thus the map $\Delta_{j}$ has second derivative of absolute value $>2 r$ in $\left[x_{j}-2 r, x_{j}+2 r\right]$ and has a unique critical point in $\left(x_{j}-r / 2, x_{j}+r / 2\right)$ which is equal to $x_{j}$ by Proposition 1.28. Note that $\Delta_{j}\left(x_{j}\right)=0$ is a local extrenum of $\Delta_{j}$. By Lemma 3.20, the maps $v_{j+1}$ and $\breve{v}_{j+1}$ are both $C^{2}$-close and $o\left(\breve{\gamma}_{j}^{2}\right)$-C $C^{0}$-close when $j$ is large. Then for $j$ large enough the following function:

$$
\breve{\Delta}_{j}: x \in I^{\unlhd_{j}} \mapsto\left(F_{1}^{๑_{j}}-\breve{v}_{j+1} \circ F_{2}^{๑_{j}}\right)\left(x, h_{j}(x)\right)
$$

has a unique critical point $\breve{x}_{j}$ in $\left(x_{j}-r, x_{j}+r\right)$ which is close to $x_{j}$. Thus by Proposition 3.11 (2), $\breve{x}_{j}$ is $\epsilon_{\mathscr{C}}$-distant to $\partial I^{®_{j}}$ for every $j$ large. Furthermore, the second derivative of $\breve{\Delta}_{j}$ is with absolute value $>r$ on $\left[-r+\breve{x}_{j}, r+\breve{x}_{j}\right]$. In particular the $C^{0}$-distance between the extremal values $\Delta_{j}\left(x_{j}\right)$ and $\breve{\Delta}_{j}\left(\breve{x}_{j}\right)$ of $\Delta_{j}$ and $\breve{\Delta}_{j}$ is small compared to $\breve{\gamma}_{j}^{2}$. Since $\Delta_{j}\left(x_{j}\right)=0$, we then have $\breve{\Delta}_{j}\left(\breve{x}_{j}\right)=o\left(\breve{\gamma}_{j}^{2}\right)$. Let us estimate the distance between $\breve{x}_{j}$ and $x_{j}$ :

$$
\left|\breve{\Delta}_{j}\left(x_{j}\right)-\breve{\Delta}_{j}\left(\breve{x}_{j}\right)\right|=\left|\int_{\breve{x}_{j}}^{x_{j}} D \breve{\Delta}_{j}(t) d t\right| \geq\left|\int_{\breve{x}_{j}}^{x_{j}} \int_{\breve{x}_{j}}^{t} D^{2} \breve{\Delta}_{j}(s) d s d t\right| \geq \frac{r}{2}\left(x_{j}-\breve{x}_{j}\right)^{2}
$$

As $\left|\breve{\Delta}_{j}\left(\breve{x}_{j}\right)-\Delta_{j}\left(x_{j}\right)\right|=o\left(\breve{\gamma}_{j}^{2}\right)$ and by Lemma $3.20\left|\breve{\Delta}_{j}\left(x_{j}\right)-\Delta_{j}\left(x_{j}\right)\right|=o\left(\breve{\gamma}_{j}^{2}\right)$, it holds $\left|\breve{x}_{j}-x_{j}\right|=$ $o\left(\breve{\gamma}_{j}\right)$.

The third step of the proof of Proposition 3.13 will be to obtain the normal form of Eq. (3.4). It is done below in Eq. (3.7). We recall that by Lemma 3.21, $\breve{x}_{j}$ is $\epsilon_{\mathscr{G}}$-distant from $\partial I^{®_{j}}$. On the other hand the curve $H^{c_{j}}$ intersects $Y^{\unlhd_{j}}$ since $c_{j} \in \mathscr{A}_{\triangleright_{j}}^{*}$ and is $\epsilon_{\mathscr{C}}$-distant to $\partial^{u} Y^{\unlhd_{j}}$ by Proposition 3.11 (3). Thus it comes:

Fact 3.22. The function $\varphi_{j}(x, y)=\left(\breve{x}_{j}+x, h_{j}\left(\breve{x}_{j}+x\right)+y\right)$ is well defined on $\left(-\epsilon_{\mathscr{C}}, \epsilon_{\mathscr{C}}\right)^{2}$ and takes its values in $Y^{\square_{j}}$.

Similarly it holds:
Fact 3.23. If $F$ is analytic, there exists $\tilde{\epsilon}_{\mathscr{C}}>0$ such that for every $j$ large enough, the map $\varphi_{j}$ is well-defined on $\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}$ and satisfies both $\varphi_{j}\left(\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}\right) \subset \tilde{Y}^{\unlhd_{j}}$ and $F^{\square_{j}} \circ \varphi_{j}\left(\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}\right) \subset \tilde{Y}^{e}$.

Proof. We recall that $\breve{x}_{j} \in I^{\square_{j}}$ is $\epsilon_{\mathscr{C}}$-distant from $\partial I^{\square_{j}}$ and so $\min \left(\rho, \epsilon_{\mathscr{C}}\right)$-distant to $\partial \tilde{I}$. Since the maps $h_{j}$ extend to $\tilde{I}$ and are uniformly $C^{1}$-bounded on $\tilde{I}$ by Definition 1.5 (2), the maps $\varphi_{j}$ are well defined on $\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}$ for $\tilde{\epsilon}_{\mathscr{C}}<\min \left(\rho, \epsilon_{\mathscr{C}}\right)$ and uniformly $C^{1}$-bounded. Thus for $\tilde{\epsilon}_{\mathscr{C}}$ small, the set $\varphi_{j}\left(\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}\right)$ has a uniformly small diameter and has real trace $\varphi_{j}\left(\left(-\epsilon_{\mathscr{C}}, \epsilon_{\mathscr{C}}\right)^{2}\right) \subset \varphi_{j}\left(\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}\right)$. So to get the two last inequalities, it suffices to use Fact 3.22 and remark that the set $\tilde{Y}^{\square_{j}} \cap\left(F^{\square_{j}}\right)^{-1}\left(\tilde{Y}^{e}\right)$ contains a uniformly large complex neighborhood of $Y^{\square_{j}}$.

Also by Lemma 3.21 there exist real numbers $b_{j}, q_{j}$ and a $C^{2}$-function $r_{j}$ such that for every sufficiently large $j$ and $(x, y) \in\left(-\epsilon_{\mathscr{C}}, \epsilon_{\mathscr{C}}\right)^{2}$ :

$$
\begin{gather*}
\left(F_{1}^{๑_{j}}-\breve{v}_{j+1} \circ F_{2}^{๑_{j}}\right) \circ \varphi_{j}(x, y)=q_{j} \cdot x^{2}+b_{j} \cdot y+r_{j}(x, y),  \tag{3.7}\\
\text { with } \quad 0=\partial_{x} r_{j}(0)=\partial_{y} r_{j}(0,0)=\partial_{x}^{2} r_{j}(0) \text { and } r_{j}(0)=o\left(\breve{\gamma}_{j}^{2}\right) .
\end{gather*}
$$

If $F$ is analytic, the map $\breve{v}_{j+1}$ extends holomorphically to $\tilde{I}$ and then, by Fact 3.23 , the function $r_{j}$ extends holomorphically to $\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}$. The next lemma achieves the proof of the proposition.

Lemma 3.24. The families $\left(\varphi_{j}\right)_{j}$ and $\left(r_{j}\right)_{j}$ are normal in the $C^{2}$-topology, and if $F$ is analytic, they are also normal on $\mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}$. Moreover, there exists $K \geq 1$ such that for $j$ large enough:

$$
K^{-1}<\left|q_{j}\right|<K \quad \text { and } \quad K^{-1}<\left|b_{j}\right|<K
$$

Proof. By finiteness of $\mathscr{C}$, by compactness of $\overrightarrow{\mathscr{A}}$ and $\overleftarrow{\mathscr{A}}_{\square}$ and by the Inclination Lemma, the following are compact subsets of $C^{2}(I, I)$ :

$$
\begin{gathered}
\left\{w^{s}: s \in \overrightarrow{\mathscr{A}}\right\} \cup\left\{\mathcal{X}^{c}(x, \cdot): c \in \mathscr{A}^{*} \& x \in I\right\} \\
\left\{w^{u}: \boxtimes \in \mathscr{C}, u \in \overleftarrow{\mathscr{A}}_{\square}\right\} \cup\left\{\mathcal{Y}^{c}(\cdot, 0): \boxtimes \in \mathscr{C}, c \in \mathscr{A}_{\square}^{*}\right\}
\end{gathered}
$$

We recall that $\breve{v}_{j+1}$ belongs to the first compact set and $h_{j}$ belongs to the second one. Thus the sequence of $C^{2}$-diffeomorphisms $\left(F^{\unlhd_{j}} \circ \varphi_{j}\right)_{j}$ is normal and $\left((x, y) \in I^{2} \mapsto\left(x-\breve{v}_{j+1}(y), y\right)\right)_{j}$ as well, and so also the family of their compositions. Thus $\left(r_{j}\right)_{j}$ is normal and $\left(\left|b_{j}\right|\right)_{j},\left(\left|q_{j}\right|\right)_{j},\left(\left\|r_{j}\right\|_{C^{2}}\right)_{j}$ are bounded. If $F$ is analytic we do the same proof using this time the inclination lemma in $\mathbb{C}^{2}$. Note also that:

where the first inequality come from the fact that $\varphi_{j}$ is conservative and the second equality comes from the fact that the determinant of $(x, y) \mapsto\left(x-\breve{v}_{j+1}(y), y\right)$ is equal to 1 . As $\left|\operatorname{det} D F^{\square_{j}}\right|$ is positive on its compact domain and $\left(F_{2}^{\square_{j}} \circ \varphi_{j}\right)_{j}$ is normal, there exists a positive bound from below for $\left|b_{j}\right|$. Finally the existence of the bound from below for $\left|q_{j}\right|$ follows from the uniform bound $r$ at the end of the proof of Lemma 3.21.

Remark 3.25. The set $B_{j}$ contains both $\breve{\zeta}_{j}:=\Phi_{j}(0)$ and $\zeta^{c_{j}}$ when $j$ is large, since $\left|\breve{x}_{j}-x_{j}\right|$ is small compared to $\breve{\gamma}_{j}$ by Lemma 3.21 and thus compared to $\gamma_{j}$.

This Remark with Lemma 3.24 implying $\left|q_{j} / b_{j}\right| \leq K^{2}$ and the normality of $\left(\varphi_{j}\right)_{j}$ give:
Fact 3.26. For $j$ large enough, the set $B_{j}$ contains both $\zeta^{c_{j}}$ and $\breve{\zeta}_{j}=\Phi_{j}(0)$, and has its diameter dominated by $\gamma_{j}$. Likewise the set $\tilde{B}_{j} \supset B_{j}$ has its diameter dominated by $\gamma_{j}$.
3.3. Proof of Theorems $\mathbf{1 . 3 2}$ and 1.34. We start by proving one by one each of the six inclusions and then we will prove the two limits of Theorems 1.32 and 1.34.

Claim 3.27. For $j$ large, the domain $B_{j}$ is included in $Y^{\unlhd_{j}}$ and if $F$ is analytic, the domain $\tilde{B}_{j}$ is included in $\tilde{Y}^{\unlhd_{j}}$. Moreover they are uniformly distant to respectively $\partial Y^{\square_{j}}$ and $\partial \tilde{Y}^{\unlhd_{j}}$ 。

Proof. By Fact 3.26 , the sets $B_{j} \subset \breve{B}_{j}$ contain $\breve{\zeta}_{j}$ and have small diameters. As the point $\breve{\zeta}_{j}$ is uniformly distant to $\partial Y^{\unrhd_{j}}$ and $\partial \tilde{Y}^{\unlhd_{j}}$ by Proposition 3.13 for its $x$-coordinate $\breve{x}_{j}$ and by Proposition 3.11 (3) for its $y$-coordinate, the claim follows.

Let us show that $F^{๑_{j}}$ sends $B_{j}$ into $Y^{c_{j+1}}$, and if $F$ is analytic, that $F^{\unlhd_{j}}$ sends $\tilde{B}_{j}$ into $\tilde{Y}^{c_{j+1}}$, as claimed in the second inclusions of the two conclusions of Theorems 1.32 and 1.34. We are going to show even more, that they are respectively included in the following sets:

$$
Y_{\mathcal{B}}^{c_{j+1}}:=\left\{\left(\mathcal{X}_{j+1}(x, y), y\right):(x, y) \in \mathcal{B}\right\} \quad \text { with } \mathcal{B}:=\operatorname{pr}_{1}\left(B_{j+1}\right) \times F_{2}^{\complement_{j}}\left(B_{j}\right)
$$

with $\mathrm{pr}_{1}$ the first coordinate projection and $F_{2}^{\complement_{j}}$ the second coordinate of $F^{\complement_{j}}$, and respectively:

$$
\tilde{Y}_{\mathcal{B}}^{c_{j+1}}:=\left\{\left(\mathcal{Z}_{j+1}(z, w), w\right):(z, w) \in \tilde{\mathcal{B}}\right\} \quad \text { with } \tilde{\mathcal{B}}:=\operatorname{pr}_{1}\left(\tilde{B}_{j+1}\right) \times F_{2}^{\complement_{j}}\left(\tilde{B}_{j}\right)
$$

Proposition 3.28. For every $j \geq 1$ sufficiently large, the map $F^{\rrbracket_{j}}$ sends $B_{j}$ into $Y_{\mathcal{B}}^{c_{j+1}} \subset Y^{c_{j+1}}$. If moreover $F$ is analytic, the map $F^{\unlhd_{j}}$ sends $\tilde{B}_{j}$ into $\tilde{Y}_{\mathcal{B}}^{c_{j+1}} \subset \tilde{Y}^{c_{j+1}}$.

Proof. We define $\psi_{j}(x, y):=\left(x-\breve{v}_{j+1}(y), y\right)$. To prove this proposition, it suffices to show:

$$
\psi_{j} \circ F^{\unlhd_{j}}\left(B_{j}\right) \subset \psi_{j}\left(Y_{\mathcal{B}}^{c_{j+1}}\right) \quad \text { and } \quad \text { resp. } \psi_{j} \circ F^{\unlhd_{j}}\left(\tilde{B}_{j}\right) \subset \psi_{j}\left(\tilde{Y}_{\mathcal{B}}^{c_{j+1}}\right)
$$

Thus it suffices to prove the two following inclusions:

$$
\begin{gather*}
\psi_{j} \circ F^{\unlhd_{j}}\left(B_{j}\right) \subset\left(-\frac{1}{4}\left|\gamma_{j+1} \sigma_{j+1}\right|, \frac{1}{4}\left|\gamma_{j+1} \sigma_{j+1}\right|\right) \times F_{2}^{\complement_{j}}\left(B_{j}\right) \quad \subset \quad \psi_{j}\left(Y_{\mathcal{B}}^{c_{j+1}}\right)  \tag{3.8}\\
\text { resp. } \psi_{j} \circ F^{\complement_{j}}\left(\tilde{B}_{j}\right) \quad \subset \mathbb{D}\left(\frac{1}{4}\left|\gamma_{j+1} \sigma_{j+1}\right|\right) \times F_{2}^{®_{j}}\left(\tilde{B}_{j}\right) \quad \subset \quad \psi_{j}\left(\tilde{Y}_{\mathcal{B}}^{c_{j+1}}\right) \tag{3.9}
\end{gather*}
$$

Proof of the first inclusions of Eq. (3.8) and Eq. (3.9). By Definition 3.17, we have $B_{j}=$ $\Phi_{j}\left((-0.3,0.3)^{2}\right)$. Thus it suffices to show that for every $(X, Y) \in(-0.3,0.3)^{2}$ the first coordinate $x$ of the following point is in $\left(-\frac{1}{4}\left|\gamma_{j+1} \sigma_{j+1}\right|, \frac{1}{4}\left|\gamma_{j+1} \sigma_{j+1}\right|\right)$ :

$$
\begin{equation*}
(x, y):=\psi_{j} \circ F^{\complement_{j}} \circ \Phi_{j}(X, Y)=\psi_{j} \circ F^{\complement_{j}} \circ \varphi_{j}\left(\gamma_{j} \cdot X, \frac{q_{j} \cdot \gamma_{j}^{2}}{2 b_{j}} \cdot Y\right) \tag{3.10}
\end{equation*}
$$

Note that $|x|$ is the horizontal distance from $F^{\unlhd_{j}} \circ \Phi_{j}(X, Y)$ to $\breve{V}^{j+1}$. By Eq. (3.4) it holds:

$$
\begin{equation*}
x=q_{j} \cdot \gamma_{j}^{2} \cdot\left(X^{2}+\frac{Y}{2}+R_{j}(X, Y)\right) \quad \text { with } R_{j}(X, Y):=\frac{1}{q_{j} \gamma_{j}^{2}} r_{j}\left(\gamma_{j} X, \frac{q_{j}}{2 b_{j}} \gamma_{j}^{2} Y\right) \tag{3.11}
\end{equation*}
$$

On the other hand, if $F$ is analytic, we have $\tilde{B}_{j}=\Phi_{j}\left(\mathbb{D}(0.3)^{2}\right)$ by Definition 3.17. Thus it suffices to show that for every $(Z, W) \in \mathbb{D}(0.3)^{2}$ the first coordinate $z$ of the following point is in $\mathbb{D}\left(-\frac{1}{4}\left|\gamma_{j+1} \sigma_{j+1}\right|\right):$

$$
\begin{equation*}
(z, w):=\psi_{j} \circ F^{\unlhd_{j}} \circ \Phi_{j}(Z, W)=\psi_{j} \circ F^{\complement_{j}} \circ \varphi_{j}\left(\gamma_{j} \cdot Z, \frac{q_{j} \cdot \gamma_{j}^{2}}{2 b_{j}} \cdot W\right) \tag{3.12}
\end{equation*}
$$

Again $|z|$ is the distance from $F^{®_{j}} \circ \Phi_{j}(Z, W)$ to the curve $\breve{V}^{j+1}$. By Eq. (3.7) it holds:

$$
\begin{equation*}
z=q_{j} \cdot \gamma_{j}^{2} \cdot\left(Z^{2}+\frac{W}{2}+R_{j}(Z, W)\right) \tag{3.13}
\end{equation*}
$$

To conclude, we need the following, whose proof is done below:
Lemma 3.29. (i) For $j$ large, the map $(X, Y) \in(-0.3,0.3)^{2} \mapsto R_{j}(X, Y)$ is $C^{0}$-small,
(ii) if moreover $F$ is analytic, the map $(Z, W) \in \mathbb{D}(0.3)^{2} \mapsto R_{j}(Z, W)$ is $C^{0}$-small.

Since $|X|^{2}+\frac{1}{2}|Y| \leq(0.3)^{2}+\frac{1}{2} \cdot 0.3<1 / 4$ with $\left|q_{j} \cdot \gamma_{j}^{2}\right|=\left|\gamma_{j+1} \cdot \sigma_{j+1}\right|$, the latter lemma implies immediately that for $j$ large enough and $(X, Y) \in(-0.3,0.3)^{2}$, the number $x$ is in $\left(-\frac{1}{4}\left|\gamma_{j+1} \sigma_{j+1}\right|, \frac{1}{4}\left|\gamma_{j+1} \sigma_{j+1}\right|\right)$, which concludes the proof of the first inclusion of Eq. (3.8).

If $F$ is analytic, since $|Z|^{2}+\frac{1}{2}|W| \leq(0.3)^{2}+\frac{1}{2} \cdot 0.3<1 / 4$ for every $(Z, W) \in \mathbb{D}(0.3)^{2}$, the same argument gives that $z$ is in $\mathbb{D}\left(\left|\gamma_{j+1} \sigma_{j+1}\right| / 4\right)$ and so gives the first inclusion of Eq. (3.9).

Proof of Lemma 3.29. Let us first prove (i). We recall that $R_{j}(X, Y):=\frac{1}{q_{j} \gamma_{j}^{2}} r_{j}\left(\gamma_{j} X, \frac{q_{j}}{2 b_{j}} \gamma_{j}^{2} Y\right)$ for every $(X, Y) \in(-0.3,0.3)^{2}$. We recall that $\gamma_{j}$ is small for $j$ large. Moreover by Eq. (3.5) it holds:

$$
0=\partial_{x} R_{j}(0)=\partial_{y} R_{j}(0)=\partial_{x}^{2} R_{j}(0)
$$

By Eq. (3.5), $r_{j}(0)$ is small compared to $\breve{\gamma}_{j}^{2}$ and thus compared to $\gamma_{j}^{2}$. Thus $R_{j}(0)$ is small and so the $C^{0}$-norm of $R_{j}$ is dominated by a small constant plus the $C^{0}$-norm of its second derivative. We recall that $\left(r_{j}\right)_{j}$ is a normal family in the $C^{2}$-topology by Lemma 3.24. Let $\mu$ be a uniform modulus of continuity of their second derivative. We observe that $\left(\partial_{Y}^{2} R_{j}\right)_{j}$ and $\left(\partial_{X} \partial_{Y} R_{j}\right)_{j}$ converge uniformly to 0 in the $C^{0}$ topology. On the other hand $\partial_{X}^{2} R_{j}(0)=0$ and so $\partial_{X}^{2} R_{j}(X, Y)=\frac{1}{q_{j}} \partial_{x}^{2} r_{j}\left(\gamma_{j} X, \frac{q_{j}}{2 b_{j}} \gamma_{j}^{2} Y\right)$ is bounded by $\frac{1}{\left|q_{j}\right|} \mu\left(\operatorname{diam}(-0.3,0.3)^{2} \cdot K^{2} \cdot \gamma_{j}\right)$ which converges uniformly to 0 since $\gamma_{j} \rightarrow 0$ and $\left|q_{j}\right|^{-1} \leq K$ (by Lemma 3.24). Thus $R_{j}$ is $C^{0}$-small for $j$ large.

Let us now prove $(i i)$. Let us notice that $\left(r_{j} \mid \mathbb{D}\left(\tilde{\epsilon}_{\mathscr{C}}\right)^{2}\right)_{j}$ is normal by Lemma 3.24. Then the proof of $(i i)$ is the same as the proof of $(i)$, replacing $(-0.3,0.3)^{2}$ by $\mathbb{D}(0.3)^{2}$.

Proof of the second inclusions of Eq. (3.8) and Eq. (3.9). We will use the following:
Lemma 3.30. The following map is $C^{1}$-close to the first coordinate projection when $j$ is large:

$$
\Upsilon:(x, y) \in \mathcal{B} \mapsto \frac{\mathcal{X}_{j+1}(x, y)-\mathcal{X}_{j+1}\left(x_{j+1}, y\right)}{\sigma_{j+1}}+x_{j+1}
$$

If $F$ is analytic, the following map is $C^{1}$-close to the first coordinate projection when $j$ is large:

$$
\tilde{\Upsilon}:(z, w) \in \tilde{\mathcal{B}} \mapsto \frac{\mathcal{Z}_{j+1}(z, w)-\mathcal{Z}_{j+1}\left(x_{j+1}, w\right)}{\sigma_{j+1}}+x_{j+1}
$$

Proof. To show this lemma it suffices to use the distortion's Proposition B.3, whose assumptions must be fulfilled. The first is that when $j$ is large, the diameter of $\mathcal{B}$ must be small and when $F$ is analytic, that the diameter of $\tilde{\mathcal{B}}$ must be small compared to $\sigma_{j+1}$. Indeed by Remark 1.33, it holds $\breve{\gamma}_{j}=o\left(\left|Y^{c_{j+1}}\right|\right)$ and so $\breve{\gamma}_{j}=o\left(\left|\sigma_{j+1}\right|\right)$ by Corollary 3.10. Thus by Eq. (3.6), we have $\gamma_{j}=o\left(\left|\sigma_{j+1}\right|\right)$. The second assumption in the real case is that the distortion of $\left(\mathcal{X}_{j}\right)_{j}$ is bounded; this is the case by Theorem 3.9. The second assumption in the complex case is that $\tilde{\mathcal{B}}$ is uniformly distant to the boundary of $\tilde{I}^{2}$. This is indeed the case since $\tilde{B}_{j} \ni \zeta^{c_{j}}$ has a small diameter by Fact 3.26, while the first coordinate of $\zeta^{c_{j}} \in Y^{\unlhd_{j}}$ and the second coordinate of $F^{\unlhd_{j}}\left(\zeta^{c_{j}}\right) \in F^{®_{j}}\left(Y^{\unlhd_{j}}\right)$ are uniformly distant from $\partial \tilde{I}$, by Proposition 3.11 (1) and (2).

Let us now use this lemma to prove the second inclusions of Eq. (3.8) and Eq. (3.9). We notice that it suffices to show:

$$
\begin{aligned}
& \left\{\mathcal{X}_{j+1}(x, y)-\breve{v}_{j+1}(y): x \in \operatorname{pr}_{1}\left(B_{j+1}\right)\right\} \supset\left(-\left|\gamma_{j+1} \sigma_{j+1}\right| / 4,\left|\gamma_{j+1} \sigma_{j+1}\right| / 4\right) \quad \forall y \in F_{2}^{®_{j}}\left(B_{j}\right) \\
& \quad \text { resp. }\left\{\mathcal{Z}_{j+1}(z, w)-\breve{v}_{j+1}(w): z \in \operatorname{pr}_{1}\left(\tilde{B}_{j+1}\right)\right\} \supset \mathbb{D}\left(\left|\gamma_{j+1} \sigma_{j+1}\right| / 4\right) \quad \forall w \in F_{2}^{®_{j}}\left(\tilde{B}_{j}\right)
\end{aligned}
$$

Now we infer that $\breve{v}_{j+1}(y)=\mathcal{X}_{j+1}\left(x_{j+1}, y\right)$ (resp. $\left.\breve{v}_{j+1}(w)=\mathcal{Z}_{j+1}\left(x_{j+1}, w\right)\right)$. So it suffices to show that:

$$
\begin{aligned}
& \left\{\Upsilon(x, y)-x_{j+1}: x \in \operatorname{pr}_{1}\left(B_{j+1}\right)\right\} \supset\left(-\left|\gamma_{j+1}\right| / 4,\left|\gamma_{j+1}\right| / 4\right) \quad \forall y \in F_{2}^{®_{j}}\left(B_{j}\right) . \\
& \quad \text { resp. }\left\{\tilde{\Upsilon}(z, w)-x_{j+1}: z \in \operatorname{pr}_{1}\left(\tilde{B}_{j+1}\right)\right\} \supset \mathbb{D}\left(\left|\gamma_{j+1}\right| / 4\right) \quad \forall w \in F_{2}^{®_{j}}\left(\tilde{B}_{j}\right) .
\end{aligned}
$$

Now we infer that $\left|\breve{x}_{j+1}-x_{j+1}\right|$ is small compared to $\breve{\gamma}_{j+1}$ by Lemma 3.21 and so compared to $\gamma_{j+1}$ by Eq. (3.6). As $\operatorname{pr}_{1}\left(B_{j+1}\right)=\breve{x}_{j+1}+\left(-0.3 \cdot \gamma_{j+1}, 0.3 \cdot \gamma_{j+1}\right)\left(\right.$ resp. $\left.\operatorname{pr}_{1}\left(\tilde{B}_{j+1}\right)=\breve{x}_{j+1}+\mathbb{D}\left(0.3 \cdot \gamma_{j+1}\right)\right)$, the result follows from Lemma 3.30 and $1 / 4<0.3$.

The following (with Proposition 3.28) shows the two remaining inclusions of Theorems 1.32 and 1.34:

Proposition 3.31. For every $j \geq 1$ large enough, the map $F^{c_{j+1}}$ sends $Y_{\mathcal{B}}^{c_{j+1}}$ into $B_{j+1}$. If moreover $F$ is analytic, the map $F^{c_{j+1}}$ sends $\tilde{Y}_{\mathcal{B}}^{c_{j+1}}$ into $\tilde{B}_{j+1}$.
Proof. Since $Y_{\mathcal{B}}^{c_{j+1}}=\left\{\left(\mathcal{X}_{j+1}(x, y), y\right):(x, y) \in \mathcal{B}\right\}$ and $\tilde{Y}_{\mathcal{B}}^{c_{j+1}}=\left\{\left(\mathcal{Z}_{j+1}(z, w), w\right):(z, w) \in \tilde{\mathcal{B}}\right\}$, by definition of the implicit representation, it holds:

Fact 3.32. The map $F^{c_{j+1}}$ sends $Y_{\mathcal{B}}^{c_{j+1}}$ to $\left\{\left(x, \mathcal{Y}_{j+1}(x, y)\right):(x, y) \in \mathcal{B}\right\}$. If moreover $F$ is analytic, the map $F^{c_{j+1}}$ sends $\tilde{Y}_{\mathcal{B}}^{c_{j+1}}$ to $\left\{\left(z, \mathcal{W}_{j+1}(z, w)\right):(z, w) \in \tilde{\mathcal{B}}\right\}$.

Thus we shall prove that $\left\{\left(x, \mathcal{Y}_{j+1}(x, y)\right):(x, y) \in \mathcal{B}\right\}$ is included in $B_{j+1}$ and respectively that $\left.\left\{\left(z, \mathcal{W}_{j+1}(z, w)\right):(z, w) \in \tilde{\mathcal{B}}\right\}\right)$ is included in $\tilde{B}_{j+1}$. By the skew product form of the sets $B_{j+1}$ and $\tilde{B}_{j+1}$ (see Definitions 3.15 and 3.17 ), it suffices to show that:

$$
\begin{aligned}
& \left|\mathcal{Y}_{j+1}(x, y)-h_{j+1}(x)\right|=\left|\mathcal{Y}_{j+1}(x, y)-\mathcal{Y}_{j+1}(x, 0)\right| \\
& \quad \quad \text { and } \quad \text { resp. } \quad\left|\mathcal{W}_{j+1}(z, w)-h_{j+1}(z)\right|=\left|\mathcal{W}_{j+1}(z, w)-\mathcal{W}_{j+1}(z, 0)\right|
\end{aligned}
$$

are small compared $\gamma_{j+1}^{2}$ for any $(x, y) \in \mathcal{B}$ and resp. $(z, w) \in \tilde{\mathcal{B}}$. They are respectively smaller than $\left|F^{c_{j+1}}\left(Y^{c_{j+1}}\right)\right|$ and $\left|F^{c_{j+1}}\left(\tilde{Y}^{c_{j+1}}\right)\right|$. By Theorem 1.32 (iii) and Theorem $1.34(\widetilde{i i i})$, we have
$\left|F^{c_{j+1}}\left(Y^{c_{j+1}}\right)\right|=o\left(\breve{\gamma}_{j+1}^{2}\right)$, and moreover $\left|F^{c_{j+1}}\left(\tilde{Y}^{c_{j+1}}\right)\right|=o\left(\breve{\gamma}_{j+1}^{2}\right)$ if $F$ is analytic. Since $\breve{\gamma}_{j+1}^{2}=$ $O\left(\gamma_{j+1}^{2}\right)$ by Eq. (3.6), the result follows.

We are now in position to prove:
Proof of Theorems 1.32 and 1.34. The six inclusions of Theorems 1.32 and 1.34 are given by Claim 3.27, Proposition 3.28 and then Proposition 3.31. Hence it remains only to show:

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \operatorname{diam} F^{k}\left(B_{j} \times\left\{a\left(\square_{j}\right)\right\}\right)=0 \quad \text { and } \quad \lim _{k \rightarrow \infty} \operatorname{diam} F^{k}\left(\tilde{B}_{j} \times\left\{a\left(\square_{j}\right)\right\}\right)=0 \quad \text { if } F \text { is analytic } \tag{3.14}
\end{equation*}
$$

when $j$ is large. To this end, it suffices to show that for every $\eta>0$, for every $j$ sufficiently large, the $\left|c_{j+1}\right|$ first iterates of $F^{\square_{j}}\left(B_{j}\right)$ (resp. $F^{\square_{j}}\left(\tilde{B}_{j}\right)$ ) have diameters smaller than $\eta$. To this end we use:

Lemma 3.33. When $N$ is large, $\max _{(p, q) \in \mathscr{B}^{*}:|p|,|q| \geq N} \operatorname{diam} Y^{p} \cap F^{q}\left(Y^{q}\right)$ is small.
Proof. For every $s, t \in I$, let $\mathcal{H}_{s}$ be the image of $\{y=s\} \cap Y^{q}$ under $F^{q}$ and $\mathcal{V}_{t}$ be the image of $\{x=t\} \cap F^{\mathcal{R}}\left(Y^{p}\right)$ under $\left(F^{p}\right)^{-1}$. We notice that $\mathcal{H}_{s}$ and $\mathcal{V}_{t}$ intersects at a unique point $\phi(s, t) \in Y^{p} \cap F^{q}\left(Y^{q}\right)$, since by the cone properties of Definition $1.4, \mathcal{H}_{s}$ is a graph over $x \in I$ with tangent spaces in $\chi_{h}$ and $\mathcal{V}_{t}$ is a graph over $y \in I$ with tangent spaces in $\chi_{v}$. It is straightforward that $\phi$ sends $I^{2}$ onto $Y^{p} \cap F^{q}\left(Y^{q}\right)$. Also by transversality, $\phi$ is differentiable. For every $(s, t) \in I^{2}$, the vector $\partial_{t} \phi(s, t)$ is tangent to $\mathcal{H}_{s}$ and so in $\chi_{h}$ and the norm of its first coordinate projection is at most $\lambda^{-|p|}$ by property (2) of Definition 1.4. Thus the norm of $\partial_{t} \phi(s, t)$ is at most $\sqrt{1+\theta^{2}} \cdot \lambda^{-|p|}$. Likewise the norm of $\partial_{s} \phi(s, t)$ is at most $\sqrt{1+\theta^{2}} \cdot \lambda^{-|q|}$. Thus the norm of $D \phi$ is small when $|p|$ and $|q|$ are large and so the diameter of $\phi\left(I^{2}\right)=Y^{p} \cap F^{q}\left(Y^{q}\right)$ is small.

By Lemma 3.33, there exists $N$ independent of $j$ such that the $k^{t h}$ iterate of $F^{\complement_{j}}\left(B_{j}\right)$ has diameter smaller than $\eta$ for every $N \leq k \leq\left|c_{j+1}\right|-N$. By Fact 3.26, the diameter of $B_{j}$ is dominated by $\gamma_{j}$ which is small when $j$ is large and likewise for $F^{c_{j+1}} \circ F^{®_{j}}\left(B_{j}\right) \subset B_{j+1}$. Thus their $N$ first iterates and preimages have also small diameter when $j$ is large. This implies the first limit in Eq. (3.14).

If $F$ is analytic, we prove the second limit in Eq. (3.14) exactly the same way, replacing Lemma 3.33 by Lemma 3.34 below.

Lemma 3.34. When $N$ is large, $\max _{(p, q) \in \mathscr{B}^{*}:|p|,|q| \geq N} \operatorname{diam} \tilde{Y}^{p} \cap F^{q}\left(\tilde{Y}^{q}\right)$ is small.
Proof. The proof is exactly the same as for Lemma 3.33, up to changing $I$ to $\tilde{I}$ and $Y$ to $\tilde{Y}$.

## 4. Parameter selection

In this Section, we will work with an unfolding $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of wild type $(\mathscr{A}, \mathscr{C})$. Under extra assumptions ( $\# \mathscr{C}=5$ and moderate dissipativeness), given any $p_{0} \in \operatorname{int} \mathcal{P}$, we will define in Section 4.3 a sequence of parameters $\left(p_{i}\right)_{i \geq 0}$ which converges to a certain $p_{\infty}$ close to $p_{0}$ such that at each $p_{i}$, the map $F_{p_{i}}$ will satisfy tangency's conditions. Then we prove Theorem 1.45 in Section 4.4. To this aim, we first show that $F_{p_{\infty}}$ satisfies the assumptions of Theorem 1.32, and those of Theorem 1.34 if $F_{p_{\infty}}$ is analytic, and finally we will show property $(\boldsymbol{\downarrow})$.
4.1. Estimates for parameter families. We set up here some analytic constants useful for the proof. For every $c \in \mathscr{A}^{*}$ and $p \in \mathcal{P}$, let $\left(\mathcal{X}_{p}^{c}, \mathcal{Y}_{p}^{c}\right)$ denote the implicit representation of $\left(Y_{p}^{c}, F_{p}^{c}\right)$. We recall that by Proposition 1.38, the family $\left(\mathcal{X}_{p}^{c}, \mathcal{Y}_{p}^{c}\right)_{p \in \mathcal{P}}$ is of class $C^{2}$. Here is a consequence of Theorem 3.9:

Corollary 4.1. There exists $\overline{\mathbf{D}}>0$ such that for all $c \in \mathscr{A}^{*}$ and $p \in \mathcal{P}: \mathbf{D}\left(Y_{p}^{c}, F_{p}^{c}\right) \leq \overline{\mathbf{D}}$.
Proof. As the domain of each $F_{p}$ is compact, and the set $\mathcal{P}$ is compact, the $C^{2}$-norm of $F_{p}$ and its inverse are bounded independently of $p \in \mathcal{P}$. Thus for every $a \in \mathscr{A}$, the distortion of each hyperbolic transformation $\left(Y_{p}^{a}, F_{p}^{a}\right)$ is bounded independently of $p \in \mathcal{P}$. By finiteness of $\mathscr{A}$, these distortions are bounded independently of $p \in \mathcal{P}$ and $a \in \mathscr{A}$. We conclude using Theorem 3.9.

An immediate consequence of this bound is:
Corollary 4.2. For every $c \in \mathscr{A}^{*}$, for every $(x, y),\left(x^{\prime}, y^{\prime}\right) \in I^{2}$, it holds for every $p \in \mathcal{P}$ :

$$
\left|\frac{\partial_{x} \mathcal{X}_{p}^{c}(x, y)}{\partial_{x} \mathcal{X}_{p}^{c}\left(x^{\prime}, y^{\prime}\right)}\right| \leq \exp \left(\overline{\mathbf{D}} \cdot\left\|\left(x-x^{\prime}, y-y^{\prime}\right)\right\|\right) \leq \overline{\mathbf{B}}, \quad \text { with } \overline{\mathbf{B}}:=\exp (2 \sqrt{2} \cdot \overline{\mathbf{D}})<\infty
$$

Definition 4.3. Given a non-empty compact subset $\mathcal{P}_{0} \subset \mathcal{P}$, we define the maximal and minimal widths of $c \in \mathscr{A}^{*}$ [among $p \in \mathcal{P}_{0}$ ] by respectively:

$$
\underline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right):=\min _{p \in \mathcal{P}_{0}, I^{2}}\left|\partial_{x_{1}} \mathcal{X}_{p}^{c}\right| \quad \text { and } \quad \overline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right):=\max _{p \in \mathcal{P}_{0}, I^{2}}\left|\partial_{x_{1}} \mathcal{X}_{p}^{c}\right|
$$

The maximal and minimal vertical height of $c$ are respectively:

$$
\underline{\mathrm{h}}\left(c, \mathcal{P}_{0}\right):=\min _{p \in \mathcal{P}_{0}, I^{2}}\left|\partial_{y_{0}} \mathcal{Y}_{p}^{c}\right| \quad \text { and } \quad \overline{\mathrm{h}}\left(c, \mathcal{P}_{0}\right):=\max _{p \in \mathcal{P}_{0}, I^{2}}\left|\partial_{y_{0}} \mathcal{Y}_{p}^{c}\right|
$$

Remark 4.4. For any $p \in \mathcal{P}_{0}$ and $c \in \mathscr{A}^{*}$, we notice that any line $\mathbb{R} \times\{y\}, y \in I$, intersects $Y_{p}^{c}$ at a segment of length in $\left[2 \underline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right), 2 \overline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right)\right]$; whereas any line $\{x\} \times \mathbb{R}, x \in I$, intersects $F_{p}^{c}\left(Y_{p}^{c}\right)$ at a segment of length in $\left[2 \underline{h}\left(c, \mathcal{P}_{0}\right), 2 \overline{\mathrm{~h}}\left(c, \mathcal{P}_{0}\right)\right]$. In particular, with the notations of Definition 1.27, is holds:

$$
2 \underline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right) \leq\left|Y_{p}^{c}\right| \leq 2 \overline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right) \quad \text { and } \quad 2 \underline{\mathrm{~h}}\left(c, \mathcal{P}_{0}\right) \leq\left|F_{p}^{c}\left(Y_{p}^{c}\right)\right| \leq 2 \overline{\mathrm{~h}}\left(c, \mathcal{P}_{0}\right)
$$

Definition $4.5\left(p_{m}\right)$. For every $m \in \mathbb{N}^{*}$, for every $s=\left(a_{j}\right)_{j \geq 0} \in \overrightarrow{\mathscr{A}}$ and $u=\left(a_{j}\right)_{j<0} \in \overleftarrow{\mathscr{A}}$, let $\mathcal{R}_{m}(s):=a_{0} \cdots a_{m-1}$ be the concatenation of the $m$ first terms of $s$ and let $p_{m}(u)=a_{-m} \cdots a_{-1}$ be the concatenation of the $m$ last terms of $u$.

Similarly as for Remark 4.4, it holds:
Lemma 4.6. For all $x_{*} \in I, p \in \mathcal{P}_{0}$ and $j>0$, it holds:
(1) For any $u \in \overleftarrow{\mathscr{A}}$, the push forward by $F_{p}^{p_{j}(u)}$ of $I \times\{0\}$ is $2 \overline{\mathrm{~h}}\left(p_{j}(u)\right)-C^{0}$ - close to $W_{p}^{u}$.
(2) For any $s \in \overrightarrow{\mathscr{A}}$, the pull back by $F_{p}^{\boldsymbol{p}_{j}(u)}$ of $\left\{x_{*}\right\} \times I$ is $2 \cdot \overline{\mathrm{w}}\left(\boldsymbol{p}_{j}(\jmath)\right)-C^{0}$-close to $W_{p}^{s}$.

From now on, we will suppose that $\left(F_{p}, \pi_{p}\right)_{p}$ is moderately dissipative and we fix a parameter $p_{0} \in \operatorname{int} \mathcal{P}$.
The following will be used many times and proved in Appendix B:

Proposition 4.7 (and definition of $\eta_{0}$ and $\mathcal{P}_{0}$ ). There exists $\eta_{0}>0$ and a compact neighborhood $\mathcal{P}_{0} \Subset \mathcal{P}$ of $p_{0} \in \mathcal{P}$ satisfying the following properties for every $c \in \mathscr{A}^{*}$ s.t. $\overline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right)<\eta_{0}$ :
(1) $\underline{\mathrm{w}}(c) \leq \overline{\mathrm{w}}(c)<\underline{\mathrm{w}}(c)^{1-\epsilon^{2}}, \quad$ with $\underline{\mathrm{w}}(c):=\underline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right)$ and $\overline{\mathrm{w}}(c):=\overline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right)$,
(2) $\overline{\mathrm{h}}(c)^{\epsilon} \leq \underline{\mathrm{w}}(c), \quad$ with $\overline{\mathrm{h}}(c):=\overline{\mathrm{h}}\left(c, \mathcal{P}_{0}\right)$.

If at some $\tilde{p} \in \mathcal{P}_{0}$ the map $F_{\tilde{p}}$ is analytic, then for $\rho$ sufficiently small, for every $c \in \mathscr{A}^{*}$, the hyperbolic transformation $\left(Y_{\tilde{p}}^{c}, F_{\tilde{p}}^{c}\right)$ extends to a $C_{\rho}^{\omega}$-hyperbolic transformation whose implicit representation $\left(\mathcal{Z}_{\tilde{p}}^{c}, \mathcal{W}_{\tilde{p}}^{c}\right)$ satisfies:
(2') $\max _{\tilde{I}^{2}}\left|\partial_{w_{0}} \mathcal{W}_{\tilde{p}}^{c}\right|^{\epsilon} \leq \min _{\tilde{I}^{2}}\left|\partial_{z_{1}} \mathcal{Z}_{\tilde{p}}^{c}\right| \quad$ if $\quad \overline{\mathrm{w}}\left(c, \mathcal{P}_{0}\right)<\eta_{0}$.
The following will be used many times.
Lemma 4.8 (and definition of $\eta_{1}$ and $\left.\epsilon^{\prime}\right)$. There exists $\eta_{1} \in\left(0, \eta_{0}\right)$ and $\epsilon^{\prime}>0$ such that for every $s \in \overrightarrow{\mathscr{A}}, u \in \overleftarrow{\mathscr{A}}$, for every $\delta \in\left(0, \eta_{1}\right)$, if $m \geq 1$ and $m^{\prime} \geq 1$ are maximal such that:

$$
\overline{\mathrm{w}}\left(p_{m}(s)\right)>\delta \quad \text { and } \quad \overline{\mathrm{h}}\left(p_{m^{\prime}}(u)\right)>\delta
$$

for every $d \in \mathscr{A}^{*}$ with $\underline{\mathrm{w}}(d) \geq \delta^{\epsilon^{4}}$ and such that $c:=p_{m+1}(s) \cdot d \cdot \mathcal{p}_{m^{\prime}+1}(u) \in \mathscr{A}^{*}$, it holds:

$$
\delta \geq \overline{\mathrm{w}}(c) \geq \underline{\mathrm{w}}(c) \geq \delta^{1+\epsilon+2 \epsilon^{2}}
$$

Moreover such $d \in \mathscr{A}^{*}$ always exist.
Proof. By maximality of $m$ and $m^{\prime}$, the maximal widths of $p_{m}(s)$ and $p_{m^{\prime}}(u)$ are small when $\eta_{1}$ is small, and sufficiently to apply Proposition 4.7 to these words. Also by maximality, it holds $\delta \geq \overline{\mathrm{w}}\left(p_{m+1}(s)\right)$ and $\delta \geq \overline{\mathrm{h}}\left(p_{m^{\prime}+1}(u)\right)$. Thus:

$$
\delta \geq \overline{\mathrm{w}}(c) \geq \underline{\mathrm{w}}(c) \quad \text { and } \quad \delta \geq \overline{\mathrm{h}}(c) \geq \underline{h}(c)
$$

It remains to prove that $\underline{\mathrm{w}}(c) \geq \delta^{1+\epsilon+2 \epsilon^{2}}$. By Lemma B. 1 we have:

$$
\underline{\mathrm{w}}(c) \geq \underline{\mathrm{w}}\left(\mathcal{p}_{m}(s)\right) \cdot \min _{a} \underline{\mathrm{w}}(a) \cdot \underline{\mathrm{w}}(d) \cdot \min _{a} \underline{\mathrm{w}}(a) \cdot \underline{\mathrm{w}}\left(\mathcal{R}_{m^{\prime}}(u)\right) \cdot \frac{1}{\left(1+\theta^{2}\right)^{4}} .
$$

We now use Proposition 4.7 (1) with $\underline{\mathrm{w}}\left(p_{m}(\jmath)\right)^{1-\epsilon^{2}}>\overline{\mathrm{w}}\left(p_{m}(\jmath)\right)>\delta$ to obtain:

$$
\underline{\mathrm{w}}(c)>\delta^{\left(1-\epsilon^{2}\right)^{-1}} \cdot\left(\min _{a} \underline{\mathrm{w}}(a)\right)^{2} \cdot \underline{\mathrm{w}}\left(p_{m^{\prime}}(u)\right) \cdot \underline{\mathrm{w}}(d) \cdot \frac{1}{\left(1+\theta^{2}\right)^{4}}
$$

By Proposition $4.7(2), \underline{\mathrm{w}}\left(p_{m^{\prime}}(u)\right) \geq \overline{\mathrm{h}}\left(p_{m^{\prime}}(u)\right)^{\epsilon}$. Also $\underline{\mathrm{w}}(d) \geq \delta^{\epsilon^{4}}$. Thus:

$$
\underline{\mathrm{w}}(c)>\delta^{\left(1-\epsilon^{2}\right)^{-1}} \cdot\left(\min _{a} \underline{\mathrm{w}}(a)\right)^{2} \cdot \delta^{\epsilon}\left(1+\theta^{2}\right)^{-4} \cdot \delta^{\epsilon^{4}}
$$

We now use the logarithm to compare this with the sought inequality:

$$
\begin{array}{r}
\log \left(\underline{\mathrm{w}}(c) \delta^{-1-\epsilon-2 \epsilon^{2}}\right)>\left(\left(1-\epsilon^{2}\right)^{-1}+\epsilon+\epsilon^{4}-1-\epsilon-2 \epsilon^{2}\right) \log \delta+2 \log \min _{a} \underline{\mathrm{w}}(a)-4 \log \left(1+\theta^{2}\right)  \tag{4.1}\\
=\left(\epsilon^{-4}\left(1-\epsilon^{2}\right)^{-1}+\epsilon^{4}-\epsilon^{2}\right) \log \delta+2 \log \min _{a} \underline{\mathrm{w}}(a)-4 \log \left(1+\theta^{2}\right)
\end{array}
$$

Note that $\left(\epsilon^{-4}\left(1-\epsilon^{2}\right)^{-1}+\epsilon^{4}-\epsilon^{2}\right)<0$, hence for $\eta_{1}>\delta$ sufficiently small, the right most sought inequality is obtained. Also, by connectedness of the graph, when $\delta$ is small, given any two vertexes in $\mathscr{V}$, there exists a word $d$ which joints them and so that $\underline{w}(d)>\delta^{\epsilon^{4}}$.

Remark 4.9. Let $0<\mathrm{W} \leq 1$ and $\mathcal{L}$ be numbers such that every pair of vertexes in $\mathscr{V}$ can be connected by a path $d$ such that $\underline{\mathrm{w}}(d)>\mathrm{W}$ and $|d| \leq \mathcal{L}$. Then under the assumptions of Lemma 4.8, for every word $\boldsymbol{d}^{\circ} \in \mathscr{A}^{*}$ satisfying

$$
\underline{\mathrm{w}}\left(d^{o}\right)>\delta^{\epsilon^{4}} \cdot \frac{1+\theta^{2}}{\mathrm{~W}} \quad \text { and } \quad a\left(d^{o}\right)=t\left(p_{m+1}(s)\right)
$$

there exists $d^{r} \in \mathscr{A}^{*}$ such that $\underline{\mathrm{w}}\left(d^{r}\right)>\mathrm{W},\left|d^{r}\right| \leq \mathcal{L}$ and we can assume that $d=d^{o} \cdot d^{r}$.
The following provides some uniform bounds on the unfolding.
Proposition 4.10 (and definition of $\vartheta$ and $L_{0}$ ). There exist $\vartheta>0$ and $L_{0}>0$ such that for every $p^{\prime} \in \mathcal{P}_{0}, m:=\left(u_{\square}, s_{\square}\right)_{\square \in \mathscr{C}} \in \prod_{\square \in \mathscr{C}} \overleftarrow{\mathscr{A}}_{\square} \times \overrightarrow{\mathscr{A}}_{\square}$, the map $\Psi_{m}: p \in \mathcal{P} \mapsto\left(\mathcal{V}_{p}^{u_{\square}, s_{\square}}\right)_{i}$ is $\frac{L_{0}}{2}$-biLipschitz from a neighborhood of $p^{\prime}$ in $\mathcal{P}$ onto the $2 \vartheta$-neighborhood of $\Psi_{m}\left(p^{\prime}\right)$.

Proof. Let $\alpha>0$ be such that for every $p^{\prime} \in \mathcal{P}$, the euclidean closed ball $B\left(p^{\prime}, \alpha\right)$ is included in int $\mathcal{P}$. Let $m=\left(u_{\square}, s_{\square}\right)_{\square \in \mathscr{C}} \in \prod_{\square \in \mathscr{C}} \overleftarrow{\mathscr{A}}_{\square} \times \overrightarrow{\mathscr{A}}_{\square}$. By $\left(\mathbf{H}_{2}\right)$, the restriction $\Psi_{m} \mid B\left(p^{\prime}, \alpha\right)$ is a diffeomorphism onto its image. By compactness of $B\left(p^{\prime}, \alpha\right)$, the restriction $\Psi_{m} \mid B\left(p^{\prime}, \alpha\right)$ is biLipschitz; let $L\left(p^{\prime}, m\right) / 2>0$ be its bi-Lipschitz constant. Also let $\rho\left(p^{\prime}, m\right)>0$ be maximal such that $\Psi_{m}\left(B\left(p^{\prime}, \alpha\right)\right)$ contains the $2 \rho\left(p^{\prime}, m\right)$-neighborhood of $\Psi_{m}\left(p^{\prime}\right)$. By Proposition 1.42 , the functions $\rho$ and $L$ are continuous on the compact set $\mathcal{P} \times \prod_{\square \in \mathscr{C}} \overleftarrow{\mathscr{A}}_{\square} \times \overrightarrow{\mathscr{A}}_{\square}$ and so bounded from below by positive constants $\vartheta$ and $L_{0}$.
4.2. Dynamics at the folding map. In this subsection, we continue to consider a moderately dissipative unfolding $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of wild type $(\mathscr{A}, \mathscr{C})$ and a fixed parameter $p_{0} \in \operatorname{int} \mathcal{P}$. The compact neighborhood $\mathcal{P}_{0} \Subset \mathcal{P}$ of $p_{0}$ was defined by Proposition 4.7, whereas $\vartheta, L_{0}$ and $\eta_{1}$ were defined in Proposition 4.10 and Lemma 4.8.

To study the folding map $F_{p}^{\square} \mid Y_{p}^{\boxminus}$ for any $\boxtimes \in \mathscr{C}$, we will use analogous functions to $\mathcal{V}_{p}^{u, s}=$ $a_{p}^{u}-b_{p}^{s}$ defined in Definition 1.40, when $u$ and/or $s$ are replaced by finite words $c \in \mathscr{A}^{*}$. The definitions of the counterparts of $a_{p}^{u}$ and $b_{p}^{s}$ use the following counterparts of $W_{p}^{u}$ or $W_{p}^{s}$ :

Definition $4.11\left(H_{p}^{c}\right.$ and $\left.V_{p}^{c}(\zeta)\right)$. For every $p \in \mathcal{P}_{0}$, let $H_{p}^{c}$ be the push forward by $F_{p}^{c}$ of $I \times\{0\}$ and given a point $\zeta \in Y^{e}$, let $V_{p}^{c}(\zeta)$ be the pullback by $F_{p}^{c}$ of the $\pi_{p}^{t(c)}$-fiber containing $\zeta$ :

$$
H_{p}^{c}=F_{p}^{c}\left(Y^{c} \cap I \times\{0\}\right) \quad \text { and } \quad V_{p}^{c}(\zeta)=\left(\pi_{p}^{t(c)} \circ F_{p}^{c}\right)^{-1}\left(\left\{\pi_{p}^{t(c)}(\zeta)\right\}\right)
$$

The following is an immediate consequence of Definition 1.23(3):
Fact 4.12. The curve $V_{p}^{c}(\zeta)$ is the fiber of $\pi_{p}^{o(c)}$ containing $\left(F_{p}^{c}\right)^{-1}(\zeta)$.
We now introduce the analogous of Proposition 1.28 for a parameter family.
Proposition 4.13 (and definition of $\breve{\mathscr{A}}_{\square}^{*}, \zeta_{p}^{c}$ and $\left.\eta_{2}\right)$. There exists $\eta_{2} \in\left(0, \eta_{1}\right)$ such that for every $\square \in \mathscr{C}, p \in \mathcal{P}_{0}$ and $c \in \mathscr{A}^{*}$ in the set:

$$
\breve{\mathscr{A}}_{\square}^{*}:=\left\{c \in \mathscr{A}^{*}: c \text { coincides with the }|c| \text { last letters of a certain } u \in \overleftarrow{\mathscr{A}}_{\square} \quad \text { and } \quad \overline{\mathrm{h}}(c)<\eta_{2}\right\}
$$

the curve $H_{p}^{c}$ intersects $Y_{p}^{\boxminus}$ but not $\partial^{u} Y_{p}^{\boxminus}$, the curve $H_{p}^{c}$ is tangent to a unique fiber of $\pi_{p}^{t(\boxminus)} \circ F_{p}^{\boxminus}$ at a unique point $\zeta_{p}^{c} \in \operatorname{int} Y_{p}^{\bullet}$, and the tangency is quadratic.

Proof. First observe that when $\eta_{2}$ is small, the minimal length of the words in $\breve{\mathscr{A}}_{\square}^{*}$ is large. Thus, the proof of this Proposition is the same as the one of Proposition 1.28 but for a parameter family this time. Indeed, note that $\overleftarrow{\mathscr{A}}_{\square}$ does not depend on $p$ for every $\boxtimes \in \mathscr{C}$ (by Definition 1.36), so it suffices to replace the Inclination Lemma and the condition on $|c|$ by respectively the parametric Inclination Lemma C. 6 Page 72 and the condition on $\overline{\mathrm{h}}(c)$. By finiteness of $\mathscr{C}$, one can take the same value $\eta_{2}$ for every $\square \in \mathscr{C}$.

Remark 4.14. Up to reducing $\eta_{2}$, we can assume that $\mathscr{A}_{\square}^{*}\left(F_{p}\right) \supset \breve{\mathscr{A}}_{\square}^{*}$ for any $p \in \mathcal{P}_{0}$, where $\mathscr{A}_{\square^{*}}^{*}\left(F_{p}\right)$ was defined in Proposition 1.28.

Definition 4.15. For every $\square \in \mathscr{C}$ and $c \in \breve{A}_{\square}^{*}$, we denote:

$$
a_{p}^{c}:=\pi_{p}^{t(\boxminus)} \circ F_{p}^{\square}\left(\zeta_{p}^{c}\right) \in I \quad \text { and } \quad b_{p}^{c}:=\pi_{p}^{o(c)} \circ\left(F_{p}^{c}\right)^{-1}\left(\zeta_{p}^{c}\right)
$$

The following compares $a_{p}^{c}$ and $b_{p}^{c}$ to the limit case $a_{p}^{u}$ and $b_{p}^{s}$ :
Proposition 4.16. For every $\square \in \mathscr{C}, c \in \breve{\mathscr{A}}_{\square}^{*}$ and $p \in \mathcal{P}_{0}$, it holds:
(1) For every $s \in \overrightarrow{\mathscr{A}}$ starting with $c$, it holds $\left|b_{p}^{c}-b_{p}^{s}\right| \leq 2 \cdot \overline{\mathrm{w}}(c)$.
(2) For every $u \in \overleftarrow{\mathscr{A}}$ ending with $c$, it holds $\left|a_{p}^{c}-a_{p}^{u}\right| \leq\left\|D\left(\pi_{p}^{t(\boxminus)} \circ F_{p}^{\boxminus}\right)\right\|_{C^{0}} \cdot 2 \cdot \overline{\mathrm{~h}}(c)$

Proof. We first prove (1). Let $s^{\prime} \in \overrightarrow{\mathscr{A}}$ be such that $s=c \cdot s^{\prime}$. We recall that $\zeta_{p}^{c} \in Y_{p}^{\square}$ for every $p$. Thus by Definition $1.23(0)$ and (1), the $\pi_{p}^{t(c)}$-fiber of $\zeta_{p}^{c}$ is a curve with tangent space in $\chi_{v}$, included in $Y^{e}$ and with endpoint in and endpoints in $I \times \partial I$. These properties also hold true for $W_{p}^{s^{\prime}}$ (by Proposition 1.25 for $W_{p}^{s^{\prime}}$ ). By Definition 1.4 (3), their respective pullbacks $V_{p}^{c}\left(\zeta_{p}^{c}\right)$ and $W_{p}^{s}$ by $F_{p}^{c}$ are then also graphs over $y \in I$ with direction in $\chi_{v}$ and endpoints in $I \times \partial I$, which are moreover included in $Y_{p}^{c}$. They are then $\left|Y_{p}^{c}\right|-C^{0}$-close and so $2 \overline{\mathrm{w}}(c)$ - $C^{0}$-close by Remark 4.4. By Definition 1.23 (3), they are also fibers of $\pi_{p}^{\circ(c)}$. By Definition 1.23-(2), their images by $\pi_{p}^{\circ(c)}$ (equal to $b_{p}^{c}$ and $b_{p}^{s}$ ) are the $x$-coordinates of the intersection points of these fibers with $I \times\{0\}$. This implies $\left|b_{p}^{c}-b_{p}^{s}\right| \leq 2 \cdot \overline{\mathrm{w}}(c)$.

Let us prove the second item. Similarly to the case of vertical curves, we can show that $H_{p}^{c} \cap$ $Y_{p}^{\boxminus} \subset H_{p}^{c}$ and $W_{p}^{u} \cap Y_{p}^{\boxminus} \subset W_{p}^{u}$ are $2 \overline{\mathrm{~h}}(c)$ - $C^{0}$-close. Their respective images by $\pi_{p}^{t(\boxminus)} \circ F_{p}^{\complement}$ are consequently $\left\|D\left(\pi_{p}^{t(\square)} \circ F_{p}^{\square}\right)\right\|_{C^{0}} \cdot 2 \overline{\mathrm{~h}}(c)$-close. Thus the endpoints of $\pi_{p}^{t(\boxminus)} \circ F_{p}^{\boxminus}\left(H_{p}^{c} \cap Y_{p}^{\square}\right)$ and $\pi_{p}^{t(\boxminus)} \circ F_{p}^{\unlhd}\left(W_{p}^{u} \cap Y_{p}^{\square}\right)$ are $2 \cdot\left\|D\left(\pi_{p}^{t(\boxminus)} \circ F_{p}^{\boxminus}\right)\right\|_{C^{0}} \cdot \overline{\mathrm{~h}}(c)$-close. We conclude by noting that $a_{p}^{c}$ and $a_{p}^{u}$ are endpoints of these segments.
Definition 4.17. For all $\boxtimes, \square^{\prime} \in \mathscr{C}, s \in \overrightarrow{\mathscr{A}}_{\square}, c \in \breve{\mathscr{A}}_{\square}^{*}$ and $c^{\prime} \in \breve{\mathscr{A}}_{\square}^{*}$, such that $t(\square)=\sigma\left(c^{\prime}\right)$, for every $p \in \mathcal{P}_{0}$, we put:

$$
\mathcal{V}_{p}^{c, c^{\prime}}:=a_{p}^{c}-b_{p}^{c^{\prime}} \quad \text { and } \quad \mathcal{V}_{p}^{c, s}:=a_{p}^{c}-b_{p}^{s}
$$

Remark 4.18. We notice that $\mathcal{V}_{p}^{c, s}=0$ if and only if $F_{p}^{\bullet}\left(H_{p}^{c} \cap Y_{p}^{\bullet}\right)$ is tangent to $W_{p}^{s}$. Similarly, $\mathcal{V}_{p}^{c, c^{\prime}}=0$ if and only if $F_{p}^{\square}\left(H_{p}^{c} \cap Y_{p}^{\boxminus}\right)$ is tangent to $V_{p}^{c^{\prime}}\left(\zeta_{p}^{c^{\prime}}\right)$. We notice that the tangency holds at $F_{p}^{\oplus}\left(\zeta_{p}^{c}\right)$.

From now on we work under the assumptions of Theorem 1.45. In other words, the unfolding $\left(F_{p}, \pi_{p}\right)_{p}$ of wild type $(\mathscr{A}, \mathscr{C})$ is moderately dissipative and $\# \mathscr{C}=5$.

Now we are ready to define the counterpart of the function involved in $\left(\mathbf{H}_{2}\right)$. We set:

$$
\mathscr{C}=:\left\{\square_{1}, \square_{2}, \square_{3}, \square_{4}, \sqcup_{5}\right\}
$$

Given $i \in \mathbb{Z}$, let $\square_{i}:=\square_{[i]}$ where $[i]$ is the equivalence class of $i$ in $\mathbb{Z} / 5 \mathbb{Z}$. In particular $\square_{i+5}=\square_{i}$ for every $i \in \mathbb{Z}$.

Definition 4.19. For all $i_{0} \geq 0$ and $\left(c_{i}\right)_{i_{0} \leq i \leq i_{0}+4} \in \prod_{i=i_{0}}^{i_{0}+4} \breve{\mathscr{A}}_{\square_{i}}^{*}$ and $s \in \overrightarrow{\mathscr{A}}_{Ð_{i_{0}+4}}$ such that $t\left(\square_{i}\right)=a\left(c_{i+1}\right)$ for every $i_{0} \leq i \leq i_{0}+3$, we put:

$$
\Psi_{\left(c_{i}\right)_{i}, s}: p \in \mathcal{P}_{0} \mapsto\left(\left(\mathcal{V}_{p}^{c_{i}, c_{i+1}}\right)_{i_{0} \leq i \leq i_{0}+3}, \mathcal{V}_{p}^{c_{i_{0}+4, s}}\right)
$$

Here is the counterpart of Proposition 4.10:
Proposition 4.20 (and definition of $\left.\eta_{3}\right)$. There exists $\eta_{3} \in\left(0, \eta_{2}\right)$ such that for every $\left(c_{i}\right)_{i_{0} \leq i \leq i_{0}+4} \in$ $\prod_{i=i_{0}}^{i_{0}+4} \breve{\mathscr{A}}_{\square_{i}}^{*}$ such that $t\left(\square_{i}\right)=\sigma\left(c_{i+1}\right)$ for every $i_{0} \leq i \leq i_{0}+3$ and $c_{i}$ satisfies $\overline{\mathrm{w}}\left(c_{i}\right)<\eta_{3}$ for every $i_{0} \leq i \leq i_{0}+4$, the following property holds true.

For every $p_{1} \in \mathcal{P}_{0}$ and $s_{i_{0}+5} \in \overrightarrow{\mathscr{A}}_{\oplus_{i_{0}+4}}$, the restriction of $\Psi_{\left(c_{i}\right)_{i}, s_{i_{0}+5}}$ to a neighborhood of $p_{1}$ is a $L_{0}$-biLipschitz map onto the $\vartheta$-neighborhood of $\left.\Psi_{\left(c_{i}\right)_{i}, s_{i_{0}+5}}, p_{1}\right)$.
Proof. Let us take for each $i_{0} \leq i \leq i_{0}+4$ a sequence $u_{i} \in \overleftarrow{\mathscr{A}}_{\square_{i}}$ such that the last letters of $u_{i}$ are equal to $c_{i}$. By Proposition 4.7, we have $\overline{\mathrm{h}}\left(c_{i}\right)<\overline{\mathrm{h}}\left(c_{i}\right)^{\epsilon} \leq \underline{\mathrm{w}}\left(c_{i}\right)<\eta_{3}<\eta_{2}<\eta_{0}$. Let $K \geq 1$ be greater than the supremum of the $C^{1}$-norms of $p \in \mathcal{P}_{0} \mapsto F_{p}^{\boxminus}\left(\zeta_{p}^{u}\right)$ among $\square \in \mathscr{C}$ and $u \in \overleftarrow{\mathscr{A}}_{\square}$ By finiteness of $\mathscr{C}$ and compactness of $\mathcal{P}_{0}$ and $\overleftarrow{\mathscr{A}}_{\square}$, the constant $K$ is finite. By the parametric Inclination Lemma C. 6 Page 72 when $\eta_{3}$ is small, the following pairs of maps are $C^{1}$-close:

$$
p \in \mathcal{P}_{0} \mapsto \zeta_{p}^{c_{i}} \quad \text { and } \quad p \in \mathcal{P}_{0} \mapsto \zeta_{p}^{u_{i}} \quad \forall i_{0} \leq i \leq i_{0}+4
$$

By Proposition 1.42 and compactness of $\bigcup_{\boxminus \in \mathscr{C}} \overleftarrow{\mathscr{A}}_{\square}$, the maps $p \in \mathcal{P}_{0} \mapsto \zeta_{p}^{u_{i}}$ are $C^{1}$-uniformly bounded among all upplets $\left(u_{i}\right)_{i_{0} \leq i \leq i_{0}+4}$ of sequences $u_{i} \in \overleftarrow{\mathscr{A}}_{\oplus_{i}}$ and so also among all $\left(u_{i}\right)_{i_{0} \leq i \leq i_{0}+4}$ defined as before. Then the maps $p \in \mathcal{P}_{0} \mapsto \zeta_{p}^{c_{i}}$ are $C^{1}$-uniformly bounded among all upplets $\left(c_{i}\right)_{i}$ satisfying the hypotheses of the proposition. Also by looking at their images by $\pi_{p}^{t\left(\varpi_{i}\right)} \circ F_{p}^{Ð_{i}}$, the following pairs of maps are $C^{1}$-close, when $\eta_{3}>0$ is small:

$$
p \in \mathcal{P}_{0} \mapsto a_{p}^{c_{i}} \quad \text { and } \quad p \in \mathcal{P}_{0} \mapsto a_{p}^{u_{i}} \quad \forall i_{0} \leq i \leq i_{0}+4
$$

Let us take for each $i_{0} \leq i \leq i_{0}+4$ a sequence $s_{i} \in \overrightarrow{\mathscr{A}}$ such that the first letters of $s_{i}$ are equal to $c_{i}$. As each map $p \mapsto \zeta_{p}^{c_{i}}$ is uniformly $C^{1}$-bounded, the family of the $\pi_{p}^{t\left(c_{i}\right)}$-fibers going through $\zeta_{p}^{c_{i}}$ indexed by $p \in \mathcal{P}_{0}$ is a $C^{1}$-bounded family, uniformly transverse to the unstable lamination of the horseshoe $\Lambda_{p}$. Applying a last time the parametric inclination Lemma C. 6 Page 72, we obtain that its pull back by $\left(F_{p}^{c_{i}}\right)_{p}$ is $C^{1}$-close to $\left(W_{p}^{s_{i}}\right)_{p \in \mathcal{P}_{0}}$, and is equal by definition to $\left(V_{p}^{c_{i}}\left(\zeta_{p}^{c_{i}}\right)\right)_{p}$. By Fact 4.12, the image by $\left(\pi_{p}^{o\left(c_{i}\right)}\right)_{p}$ of this family of curves is $\left(\pi_{p}^{o\left(c_{i}\right)} \circ\left(F^{c_{i}}\right)^{-1}\left(\zeta_{p}^{c_{i}}\right)\right)_{p}=:\left(b_{p}^{c_{i}}\right)_{p}$ which is therefore uniformly $C^{1}$-close to $\left(b_{p}^{s_{i}}\right)_{p}$. In other words, the following pairs of maps are $C^{1}$-close, when $\eta_{3}>0$ is small:

$$
p \in \mathcal{P}_{0} \mapsto b_{p}^{c_{i}} \quad \text { and } \quad p \in \mathcal{P}_{0} \mapsto b_{p}^{s_{i}} \quad \forall i_{0} \leq i \leq i_{0}+4
$$

Consequently when $\eta_{3}>0$ is small, the map $\Psi_{\left(c_{i}\right)_{i}, j_{i}+5}$ is uniformly $C^{1}$-close to $p \in \mathcal{P}_{0} \mapsto$ $\left(\mathcal{V}_{p}^{u_{i}, \dot{J}_{i+1}}\right)_{i}$ among all the $\left(c_{i}\right)_{i}$. Thus Proposition 4.10 implies the proposition by the inverse local theorem and compactness of the spaces involved.
4.3. Infinite chain of nearly heteroclinic tangencies. In this subsection, we work under the assumptions of Theorem 1.45. We consider a moderately dissipative unfolding $\left(F_{p}, \pi_{p}\right)_{p}$ of wild type $(\mathscr{A}, \mathscr{C})$ such that $\mathscr{C}$ is formed by five elements $\left\{\square_{1}, \square_{2}, \square_{3}, \square_{4}, \square_{5}\right\}$. We fix $p_{0} \in \operatorname{int} \mathcal{P}$ as in the previous subsection.

The parameter selection is given by Proposition 4.23 below and proved along this subsection. Its statement involves the following constants which are positive and finite:

$$
\begin{equation*}
C_{1}:=2 \sqrt{5} \cdot L_{0} \cdot \sup _{p \in \mathcal{P}_{0}}\left(\left\|D\left(\pi_{p} \circ F_{p}^{\mathscr{C}}\right)\right\|+1\right) \quad \text { and } \quad \eta_{4}:=\min \left(L_{0} \vartheta / C_{1}, \eta_{3}\right) \tag{4.2}
\end{equation*}
$$

where the compact subset $\mathcal{P}_{0} \subset \mathcal{P}$ was defined by Proposition 4.7 , the constant $\eta_{3}$ was defined in Proposition 4.20, whereas the constants $\vartheta$ and $L_{0}$ were defined in Proposition 4.10.

We are going to define a sequence of parameters $\left(p_{i}\right)_{i \geq 1}$ which converges to a certain $p_{\infty}$ such that each $p_{i}$ is close to $p_{i-1}$ and at $p_{i}$ the map $F_{p_{i}}$ satisfy some tangency's conditions. To this end we will ask that the distance between $p_{i}$ and $p_{i+1}$ is at most of the order $\delta^{\beta^{i}}$ for a certain $\delta$ small and a fixed $\beta>1$. Hence the following immediate fact will be useful to ensure that all the $p_{i}$ remain in $\mathcal{P}_{0}$ (at which all our bounds were done):

Fact 4.21. For every $\beta>1$, there exists $M_{\beta}>0$ independent of $1 / 2>\delta>0$ such that $\sum_{j \geq 0} \delta^{\beta^{j}} \leq$ $M_{\beta} \cdot \delta$ for every $1 / 2>\delta>0$.

This fact enables us to fix the following constant:
Definition $4.22\left(\beta\right.$ and $\left.\left(\delta_{j}\right)_{j}\right)$. Let $\beta=3 / 2$ and let $1 / 2>\delta_{0}>0$ be any number smaller than $\eta_{4}$ such that the ball centered at $p_{0}$ with radius $C_{1} \cdot M_{\beta} \cdot \delta_{0}$ is included in $\mathcal{P}_{0}$. We put

$$
\begin{equation*}
\delta_{j}:=\delta_{0}^{\beta^{j}} \quad \text { with } \beta=\frac{3}{2} . \tag{4.3}
\end{equation*}
$$

We notice that the sequence $\left(\delta_{j}\right)_{j}$ decreases super exponentially fast to 0 . Also it satisfies the following for every $i, m \geq 0$ :

$$
\begin{equation*}
\delta_{i+m}=\left(\delta_{i}\right)^{\beta^{m}} \quad \text { and } \quad \sum_{j \geq i} \delta_{j} \leq M_{\beta} \cdot \delta_{i} \tag{4.4}
\end{equation*}
$$

We recall that for every $i \in \mathbb{Z}$, we put $\square_{i}=\square_{[i]}$, where $[i]$ is the equivalent class of $i$ in $\mathbb{Z} / 5 \mathbb{Z}$. The following involves the curves $H_{p}^{c}$ and $V_{p}^{c}(\zeta)$ introduced in Definition 4.11.

Proposition 4.23 (Main). There exist a sequence of words $\left(c_{i}\right)_{i \geq 1} \in \prod_{i \geq 1} \breve{\mathscr{A}}_{\square_{i}}^{*}$ and a sequence of parameters $\left(p_{i}\right)_{i \geq 1} \in \mathcal{P}_{0}^{\mathbb{N}}$, such that it holds $t\left(\square_{i}\right)=0\left(c_{i+1}\right)$ and the following for any $i \geq 1$ :
$\left(\mathbf{C}_{1}\right) \delta_{i}^{1+\epsilon+2 \epsilon^{2}} \leq \underline{\mathrm{w}}\left(c_{i}\right) \leq \overline{\mathrm{w}}\left(c_{i}\right) \leq \delta_{i}$,
$\left(\mathbf{C}_{2}\right)$ the map $F_{p_{i}}^{\square_{i}}$ sends $H_{p_{i}}^{c_{i}} \cap Y_{p_{i}}^{\square_{i}}$ tangent to $V_{p_{i}}^{c_{i+1}}\left(\zeta_{p_{i}}^{c_{i+1}}\right)$ at the point $F_{p_{i}}^{\unlhd_{i}{ }_{i}}\left(\zeta_{p_{i}}^{c_{i}}\right)$,
$\left(\mathbf{C}_{3}\right)\left\|p_{0}-p_{1}\right\|<C_{1} \delta_{1}$ and $\left\|p_{i}-p_{i-1}\right\|<C_{1} \delta_{i+4}$ whenever $i \geq 2$.
We will prove this proposition below. From this we deduce the following result involving the constants $L_{0}, M_{\beta}$ and $C_{1}$ defined respectively in Proposition 4.10, Fact 4.21 and Eq. (4.2) Page 57:

Corollary 4.24. There exist $p_{\infty} \in \mathcal{P}_{0}$ which is $C_{1} M_{\beta} \delta_{0}$-close to $p_{0}$ and a sequence of words $\left(c_{i}\right)_{i \geq 1} \in \prod_{i \geq 1} \breve{\mathscr{A}}_{\square_{i}}^{*}$ such that for any $i \geq 1$ it holds $t\left(\square_{i}\right)=\sigma\left(c_{i+1}\right)$ and:
$\left(\mathbf{C}_{1}\right) \quad \delta_{i}^{1+\epsilon+2 \epsilon^{2}} \leq \underline{\mathrm{w}}\left(c_{i}\right) \leq \overline{\mathrm{w}}\left(c_{i}\right) \leq \delta_{i}$,
$\left(\mathbf{C}_{2}^{\prime}\right) \quad\left|\mathcal{V}_{p_{\infty}}^{c_{i}, c_{i+1}}\right| \leq C_{1} L_{0} M_{\beta} \delta_{i+5}$.
Proof. We take the sequence of words $\left(c_{i}\right)_{i \geq 1}$ and the sequence of parameters $\left(p_{i}\right)_{i \geq 1}$ given by Proposition 4.23. In particular, they satisfy $\left(\mathbf{C}_{1}\right)$. By Conclusion $\left(\mathbf{C}_{3}\right)$ of Proposition 4.23, we have $\left\|p_{i+1}-p_{i}\right\|<C_{1} \delta_{i+5}$ whenever $i \geq 1$ with $\left(\delta_{i}\right)_{i \geq 1}$ decreasing super exponentially fast to 0 . Thus the sequence $\left(p_{i}\right)_{i \geq 1}$ converges to a parameter $p_{\infty} \in \mathbb{R}^{5}$. By Eq. (4.4), it holds:

$$
\left\|p_{0}-p_{\infty}\right\| \leq\left\|p_{1}-p_{0}\right\|+\sum_{i \geq 1}\left\|p_{i+1}-p_{i}\right\|<C_{1} \delta_{1}+C_{1} \sum_{i \geq 1} \delta_{i+5} \leq C_{1} M_{\beta} \delta_{0}
$$

By Definition 4.22, this implies that $p_{\infty}$ is in $\mathcal{P}_{0}$. Still using Conclusion $\left(\mathbf{C}_{3}\right)$ of Proposition 4.23 and then Eq. (4.4), we have $\left\|p_{i}-p_{\infty}\right\| \leq C_{1} M_{\beta} \delta_{i+5}$ for $i \geq 1$. We notice that $\mathcal{V}_{p_{i}}^{c_{i}, c_{i+1}}=0$ by Conclusion $\left(\mathbf{C}_{2}\right)$ of Proposition 4.23. Then $\left(\mathbf{C}_{2}^{\prime}\right)$ follows from Proposition 4.20.

In Section 4.4, we will show that $F_{p_{\infty}}$ displays a stable wandering domain by showing that the conclusions of Corollary 4.24 imply the assumptions of Theorem 1.32. Proposition 4.23 will be proved by induction. The first step will be implied by the following:

Lemma 4.25. There exist $\left(c_{i}\right)_{1 \leq i \leq 5} \in \prod_{i=1}^{5} \breve{\mathscr{A}}_{\square_{i}}^{*}$ with $t\left(\square_{i}\right)=a\left(c_{i+1}\right)$ for any $1 \leq i \leq 4$, $s_{6} \in \overrightarrow{\mathscr{A}}_{\square_{5}}$ and a parameter $p_{1} \in \mathcal{P}_{0}$ which is $C_{1} \delta_{1}$-close to $p_{0}$ such that:
$\left(\mathbf{C}_{1}\right)$ For every $1 \leq i \leq 5$, it holds $\delta_{i}^{1+\epsilon+2 \epsilon^{2}} \leq \underline{w}\left(c_{i}\right) \leq \overline{\mathrm{w}}\left(c_{i}\right) \leq \delta_{i}$.
$\left(\mathbf{C}_{2}^{\prime \prime}\right)$ For any $1 \leq i \leq 4$, it holds $\mathcal{V}_{p_{1}}^{c_{i}, c_{i+1}}=0$ and $\mathcal{V}_{p_{1}}^{c_{5}, s_{6}}=0$.
Proof. By $\left(\mathbf{H}_{\mathbf{1}}\right)$ of Definition 1.43, there exist

$$
\begin{aligned}
& u_{1} \in \overleftarrow{\mathscr{A}}_{\oplus_{1}}, u_{2} \in \overleftarrow{\mathscr{A}}_{Ð_{2}}, u_{3} \in \overleftarrow{\mathscr{A}}_{Ð_{3}}, u_{4} \in \overleftarrow{\mathscr{A}}_{\oplus_{4}}, u_{5} \in \overleftarrow{\mathscr{A}}_{\square_{5}}
\end{aligned}
$$

such that:
$\left(T^{\prime}\right)$ for every $1 \leq j \leq 5$, the curve $F_{p_{0}}^{\square_{j}}\left(W_{p_{0}}^{u_{j}} \cap Y_{p_{0}}^{\square_{j}}\right)$ is tangent to $W_{p_{0}}^{\jmath_{j+1}}$.
We also pick any $s_{1} \in \overrightarrow{\mathscr{A}}_{\oplus_{0}}$. For every $1 \leq i \leq 5$, let $m_{i} \geq 1$ and $m_{i}^{\prime} \geq 1$ be maximal such that:

$$
\overline{\mathrm{w}}\left(p_{m_{i}}\left(s_{i}\right)\right)>\delta_{i} \quad \text { and } \quad \overline{\mathrm{h}}\left(p_{m_{i}^{\prime}}\left(u_{i}\right)\right)>\delta_{i}
$$

As $\delta_{i}<\eta_{4}<\eta_{1}$, by Lemma 4.8 , for $1 \leq i \leq 5$, there exists $d_{i} \in \mathscr{A}^{*}$ such that $c_{i}:=\mathcal{p}_{1+m_{i}}\left(\mathcal{s}_{i}\right)$. $d_{i} \cdot p_{1+m_{i}^{\prime}}\left(u_{i}\right) \in \mathscr{A}^{*}$ is well-defined and satisfies:

$$
\begin{equation*}
\delta_{i} \geq \overline{\mathrm{h}}\left(c_{i}\right) \geq \underline{h}\left(c_{i}\right) \quad \text { and } \quad \delta_{i} \geq \overline{\mathrm{w}}\left(c_{i}\right) \geq \underline{\mathrm{w}}\left(c_{i}\right) \geq \delta_{i}^{1+\epsilon+2 \epsilon^{2}} \tag{4.5}
\end{equation*}
$$

Moreover $c_{i} \in \breve{\mathscr{A}}_{\square_{i}}^{*}$ since $u_{i} \in \overleftarrow{\mathscr{A}}_{\square_{i}}$ and since $\bar{h}\left(c_{i}\right) \leq \delta_{i}<\eta_{4}<\eta_{2}$. We remark that $t\left(\square_{i}\right)=$ $a\left(s_{i+1}\right)=a\left(c_{i+1}\right)$. This proves $\left(\mathbf{C}_{1}\right)$.

Let us prove $\left(\mathbf{C}_{2}^{\prime \prime}\right)$. In Definition 4.19, we introduced:

$$
\left.\Psi_{\left(c_{i}\right)_{i}, \Delta_{6}}: p \in \mathcal{P}_{0} \mapsto\left(\left(\mathcal{V}_{p}^{c_{i}, c_{i+1}}\right)_{1 \leq i \leq 4}, \mathcal{V}_{p}^{c_{5}, s_{6}}\right)\right)
$$

By Proposition 4.16, the latter function is $2 \sqrt{5} \cdot \sup _{p \in \mathcal{P}_{0}}\left(\left\|D\left(\pi_{p} \circ F_{p}^{\mathscr{C}}\right)\right\|+1\right) \cdot \delta_{1}=C_{1} \delta_{1} / L_{0}-C^{0}-$ close to the following:

$$
\Psi_{\left(u_{i}, s_{i+1}\right)_{i}}: p \in \mathcal{P}_{0} \mapsto\left(\mathcal{V}_{p}^{u_{1}, \hat{s}_{2}}, \mathcal{V}_{p}^{u_{2}, \delta_{3}}, \mathcal{V}_{p}^{u_{3}, s_{4}}, \mathcal{V}_{p}^{u_{4}, \hat{s}_{5}}, \mathcal{V}_{p}^{u_{5}, \hat{s}_{6}}\right)
$$

where $\mathcal{V}_{p}^{u_{i}, s_{j}}$ was defined in Definition 1.40. By definition of $\left(u_{i}, s_{i+1}\right)_{i}$, this function vanishes at $p=p_{0}$. Thus it holds:

$$
\left\|\Psi_{\left(e_{i}\right)_{i}, ग_{6}}\left(p_{0}\right)\right\| \leq C_{1} \delta_{1} / L_{0}
$$

By Proposition 4.20, a restriction of $\Psi_{\left(c_{i}\right)_{i}, s_{6}}$ to a neighborhood of $p_{0}$ is $L_{0}$-biLipschitz onto the $\vartheta$-ball about $\Psi_{\left(c_{i}\right)_{i}, \delta_{6}}\left(p_{0}\right)$. As $C_{1} \delta_{1} / L_{0} \leq C_{1} \eta_{4} / L_{0} \leq \vartheta$, this ball contains 0 . Its preimage $p_{1}$ is $C_{1} \delta_{1}$-small. Thus Conclusion $\left(\mathbf{C}_{2}^{\prime \prime}\right)$ holds true.

Here is the proposition which will give the step $i_{0} \rightarrow i_{0}+1$ of the aforementioned induction.
Lemma 4.26. If $p_{i_{0}} \in \mathcal{P}_{0},\left(c_{i}\right)_{i_{0} \leq i \leq i_{0}+4} \in \prod_{i=i_{0}}^{i_{0}+4} \breve{\mathscr{A}}_{\square_{i}}^{*}$ with $t\left(\square_{i}\right)=\sigma\left(c_{i+1}\right)$ when $i \leq i_{0}+3$ and $s_{i_{0}+5} \in \overrightarrow{\mathscr{A}}_{\square_{i_{0}+4}}$ satisfy:
$\left(h_{1}\right) \overline{\mathrm{w}}\left(c_{i}\right)$ is smaller than $\eta_{4}$ for any $i_{0} \leq i \leq i_{0}+4$.
$\left(h_{2}\right)$ For every $i_{0} \leq i \leq i_{0}+3$, it holds $\mathcal{V}_{p_{i_{0}}}^{c_{i}, c_{i+1}}=0$ and $\mathcal{V}_{p_{i_{0}}}^{c_{i_{0}+4}, s_{i_{0}+5}}=0$.
Then there exist $c_{i_{0}+5} \in \breve{\mathscr{A}}_{\square_{i_{0}+5}}^{*}, s_{i_{0}+6} \in \overrightarrow{\mathscr{A}}_{\oplus_{i_{0}+5}}$ with $t\left(\square_{i_{0}+4}\right)=\sigma\left(c_{i_{0}+5}\right)$ and a parameter $p_{i_{0}+1} \in \mathcal{P}_{0}$ which is $C_{1} \delta_{i_{0}+5 \text {-close }}$ to $p_{i_{0}}$ such that:
$\left(\mathbf{C}_{1}\right) \delta_{i_{0}+5}^{1+\epsilon+2 \epsilon^{2}} \leq \underline{\mathrm{w}}\left(c_{i_{0}+5}\right) \leq \overline{\mathrm{w}}\left(c_{i_{0}+5}\right) \leq \delta_{i_{0}+5}$.
$\left(\mathbf{C}_{2}^{\prime \prime}\right)$ It holds $\mathcal{V}_{p_{i_{0}+1}}^{c_{i}, c_{i+1}}=0$ for every $i_{0}+1 \leq i \leq i_{0}+4$ and $\mathcal{V}_{p_{i_{0}+1}}^{c_{i_{0}+5}, s_{i_{0}+6}}=0$.
Proof. The proof is similar to the one of Lemma 4.25. By $\left(\mathbf{H}_{\mathbf{1}}\right)$ of Definition 1.43, there exist $u_{i_{0}+5} \in \overleftarrow{\mathscr{A}}_{Ð_{i_{0}+5}}$ and $s_{i_{0}+6} \in \overrightarrow{\mathscr{A}}_{Ð_{i_{0}+5}}$ such that the curve $W_{p_{i_{0}}}^{u_{i_{0}+5}} \cap Y_{p_{i_{0}}}^{\complement_{i_{0}+5}}$ is sent tangent to $W_{p_{i_{0}}}^{s_{i_{0}}+6}$ by $F_{p_{i_{0}}}^{\square_{i_{0}+5}}$. Let $m_{i_{0}+5} \geq 1$ and $m_{i_{0}+5}^{\prime} \geq 1$ be maximal such that:

$$
\overline{\mathrm{w}}\left(p_{m_{i_{0}+5}}\left(s_{i_{0}+5}\right)\right)>\delta_{i_{0}+5} \quad \text { and } \quad \overline{\mathrm{h}}\left(p_{m_{i_{0}+5}^{\prime}}\left(u_{i_{0}+5}\right)\right)>\delta_{i_{0}+5}
$$

As $\delta_{i_{0}+5}<\eta_{4}<\eta_{1}$, by Lemma 4.8 there exists $\boldsymbol{d}_{i_{0}+5} \in \mathscr{A}^{*}$ such that the word $\boldsymbol{c}_{i_{0}+5}:=$ $p_{1+m_{i_{0}+5}}\left(s_{i_{0}+5}\right) \cdot d_{i_{0}+5} \cdot p_{1+m_{i_{0}+5}^{\prime}}\left(u_{i_{0}+5}\right) \in \mathscr{A}^{*}$ is well-defined and satisfies $t\left(\square_{i_{0}+4}\right)=0\left(s_{i_{0}+5}\right)=$ $a\left(c_{i_{0}+5}\right)$ with:

$$
\begin{equation*}
\delta_{i_{0}+5} \geq \overline{\mathrm{h}}\left(c_{i_{0}+5}\right) \geq \underline{h}\left(c_{i_{0}+5}\right) \quad \text { and } \quad \delta_{i_{0}+5} \geq \overline{\mathrm{w}}\left(c_{i_{0}+5}\right) \geq \underline{\mathrm{w}}\left(c_{i_{0}+5}\right) \geq \delta_{i_{0}+5}^{1+\epsilon+2 \epsilon^{2}} \tag{4.6}
\end{equation*}
$$

Moreover $c_{i_{0}+5} \in \breve{\mathscr{A}}_{\square_{i_{0}+5}}^{*}$ since $u_{i_{0}+5} \in \overleftarrow{\mathscr{A}}_{\square_{i_{0}+5}}$ and since $\bar{h}\left(c_{i_{0}+5}\right) \leq \delta_{i_{0}+5}<\eta_{4}<\eta_{2}$. This proves $\left(\mathbf{C}_{1}\right)$. Let us prove $\left(\mathbf{C}_{2}^{\prime \prime}\right)$. We recall:

$$
\Psi_{\left(c_{i}\right)_{i}, s_{i_{0}+6}}: p \in \mathcal{P}_{0} \mapsto\left(\left(\mathcal{V}_{p}^{c_{i}, c_{i+1}}\right)_{i_{0}+1 \leq i \leq i_{0}+4}, \mathcal{V}_{p}^{c_{i_{0}+5}, s_{i_{0}}+6}\right)
$$

By Proposition 4.16, the latter function is $C_{1} \delta_{i_{0}+5} / L_{0}-C^{0}$-close to the following:

$$
\stackrel{\circ}{\Psi}: p \in \mathcal{P}_{0} \mapsto\left(\left(\mathcal{V}_{p}^{c_{i}, c_{i+1}}\right)_{i_{0}+1 \leq i \leq i_{0}+3}, \mathcal{V}_{p}^{c_{i_{0}+4, s_{5}}}, \mathcal{V}_{p}^{u_{5}, s_{i_{0}+6}}\right)
$$

By $\left(h_{2}\right)$ and definition of $\left(u_{i}, \jmath_{i+1}\right)_{i}$, this function vanishes at $p=p_{i_{0}}$. Thus it holds:

$$
\left\|\Psi_{\left(c_{i}\right)_{i}, s_{i}+6}\left(p_{i_{0}}\right)\right\| \leq C_{1} \delta_{i_{0}+5} / L_{0}
$$

By Proposition 4.20, a restriction of $\Psi_{\left(c_{i}\right)_{i}, s_{i_{0}+6}}$ to a neighborhood of $p_{i_{0}}$ is $L_{0}$-biLipschitz onto the $\vartheta$-ball about $\Psi_{\left(c_{i}\right)_{i}, s_{i_{0}+6}}\left(p_{1}\right)$. As $C_{1} \delta_{i_{0}+5} / L_{0} \leq C_{1} \eta_{4} / L_{0} \leq \vartheta$, this ball contains 0 . Its preimage $p_{i_{0}+1}$ is $C_{1} \delta_{i_{0}+5}$-close to $p_{i_{0}}$. Thus conclusion $\left(\mathbf{C}_{2}^{\prime \prime}\right)$ holds true.

We are now ready to:
Proof of Proposition 4.23. We prove by induction on $i_{0} \geq 1$ the existence of parameters $\left(p_{i}\right)_{1 \leq i \leq i_{0}}$, words $\left(c_{i}\right)_{i \leq i_{0}+4} \in \prod_{i=1}^{i_{0}+4} \breve{\mathscr{A}}_{\square_{i}}^{*}$ with $t\left(\square_{i}\right)=\sigma\left(c_{i+1}\right)$ for $i \leq i_{0}+3$ and of $s_{i_{0}+5} \in \overrightarrow{\mathscr{A}}_{\square_{i_{0}+4}}$ satisfying: $\left(\mathbf{C}_{1}\right) \delta_{i}^{1+\epsilon+2 \epsilon^{2}} \leq \underline{\mathrm{w}}\left(c_{i}\right) \leq \overline{\mathrm{w}}\left(c_{i}\right) \leq \delta_{i}$ for every $i \leq i_{0}+4$.
$\left(\mathbf{C}_{2}^{\prime \prime}\right) \mathcal{V}_{p_{i_{0}}}^{c_{i}, c_{i+1}}=0$, for every $i_{0} \leq i \leq i_{0}+3$ and $\mathcal{V}_{p_{i_{0}}}^{c_{i_{0}+4}, \dot{s}_{0}+5}=0$.
$\left(\mathbf{C}_{3}\right)$ The parameters $p_{i_{0}-1}$ and $p_{i_{0}}$ are $<C_{1} \delta_{i_{0}+4}$ distant whenever $i_{0} \geq 2$. The parameters $p_{0}$ and $p_{1}$ are $C_{1} \delta_{1}$ distant.
Observe that by Remark 4.18, this statement implies Proposition 4.23.
The proof of this statement is done by induction on $i_{0} \geq 1$ using Lemmas 4.25 and 4.26. We notice that step $i_{0}=1$ is an immediate consequence of Lemma 4.25. Let $i_{0} \geq 1$ and assume $\left(c_{i}\right)_{i \leq i_{0}+4} \in\left(\mathscr{A}^{*}\right)^{i_{0}+4}$ constructed such that $\left(\mathbf{C}_{1}-\mathbf{C}_{2}^{\prime \prime}-\mathbf{C}_{3}\right)$ are satisfied at every $i \leq i_{0}$. By Fact 4.21, the distance between the parameter $p_{i_{0}}$ and $p_{0}$ is at most $C_{1} \sum_{i=1}^{i_{0}+4} \delta_{i}<C_{1} \cdot M_{\beta} \cdot \delta_{0}$, and so by Definition 4.22 the parameter $p_{i_{0}}$ belongs to $\mathcal{P}_{0}$. Thus we can use Lemma 4.26 which implies $\left(\mathbf{C}_{1}-\mathbf{C}_{2}^{\prime \prime}-\mathbf{C}_{3}\right)$ at step $i_{0}+1$.

Here is a consequence of the above proof which will be useful to obtain Conclusions ( $\boldsymbol{\bullet}$ ) of Theorem 1.45. We recall that $(\mathscr{A} \sqcup \mathscr{C})^{*}$ denote the set of words $\mathscr{g}=a_{1} \cdots a_{m}$ formed by letters $a_{i} \in \mathscr{A} \sqcup \mathscr{C}$ such that $t\left(a_{i}\right)=a\left(a_{i+1}\right)$ for every $i<m$.

Corollary 4.27. Under the assumptions of Corollary 4.24, there exists $\nu>0$ such that given any function $\mathrm{C}^{\circ}:(\mathscr{A} \sqcup \mathscr{C})^{*} \mapsto \mathscr{A}^{*}$ satisfying

$$
\left|C^{o}(c)\right| \leq \nu|c| \quad \text { and } \quad t(c)=a\left(C^{o}(c)\right) \quad \forall c \in(\mathscr{A} \sqcup \mathscr{C})^{*}
$$

then Conclusions $\left(\mathbf{C}_{1}\right)$ and $\left(\mathbf{C}_{2}^{\prime}\right)$ of Corollary 4.24 hold true with words $\left(c_{i}\right)_{i}$ satisfying:
$\left(\mathbf{C}_{0}\right)$ the sequence $\mathcal{f}:=c_{1} \cdot \square_{1} \cdot c_{2} \cdot \square_{2} \cdots c_{i} \cdot \square_{i} \cdots$ has infinitely many $m \geq m^{\prime}$ such that:

$$
p_{m}(f)=p_{m^{\prime}}(f) \cdot C^{o}\left(p_{m^{\prime}}(f)\right)
$$

Proof. In the proof of Lemma 4.26, the word $c_{i_{0}+5}$ is chosen in function of $\left(u_{i_{0}+5}, \jmath_{i_{0}+5}\right)$ and of the form:

$$
c_{i_{0}+5}=p_{1+m_{i_{0}+5}}\left(s_{i_{0}+5}\right) \cdot d_{i_{0}+5} \cdot p_{1+m_{i_{0}+5}^{\prime}}\left(u_{i_{0}+5}\right)
$$

By using Remark 4.9 instead of Lemma 4.8, the word $d_{i_{0}+5}$ can be chosen of the form

$$
d_{i_{0}+5}=d_{i_{0}+5}^{o} \cdot d_{i_{0}+5}^{r}
$$

for any word $\mathscr{d}_{i_{0}+5}^{o}$ in $\mathscr{A}^{*}$ such that:

$$
\begin{equation*}
\underline{\mathrm{w}}\left(d_{i_{0}+5}^{o}\right)>\delta_{i_{0}+5}^{\epsilon^{4}} \cdot \frac{1+\theta^{2}}{\mathrm{~W}} \quad \text { and } \quad a\left(d_{i_{0}+5}^{o}\right)=t\left(p_{1+m_{i_{0}+5}}\left(\jmath_{i_{0}+5}\right)\right) \tag{4.7}
\end{equation*}
$$

and where $\mathscr{d}_{i_{0}+5}^{r}$ is a word in $\mathscr{A}^{*}$ such that $\underline{\mathrm{w}}\left(\mathcal{d}_{i_{0}+5}^{r}\right)>\mathrm{W}$ and such that $c_{i_{0}+5}$ is in $\mathscr{A}^{*}$. By definition of W , the word $d_{i_{0}+5}^{r}$ always exists whatever are $d_{i_{0}+5}^{o}$ and $p_{1+m_{i_{0}+5}^{\prime}}^{\prime}\left(u_{i_{0}+5}\right)$. Now we
recall that $p_{1+m_{i_{0}+5}}\left(\Im_{i_{0}+5}\right)$ depends only on $p_{i_{0}}$ and $i_{0}$ which do not depend on $d_{i_{0}+5}^{o}$. So really $d_{i_{0}+5}^{o}$ can be any word satisfying Eq. (4.7). For $j \geq 1$, let

$$
g_{j}:=c_{1} \cdot \square_{1} \cdot c_{2} \cdots c_{j-1} \cdot \square_{j-1} \cdot p_{1+m_{j}}\left(\jmath_{j}\right)
$$

With $j=i_{0}+5$, Eq. (4.7), is equivalent to:

$$
\begin{equation*}
\log \underline{\mathrm{w}}\left(d_{j}^{o}\right)>\epsilon^{4} \log \delta_{j}+\log \left(\left(1+\theta^{2}\right) / \mathrm{W}\right) \quad \text { and } \quad a\left(d_{j}^{o}\right)=t\left(p_{1+m_{j}}\left(s_{j}\right)\right) \tag{4.8}
\end{equation*}
$$

Put:

$$
\mu:=-\inf _{d \in \mathscr{A}^{*}} \frac{\log \underline{w}(d)}{|d|}
$$

We observe that $0<\mu<\infty$ since $\mathscr{A}$ is finite. Then Eq. (4.8) is implied by:

$$
\begin{equation*}
\left|d_{j}^{o}\right|<-\frac{\epsilon^{4}}{\mu} \log \delta_{j}-\frac{\log \left(\left(1+\theta^{2}\right) / \mathrm{W}\right)}{\mu} \quad \text { and } \quad a\left(d_{j}^{o}\right)=t\left(p_{1+m_{j}}\left(s_{j}\right)\right) \tag{4.9}
\end{equation*}
$$

Now let us replace the latter inequality by a stronger one involving $\left|g_{j}\right|$ instead of $\delta_{j}$. It holds:

$$
\left|g_{j}\right|=j+\left|c_{1}\right|+\cdots+\left|c_{j-1}\right|+m_{j} \leq 2\left(\left|c_{1}\right|+\cdots+\left|c_{j}\right|\right)
$$

For every $i$ we have by $\left(\mathbf{C}_{1}\right)$ :

$$
\delta_{i}^{1+\epsilon+2 \epsilon^{2}} \leq \overline{\mathrm{w}}\left(c_{i}\right) \leq \lambda^{-\left|c_{i}\right|} \Rightarrow\left|c_{i}\right| \log \lambda \leq-\left(1+\epsilon+2 \epsilon^{2}\right) \log \delta_{i}=-\left(1+\epsilon+2 \epsilon^{2}\right) \beta^{i-j} \log \delta_{j}
$$

Thus:

$$
\left|g_{j}\right| \leq-2\left(1+\epsilon+2 \epsilon^{2}\right) \sum_{i=1}^{j} \beta^{i-j} \log _{\lambda} \delta_{j} \leq-2\left(1+\epsilon+2 \epsilon^{2}\right) \frac{\beta}{\beta-1} \log _{\lambda} \delta_{j} \leq-10 \cdot \log \delta_{j}
$$

Thus Eq. (4.9) and so Eq. (4.7) is implied by

$$
\begin{equation*}
\left|d_{j}^{o}\right|<\frac{\epsilon^{4}}{10 \cdot \mu}\left|q_{j}\right|-\frac{\log \left(\left(1+\theta^{2}\right) / \mathrm{W}\right)}{\mu} \quad \text { and } \quad a\left(d_{j}^{o}\right)=t\left(p_{m_{j}}\left(s_{j}\right)\right) \tag{4.10}
\end{equation*}
$$

As $\left|g_{j}\right| \geq 1$ for every $j$, we obtain the existence of $\nu>0$ so that the corollary holds true.
4.4. Proof of Theorem 1.45. To show the first part of Theorem 1.45 regarding the existence of a family of sets $B_{j} \times\left\{o\left(\square_{j}\right)\right\}$ satisfying the conclusions of Theorem 1.32 (resp. Theorem 1.34), it suffices to show that the assumptions of these theorems are verified at a dense set of parameters.

Proof of the first part of Theorem 1.45. Let $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ be a moderately dissipative unfolding of wild type $(\mathscr{A}, \mathscr{C})$ with $\# \mathscr{C}=5$. In Section 4.3 , we showed that for every $p_{0} \in \mathcal{P}$ and every sufficiently small $\delta_{0}$, we can find $p_{\infty} \in \mathcal{P}$ which is $C_{1} M_{\beta} \delta_{0}$-close to $p_{0}$, a sequence $\left(\square_{j}\right)_{j \geq 1} \in \mathscr{C}^{\mathbb{N}^{*}}$ and a sequence of words $\left(c_{j}\right)_{j \geq 1} \in \prod_{j \geq 1} \breve{\mathscr{A}}_{\square_{j}}^{*}$ such that $t\left(\square_{j}\right)=\sigma\left(c_{j+1}\right)$ for every $j \geq 1$ and satisfying the conclusions $\left(\mathbf{C}_{\mathbf{1}}\right)$ and $\left(\mathbf{C}_{\mathbf{2}}^{\prime}\right)$ of Corollary 4.24. We are going to apply Theorem 1.32 , or respectively Theorem 1.34 if $F_{p_{\infty}}$ is analytic.

Proposition 4.28. The system $F_{p_{\infty}}$ of type $(\mathscr{A}, \mathscr{C})$ endowed with the adapted projection $\pi_{p_{\infty}}$ and the sequences $\left(\square_{j}\right)_{j \geq 1},\left(c_{j}\right)_{j \geq 1}$ satisfies the assumptions of Theorem 1.32, and moreover those of Theorem 1.34 if $F_{p_{\infty}}$ is analytic.

The proof of Proposition 4.28 is given below. Together with Theorem 1.32 (resp. Theorem 1.34), it implies the existence of families of non-empty subsets $B_{j}$ of $Y^{e}$ satisfying the conclusions of Theorem 1.32 (resp. Theorem 1.34).

The rest of this subsection is devoted to the proof of Proposition 4.28. In the following, we fix a parameter $p_{\infty}$ as in Corollary 4.24. We are going to show that $F_{p_{\infty}}$ endowed with $\pi_{p_{\infty}},\left(\square_{j}\right)_{j \geq 1}$ and $\left(c_{j}\right)_{j \geq 1}$ satisfies the assumptions of Theorem 1.32, by checking these assumptions one by one.

For simplicity and since we work at fixed parameter, we will omit to note $p_{\infty}$ in index in the following. We begin by proving the following crucial estimates between the coefficients $\delta_{j}$ (defined in Definition 4.22), $\left|Y^{c_{j}}\right|$ (defined in Definition 1.27) and $\breve{\gamma}_{j}$ (defined in Theorem 1.32):

Lemma 4.29. For every $j \geq 1$, the coefficient $\breve{\gamma}_{j}$ is non zero and it holds:

$$
2 \delta_{j}^{1+\epsilon+2 \epsilon^{2}} \leq\left|Y^{c_{j}}\right| \leq 2 \delta_{j} \quad \text { and } \quad 2 \delta_{j}^{3\left(1+\epsilon+2 \epsilon^{2}\right)} \leq \breve{\gamma}_{j} \leq 2 \delta_{j}^{3}
$$

Proof. To prove the first inequalities, we recall that by Remark 4.4 we have $2 \underline{\mathrm{w}}\left(c_{j}\right) \leq\left|Y^{c_{j}}\right| \leq$ $2 \overline{\mathrm{w}}\left(c_{j}\right)$ for every $j \geq 1$. Thus the inequalities follow immediately from Corollary $4.24\left(\mathbf{C}_{1}\right)$.

We now prove the second ones. We first recall that:

$$
\breve{\gamma}_{j}=\left|Y^{c_{j+1}}\right|^{1 / 2} \cdot\left|Y^{c_{j+2}}\right|^{1 / 2^{2}} \cdots\left|Y^{c_{j+k}}\right|^{1 / 2^{k}} \cdots
$$

if it is well-defined. To show that $\breve{\gamma}_{j}$ is indeed well-defined and positive, we compute the logarithm of the latter infinite product (with $\log 0=-\infty$ by convention) and we show that it is finite:

$$
\log \breve{\gamma}_{j}=\sum_{k>j} \log \left|Y^{c_{k}}\right|^{1 / 2^{k-j}}=\sum_{k>j} \log \left|Y^{c_{k}}\right|^{2^{j-k}}=\sum_{k>j} 2^{j-k} \log \left|Y^{c_{k}}\right|=\sum_{k>0} 2^{-k} \log \left|Y^{c_{j+k}}\right|
$$

We just proved that $2 \delta_{k}^{1+\epsilon+2 \epsilon^{2}} \leq\left|Y^{c_{k}}\right| \leq 2 \delta_{k}$ for every $k>0$. Thus we have:

$$
\log (2)+\left(1+\epsilon+2 \epsilon^{2}\right) \sum_{k>0} 2^{-k} \log \delta_{k+j} \leq \log \breve{\gamma}_{j} \leq \log (2)+\sum_{k>0} 2^{-k} \log \delta_{k+j}
$$

We recall that $\delta_{k+j}=\delta_{j}^{\beta^{k}}$ by Eq. (4.4). Thus it holds:

$$
\sum_{k>0} 2^{-k} \log \delta_{k+j}=\sum_{k>0}\left(\frac{\beta}{2}\right)^{k} \log \delta_{j}=\left(\frac{\beta}{2-\beta}\right) \log \delta_{j}
$$

Then $\log (2)+\left(\frac{\beta}{2-\beta}\right)\left(1+\epsilon+2 \epsilon^{2}\right) \log \delta_{j} \leq \log \breve{\gamma}_{j} \leq \log (2)+\left(\frac{\beta}{2-\beta}\right) \log \delta_{j}$. In particular, $\breve{\gamma}_{j}$ is positive and:

$$
\begin{equation*}
2 \delta_{j}^{\frac{\beta}{2-\beta}\left(1+\epsilon+2 \epsilon^{2}\right)} \leq \breve{\gamma}_{j} \leq 2 \delta_{j}^{\frac{\beta}{2-\beta}} \tag{4.11}
\end{equation*}
$$

which concludes the proof of the second inequalities since $\frac{\beta}{2-\beta}=3$.
Here is an important consequence of the latter lemma:
Fact 4.30. When $j$ is large, $\delta_{j+5}$ is small compared to $\breve{\gamma}_{j}^{2}$.
Proof. By Eq. (4.4), we have $\delta_{j+5}=\delta_{j}^{\beta^{5}}=\delta_{j}^{7.59375}$. By Eq. (4.11), we have $\breve{\gamma}_{j} \geq 2 \delta_{j}^{\frac{\beta}{2-\beta}\left(1+\epsilon+2 \epsilon^{2}\right)}$. Since $\frac{2 \beta}{2-\beta}\left(1+\epsilon+2 \epsilon^{2}\right)=6\left(1+\epsilon+2 \epsilon^{2}\right)<7.59375$, the fact follows.

Remark 4.31. The proof of Fact 4.30 enables to see why we needed 5 parameters in the proof of the main theorems of this article. Indeed, for a $k$-parameter family, we would have to find $a$ $2>\beta>1$ which statisfies:

$$
\frac{\beta}{2-\beta}\left(1+\epsilon+2 \epsilon^{2}\right)<\beta^{k}
$$

By taking $\epsilon$ small this gives the inequality:

$$
\frac{\beta}{2-\beta}<\beta^{k}
$$

which does not have any solution for $2>\beta>1$ and $k<5$.
Proof of Proposition 4.28. Let us show that $F$ endowed with $\pi,\left(\square_{j}\right)_{j},\left(c_{j}\right)_{j}$ satisfies the assumptions of Theorem 1.32. First, by Lemma 4.29, the coefficient $\breve{\gamma}_{j}$ is positive. This shows $(i)$.

We now show (ii). We remark that $t\left(\square_{j}\right)=\sigma\left(c_{j+1}\right)$ by Corollary 4.24. We want to prove that the $\pi^{t\left(\complement_{j}\right)}$-fiber of $F^{\unlhd_{j}}\left(\zeta^{c_{j}}\right)$ is $o\left(\breve{\gamma}_{j}^{2}\right)-C^{0}$-close to the $\pi^{0\left(c_{j+1}\right)}$-fiber of $\left(F^{c_{j+1}}\right)^{-1}\left(\zeta^{c_{j+1}}\right)$. By Corollary $4.24\left(\mathbf{C}_{2}^{\prime}\right)$, we have $\left|\mathcal{V}^{c_{j}, c_{j+1}}\right|=O\left(\delta_{j+5}\right)$, where, by Definitions 4.15 and 4.17 , it holds:

$$
\mathcal{V}^{c_{j}, c_{j+1}}=a^{c_{j}}-b^{c_{j+1}}=\pi^{t\left(\oplus_{j}\right)} \circ F^{®_{j}}\left(\zeta^{c_{j}}\right)-\pi^{0\left(c_{j+1}\right)} \circ\left(F^{c_{j+1}}\right)^{-1}\left(\zeta^{c_{j+1}}\right)
$$

By Fact 4.30, we have $\delta_{j+5}=o\left(\breve{\gamma}_{j}^{2}\right)$. Since $F^{c_{j+1}}$ expands the horizontal distance is expanded by a factor dominated by $\left|Y^{c_{j+1}}\right|^{-1}$ and since $\breve{\gamma}_{j}^{2}=\left|Y^{c_{j+1}}\right| \cdot \breve{\gamma}_{j+1}$ we obtain (ii).

By Lemma 4.29, $\left|Y^{c_{j}}\right| \leq 2 \delta_{j}$ and so $\left|Y^{c_{j}}\right|$ is small when $j$ is large.
This shows the first statement of (iii). Finally, by Remark 4.4 and Proposition 4.7 (2), we have:

$$
\left|F^{c_{j}}\left(Y^{c_{j}}\right)\right| \leq 2 \overline{\mathrm{~h}}\left(c_{j}\right) \leq 2 \underline{\mathrm{w}}\left(c_{j}\right)^{1 / \epsilon} \leq 2\left|Y^{c_{j}}\right|^{10}
$$

By Lemma 4.29, we have $\left|Y^{c_{j}}\right|^{6\left(1+\epsilon+2 \epsilon^{2}\right)}=O\left(\breve{\gamma}_{j}^{2}\right)$. As $6\left(1+\epsilon+2 \epsilon^{2}\right)=6 \times 1.12<10$, we obtain $\left|F^{c_{j}}\left(Y^{c_{j}}\right)\right|=o\left(\breve{\gamma}_{j}^{2}\right)$ and the proof of (iii) is complete.

If moreover $F$ is analytic, by Proposition $4.7\left(2^{\prime}\right)$, we have $\max _{\tilde{I}^{2}}\left|\partial_{w_{0}} \mathcal{W}^{c_{j}}\right|^{\epsilon} \leq \min _{\tilde{I}^{2}}\left|\partial_{z_{1}} \mathcal{Z}^{c_{j}}\right|$ when $j$ is large, after reducing $\rho$ if necessary, and we conclude the same. This shows assumption $(\widetilde{i i i})$ of Theorem 1.34.

We are now going to prove the second part of Theorem 1.45 regarding Property $(\boldsymbol{*})$. This will be done using a new symbolic algorithm. We recall that $\overrightarrow{\mathscr{A}} \subset \overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$ denote the sets of infinite words $s=s_{1} \cdots s_{m} \cdots$ in respectively the alphabets $\mathscr{A}$ and $\mathscr{A} \sqcup \mathscr{C}$ which are admissible: $t\left(a_{i}\right)=a\left(a_{i+1}\right)$, for every $i$. The space $\overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$ is compact endowed with the following distance:

$$
d\left(s, s^{\prime}\right)=2^{-\min \left(i: s_{i} \neq s_{i}^{\prime}\right)} \quad \text { for } s=s_{1} s_{2} \cdots, s^{\prime}=s_{1}^{\prime} s_{2}^{\prime} \cdots \in \overrightarrow{\mathscr{A} \sqcup \mathscr{C}}
$$

Let $g$ and $g^{\mathscr{A}}$ be the shift dynamics on $\overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$ and $\overrightarrow{\mathscr{A}}$ :

$$
g: s_{1} \cdots s_{m} \cdots \in \overrightarrow{\mathscr{A} \sqcup \mathscr{C}} \mapsto s_{2} \cdots s_{m+1} \cdots \in \overrightarrow{\mathscr{A} \sqcup \mathscr{C}} \quad \text { and } \quad g^{\mathscr{A}}:=g \mid \overrightarrow{\mathscr{A}}
$$

Let $\mathcal{M}(g)$ be the space of $g$-invariant measure on $\overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$ and let $\mathcal{M}\left(g^{\mathscr{A}}\right)$ be the space of $g$-invariant measures supported by $\overrightarrow{\mathscr{A}}$. We endow $\mathcal{M}(g) \supset \mathcal{M}\left(g^{\mathscr{A}}\right)$ with the Wasserstein distance. Given $f \in \overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$, let $\mathrm{e}_{n}(f)$ be the probability measure on $\overrightarrow{\mathscr{A} \sqcup \vec{C}}$ which is equidistributed on the $n$ first $g$-iterates of $\mathcal{f}$. For every $c \in \mathscr{A} \sqcup \mathscr{C}^{*}$, let us chose $\tilde{c} \in \overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$ such that $\mathcal{R}_{|c|}(\tilde{c})=c$ and put $\mathrm{e}_{n}(c):=\mathrm{e}_{n}(\tilde{c})$ for every $n \leq|c|$. Here is the algorithm, it uses $\nu>0$ defined in Corollary 4.27.

Definition 4.32. Let $\mathrm{C}^{o}:(\mathscr{A} \sqcup \mathscr{C})^{*} \rightarrow \mathscr{A}^{*}$ be a map sending $c \in(\mathscr{A} \sqcup \mathscr{C})^{*}$ to a word $d$ such that:
(1) $t(c)=a(d)$ and $|d| \leq \nu|c|$.
(2) for every word $d^{\prime} \in \mathscr{A}^{*}$ such that $t(c)=\sigma\left(d^{\prime}\right)$ and $\left|d^{\prime}\right| \leq \nu|c|$, it holds:

$$
d\left(\mathrm{e}_{\left|c \cdot d^{\prime}\right|}\left(c \cdot d^{\prime}\right),\left\{\mathrm{e}_{n}(c): 1 \leq n \leq|c|\right\}\right) \leq d\left(\mathrm{e}_{|c \cdot d|}(c \cdot d),\left\{\mathrm{e}_{n}(c): 1 \leq n \leq|c|\right\}\right)
$$

Corollary 4.33. Under the assumptions of Corollary 4.24, there is a parameter $p_{\infty} \in \mathcal{P}_{0}$ which is $C_{1} M_{\beta} \delta_{0}$-close to $p_{0}$ and a sequence of words $\left(c_{i}\right)_{i}$ satisfying Conclusion $\left(\mathbf{C}_{1}\right)$ and $\left(\mathbf{C}_{2}^{\prime}\right)$ of Corollary 4.24 and moreover the point $\mathcal{f}:=c_{1} \cdot \square_{1} \cdot c_{2} \cdot \square_{2} \cdots c_{i} \cdot \square_{i} \cdots \in \vec{A} \sqcup \mathscr{C}$ satisfies:
$\left(\mathbf{C}_{0}\right)$ the limit set $L$ of $\left(\mathrm{e}_{n}(\mathcal{F})\right)_{n \geq 0}$ contains $\frac{1}{1+\nu} \cdot \check{\mu}+\frac{\nu}{1+\nu} \cdot \mathcal{M}\left(g^{\text {gl }}\right)$ for a certain $\check{\mu} \in \mathcal{M}\left(g^{\text {sl }}\right)$.
Proof. We recall that for every $m \geq 0$ and $s=s_{0} \cdots s_{j} \cdots \in \overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$, we defined $p_{m}(s)=s_{0} \cdots s_{m-1}$. We apply Corollary 4.27 with the function $C^{o}$ of Definition 4.32. It gives a parameter $p_{\infty} \in \mathcal{P}_{0}$ which is $C_{1} M_{\beta} \delta_{0}$-close to $p_{0}$ and a sequence of words $\left(c_{i}\right)_{i}$ satisfying Conclusions $\left(\mathbf{C}_{1}\right)$ and $\left(\mathbf{C}_{2}^{\prime}\right)$ of Corollary 4.24 and moreover $\mathcal{K}:=c_{1} \cdot \square_{1} \cdot c_{2} \cdot \square_{2} \cdots c_{i} \cdot \square_{i} \cdots$ satisfies:
$\left(\mathbf{C}_{0}^{\prime}\right)$ there are infinitely many $m \geq 0$ such that with $K_{m}:=\left\{\mathrm{e}_{n}\left(\mathcal{p}_{n}(\mathcal{f})\right): 1 \leq n \leq m\right\}$ and $m^{\prime}:=$ $m+\left|\mathrm{C}^{o}\left(p_{m}(f)\right)\right|$, it holds:

$$
d\left(\mathrm{e}_{m^{\prime}}\left(\mathcal{p}_{m^{\prime}}(f), K_{m}\right) \geq \max _{\left\{d \in \mathscr{A}^{*}: t\left(p_{m}(f)\right)=o(d) \text { and }|d| \leq \nu m\right\}} d\left(\mathrm{e}_{m+|d|}\left(\mathcal{p}_{m}(f) \cdot d\right), K_{m}\right)\right.
$$

Now we use:
Lemma 4.34. For every $s \in \overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$, the measures $\mathrm{e}_{n}(\jmath)$ and $\mathrm{e}_{n}\left(\mathfrak{p}_{n}(\jmath)\right)$ are $2 / n$-close.
Proof. The measures $\mathrm{e}_{n}(\mathrm{~s})$ and $\mathrm{e}_{n}\left(p_{n}(\mathrm{~s})\right)$ are equidistributed on $n$ atoms which pairwise $2^{-n+k_{-}}$ close for $0 \leq k \leq n$. Thus the transport cost is at most $\frac{1}{n} \sum 2^{-n+k}=2 / n$.

Thus $\mathrm{e}_{m^{\prime}}\left(p_{m^{\prime}}(f)\right)$ is close to $\mathrm{e}_{m^{\prime}}(f)$ when $m$ is large. Also $\mathrm{e}_{m+|d|}\left(p_{m}(f) \cdot d\right)$ is close to $\frac{m}{m+|d|} \mathrm{e}_{m}(\mathcal{\ell})+\frac{|d|}{m+|d|} \mathrm{e}_{|d|}(d)$ when $m$ is large. Consequently when $m$ is large:
$d\left(\mathrm{e}_{m^{\prime}}(f), K_{m}\right) \geq_{\left\{d \in \mathscr{A}^{*}: t\left(p_{m}(f)\right)=o(d) \text { and }|d| \leq \nu m\right\}} d\left(\frac{m}{m+|d|} \mathrm{e}_{m}(f)+\frac{|d|}{m+|d|} \mathrm{e}_{|d|}(d), K_{m}\right)+o(1)$.
As $(\mathscr{V}, \mathscr{A})$ is a finite and connected graph, we do not need to assume that $p_{m}(f) \cdot d$ is admissible in the above inequality, also we can restrict this inequality to the case $|d|=[\nu \cdot m]$ :

$$
d\left(\mathrm{e}_{m^{\prime}}(\mathcal{f}), K_{m}\right) \geq \max _{\{d \in \mathscr{A} *:|d|=[\nu \cdot m]\}} d\left(\frac{m}{m+|d|} \mathrm{e}_{m}(f)+\frac{|d|}{m+|d|} \mathrm{e}_{|d|}(d), K_{m}\right)+o(1)
$$

Thus with $K_{\infty}=c l\left(\bigcup_{m \geq 0} K_{m}\right)$ the limit of $\left(K_{m}\right)_{m}$ in the Hausdorff topology, we have:

$$
0=d\left(\mathrm{e}_{m^{\prime}}(f), K_{\infty}\right)+o(1) \geq \max _{\left\{d \in \mathscr{A}^{*}:|d|=[\nu \cdot m]\right\}} d\left(\frac{1}{1+\nu} \mathrm{e}_{m}(f)+\frac{\nu}{1+\nu} \mathrm{e}_{|d|}(d), K_{\infty}\right)+o(1)
$$

Hence for any accumulation point $\check{\mu}$ of such $\mathrm{e}_{m}(f)$, as $\mathrm{e}_{|d|}(d)$ can approximate any measure $\mu^{\prime}$ of $\mathcal{M}\left(g^{\text {d }}\right)$, it holds:

$$
\max _{\mu^{\prime} \in \mathcal{M}\left(g^{s l}\right)} d\left(\frac{1}{1+\nu} \check{\mu}+\frac{\nu}{1+\nu} \mu^{\prime}, K_{\infty}\right)=0+o(1)
$$

Hence we obtained that the connected set $\frac{1}{1+\nu} \cdot \check{\mu}+\frac{\nu}{1+\nu} \cdot \mathcal{M}\left(g^{\mathscr{A}}\right)$ is included in $K_{\infty}$. Finally we notice that the proportion of letters in $\mathscr{C}$ is asymptotically small by $\left(\mathbf{C}_{1}\right)$. Thus $\check{\mu}$ belongs to $\mathcal{M}\left(g^{\mathscr{A}}\right)$. As $K_{\infty}$ is equal to the union of the limit set $L$ of $\left(e_{m}(f)\right)$ with the discrete set $\bigcup_{m} K_{m}$, the limit set $L$ contains $\frac{1}{1+\nu} \cdot \check{\mu}+\frac{\nu}{1+\nu} \cdot \mathcal{M}\left(g^{d}\right)$.

Proof of property $(\boldsymbol{*})$ of Theorem 1.45. We still consider a moderately dissipative unfolding $\left(F_{p}, \pi_{p}\right)_{p \in \mathcal{P}}$ of wild type $(\mathscr{A}, \mathscr{C})$, such that $\mathscr{C}$ is formed by five elements $\left\{\square_{1}, \square_{2}, \square_{3}, \square_{4}, \square_{5}\right\}$, and a fixed parameter $p_{0} \in \operatorname{int} \mathcal{P}$. We define $\mathcal{P}_{0}$ as in Proposition 4.7. We apply Corollary 4.33 which gives the existence of a parameter $p_{\infty} \in \mathcal{P}_{0}$ which is $C_{1} M_{\beta} \delta_{0}$-close to $p_{0}$ and a sequence of words $\left(c_{i}\right)_{i}$ satisfying Conclusions $\left(\mathbf{C}_{1}\right)$ and $\left(\mathbf{C}_{2}^{\prime}\right)$ of Corollary 4.24 and moreover:
$\left(\mathbf{C}_{0}\right)$ the point $f:=c_{1} \cdot \boxtimes_{1} \cdot c_{2} \cdot \square_{2} \cdots c_{i} \cdot \boxtimes_{i} \cdots \in \overrightarrow{\mathscr{A} \sqcup \mathscr{C}}$ satisfies that the limit set $L$ of $\left(\mathrm{e}_{n}(f)\right)_{n \geq 0}$ contains $\frac{1}{1+\nu} \cdot \check{\mu}+\frac{\nu}{1+\nu} \cdot \mathcal{M}\left(g^{g l}\right)$ for a certain $\check{\mu} \in \mathcal{M}\left(g^{d d}\right)$.
Let $B:=\bigcap_{n \geq 0} F_{p_{\infty}}^{-n}\left(Y_{e} \times \mathscr{V}\right)$. The map $F_{p_{\infty}}$ sends $B$ into itself, and the restriction $F_{p_{\infty}} \mid B$ is semi-conjugate to $g$ via the following map:

$$
h:(z, v) \in B \mapsto\left(d_{i}\right)_{i \geq 0} \quad \text { with } d_{i} \in \mathscr{A} \sqcup \mathscr{C} \text { such that } F_{p_{\infty}}^{i}(z, v) \in Y_{p_{\infty}}^{d_{i}} \times\left\{a\left(d_{i}\right)\right\}
$$

We are going to use the following proved below:
Lemma 4.35. The set of invariant probability measures $\mathcal{M}\left(F_{p^{\infty}}^{g l}\right)$ of $F_{p^{\infty}}^{g l}$ is homeomorphic to $\mathcal{M}\left(g^{\text {d }}\right)$ via the pushforward $h_{*}$ induced by $h$.

By Proposition 4.28 and Theorem 1.32, there exists $J \geq 0$ and a stable domain $B_{J} \subset B$ which is sent by $h$ to $\left\{g^{J}(\boldsymbol{f})\right\}$. By $\left(\mathbf{C}_{0}\right)$ and Lemma 4.35, the limit set of the empirical measures of points in $B_{J}$ contains:
$h_{*}^{-1}\left(\frac{1}{1+\nu} \cdot \check{\mu}+\frac{\nu}{1+\nu} \cdot \mathcal{M}\left(g^{\mathscr{A}}\right)\right)=\frac{1}{1+\nu} \cdot h_{*}^{-1} \check{\mu}+\frac{\nu}{1+\nu} \cdot h_{*}^{-1}\left(\mathcal{M}\left(g^{\mathscr{A}}\right)\right)=\frac{1}{1+\nu} \cdot \mu+\frac{\nu}{1+\nu} \cdot \mathcal{M}\left(F_{p^{\infty}}^{\mathscr{D}}\right)$, with $\mu=h_{*}^{-1} \check{\mu} \in \mathcal{M}\left(F_{p \infty}^{g l}\right)$.
Proof of Lemma 4.35. The map $\mu^{\prime} \in \mathcal{M}\left(F_{p}^{d}\right) \mapsto h_{*} \mu^{\prime} \in \mathcal{M}\left(g^{s d}\right)$ is continuous on its domains which is compact. Thus it suffices to show that it is a bijection. In order to show this, it suffices to recall that the space of invariant measures of a continuous map are homeomorphic to those of its inverse limit, and that the dynamics on the inverse limit $\overleftrightarrow{\mathscr{A}}$ of $\overrightarrow{\mathscr{A}}$ is conjugate to the dynamics on the inverse limit $\overleftrightarrow{\Lambda}$ of $\Lambda$ by Remark 1.14

## Appendix A. Some results on $C_{\rho}^{\omega}$-Hyperbolic transformations

Proof of Proposition 1.7. Let us just do the proof for $C_{\rho}^{\omega}$-hyperbolic transformations, the one for (real) hyperbolic transformations is similar. We assume that $(\tilde{Y}, F) \neq\left(\tilde{Y}^{e}, i d\right) \neq\left(\tilde{Y}^{\prime}, F^{\prime}\right)$, otherwise the proof is obvious. Put:

$$
\left(\tilde{Y}^{\prime \prime}, F^{\prime \prime}\right):=(\tilde{Y}, F) \star\left(\tilde{Y}^{\prime}, F^{\prime}\right)
$$

The proposition is the consequence of the two following Lemmas A. 1 and A.2.
Lemma A.1. The set $\tilde{Y}^{\prime \prime}$ is a $C_{\rho}^{\omega}$-box.
Proof. Note first that $\tilde{Y}^{\prime \prime} \subset \tilde{Y} \subset \tilde{Y}^{e}$. Let $\zeta:(z, w) \in \tilde{I}^{2} \rightarrow\left(\mathcal{Z}_{1}(z, w), w\right) \in \tilde{Y}$ and $\zeta^{\prime}:(z, w) \in$ $\tilde{I}^{2} \rightarrow\left(\mathcal{Z}_{2}(z, w), w\right) \in \tilde{Y}^{\prime}$ be two biholomorphisms given by Definition 1.2 for the $C_{\rho}^{\omega}$-boxes $\tilde{Y}$ and $\tilde{Y}^{\prime}$. By the cone property (2) of Definition 1.5 , at $w \in \tilde{I}$ fixed, the (complex) curve $\mathbb{D}_{1 w}:=$ $\left\{\left(\mathcal{Z}_{1}(z, w), w\right): z \in \tilde{I}\right\}$ is sent by $F$ to a curve with tangent spaces in $\tilde{\chi}_{h}$. By property (1) of Definition 1.5 , the curve $F\left(\mathbb{D}_{1 w}\right)$ is disjoint from $\partial^{u} \tilde{Y}^{e}$ and with boundary in $\partial^{s} \tilde{Y}^{e}$. Still by
property (1) of Definition 1.5 , for every $z_{2} \in \tilde{I}$, the curve $\left\{\left(\mathcal{Z}_{2}\left(z_{2}, w\right), w\right): w \in \tilde{I}\right\}$ is disjoint from $\partial^{s} \tilde{Y}^{e}$. By Definition 1.2, it has its tangent spaces in $\tilde{\chi}_{v}$ and its boundary is in $\partial^{u} \tilde{Y}^{e}$. Thus $\left\{\left(\mathcal{Z}_{2}\left(z_{2}, w\right), w\right): w \in \tilde{I}\right\}$ intersects $F\left(\mathbb{D}_{1 w}\right)$ transversally at a unique point. We denote by $\zeta^{\prime \prime}\left(z_{2}, w\right)=\left(\mathcal{Z}_{12}\left(z_{2}, w\right), w\right)$ its preimage by $F$. By transversality, $\zeta^{\prime \prime}$ is a holomorphic map from $\tilde{I}^{2}$ onto $\tilde{Y}^{\prime \prime}$. Since $\left|\partial_{w} \mathcal{Z}_{2}\right|<\theta$ and by the cone property (3) of Definition 1.5, it holds $\left|\partial_{w} \mathcal{Z}_{12}\right|<\theta$. Also $\left(\left\{\left(\mathcal{Z}_{2}\left(z_{2}, w\right), w\right): w \in \tilde{I}\right\}\right)_{z_{2} \in \tilde{I}}$ are the leaves of a foliation. Each leaf intersects transversally the curve $F\left(\mathbb{D}_{1 w}\right)$ at the unique point $F \circ \zeta^{\prime \prime}\left(z_{2}, w\right)$, so the derivative of $F \circ \zeta^{\prime \prime}$ w.r.t. $z_{2}$ is non-zero. Thus, the derivative of $\zeta^{\prime \prime}\left(z_{2}, w\right)$ w.r.t. $z_{2}$ is non-zero. By definition, $\zeta^{\prime \prime}\left(z_{2}, w\right)$ belongs to $\tilde{I} \times\{w\}$ and so $\partial_{z_{2}} \zeta^{\prime \prime}$ is horizontal and non zero. Thus $D \zeta^{\prime \prime}$ is invertible and so $\zeta^{\prime \prime}$ is a biholomorphism.

Let us denote by $Y, Y^{\prime}, Y^{\prime \prime}$ the real traces $\tilde{Y} \cap Y^{e}, \tilde{Y}^{\prime} \cap Y^{e}, \tilde{Y}^{\prime \prime} \cap Y^{e}$ of $\tilde{Y}, \tilde{Y}^{\prime}, \tilde{Y}^{\prime \prime}$. It remains to show that $Y^{\prime \prime}$ is a (real) box. It is enough to show that $Y^{\prime \prime}=\breve{Y}^{\prime \prime}$ where we set $\left(\breve{Y}^{\prime \prime}, F^{\prime \prime}\right):=(Y, F) \star\left(Y^{\prime}, F^{\prime}\right)$. Indeed both $Y$ and $Y^{\prime}$ are boxes by Definition 1.2 and so by [PY01, $\S 3.2 .1]$ it is also the case for $\breve{Y}^{\prime \prime}$. Then one just has to remark that:

$$
Y^{\prime \prime}=\tilde{Y} \cap F^{-1}\left(\tilde{Y}^{\prime}\right) \cap Y^{e}=Y \cap F^{-1}\left(\tilde{Y}^{\prime}\right)=Y \cap F^{-1}\left(Y^{\prime}\right)=\breve{Y}^{\prime \prime}
$$

where the first equality is true by definition of $\tilde{Y}^{\prime \prime}$, the second equality is true by definition of $Y$, the fourth one is given by definition of the real $\star$-product and the third one uses that $F \mid Y$ is a bijection onto $F(Y)$ whose inverse is conjugate to a hyperbolic transformation via the involution $(z, w) \mapsto(w, z)$. This shows that $\tilde{Y}^{\prime \prime}$ is a $C_{\rho}^{\omega}$-box.

Lemma A.2. The pair $\left(\tilde{Y}^{\prime \prime}, F^{\prime \prime}\right)$ is a $C_{\rho}^{\omega}$-hyperbolic transformation.
Proof. The map $F^{\prime \prime}=F^{\prime} \circ F$ is a biholomorphism on $\tilde{Y}^{\prime \prime}$ as a composition of two biholomorphisms. We now continue with the same notations as in the previous lemma. Let us prove (1) of Definition 1.5. Note that when $z_{2} \in \partial \tilde{I}$, then $\left\{\left(\mathcal{Z}_{12}\left(z_{2}, w\right), w\right): w \in \tilde{I}\right\}$ is sent by $F$ into $\left\{\left(\mathcal{Z}_{2}\left(z_{2}, w\right), w\right): w \in \tilde{I}\right\} \subset \partial^{s} \tilde{Y}^{\prime}$. Thus by this property $(1)$ for $\left(\tilde{Y}^{\prime}, F\right)$, the set $\left\{\left(\mathcal{Z}_{12}\left(z_{2}, w\right), w\right)\right.$ : $w \in \tilde{I}\}$ is sent by $F^{\prime \prime}$ into $\partial^{s} \tilde{Y}^{e}: F^{\prime \prime}\left(\partial^{s} \tilde{Y}^{\prime \prime}\right) \subset \partial^{s} \tilde{Y}^{e}$. By definition, the set $\tilde{Y}^{\prime \prime}$ is included in the definition domain of $F^{\prime \prime}$ which sends it into $F^{\prime}\left(\tilde{Y}^{\prime}\right) \subset \tilde{Y}^{e}$. As $(\tilde{Y}, F) \neq\left(\tilde{Y}^{e}, i d\right) \neq\left(\tilde{Y}^{\prime}, F^{\prime}\right)$, it holds $\tilde{Y}^{\prime \prime} \cap \partial^{s} \tilde{Y}^{e} \subset \tilde{Y} \cap \partial^{s} \tilde{Y}^{e}=\emptyset$ and $F^{\prime \prime}\left(\tilde{Y}^{\prime \prime}\right) \cap \partial^{u} \tilde{Y}^{e} \subset F^{\prime}\left(\tilde{Y}^{\prime}\right) \cap \partial^{u} \tilde{Y}^{e}=\emptyset$. Properties (2), (3) and (4) of Definition 1.5 for $\left(\tilde{Y}^{\prime \prime}, F^{\prime \prime}\right)$ come from the same properties for $(\tilde{Y}, F)$ and $\left(\tilde{Y}^{\prime}, F^{\prime}\right)$.

The remaining of this section is devoted to show the following given a finite graph $(\mathscr{V}, \mathscr{A})$ :
Proposition A.3. If $F^{\mathscr{A}}$ is a hyperbolic map of type $\mathscr{A}$ and if $F^{\mathscr{A}}$ is real analytic, then there exists $\rho>0$ such that for every $c \in \mathscr{A}^{*}$, there exists a $C_{\rho}^{\omega}$-box $\tilde{Y}^{c}$ such that $Y^{c}$ is the real trace $\tilde{Y}^{c} \cap \mathbb{R}^{2}$ of $\tilde{Y}^{c}, F^{c}$ extends to a biholomorphism from $\tilde{Y}^{c}$ and $\left(\tilde{Y}^{c}, F^{c}\right)$ is a $C_{\rho}^{\omega}$-hyperbolic transformation.

Proof. We prove below the following:
Lemma A.4. If $(Y, F)$ is a hyperbolic transformation and $F$ is real analytic, then for every $\rho>0$ sufficiently small, there exists a $C_{\rho}^{\omega}$-box $\tilde{Y}$ such that $Y$ is the real trace $\tilde{Y} \cap \mathbb{R}^{2}$ of $\tilde{Y}$, $F$ extends to a biholomorphism from $\tilde{Y}$ and $(\tilde{Y}, F)$ is a $C_{\rho}^{\omega}$-hyperbolic transformation.

Thus by finiteness of $\mathscr{A}$ and by Lemma A.4, there exist $\rho>0$ such that each hyperbolic transformation $\left(Y^{a}, F^{a}\right)$ with $a \in \mathscr{A}$ admits an extension as a $C_{\rho}^{\omega}$-hyperbolic transformation $\left(\tilde{Y}^{a}, F^{a}\right)$.

Then by Proposition 1.7, for every $c \in \mathscr{A}^{*}$, we can define the $C_{\rho}^{\omega}$-hyperbolic transformation:

$$
\left(\tilde{Y}^{c}, F^{c}\right):=\left(\tilde{Y}^{a_{1}}, F^{a_{1}}\right) \star \cdots \star\left(\tilde{Y}^{a_{k}}, F^{a_{k}}\right), \quad \forall c=a_{1} \cdots a_{k} \in \mathscr{A}^{*}
$$

By Proposition 1.7, the pair $\left(\tilde{Y}^{c}, F^{c}\right)$ is a $C_{\rho}^{\omega}$-hyperbolic transformation extending $\left(Y^{c}, F^{c}\right)$.
Proof of Lemma A.4. Let us assume that $F$ is not the identity (otherwise we take the holomorphic extension $\left.\left(\tilde{Y}^{e}, i d\right)\right)$.

As $F \mid Y$ is a real analytic diffeomorphism onto its image, it extends to a biholomorphism from a neighborhood $\hat{Y}$ of $Y$ in $\mathbb{C}^{2}$ onto its image $F(\hat{Y})$. If $\hat{Y}$ is small enough, from cone properties (2) and (3) of Definition 1.4, at every point in $\hat{Y}$, every non-zero vector in the complement of $\tilde{\chi}_{v}$ is sent into $\tilde{\chi}_{h}$ by $D F$ and has its first coordinate's modulus which is more than $\lambda$-expanded by $D F$. Also at every point in $F(\hat{Y})$, every non-zero vector in the complement of $\tilde{\chi}_{h}$ is sent into $\tilde{\chi}_{v}$ by $D F^{-1}$ and has its second coordinate modulus which is more than $\lambda$-expanded by $D F^{-1}$.

Let $\mathcal{X}: I^{2} \rightarrow \operatorname{int} I$ and $\mathcal{Y}: I^{2} \rightarrow \operatorname{int} I$ be the implicit representation of the hyperbolic transformation $(F, Y)$ as defined in Proposition 3.1. By the transversality used in the proof of this proposition, the maps $\mathcal{X}$ and $\mathcal{Y}$ are real analytic. Thus for $\rho$-small enough they extend to holomorphic functions $\mathcal{Z}$ and $\mathcal{W}$ on $\tilde{I}^{2}$, with $\tilde{I}=I+i[-\rho, \rho]$. By [PY01, Lem. 3.2], it holds:

$$
\left|\partial_{x} \mathcal{X}\right|+\left|\partial_{y} \mathcal{X}\right|<\theta \quad \text { and } \quad\left|\partial_{x} \mathcal{Y}\right|+\left|\partial_{y} \mathcal{Y}\right|<\theta, \quad \text { on } I^{2} .
$$

Thus for $\rho$ sufficiently small:

$$
\left|\partial_{z} \mathcal{Z}\right|+\left|\partial_{w} \mathcal{Z}\right|<\theta \quad \text { and } \quad\left|\partial_{z} \mathcal{W}\right|+\left|\partial_{w} \mathcal{W}\right|<\theta, \quad \text { on } \tilde{I}^{2}
$$

In particular $\mathcal{Z}$ and $\mathcal{W}$ are contracting. As they map $I^{2}$ into the interior of $I$, for $\rho$-sufficiently small $\mathcal{Z}$ and $\mathcal{W}$ map $\tilde{I}^{2}$ into the interior of $\tilde{I}$. Hence the following set is included in $\tilde{Y}^{e} \backslash \partial^{s} \tilde{Y}^{e}$ :

$$
\tilde{Y}:=\left\{\left(\mathcal{Z}\left(z_{1}, w_{0}\right), w_{0}\right):\left(z_{1}, w_{0}\right) \in \tilde{I}^{2}\right\}
$$

Also $\left|\partial_{w} \mathcal{Z}\right|<\theta$. Since $\partial_{x} \mathcal{X} \neq 0$, reducing $\rho$ if necessary, we also have $\partial_{z} \mathcal{Z} \neq 0$ and so the map $(z, w) \in \tilde{I}^{2} \mapsto(\mathcal{Z}(z, w), w) \in \tilde{Y}$ is a biholomorphism. Thus $\tilde{Y}$ is $C_{\rho}^{\omega}$-box. Note that for $\rho$ small, $\tilde{Y}$ is close to $Y$ and so included in $\hat{Y}$. Thus cone properties (2) and (3) of Definition 1.5 are satisfied. Put:

$$
\partial^{u} \tilde{Y}:=\zeta(\tilde{I} \times \partial \tilde{I}) \quad \text { and } \quad \partial^{s} \tilde{Y}:=\zeta((\partial \tilde{I}) \times \tilde{I})
$$

As for every $(x, y) \in I^{2}$, it holds $F(\mathcal{X}(x, y), y)=(x, \mathcal{Y}(x, y))$ by Proposition 3.1 and by analyticity, we have:

$$
F(\mathcal{Z}(z, w), w)=(z, \mathcal{W}(z, w)) \quad \forall(z, w) \in \tilde{I}^{2}
$$

from which property (1) of Definition 1.5 is immediate. Finally note that if for $w_{0} \in I$ and $z_{1} \in \tilde{I}$, the number $\mathcal{Z}\left(z_{1}, w_{0}\right)$ is real, then its image $\left(z_{1}, \mathcal{W}\left(z_{1}, w_{0}\right)\right)$ by $F$ is real. Thus $z_{1}$ is real and so in $I$. Thus

$$
\tilde{Y} \cap \mathbb{R}^{2}=\left\{\left(\mathcal{Z}\left(z_{1}, w_{0}\right), w_{0}\right):\left(z_{1}, w_{0}\right) \in I^{2}\right\}=Y
$$

This achieves the proof of Definition 1.5 (4).

## Appendix B. Bounds on implicit representations and distortion results

Let us recall:
Lemma B. 1 ([PY09, Ineq. (A.8) p.195]). Let $(Y, F)$ and $\left(Y^{\prime}, F^{\prime}\right)$ be hyperbolic transformations with implicit representations $\left(\mathcal{X}_{0}, \mathcal{Y}_{1}\right)$ and $\left(\mathcal{X}_{1}, \mathcal{Y}_{2}\right)$. Then the implicit representation $\left(\tilde{\mathcal{X}}_{0}, \tilde{\mathcal{Y}}_{2}\right)$ of $(Y, F) \star\left(Y^{\prime}, F^{\prime}\right)$ satisfies:

$$
\frac{1}{1+\theta^{2}} \leq \frac{\left|\partial_{x_{2}} \tilde{\mathcal{X}}_{0}\left(x_{2}, y_{0}\right)\right|}{\left|\partial_{x_{1}} \mathcal{X}_{0}\left(x_{1}, y_{0}\right)\right| \cdot\left|\partial_{x_{2}} \mathcal{X}_{1}\left(x_{2}, y_{1}\right)\right|} \leq \frac{1}{1-\theta^{2}}
$$

where $F$ sends $\left(x_{0}, y_{0}\right) \in Y$ to $\left(x_{1}, y_{1}\right) \in Y^{\prime}$ and $\left(x_{1}, y_{1}\right)$ is sent to $\left(x_{2}, y_{2}\right)$ by $F^{\prime}$.
We will also need the following complex analogous of the latter result:
Lemma B.2. Let $(\tilde{Y}, F)$ and $\left(\tilde{Y}^{\prime}, F^{\prime}\right)$ be $C_{\rho}^{\omega}$-hyperbolic transformations with implicit representations $\left(\mathcal{Z}_{0}, \mathcal{W}_{1}\right)$ and $\left(\mathcal{Z}_{1}, \mathcal{W}_{2}\right)$. Then the implicit representation $\left(\tilde{\mathcal{Z}}_{0}, \tilde{\mathcal{W}}_{2}\right)$ of $(\tilde{Y}, F) \star\left(\tilde{Y}^{\prime}, F^{\prime}\right)$ satisfies:

$$
\begin{aligned}
& \frac{1}{1+\theta^{2}} \leq \frac{\left|\partial_{z_{2}} \tilde{\mathcal{Z}}_{0}\left(z_{2}, w_{0}\right)\right|}{\left|\partial_{z_{1}} \mathcal{Z}_{0}\left(z_{1}, w_{0}\right)\right| \cdot\left|\partial_{z_{2}} \mathcal{Z}_{1}\left(z_{2}, w_{1}\right)\right|} \leq \frac{1}{1-\theta^{2}} \\
& \frac{1}{1+\theta^{2}} \leq \frac{\left|\partial_{w_{0}} \tilde{\mathcal{W}}_{2}\left(z_{2}, w_{0}\right)\right|}{\left|\partial_{w_{0}} \mathcal{W}_{1}\left(z_{1}, w_{0}\right)\right| \cdot\left|\partial_{w_{1}} \mathcal{W}_{2}\left(z_{2}, w_{1}\right)\right|} \leq \frac{1}{1-\theta^{2}}
\end{aligned}
$$

where $F$ sends $\left(z_{0}, w_{0}\right) \in \tilde{Y}$ to $\left(z_{1}, w_{1}\right) \in \tilde{Y}^{\prime}$ and $\left(z_{1}, w_{1}\right)$ is sent to $\left(z_{2}, w_{2}\right)$ by $F^{\prime}$.
Proof. We just prove the inequalities on $\tilde{\mathcal{Z}}_{0}$, those on regarding $\tilde{\mathcal{W}}_{2}$ are obtained similarly. We consider the map $\Psi:\left(z_{1}, w_{1}, z_{2}, w_{0}\right) \in \tilde{I}^{4} \rightarrow\left(z_{1}-\mathcal{Z}_{1}\left(z_{2}, w_{1}\right), w_{1}-\mathcal{W}_{1}\left(z_{1}, w_{0}\right)\right)$, which is holomorphic. For all $\left(z_{0}, w_{0}\right) \in \tilde{Y},\left(z_{1}, w_{1}\right) \in \tilde{Y}^{\prime}$ and $\left(z_{2}, w_{2}\right) \in \tilde{Y}^{e}$ such that $F$ sends $\left(z_{0}, w_{0}\right)$ to $\left(z_{1}, w_{1}\right)$ and $F^{\prime}$ sends $\left(z_{1}, w_{1}\right)$ to $\left(z_{2}, w_{2}\right)$, it holds $\Psi\left(z_{1}, w_{1}, z_{2}, w_{0}\right)=(0,0)$. By the cone properties of $C_{\rho}^{\omega}$-hyperbolic transformations, we have $\left|\partial_{w_{1}} \mathcal{Z}_{1}\right|<\theta$ and $\left|\partial_{z_{1}} \mathcal{W}_{1}\right|<\theta$. Then $\Delta:=\left|\operatorname{det}\left(\partial_{\left(z_{1}, w_{1}\right)} \Psi\right)\right|=\left|1-\partial_{w_{1}} \mathcal{Z}_{1} \cdot \partial_{z_{1}} \mathcal{W}_{1}\right|$ satisfies:

$$
\begin{equation*}
1+\theta^{2}>\Delta>1-\theta^{2}>0 \tag{B.1}
\end{equation*}
$$

and so $\partial_{\left(z_{1}, w_{1}\right)} \Psi$ is invertible. By the implicit function Theorem, there exist implicitly defined holomorphic maps $\tilde{\mathcal{Z}}_{1}$ and $\tilde{\mathcal{W}}_{1}$ such that:

$$
\begin{equation*}
\Psi\left(z_{1}, w_{1}, z_{2}, w_{0}\right)=0 \Longleftrightarrow\left(z_{1}=\tilde{\mathcal{Z}}_{1}\left(z_{2}, w_{0}\right) \text { and } w_{1}=\tilde{\mathcal{W}}_{1}\left(z_{2}, w_{0}\right)\right) \tag{B.2}
\end{equation*}
$$

Furthermore, this theorem gives:
$\left(\begin{array}{ccc}\partial_{z_{2}} \tilde{\mathcal{Z}}_{1} & \partial_{w_{0}} \tilde{\mathcal{Z}}_{1} \\ \partial_{z_{2}} & \tilde{\mathcal{W}}_{1} & \partial_{w_{0}} \tilde{\mathcal{W}}_{1}\end{array}\right)=-\partial_{\left(z_{1}, w_{1}\right)} \Psi^{-1} \cdot \partial_{z_{2}, w_{0}} \Psi=-\left(\begin{array}{cc}1 & -\partial_{w_{1}} \mathcal{Z}_{1} \\ -\partial_{z_{1}} \mathcal{W}_{1} & 1\end{array}\right)^{-1} \cdot\left(\begin{array}{cc}-\partial_{z_{2}} \mathcal{Z}_{1} & 0 \\ 0 & -\partial_{w_{0}} \mathcal{W}_{1}\end{array}\right)$.
Thus, it holds $\partial_{z_{2}} \tilde{\mathcal{Z}}_{1}=\Delta^{-1} \cdot \partial_{z_{2}} \mathcal{Z}_{1}$. As $\tilde{\mathcal{Z}}_{0}\left(z_{2}, w_{0}\right)=\mathcal{Z}_{0}\left(\tilde{\mathcal{Z}}_{1}\left(z_{2}, w_{0}\right), w_{0}\right)$, we then obtain $\partial_{z_{2}} \tilde{\mathcal{Z}}_{0}=$ $\partial_{z_{1}} \mathcal{Z}_{0} \cdot \partial_{z_{2}} \tilde{\mathcal{Z}}_{1}=\Delta^{-1} \cdot \partial_{z_{1}} \mathcal{Z}_{0} \cdot \partial_{z_{2}} \mathcal{Z}_{1}$ and the sought inequalities follow from Eq. (B.1).

We are now in position to prove:

Proof of Proposition 4.7. Proof of (1). The left inequality being obvious, let us show the right inequality. Let $c \in \mathscr{A}^{*}$ and $\left(\mathcal{X}_{p}^{c}, \mathcal{Y}_{p}^{c}\right)$ be the implicit representation of $\left(Y_{p}^{c}, F_{p}^{c}\right)$ for $p \in \mathcal{P}$. According to the cone condition (2) of Definition 1.4, we have $\left|\partial_{x} \mathcal{X}_{p}^{c}\right| \leq \lambda^{-|c|}$ and so:

$$
\left(1+\epsilon^{2}\right) \max _{I^{2}} \log \left|\partial_{x} \mathcal{X}_{p}^{c}\right| \leq \max _{I^{2}} \log \left|\partial_{x} \mathcal{X}_{p}^{c}\right|-\epsilon^{2} \cdot|c| \cdot \log \lambda
$$

Then using Corollary 4.2, we obtain:

$$
\left(1+\epsilon^{2}\right) \max _{I^{2}} \log \left|\partial_{x} \mathcal{X}_{p}^{c}\right| \leq \min _{I^{2}} \log \left|\partial_{x} \mathcal{X}_{p}^{c}\right|+\log \overline{\mathbf{B}}-\epsilon^{2} \cdot|c| \cdot \log \lambda
$$

Let $N_{0}>0$ be such that

$$
\begin{equation*}
\log \overline{\mathbf{B}}-\epsilon^{2} \cdot N_{0} \cdot \log \lambda<-\log \left(1+\theta^{2}\right)+\left(1+\epsilon^{2}\right) \log \left(1-\theta^{2}\right) \tag{B.3}
\end{equation*}
$$

Then:

$$
\left(1+\epsilon^{2}\right) \max _{I^{2}} \log \left|\left(1-\theta^{2}\right)^{-1} \cdot \partial_{x} \mathcal{X}_{p}^{c}\right|<\min _{I^{2}} \log \left|\left(1+\theta^{2}\right)^{-1} \cdot \partial_{x} \mathcal{X}_{p}^{c}\right| \quad \forall c \in \mathscr{A}^{*} \text { with }|c| \geq N_{0}
$$

As the latter inequality is strict, given any $p_{0} \in \mathcal{P}$, there exists a compact neighborhood $\mathcal{P}_{0} \subset \mathcal{P}$ of $p_{0}$ such that for every word $c \in \mathscr{A}^{*}$ with at least $N_{0}$ and at most $2 N_{0}$ letters, for every $p, p^{\prime} \in \mathcal{P}_{0}$, we obtain:

$$
\begin{equation*}
\left(1-\theta^{2}\right)^{-1-\epsilon^{2}} \cdot \max _{I^{2}}\left|\partial_{x} \mathcal{X}_{p}^{c}\right|^{1+\epsilon^{2}}<\left(1+\theta^{2}\right)^{-1} \cdot \min _{I^{2}}\left|\partial_{x} \mathcal{X}_{p^{\prime}}^{c}\right| \tag{B.4}
\end{equation*}
$$

Consequently for every word $c \in \mathscr{A}^{*}$ with at least $N_{0}$ and at most $2 N_{0}$ letters,

$$
\begin{equation*}
\left(1-\theta^{2}\right)^{-1-\epsilon^{2}} \cdot \overline{\mathrm{w}}(c)^{1+\epsilon^{2}}<\left(1+\theta^{2}\right)^{-1} \cdot \underline{\mathrm{w}}(c) \tag{B.5}
\end{equation*}
$$

Let $c \in \mathscr{A}^{*}$ be such that $|c| \geq N_{0}$. It can be written as a concatenation $c=c_{1} \cdots c_{m}$ where each $c_{i}$ has at least $N_{0}$ and at most $2 N_{0}$ letters. According to Lemma B.1, we have:

$$
\begin{equation*}
\overline{\mathrm{w}}(c) \leq\left(1-\theta^{2}\right)^{-m} \cdot \overline{\mathrm{w}}\left(c_{1}\right) \cdots \overline{\mathrm{w}}\left(c_{m}\right) \quad \text { and } \quad\left(1+\theta^{2}\right)^{-m} \cdot \underline{\mathrm{w}}\left(c_{1}\right) \cdots \underline{\mathrm{w}}\left(c_{m}\right) \leq \underline{\mathrm{w}}(c) . \tag{B.6}
\end{equation*}
$$

Together with Eq. (B.5), this yields:

$$
\begin{equation*}
\overline{\mathrm{w}}(c)^{1+\epsilon^{2}} \leq \underline{\mathrm{w}}(c) \tag{B.7}
\end{equation*}
$$

Since $1-\epsilon^{2}<\left(1+\epsilon^{2}\right)^{-1}$, this proves (1) for every word with more than $N_{0}$ letters. Thus it suffices to assume $\eta_{0}>0$ small enough such that every $c \in \mathscr{A}^{*}$ such that $\bar{w}(c)<\eta_{0}$ has at least $N_{0}$ letters.

Proof of (2). Let $c \in \mathscr{A}^{*}$. By Proposition 3.6 and since $F_{p}^{\mathscr{d}}$ is moderately dissipative for every $p \in \mathcal{P}$, for every $z=\left(x_{0}, y_{0}\right) \in Y_{p}^{c}$ sent to $\left(x_{1}, y_{1}\right)$ by $F_{p}^{c}$, it holds:

$$
\left|\partial_{y_{0}} \mathcal{Y}_{p}^{c}\left(x_{1}, y_{0}\right)\right|^{\epsilon}=\left|\operatorname{det} D_{\left(x_{0}, y_{0}\right)} F_{p}^{c}\right|^{\epsilon} \cdot\left|\partial_{x_{1}} \mathcal{X}_{p}^{c}\left(x_{1}, y_{0}\right)\right|^{\epsilon} \leq \min _{Y_{p}^{c}}\left\|D F_{p}^{c}\right\|^{-1} \cdot\left|\partial_{x_{1}} \mathcal{X}_{p}^{c}\left(x_{1}, y_{0}\right)\right|^{\epsilon}
$$

As $\left\|D F_{p}^{c}\right\|^{-1} \leq\left|\partial_{x_{1}} \mathcal{X}_{p}^{c}\right|$, it comes:

$$
\begin{equation*}
\max _{Y_{p}^{\bullet}}\left|\partial_{y_{0}} \mathcal{Y}_{p}^{c}\right|^{\epsilon} \leq \min _{Y_{p}^{c}}\left|\partial_{x_{1}} \mathcal{X}_{p}^{c}\right| \cdot \max _{Y_{p}^{\bullet}}\left|\partial_{x_{1}} \mathcal{X}_{p}^{c}\right|^{\epsilon} \leq \max _{Y_{p}^{\bullet}}\left|\partial_{x_{1}} \mathcal{X}_{p}^{c}\right|^{1+\epsilon} \tag{B.8}
\end{equation*}
$$

Thus $\overline{\mathrm{h}}(c)^{\epsilon} \leq \overline{\mathrm{w}}(c)^{1+\epsilon}$. Also by Eq. (B.7) and since $0<\epsilon<1$, if $\overline{\mathrm{w}}(c) \leq \eta_{0}$ then $\overline{\mathrm{h}}(c)^{\epsilon}<\underline{\mathrm{w}}(c)$.

Proof of (2') (and definition of $\rho$ ). Let $F_{p}$ be real analytic for a certain $p \in \mathcal{P}_{0}$. By Proposition A.3, let $\rho>0$ be such that for every $c \in \mathscr{A}^{*}$, the hyperbolic transformation $\left(Y_{p}^{c}, F_{p}^{c}\right)$ admits a $C_{\rho}^{\omega}$ extension $\left(\tilde{Y}_{p}^{c}, F_{p}^{c}\right)$. Then by Eq. (B.8), reducing $\rho$ if necessary, we have:

$$
\begin{equation*}
\max _{\tilde{Y}_{p}^{c}}\left|\partial_{w_{0}} \mathcal{W}_{p}^{c}\right|^{\epsilon} \leq \min _{\tilde{Y}_{p}^{c}}\left|\partial_{z_{1}} \mathcal{Z}_{p}^{c}\right| \cdot \max _{\tilde{Y}_{p}^{c}}\left|\partial_{z_{1}} \mathcal{Z}_{p}^{c}\right|^{\epsilon} \leq \lambda^{-N_{0} \cdot \epsilon} \cdot \min _{\tilde{Y}_{p}^{c}}\left|\partial_{z_{1}} \mathcal{Z}_{p}^{c}\right| \quad \text { if } N_{0} \leq|c| \leq 2 N_{0} \tag{B.9}
\end{equation*}
$$

By Eq. (B.3), $-\epsilon \cdot N_{0} \cdot \log \lambda<\epsilon \log \left(1-\theta^{2}\right)-\log \left(1+\theta^{2}\right)$ and so it holds:

$$
\begin{equation*}
\left(1-\theta^{2}\right)^{-\epsilon} \max _{\tilde{I}^{2}}\left|\partial_{w_{0}} \mathcal{W}_{p}^{c}\right|^{\epsilon}<\left(1+\theta^{2}\right)^{-1} \cdot \min _{\tilde{I}^{2}}\left|\partial_{z_{1}} \mathcal{Z}_{p}^{c}\right| \quad \forall c \in \mathscr{A}^{*} \text { with } N_{0} \leq|c| \leq 2 N_{0} \tag{B.10}
\end{equation*}
$$

Then, given any $c \in \mathscr{A}^{*}$ with $|c| \geq N_{0}$, we split $c$ into words of lengths in $\left[N_{0}, 2 N_{0}\right.$ ] and proceed as in the proof of Eq. (B.6) and Eq. (B.7), using Lemma B. 2 instead of Lemma B.1, to obtain:

$$
\begin{equation*}
\max _{\tilde{I}^{2}}\left|\partial_{z_{1}} \mathcal{W}_{p}^{c}\right|^{\epsilon} \leq \min _{\tilde{I}^{2}}\left|\partial_{z_{1}} \mathcal{Z}_{p}^{c}\right| \tag{B.11}
\end{equation*}
$$

for every $c \in \mathscr{A}^{*}$ such that $|c| \geq N_{0}$ and so every $c \in \mathscr{A}^{*}$ such that $\bar{w}(c) \leq \eta_{0}$.
In Section 3.3, we needed the following distortion result:
Proposition B.3. (i) Let $\left(\mathcal{X}_{j}\right)_{j}$ be a sequence of $C^{2}$-functions $\mathcal{X}_{j}: I^{2} \rightarrow \mathbb{R}$ such that $\partial_{x} \mathcal{X}_{j}$ does not vanish. Let $\left(\mathcal{B}_{j}\right)_{j}$ be a sequence of subsets of $I^{2}$ whose diameters are small when $j$ is large and let $\left(x_{j}^{0}, y_{j}^{0}\right)_{j}$ be a sequence of points in $\mathcal{B}_{j}$. If the distortion $\left(\left\|\partial_{x} \log \left|\partial_{x} \mathcal{X}_{j}\right|\right\|_{C^{0}},\left\|\partial_{y} \log \left|\partial_{x} \mathcal{X}_{j}\right|\right\|_{C^{0}}\right)_{j}$ is bounded, then the following map is $C^{1}$-close to the first coordinate projection when $j$ is large:

$$
\Upsilon_{j}:(x, y) \in \mathcal{B}_{j} \mapsto \frac{\mathcal{X}_{j}(x, y)-\mathcal{X}_{j}\left(x_{j}^{0}, y\right)}{\partial_{x} \mathcal{X}_{j}\left(x_{j}^{0}, y_{j}^{0}\right)}+x_{j}^{0}
$$

(ii) Let $\left(\mathcal{Z}_{j}\right)_{j}$ be a sequence of holomorphic functions $\mathcal{Z}_{j}: \tilde{I}^{2} \rightarrow \tilde{I}$ such that $0<\left|\partial_{z} \mathcal{Z}_{j}\right|<1$. Let $\left(\mathcal{B}_{j}\right)_{j}$ be a sequence of subsets of $\tilde{I}^{2}$ and $\left(z_{j}, w_{j}\right)_{j}$ be a sequence of points such that $\left(z_{j}, w_{j}\right) \in \mathcal{B}_{j}$ for every $j$. We suppose that the diameter of $\mathcal{B}_{j}$ is small compared to $\left|\partial_{z} \mathcal{Z}_{j}\left(z_{j}, w_{j}\right)\right|$ when $j$ is large and that the points $\left(z_{j}, w_{j}\right)$ are uniformly distant to $\partial \tilde{I}^{2}$. Then the following map is $C^{1}$-close to the first coordinate projection when $j$ is large:

$$
\tilde{\Upsilon}_{j}:(z, w) \in \mathcal{B}_{j} \mapsto \frac{\mathcal{Z}_{j}(z, w)-\mathcal{Z}_{j}\left(z_{j}, w\right)}{\partial_{z} \mathcal{Z}_{j}\left(z_{j}, w_{j}\right)}+z_{j}
$$

Proof of Proposition B.3. Let us first prove (i). Up to replacing $\mathcal{B}_{j}$ by $\operatorname{pr}_{1}\left(\mathcal{B}_{j}\right) \times \operatorname{pr}_{2}\left(\mathcal{B}_{j}\right)$ (whose diameter is still small when $j$ is large), we can suppose that $\mathcal{B}_{j}$ is a product. Note that $\Upsilon_{j}\left(x_{j}^{0}, y_{j}^{0}\right)=$ $x_{0}^{j}, \partial_{x} \Upsilon_{j}\left(x_{j}^{0}, y_{j}^{0}\right)=1$ and $\partial_{y} \Upsilon_{j}\left(x_{j}^{0}, y\right)=0$ for every $y \in \operatorname{pr}_{2}\left(\mathcal{B}_{j}\right)$. Thus by the mean value inequality, it suffices to show that the map $\partial_{x}^{2} \Upsilon_{j}=\partial_{x}^{2} \mathcal{X} / \partial_{x} \mathcal{X}\left(x_{j}^{0}, y_{j}^{0}\right)$ and $\partial_{x} \partial_{y} \Upsilon_{j}=\partial_{y} \partial_{x} \mathcal{X} / \partial_{x} \mathcal{X}\left(x_{j}^{0}, y_{j}^{0}\right)$ are bounded since the diameter of $\mathcal{B}_{j}$ is small. This is a direct consequence of the distortion bound:

$$
\begin{aligned}
& \left|\frac{\partial_{x}^{2} \mathcal{X}(x, y)}{\partial_{x} \mathcal{X}\left(x_{j}^{0}, y_{j}^{0}\right)}\right| \leq\left|\partial_{x} \log \right| \partial_{x} \mathcal{X}_{j}(x, y)\left\|\cdot\left|\frac{\partial_{x} \mathcal{X}(x, y)}{\partial_{x} \mathcal{X}\left(x_{j}^{0}, y_{j}^{0}\right)}\right| \leq\right\| D \log \left|\partial_{x} \mathcal{X}_{j}\right| \| \cdot e^{\left|\partial_{x} \log \right| \partial_{x} \mathcal{X}_{j}(x, y) \| \cdot \operatorname{diam} \mathcal{B}_{j}} \\
& \left|\frac{\partial_{y} \partial_{x} \mathcal{X}(x, y)}{\partial_{x} \mathcal{X}\left(x_{j}^{0}, y_{j}^{0}\right)}\right| \leq\left|\partial_{y} \log \right| \partial_{x} \mathcal{X}_{j}(x, y)\left\|\cdot\left|\frac{\partial_{x} \mathcal{X}(x, y)}{\partial_{x} \mathcal{X}\left(x_{j}^{0}, y_{j}^{0}\right)}\right| \leq\right\| D \log \left|\partial_{x} \mathcal{X}_{j}\right| \| \cdot e^{\left|\partial_{y} \log \right| \partial_{x} \mathcal{X}_{j}(x, y)| | \cdot \operatorname{diam} \mathcal{B}_{j}} .
\end{aligned}
$$

We now prove (ii). Let us denote $\sigma_{j}:=\partial_{z} \mathcal{Z}_{j}\left(z_{j}, w_{j}\right)$ for every $j \in \mathbb{N}$ and recall that $0<\left|\sigma_{j}\right|<1$. We notice that for any $(z, w) \in \mathcal{B}_{j}$, it holds:

$$
\partial_{z} \tilde{\Upsilon}_{j}(z, w)=\frac{1}{\sigma_{j}} \partial_{z} \mathcal{Z}_{j}(z, w) \quad \text { and } \quad \partial_{w} \tilde{\Upsilon}_{j}(z, w)=\frac{1}{\sigma_{j}}\left(\partial_{w} \mathcal{Z}_{j}(z, w)-\partial_{w} \mathcal{Z}_{j}\left(z_{j}, w\right)\right)
$$

In particular, we have $\partial_{z} \tilde{\Upsilon}_{j}\left(z_{j}, w_{j}\right)=1$ and $\partial_{w} \tilde{\Upsilon}_{j}\left(z_{j}, w_{j}\right)=0$. Since $\tilde{\Upsilon}_{j}\left(z_{j}, w_{j}\right)=z_{j}$, the tangent map of $\tilde{\Upsilon}_{j}$ at $\left(z_{j}, w_{j}\right)$ is simply the first coordinate projection.

On the other hand, the maps $\mathcal{Z}_{j}: \tilde{I}^{2} \rightarrow \tilde{I}$ are holomorphic on the interior of $\tilde{I}^{2}$ and the subsets $\mathcal{B}_{j}$ of $\tilde{I}^{2}$ are uniformly distant to $\partial \tilde{I}^{2}$ (since by assumption the points $\left(z_{j}, w_{j}\right) \in \mathcal{B}_{j}$ are uniformly distant to the boundary and $\operatorname{diam}\left(\mathcal{B}_{j}\right)$ is small). By the Cauchy integral formula, the second differential of $\mathcal{Z}_{j}$ is then uniformly $C^{0}$-bounded on $\mathcal{B}_{j}$ among $j \in \mathbb{N}$ by a constant $C$. Thus the second differential of $\tilde{\Upsilon}_{j}$ is $C^{0}$-bounded by $2 C \cdot\left|\sigma_{j}\right|^{-1}$ on $\mathcal{B}_{j}$. Since $\operatorname{diam}\left(\mathcal{B}_{j}\right) \cdot 2 C \cdot\left|\sigma_{j}\right|^{-1}$ is small when $j$ is large (by assumption) and $\left(z_{j}, w_{j}\right) \in \mathcal{B}_{j}$, the mean value inequality implies that $\tilde{\Upsilon}_{j}$ is $C^{1}$-close on $\mathcal{B}_{j}$ to its tangent map at $\left(z_{j}, w_{j}\right)$.

## Appendix C. Standard results from hyperbolic theory

Let $M$ be a manifold and $1 \leq r \leq \infty$. We recall that a compact subset $K \subset U$ is hyperbolic for a $C^{r}$-endomorphism $f$ of $M$ if it is invariant $(K=f(K))$ and there exists $n \geq 1$ and open cone fields $\chi_{s}$ and $\chi_{u}$ such that for every $z \in K$, it holds:
(1) the union of $\chi_{s}(z)$ and $\chi_{u}(z)$ is $T_{z} M \backslash\{0\}$,
(2) $D_{z} f$ sends the closure of $\chi_{u}(z)$ into $\chi_{u}(f(z)) \cup\{0\}$ and the preimage by $D_{z} f$ of the closure of $\chi_{s}(f(z))$ is in $\chi_{s}(z) \cup\{0\}$,
(3) every vector in $\chi_{u}(z)$ is expanded by $D_{z} f^{n}$, every vector in $D_{f^{n}(z)} f^{-n}\left(\chi_{s}\left(f^{n}(z)\right)\right)$ is contracted by $D_{z} f^{n}$.
We recall that the inverse limit $\overleftrightarrow{K}:=\left\{\left(x_{i}\right)_{i \in \mathbb{Z}} \in K^{\mathbb{Z}}: f\left(x_{i-1}\right)=x_{i}\right\}$ is endowed with the topology induced by the product one. The map $\overleftrightarrow{f}:=\left(x_{i}\right)_{i \in \mathbb{Z}} \in \overleftrightarrow{K} \mapsto\left(f\left(x_{i}\right)\right)_{i \in \mathbb{Z}} \in \overleftrightarrow{K}$ is semi-conjugate to $f \mid K$ via $h_{0}:\left(x_{i}\right)_{i \in \mathbb{Z}} \mapsto x_{0}$. Note that if $f \mid K$ is a homeomorphism then $h_{0}$ is a homeomorphism.

Theorem C. 1 (Anosov, Quandt [Qua88, Prop 1]). For every $f^{\prime} C^{1}$-close to $f$, there exists a unique continuous map $h_{f^{\prime}}: \overleftrightarrow{K} \rightarrow M$ which is $C^{0}$-close to $h_{0}$ and such that:
$-h_{f^{\prime}} \circ \overleftrightarrow{f}=f^{\prime} \circ h_{f^{\prime}}$,

- $K_{f^{\prime}}:=h_{f^{\prime}}(\overleftrightarrow{K})$ is hyperbolic for $f^{\prime}$, and called the hyperbolic continuation of $K$,
- the map $f^{\prime} \mapsto h_{f^{\prime}}(\overleftrightarrow{k})$ is of class $C^{\infty}$ and depends continuously on $\overleftrightarrow{k} \in \overleftrightarrow{K}$ for the $C^{1}$ topology,
- if $f^{\prime} \mid K_{f^{\prime}}$ is a homeomorphism, then $h_{f^{\prime}}$ is a homeomorphism.

Remark C.2. Quandt did not state that $f^{\prime} \mapsto h_{f^{\prime}}(\overleftrightarrow{k})$ is of class $C^{\infty}$ but he uses the fact that this map is the fixed point of a hyperbolic operator, and so the implicit function theorem implies that $f^{\prime} \mapsto h_{f^{\prime}}(\overleftrightarrow{k})$ is of class $C^{\infty}$.

Remark C.3. If $f$ is a diffeomorphism, then $h_{f^{\prime}}$ is bi-Hölder with exponent close to 1 when $f^{\prime}$ is close to $f$.

For a proof see for instance [Yoc95, Theorem $\S 2.6$ ]. For every $k \in K$ and $\overleftrightarrow{k} \in \overleftrightarrow{K}$, we define their respective stable and unstable manifolds by:

$$
\begin{align*}
W^{s}(k, f)=\left\{k^{\prime} \in M:\right. & \left.\lim _{i \rightarrow \infty} d\left(f^{i}(k), f^{i}\left(k^{\prime}\right)\right)=0\right\} \text { and }  \tag{C.1}\\
& W^{u}(\overleftrightarrow{k}, f)=\left\{k_{0}^{\prime} \in M: \exists\left(k_{i}^{\prime}\right)_{i<0}, f\left(k_{i-1}^{\prime}\right)=k_{i}^{\prime}, \lim _{i \rightarrow-\infty} d\left(k_{i}, k_{i}^{\prime}\right)=0\right\}
\end{align*}
$$

and for $\eta>0$ let:

$$
\begin{align*}
& W_{\eta}^{s}(k, f)=\left\{k^{\prime} \in M: \eta>d\left(f^{i}(k), f^{i}\left(k^{\prime}\right)\right) \xrightarrow{i \rightarrow \infty} 0\right\} \text { and }  \tag{C.2}\\
& \qquad W_{\eta}^{u}(\overleftrightarrow{k}, f)=\left\{k_{0}^{\prime} \in M: \exists\left(k_{i}^{\prime}\right)_{i<0}, f\left(k_{i-1}^{\prime}\right)=k_{i}^{\prime}, \eta>d\left(k_{i}, k_{i}^{\prime}\right) \xrightarrow{i \rightarrow-\infty} 0\right\}
\end{align*}
$$

These sets are properly embedded $C^{r}$ manifolds which depend continuously on $k \in K$ and $\overleftrightarrow{k} \in \overleftrightarrow{K}$ for the $C^{r}$-topology (see for instance [Yoc95, Thm §3.6] or [Ber10, Prop 9.1]). For the sake of simplicity let us assume that all the unstable manifolds have the same dimension $d_{u}$ and all the stable manifolds have the same dimension $d_{s}$. By continuity and invariance, this is the case when $K$ is transitive. Let us recall the following that we are going to generalize and prove in Lemma C.6:

Lemma C. 4 (Inclination Lemma for endomorphism). Let $\Gamma$ be a $d_{u}$-submanifold which is transverse to $W_{\eta}^{s}(k, f)$ for $k \in K$ at a point $z$. Then for every $n$ sufficiently large, there exists a neighborhood $\Gamma^{\prime}$ of $z$ in $\Gamma$ such that $f^{n}\left(\Gamma^{\prime}\right)$ is a submanifold which is uniformly $C^{r}$-close to $W_{\eta}^{u}\left(\overleftrightarrow{k^{n}}, f\right)$, among any $\overleftrightarrow{k^{n}} \in \overleftrightarrow{K}$ such that $h_{0}\left(\overleftrightarrow{k^{n}}\right)=f^{n}(k)$.

For $d \geq 0$ and let $\mathcal{P}$ be a regular compact subset of $\mathbb{R}^{d}$ : it is nonempty and equal to the closure of its interior. We recall that a family $\left(f_{p}\right)_{p \in \mathcal{P}}$ of maps $f_{p} \in C^{r}(M, M)$ is of class $C^{r}$ if the map $(p, x) \in \mathcal{P} \times M \mapsto f_{p}(x) \in M$ extends to a $C^{r}$ map from a neighborhood $\hat{\mathcal{P}} \times M$ of $\mathcal{P} \times M$. By regularity of $\mathcal{P} \times M$, the $C^{r}$-jet of this map is uniquely defined at every point in $\mathcal{P} \times M$. We say that two families are $C^{r}$-close if their $C^{r}$-jet functions are close for the $C^{0}$-compact-open topology.

Assume that $0 \in \mathcal{P}$ and a compact hyperbolic set $K_{0}$ for $f_{0}$ persists by Theorem C. 1 for every $p \in \mathcal{P}$. We denote $h_{p}:=h_{f_{p}}$ the semi-conjugacy and $K_{p}:=h_{p}\left(K_{0}\right)$ its hyperbolic continuation at the parameter $p \in \mathcal{P}$. For $p \in \mathcal{P}$ and ${\overleftrightarrow{k_{0}}}_{\in}^{\overleftrightarrow{K}_{0}}$, we put $k_{p}:=h_{p}\left(\overleftrightarrow{k_{0}}\right)$ and $\overleftrightarrow{k}_{p}:=\left(h_{p} \circ \overleftrightarrow{f_{0}}{ }^{n}\left(\overleftrightarrow{k_{0}}\right)\right)_{n \in \mathbb{Z}}$
Theorem C.5. The families $\left(W_{\eta}^{s}\left(k_{p}, f_{p}\right)\right)_{p \in \mathcal{P}}$ and $\left(W_{\eta}^{u}\left(\overleftrightarrow{k_{p}}, f_{p}\right)\right)_{p \in \mathcal{P}}$ of $C^{r}$-submanifolds are of class $C^{r}$ and depend continuously on respectively $k_{0} \in K_{0}$ and $\overleftrightarrow{k_{0}} \in \overleftrightarrow{K}_{0}$.

We will prove this result below. The following is a generalization of Lemma C. 4 for families.
Lemma C. 6 (Parametric inclination Lemma [Ber16, Ber19, Lemm. 1.7]). Let $\left(\Gamma_{p}\right)_{p \in \mathcal{P}}$ be a $C^{r}$-family of $d_{u}$-submanifolds such that for some $k_{0} \in K_{0}$, the manifold $\Gamma_{p}$ intersects uniformly transversally $W_{\eta}^{s}\left(k_{p}, f_{p}\right)$ at a point $z_{p}$ depending continuously on $p \in \mathcal{P}$. Then when $n \geq 0$ is large, there exists a neighborhood $\Gamma_{p}^{\prime}$ of $z_{p}$ in $\Gamma_{p}$ such that $f_{p}^{n}\left(\Gamma_{p}^{\prime}\right)$ is a submanifold $C^{r}$-close to $W_{\eta}^{u}\left(\overleftrightarrow{k_{p}^{n}}, f_{p}\right)$ and $\left(f_{p}^{n}\left(\Gamma_{p}^{\prime}\right)\right)_{p \in \mathcal{P}}$ is $C^{r}$-close to $\left(W_{\eta}^{u}\left(\overleftrightarrow{k_{p}^{n}}, f_{p}\right)\right)_{p}$ among any $\overleftrightarrow{k_{0}^{n}} \in \overleftrightarrow{K}_{0}$ such that $h_{0}\left(\overleftrightarrow{k_{0}^{n}}\right)=f_{0}^{n}\left(k_{0}\right)$.

Proof. The result being semi-local, we may assume that $f_{p}$ is a map from a small neighborhood $U_{p}$ of $K_{p}$ onto its image for every $p$. Then we extend $f_{p}$ to a larger open set $U_{p}^{\prime} \ni U_{p}$ containing a hyperbolic fixed point $H_{p}$ such that $\Gamma_{p}$ is a slice of the unstable manifold $W^{u}\left(H_{p}, f_{p}\right)$ and such that
$f_{p}^{-n}\left(\Gamma_{p}\right)$ is disjoint from $U_{p}$ for every $n \geq 0$. Then we consider the compact set $\hat{K}_{p}:=K_{p} \cup\left\{f^{n}\left(z_{p}\right)\right.$ : $n \in \mathbb{Z}\} \cup\left\{H_{p}\right\}$. We notice that $\hat{K}_{p}$ is a hyperbolic compact set. Also $\Gamma_{p}^{n}=W_{\eta}^{u}\left(f_{p}^{n}\left(z_{p}\right), f_{p}\right)$. Thus the lemma follows from Theorem C.5.

Proof of Theorem C.5. Let $\hat{\mathcal{P}}$ be an open set such that there exists a $C^{r}$-extension $\left(f_{p}\right)_{p \in \hat{\mathcal{P}}}$ of $\left(f_{p}\right)_{p \in \mathcal{P}}$ for which the compact hyperbolic set persists to $\left(K_{p}\right)_{p \in \hat{\mathcal{P}}}$. Take $p_{0} \in \mathcal{P}$. Let $r>0$ be small such that the balls $B_{r}$ and $B_{2 r}$ centered at $p_{0}$ and of radii $r$ and $2 r$ respectively are included in $\hat{\mathcal{P}}$. As the statement to be proved is local, it suffices to show it for the family $\left(f_{p}\right)_{p \in B_{r}}$. Let

$$
\rho:\left(x_{1}, \ldots, x_{d+1}\right) \in \mathbb{S}^{d}=\left\{\left(x_{1}, \ldots, x_{d+1}\right) \in \mathbb{R}^{d+1}: \sum x_{i}^{2}=1\right\} \mapsto p_{0}+2 r\left(x_{1}, \ldots, x_{d}\right) \in B_{2 r}
$$

Note that $\rho^{-1}\left(B_{r}\right)$ is formed by two components which are diffeomorphically mapped to $B_{r}$ by $\rho$. Thus it suffices to show the proposition for the family $\left(f_{\rho(s)}\right)_{s \in \mathbb{S}^{d}}$. For every $\overleftrightarrow{k} \in \overleftrightarrow{K}_{p_{0}}$, let

$$
\mathcal{L}(\overleftrightarrow{k}):=\left\{\left(s, h_{\rho(s)}(\overleftrightarrow{k})\right): s \in \mathbb{S}^{d}\right\} \subset \mathbb{S}^{d} \times M
$$

By Theorem C.1, the sphere $\mathcal{L}(\overleftrightarrow{k})$ are of class $C^{r}$ and depends continuously on $\overleftrightarrow{k} \in \overleftrightarrow{K_{p_{0}}}$ for this topology. Hence $\bigcup \overleftrightarrow{k} \in \overleftrightarrow{K} \mathcal{L}(\overleftrightarrow{k})$ is the image of $\mathbb{S}^{d} \times \overleftrightarrow{K}$ by a $C^{r}$-immersion of this lamination. Moreover this lamination is $r$-normally hyperbolic in the sense of [Ber10, Def 2.1]. Consequently, we can apply [Ber10, Prop 9.1] which states that the strong stable and unstable local manifolds of a leaf form $C^{r}$-immersed laminations. This is equivalent to say that $\left(W_{\eta}^{s}\left(k_{\rho(s)} ; f_{\rho(s)}\right)_{s \in \mathbb{S}^{d}}\right.$ and $\left(W_{\eta}^{u}\left(\overleftrightarrow{k}_{\rho(s)} ; f_{\rho(s)}\right)_{s \in \mathbb{S}^{d}}\right.$ are $C^{r}$-families which depends continuously on $k_{p_{0}} \in K_{p_{0}}$ and $\overleftrightarrow{k}_{p_{0}} \in$ $\overleftrightarrow{K}_{p_{0}}$

The following enables to extend the stable lamination to a $C^{1}$-foliation, by integrating the $C^{1}$-line field $\ell(p, \cdot)$.

Theorem C.7. If $\left(f_{p}\right)_{p \in \mathcal{P}}$ is a $C^{2}$-family of surface diffeomorphisms leaving invariant the continuation $\left(K_{p}\right)_{p}$ of a hyperbolic set, there exists a neighborhood $\hat{V}=\bigcup_{p \in \mathcal{P}}\{p\} \times V_{p}$ of $\bigcup_{p \in \mathcal{P}}\{p\} \times K_{p}$ and a $C^{1}$-function $\ell$ from $\hat{V}$ into $\mathbb{P}(\mathbb{R})$ satisfying for every $p \in \mathcal{P}$ :
$\ell(p, z)=T_{z} W_{\eta}^{s}(k) \quad$ if $z \in W_{\eta}^{s}(k)$ with $k \in K_{p} \quad$ and $\quad D_{z} f_{p}(\ell(p, z))=\ell\left(p, f_{p}(z)\right) \quad$ if $z \in V_{p} \cap f_{p}^{-1}\left(V_{p}\right)$.
Proof. The proof is the parametric counterpart of [BC16, Coro 1.11]. Let $\mathbb{P}(T M)$ be the bundle over $M$ whose fibers at $z \in M$ is the Grassmanian $\mathbb{P}\left(T_{z} M\right)$ of lines in $T_{z} M$. We consider:

$$
\hat{f}:(p, z, \ell) \in \hat{\mathcal{P}} \times \mathbb{P}(T M) \mapsto\left(p, f_{p}(z), D_{z} f_{p}(\ell)\right) \in \hat{\mathcal{P}} \times \mathbb{P}(T M)
$$

where $\hat{\mathcal{P}}$ is the open neighborhood of $\mathcal{P}$ to which the family extends. We remark that $\hat{f}$ is of class $C^{1}$ on $\hat{M}:=\hat{\mathcal{P}} \times \mathbb{P}(T M)$. We denote $E_{p}^{s}(k):=T_{k} W^{s}\left(k, f_{p}\right)$ and $E_{p}^{u}(k):=T_{k} W^{u}\left(k, f_{p}\right)$ for every $k \in K_{p}$ and $p \in \mathcal{P}$. Let $\hat{K}:=\left\{\left(p, k, E_{p}^{s}(k)\right): p \in \mathcal{P} \quad\right.$ and $\left.\quad k \in K_{p}\right\}$. It is a partially hyperbolic compact set for $\hat{f}$ with strong unstable direction at $\left(p, k, E_{p}^{s}(k)\right)$ equal to $\{(0,0)\} \times T_{E_{p}^{s}(k)} \mathbb{P}\left(T_{k} M\right)$. The strong unstable manifold of a point $\left(p, k, E_{p}^{s}(k)\right)$ is $\left(p, k, \mathbb{P}\left(T_{k} M\right) \backslash\left\{E_{p}^{u}(k)\right\}\right)$. It intersects $\hat{K}$ only at one point. Thus we can apply the following:
Theorem C. 8 ([BC16, Main result]). Let $\hat{f}$ be a $C^{1}$-diffeomorphism of a manifold $\hat{M}$ and $\hat{K}$ a partially hyperbolic compact invariant set such that $T \hat{M} \mid \hat{K}=E^{c} \oplus E^{u u}$. Then, the next two properties are equivalent.
(1) There exists a compact $C^{1}$-submanifold $S$ with boundary which:

- contains $\hat{K}$ in its interior,
- is tangent to $E^{c}$ at each point of $\hat{K}$ (i.e. $T_{x} S=E^{c}(x)$ for each $\left.x \in \hat{K}\right)$,
- is locally invariant: $\hat{f}(S)$ contains a neighborhood of $\hat{K}$ in $S$,
(2) The strong unstable manifold of any $x \in \hat{K}$ intersect $\hat{K}$ only at $x$ (i.e. $W^{u u}(x) \cap \hat{K}=\{x\}$ ).

Remark C.9. By [BC16, Coro. 1.5.], if $\hat{f}$ is of class $C^{1+}$, then $S$ can be chosen of class $C^{1+}$.
As the strong unstable direction is the last coordinate, up to restricting $S$, there exists a function from a neighborhood $\hat{V}$ of $\hat{K}$ and a $C^{1}$-function $\ell$ from $\hat{V}$ into $\mathbb{P}(\mathbb{R})$ whose graph is equal to $S$. Thus by invariance and since the action of $\hat{f}$ on the two first coordinates is $(p, z) \mapsto\left(p, f_{p}(z)\right)$, we have:

$$
\hat{f}(p, z, \ell(p, z))=\left(p, f_{p}(z), \ell\left(p, f_{p}(z)\right)\right), \quad \forall p \in \mathcal{P} \text { and } z \in V_{p} \cap f_{p}^{-1}\left(V_{p}\right)
$$

As by definition $\hat{f}(p, z, \ell(p, z))=\left(p, f_{p}(z), D_{z} f_{p}(\ell(p, z))\right)$, we obtain $D_{z} f_{p}(\ell(p, z))=\ell\left(p, f_{p}(z)\right)$. Given $p \in \mathcal{P}, k \in K_{p}$ and $z \in W_{\eta}^{s}\left(k, f_{p}\right)$, if $\ell(p, z) \neq T_{z} W_{\eta}^{s}\left(k, f_{p}\right)$, then by invariance $\ell\left(p, f_{p}^{n}(z)\right)=$ $D_{z} f_{p}^{n}(\ell(p, z)) \neq D_{z} f_{p}^{n}\left(T_{z} W_{\eta}^{s}\left(k, f_{p}\right)\right)=T_{f_{p}^{n}(z)} W_{\eta}^{s}\left(f_{p}^{n}(k), f_{p}\right)$ for every $n \geq 0$. Moreover, by hyperbolicity the angle between $D_{z} f_{p}^{n}(\ell(p, z))$ and $D_{z} f_{p}^{n}\left(T_{z} W_{\eta}^{s}\left(k, f_{p}\right)\right)$ is bounded from below. This contradicts that the angle between $\ell\left(p, f_{p}^{n}(z)\right), \ell\left(p, f_{p}^{n}(k)\right)$ and $T_{f_{p}^{n}(z)} W_{\eta}^{s}\left(f_{p}^{n}(k), f_{p}\right)$ goes to 0 as $d\left(f_{p}^{n}(z), f_{p}^{n}(k)\right) \rightarrow 0$.

Using Remark C. 9 we obtain:
Corollary C.10. If under the assumptions of Theorem C.7, the family $\left(f_{p}\right)_{p \in \mathcal{P}}$ is of class $C^{2+}$, then the function $\ell$ given by Theorem $C .7$ can be assumed of class $C^{1+}$.

Appendix D. Modeling a horseshoe by a hyperbolic map of type $\mathscr{A}$
Proof of Proposition 2.25. Construction of the maps $\breve{H}_{p}$ and verification of property (0). Let $\mathcal{P}$ be a closed ball centered at 0 included in $\hat{\mathcal{P}}$. Hyperbolicity implies the existence of a splitting $T M \mid K_{p}=E_{p}^{s} \oplus E_{p}^{u}$ of two $D f_{p}$-invariant directions which are respectively contracted and expanded by $D f_{p}$. Let $\hat{K}=\bigcup_{p \in \mathcal{P}}\{p\} \times K_{p}$. By Theorem C. 7 Page 73, we have:

Theorem D.1. There exist a neighborhood $\hat{U}=\bigcup_{p \in \mathcal{P}}\{p\} \times U_{p}$ of $\hat{K}$ and two $C^{1}$-families $\left(e_{p}^{s}\right)_{p}$ and $\left(e_{p}^{u}\right)_{p}$ of unit vector fields $e_{p}^{s}$ and $e_{p}^{u}$ on $U_{p}$ such that:
(1) for every $z \in K_{p}$, the restriction $e_{p}^{s} \mid W_{\text {loc }}^{s}\left(z, f_{p}\right) \cap U_{p}$ is a section of $T W_{l o c}^{s}\left(z, f_{p}\right) \cap U_{p}$ and $e_{p}^{u} \mid W_{\text {loc }}^{u}\left(z, f_{p}\right) \cap U_{p}$ is a section of $T W_{l o c}^{u}\left(z, f_{p}\right) \cap U_{p}$,
(2) for every $z \in U_{p} \cap f_{p}^{-1}\left(U_{p}\right)$, the vector $D_{z} f_{p}\left(e_{p}^{s}(z)\right)$ is co-linear to $e_{p}^{s}\left(f_{p}(z)\right)$ and the vector $D_{z} f_{p}\left(e_{p}^{u}(z)\right)$ is co-linear to $e_{p}^{u}\left(f_{p}(z)\right)$.
Let $\phi_{p}^{s}$ and $\phi_{p}^{u}$ be the flows associated to these vector fields for every $p \in \mathcal{P}$. By compactness of $\mathcal{P}$, there exists $\eta>0$ such that for every $z \in K_{p}$ and $p \in \mathcal{P}$, the following is a diffeomorphism onto its image included in $U_{p} \cap f_{p}^{-1}\left(U_{p}\right)$ :

$$
\left(t, t^{\prime}\right) \in(-\eta, \eta)^{2} \mapsto \phi_{p}^{s}\left(\phi_{p}^{u}(z, t), t^{\prime}\right)
$$

By [Bow08, Thm. 3.12], there is a Markov partition of $K_{0}$ by rectangle of small diameter. This is a finite partition $\left(R_{0}^{v}\right)_{v \in \mathscr{V}}$ by clopen subsets $R_{0}^{v} \subset K_{0}$ satisfying:

- Each subset $R_{0}^{v}$ is a rectangle: there exist two segments $I_{0}^{v}, J_{0}^{v} \subset(-\eta, \eta)$ satisfying
(D.1) $\quad R_{0}^{v}=\breve{\phi}_{0}^{u}\left(I_{0}^{v} \times J_{0}^{v}\right) \cap K_{0} \quad$ with $\quad \breve{\phi}_{0}^{u}:\left(t, t^{\prime}\right) \in(-\eta, \eta)^{2} \mapsto \phi_{0}^{s}\left(\phi_{0}^{u}\left(z_{0}^{u}, t\right), t^{\prime}\right)$ and $z_{0}^{u} \in R_{0}^{v}$.
- For every pair $(v, w) \in \mathscr{V}^{2}$ such that $R_{0}^{v w}:=R_{0}^{v} \cap f_{0}^{-1}\left(R_{0}^{w}\right) \neq \emptyset$, there are two ${ }^{11}$ segments $I_{0}^{u \omega} \subset I_{0}^{u}$ and $J_{0}^{u \omega} \subset J_{0}^{w}$ such that:

$$
\begin{equation*}
R_{0}^{v w}=\breve{\phi}_{0}^{v}\left(I_{0}^{v w} \times J_{0}^{v}\right) \cap K_{0} \quad \text { and } \quad f_{0} \circ \breve{\phi}_{0}^{v}\left(I_{0}^{v w} \times J_{0}^{v}\right)=\breve{\phi}_{0}^{w}\left(I_{0}^{w} \times J_{0}^{v w}\right) \tag{D.2}
\end{equation*}
$$

We assume the segments $I_{0}^{v}$ and $J_{0}^{u}$ minimal such that Eq. (D.1) holds true. Then $I_{0}^{v w}$ and $J_{0}^{v w}$ are uniquely defined by Eq. (D.2). Let $\left(R_{p}^{v}\right)_{p \in \mathcal{P}},\left(R_{p}^{v w}\right)_{p \in \mathcal{P}}$ and $\left(z_{p}^{v}\right)_{p \in \mathcal{P}}$ be the hyperbolic continuations of $R_{0}^{v}, R_{0}^{v \omega}$ and $z_{0}^{v}$ (see Theorem C. 1 Page 71). For $\mathcal{P}$ sufficiently small, for every $p \in \mathcal{P}$, there are minimal segments $I_{p}^{v}$ and $J_{p}^{v}$ included in $(-\eta, \eta)$ satisfying

$$
\begin{equation*}
R_{p}^{v}=\breve{\phi}_{p}^{v}\left(I_{p}^{v} \times J_{p}^{v}\right) \cap K_{p} \quad \text { with } \quad \breve{\phi}_{p}^{v}:\left(t, t^{\prime}\right) \in(-\eta, \eta)^{2} \mapsto \breve{\phi}_{p}^{s}\left(\phi_{p}^{u}\left(z_{p}^{v}, t\right), t^{\prime}\right) \tag{D.3}
\end{equation*}
$$

and there are segments $I_{p}^{v \omega} \subset I_{p}^{v}$ and $J_{p}^{v \omega} \subset J_{p}^{w}$ such that:

$$
\begin{equation*}
R_{p}^{v w}=\breve{\phi}_{p}^{v}\left(I_{p}^{v w} \times J_{p}^{v}\right) \cap K_{p} \quad \text { and } \quad f_{p} \circ \breve{\phi}_{p}^{u}\left(I_{p}^{v w} \times J_{p}^{v}\right)=\breve{\phi}_{p}^{w}\left(I_{p}^{w} \times J_{p}^{v w}\right) \tag{D.4}
\end{equation*}
$$

By minimality, the segments $I_{p}^{v}, J_{p}^{v}, I_{p}^{v \omega}$ and $J_{p}^{v \omega}$ vary continuously with $p \in \mathcal{P}$. Let:

$$
Q_{p}^{v}:=\breve{\phi}_{p}^{v}\left(I_{p}^{v} \times J_{p}^{v}\right), \quad \partial^{s} Q_{p}^{v}:=\breve{\phi}_{p}^{v}\left(\partial\left(I_{p}^{v}\right) \times J_{p}^{v}\right) \quad \text { and } \quad \partial^{u} Q_{p}^{v}:=\breve{\phi}_{p}^{v}\left(I_{p}^{v} \times \partial\left(J_{p}^{v}\right)\right)
$$

Fact D.2. The compact sets $\left(Q_{p}^{v}\right)_{v}$ are disjoint for every $p \in \mathcal{P}$.
Proof. For every $v \in \mathscr{V}$, by minimality of $I_{p}^{v}$ and Eq. (D.3), each of the two components of $\partial^{s} Q_{p}^{v}$ contains a point of $K_{p}$. Hence, by Theorem D. 1 (1), each of the two components of $\partial^{s} Q_{p}^{u}$ is equal to a local stable manifold of it. Likewise by minimality of $J_{p}^{v}$, Eq. (D.3) and Theorem D. 1 (1), each of the two components of $\partial^{u} Q_{p}^{v}$ is equal to a local unstable manifold of a point of $K_{p}$. Hence by local maximality of $K_{p}$, the corners $\partial^{u} Q_{p}^{u} \cap \partial^{s} Q_{p}^{v}$ of $Q_{p}^{v}$ are in $K_{p}$. Hence for $v, w \in \mathscr{V}$, if $Q_{p}^{u}$ intersects $Q_{p}^{w}$ then the sets $Q_{p}^{v}$ and $Q_{p}^{w}$ are nested or a stable boundary of $Q_{p}^{u}$ (resp. $Q_{p}^{w}$ ) intersects an unstable boundary of $Q_{p}^{w}$ (resp. $Q_{p}^{v}$ ). Thus if $Q_{p}^{v}$ intersects $Q_{p}^{w}$, then $R_{p}^{v}=Q_{p}^{v} \cap K_{p}$ intersects $R_{p}^{w}=Q_{p}^{w} \cap K_{p}$ and so $v=w$.

Let $k \geq 1$ be a large integer so that $f_{p}^{k} \mid K_{p}$ is still transitive. Put:
$\mathscr{A}:=\left\{\left(v_{1}, \ldots, v_{k}\right) \in \mathscr{V}^{k}: \bigcap_{j=1}^{k} f_{p}^{-j}\left(R_{p}^{v_{j}}\right) \neq \emptyset\right\}, \quad a\left(v_{1}, \ldots, v_{k}\right)=v_{1} \quad$ and $\quad t\left(v_{1}, \ldots, v_{k}\right)=v_{k}$.
Put $R_{p}^{a}:=\bigcap_{j=1}^{k} f_{p}^{-j}\left(R_{p}^{v_{j}}\right)$ for every $a \in \mathscr{A}$ and $p \in \mathcal{P}$. We recall that each $\breve{\phi}_{p}^{w}$ has its range in $U_{p} \cap f_{p}^{-1}\left(U_{p}\right)$, and so each map $\left(\breve{\phi}_{p}^{w}\right)^{-1} \circ f_{p} \circ \breve{\phi}_{p}^{v}$ preserves the horizontal and vertical directions, and respectively expands and contracts them. Thus by Eq. (D.4), for every $a \in \mathscr{A}$, there exist unique $I_{p}^{a} \subset I_{p}^{o(a)}$ and $J_{p}^{a} \subset J_{p}^{t(a)}$ such that:
(D.5) $\quad R_{p}^{a}=\breve{\phi}_{p}^{u}\left(I_{p}^{a} \times J_{p}^{a(a)}\right) \cap K_{p} \quad$ and $\quad f_{p}^{k} \circ \breve{\phi}_{p}^{a(a)}\left(I_{p}^{a} \times J_{p}^{a(a)}\right)=\breve{\phi}_{p}^{t(a)}\left(I_{p}^{t(a)} \times J_{p}^{a}\right)$.

For every $a \in \mathscr{A}$ and $p \in \mathcal{P}$, the following map preserves again the horizontal and vertical directions

$$
\begin{equation*}
\breve{f}_{p}^{a}:=\left(\breve{\phi}_{p}^{t(a)}\right)^{-1} \circ f_{p}^{k} \circ \breve{\phi}_{p}^{o(a)} \mid I_{p}^{a} \times J_{p}^{a(a)} . \tag{D.6}
\end{equation*}
$$

11. When the diameters of the rectangle are not small one has to consider more segments.

We assume $k$ large enough so that it satisfies moreover:

$$
\begin{equation*}
\min \left|\partial_{x} \breve{f}_{p}^{a}\right|>\frac{C \cdot \lambda}{\theta^{2}}, \quad \max \left|\partial_{y} \breve{f}_{p}^{a}\right|<\frac{\theta^{2}}{C \cdot \lambda} \quad \text { where } C:=\max _{(v, w) \in \mathscr{V}^{2}}\left\{\frac{\left|I_{p}^{v}\right|}{\left|I_{p}^{w}\right|}, \frac{\left|J_{p}^{v}\right|}{\left|J_{p}^{w}\right|}\right\} \tag{D.7}
\end{equation*}
$$

Thus there are $\theta^{2} /(C \lambda)$-contracting maps $\breve{x}_{p}^{a}: I_{p}^{t(a)} \rightarrow I_{p}^{a}$ and $\breve{y}_{p}^{a}: J_{p}^{a(a)} \rightarrow J_{p}^{a}$ such that:

$$
\begin{equation*}
\breve{f}_{p}^{a}\left(\breve{x}_{p}^{a}\left(x_{1}\right), y_{0}\right)=\left(x_{1}, \breve{y}_{p}^{a}\left(y_{0}\right)\right), \tag{D.8}
\end{equation*}
$$

for any $\left(x_{1}, y_{0}\right) \in I_{p}^{t(a)} \times J_{p}^{a(a)}$. Note that $\breve{x}_{p}^{a}$ and $\breve{y}_{p}^{a}$ are defined via Eqs. (D.6) and (D.8) on some $\eta^{\prime}$-neighborhood $V_{p}^{t(a)}$ and $W^{a(a)}$ of $I_{p}^{t(a)}$ and $J_{p}^{a(a)}$ in such a way that Eq. (D.8) remains valid. If $\eta^{\prime}$ is sufficiently small, these maps are still $\frac{\theta^{2}}{C \cdot \lambda}$-contracting by Eq. (D.7). Let $\eta^{\prime}>0$ be small enough so that these properties are valid for every $a$ in the finite set $\mathscr{A}$ and $p$ in the compact set $\mathcal{P}$. By contraction, it holds:

$$
\begin{equation*}
\breve{x}_{p}^{a}\left(V_{p}^{t(a)}\right) \Subset V_{p}^{o(a)} \quad \text { and } \quad \breve{y}_{p}^{a}\left(W_{p}^{o(a)}\right) \Subset W_{p}^{t(a)} \tag{D.9}
\end{equation*}
$$

For every $v \in \mathscr{V}$, let $\Delta_{p}^{v}$ be the affine map with diagonal and positive linear part which sends $Y^{e}$ to $V_{p}^{v} \times W_{p}^{u}$. We define the $C^{1}$-family $\left(\breve{H}_{p}\right)_{p \in \mathcal{P}}$ by:

$$
\breve{H}_{p}:(z, v) \in Y^{e} \times \mathscr{V} \mapsto \breve{H}_{p}^{v}(z) \in M, \quad \text { with } \breve{H}_{p}^{v}: z \in Y^{e} \mapsto \breve{\phi}_{p}^{v} \circ \Delta_{p}^{v}(z)
$$

Let $\breve{Y}_{p}^{a}$ be the preimage by $\Delta_{p}^{a(a)}$ of $\breve{x}_{p}^{a}\left(V_{p}^{t(a)}\right) \times W^{a(a)}$. By Eq. (D.9), $\breve{Y}_{p}^{a}$ is included in $Y^{e} \backslash \partial^{s} Y^{e}$. Also, $\breve{Y}_{p}^{a}$ is sent by $\breve{f}_{p}^{a} \circ \Delta_{p}^{a(a)}$ onto $V_{p}^{t(a)} \times \breve{y}_{p}^{a}\left(W^{o(a)}\right) \subset \Delta_{p}^{t(a)}\left(Y^{e} \backslash \partial^{u} Y^{e}\right)$. Put:

$$
\breve{F}_{p}^{a}:=\left(\Delta_{p}^{t(a)}\right)^{-1} \circ \breve{f}_{p}^{a} \circ \Delta_{p}^{a(a)}: \breve{Y}_{p}^{a} \rightarrow Y^{e} \backslash \partial^{u} Y^{e}
$$

Let $\left(\breve{F}_{p}^{\mathscr{A}}\right)_{p \in \mathcal{P}}$ be the $C^{1}$-family of maps $\breve{F}_{p}^{\mathscr{A}}$ from $\breve{D}_{p}(\mathscr{A}):=\bigsqcup_{a \in \mathscr{A}} \breve{Y}_{p}^{a} \times\{a(a)\}$ equal to $\left(\breve{F}_{p}^{a}, t(a)\right)$ at $\breve{Y}_{p}^{a} \times\{a(a)\}$. Then Proposition $2.25(0)$ holds true since $\breve{F}_{p}^{\mathscr{A}}$ leaves invariant the vertical and horizontal directions, and by Eq. (D.6):

$$
\begin{equation*}
f_{p}^{k} \circ \breve{H}_{p}^{o(a)} \mid \breve{Y}_{p}^{a}=\breve{\phi}_{p}^{t(a)} \circ \breve{f}_{p}^{a} \circ \Delta_{p}^{o(a)}=\breve{H}_{p}^{t(a)} \circ \breve{F}_{p}^{a} \tag{D.10}
\end{equation*}
$$

Verification of property (1). Let $\left(H_{p}\right)_{p \in \mathcal{P}}$ be a $C^{r}$-family which is $C^{1}$-close to $\left(\breve{H}_{p}\right)_{p \in \mathcal{P}}$. For every $v \in \mathscr{V}$ and $z \in Y^{e}$, let $H_{p}^{v}(z):=H_{p}(z, v)$.

Fact D.3. Every nonzero vector in the complement of $\chi_{v}$ is sent by $D \breve{F}_{p}^{\text {sl }}$ into $\chi_{h}$ and $\lambda$-expands its horizontal component, also $D\left(\breve{F}_{p}^{\text {sl }}\right)^{-1}$ sends every nonzero vector in the complement of $\chi_{h}$ into $\chi_{v}$ and expands its vertical component.

Proof. This is a consequence of Eq. (D.7) which implies $\left|\partial_{x} \breve{F}_{p}^{\text {sd }}\right|>\lambda / \theta^{2}$ and $\left|\partial_{y} \breve{F}_{p}^{\text {sd }}\right|<\theta^{2} / \lambda$.
Also by Eq. (D.10) and then Eqs. (D.8) and (D.9), for every $\left(x_{1}, y_{0}\right) \in I^{2}$, there exist $\breve{\mathcal{X}}_{p}^{a}\left(x_{1}\right)$ and $\breve{\mathcal{Y}}_{p}^{a}\left(y_{0}\right)$ in the interior of $I$ such that:

$$
f_{p}^{k} \circ \breve{H}_{p}^{a(a)}\left(\breve{\mathcal{X}}_{p}^{a}\left(x_{1}\right), y_{0}\right)=\breve{H}_{p}^{t(a)} \circ \breve{F}_{p}^{a}\left(\breve{\mathcal{X}}_{p}^{a}\left(x_{1}\right), y_{0}\right)=\breve{H}_{p}^{t(a)}\left(x_{1}, \breve{\mathcal{Y}}_{p}^{a}\left(y_{0}\right)\right)
$$

Thus, as $H_{p}^{a(a)}$ and $H_{p}^{t(a)}$ are $C^{1}$-close to $\breve{H}_{p}^{a(a)}$ and $\breve{H}_{p}^{t(a)}$, by hyperbolicity of $\breve{F}_{p}^{a}$ and the implicit function theorem, there exist unique points $\mathcal{X}_{p}^{a}\left(x_{1}, y_{0}\right)$ and $\mathcal{Y}_{p}^{a}\left(x_{1}, y_{0}\right)$ close to $\breve{\mathcal{X}}_{p}^{a}\left(x_{1}\right)$ and $\breve{\mathcal{Y}}_{p}^{a}\left(y_{0}\right)$ satisfying :

$$
f_{p}^{k} \circ H_{p}^{o(a)}\left(\mathcal{X}_{p}^{a}\left(x_{1}, y_{0}\right), y_{0}\right)=H_{p}^{t(a)}\left(x_{1}, \mathcal{Y}_{p}^{a}\left(x_{1}, y_{0}\right)\right)
$$

Still by the implicit function theorem, the map $\left(p, x_{1}, y_{0}\right) \mapsto\left(\mathcal{X}_{p}^{a}, \mathcal{Y}_{p}^{a}\right)\left(x_{1}, y_{0}\right)$ is of class $C^{r}$ and is $C^{1}$-close to $\left(p, x_{1}, y_{0}\right) \mapsto\left(\breve{\mathcal{X}}_{p}^{a}\left(x_{1}\right), \mathcal{Y}_{p}^{a}\left(y_{0}\right)\right)$. Hence the following set:

$$
Y_{p}^{a}:=\left\{\left(\mathcal{X}_{p}^{a}\left(x_{1}, y_{0}\right), y_{0}\right):\left(x_{1}, y_{0}\right) \in I^{2}\right\}
$$

is a box included in $Y^{e} \backslash \partial^{s} Y^{e}$. Moreover it is sent by $F_{p}^{a}:=\left(H_{p}^{t(a)}\right)^{-1} \circ f^{k} \circ H_{p}^{a(a)} \mid Y_{p}^{a}$ into $\left\{\left(x_{1}, \mathcal{Y}_{p}^{a}\left(x_{1}, y_{0}\right)\right):\left(x_{1}, y_{0}\right) \in I^{2}\right\} \subset Y^{e} \backslash \partial^{u} Y^{e}$. Hence by Fact D.3, $\left(F_{p}^{a}, Y_{p}^{a}\right)$ is a hyperbolic transformation. Also the $C^{r}$-family of maps $F_{p}^{\mathscr{A}}$ from $D_{p}(\mathscr{A}):=\bigsqcup_{a \in \mathscr{A}} Y_{p}^{a} \times\{a(a)\}$ equal to $\left(F_{p}^{a}, t(a)\right)$ at $Y_{p}^{a} \times\{a(a)\}$ satisfies Proposition 2.25 (1).

Verification of property (2). If $\lim \sup _{n \rightarrow \infty}\left\|\operatorname{det} D f_{0}^{n}\left|K_{0}\left\|^{1 / \epsilon} \cdot\right\| D f_{0}^{n}\right| K_{0}\right\|<1$, then for $k$ large enough, $\breve{F}_{0}^{s d}$ and so $F_{0}^{d}$ are moderately dissipative. So are any $F_{p}^{d}$ for $p$ in a sufficiently small $\mathcal{P}$.

Verification of property (3). We notice that the pull back by $H_{p}$ of the vector field $e_{p}^{s}$ is $C^{1}$ close to be tangent to $\{0\} \times \mathbb{R}$ and is defined on $Y^{e} \times \mathscr{V}$. Thus we can extend it to a $C^{1}$-nonzero-vector field on $\mathbb{R} \times I \times \mathscr{V}$ with values in $\chi_{v}$. Then the holonomy $\pi_{p}$ of this vector field restricted to the transverse section $\mathbb{R} \times\{0\} \times \mathscr{V}$ defines a family of projections $\left(\pi_{p}\right)_{p \in \mathcal{P}}$ satisfying (0)-(1)-(2)-(3) of Definition 1.23 , which shows property (3) of Proposition 2.25 and concludes the proof.

Proof of Corollary 2.26. We use Corollary C. 10 while using Theorem C. 7 in the above proof.

## Appendix E. A general result on Emergence

The aim of this section is to prove Theorem 0.4 Page 5. Let $K$ is a horseshoe of a $C^{1+\alpha}$-surface diffeomorphism $f$. We want to show that the covering number $\mathcal{N}(\epsilon)$ at scale $\epsilon$ of $\mathcal{M}_{f}(K)$ for the Wasserstein distance (see Section 0.2) has order at least the unstable dimension $d_{u}$ of $K$ :

$$
\liminf _{\epsilon \rightarrow 0} \frac{\log \log \mathcal{N}(\epsilon)}{-\log \epsilon} \geq d_{u}
$$

Proof of Theorem 0.4. First we can assume $f$ of class $C^{\infty}$. Indeed, we can consider a smooth approximation $\tilde{f}$ of $f$. Then by Remark C.3, the restriction of $\tilde{f}$ to the continuation $\tilde{K}$ of $K$ is conjugated to $f \mid K$ via a bi-Hölder conjugacy with exponent close to 1 . Thus the unstable dimension of $\tilde{K}$ is close to the one of $K$ and the space of invariant probability measures $\mathcal{M}_{\tilde{f}}(\tilde{K})$ of $\tilde{f} \mid K$ is in bijection to $\mathcal{M}_{f}(K)$ via a bi-Hölder map with exponent close to 1 . So the orders of their covering numbers are close.

Thus we assume that $f$ is of class $C^{\infty}$. Then by Corollary 2.26 of Proposition 2.25 (for $d=0$ ), the map $f \mid K$ is $C^{1+}$-semi-conjugate (via $\pi$ ) to a $C^{1+}$-function $g$ on a disjoint union of segments of the real line, which leaves invariant an expanding Cantor set $\Lambda_{g}$ of Hausdorff dimension equal to the unstable dimension $d^{u}$ of $K$.

Then by [BB19, Thm A] the covering number $\mathcal{N}_{g}$ of the space of invariant measures of $g \mid \Lambda_{g}$ has order $d^{u}$ :

$$
\log \log \mathcal{N}_{g}(\eta) \sim-d^{u} \cdot \log \eta, \quad \text { as } \eta \rightarrow 0
$$

The push forward by the Lipshitz semi-conjugacy is a Lipschitz map from the space of invariant measure of $F^{g}$ to those of $g$. By the same reasoning as in Lemma 4.35, this pushforward is a bijection between the spaces of invariant measures of $g$ and $F^{s}$. Thus the covering number of the space of invariant measure of $F^{\mathscr{A}}$ has the sought property.

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Pierre Berger
Sébastien Biebler
pierre.berger@imj-prg.fr sebastien.biebler@imj-prg.fr
Sorbonne Université, Université de Paris, CNRS,
Institut de Mathématiques de Jussieu-Paris Rive Gauche
F-75005 Paris, France


[^0]:    3. The argument is contained in the proof of Corollary 1.50 .
[^1]:    6. See Theorem C. 1 Page 71.
    7. This assumption is actually not necessary using Duarte's theorem [Dua08], but sufficient for all our applications needing the moderate dissipativeness assumption.
[^2]:    8. a $q$-periodic point $P$ satisfying $\left|\operatorname{det} D f^{q}(P)\right|<1$.
