# Toward a theory of rank one attractors 

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## Introduction

This paper is about a class of strange attractors that have the dual property of occurring naturally and being amenable to analysis. Roughly speaking, a rank one attractor is an attractor that has some instability in one direction and strong contraction in $m-1$ directions, $m$ here being the dimension of the phase space.

The results of this paper can be summarized as follows. Among all maps with rank one attractors, we identify, inductively, subsets $\mathcal{G}_{n}, n=1,2,3, \cdots$,

[^0]consisting of maps that are "well-behaved" up to the $n$th iterate. The maps in $\mathcal{G}:=\cap_{n>0} \mathcal{G}_{n}$ are then shown to be nonuniformly hyperbolic in a controlled way and to admit natural invariant measures called SRB measures. This is the content of Part II of this paper. The purpose of Part III is to establish existence and abundance. We show that for large classes of 1-parameter families $\left\{T_{a}\right\}$, $T_{a} \in \mathcal{G}$ for positive measure sets of $a$.

Leaving precise formulations to Section 1, we first put our results into perspective.
A. In relation to hyperbolic theory. Axiom A theory, together with its extension to the theory of systems with invariant cones and discontinuities, has served to elucidate a number of important examples such as geodesic flows and billiards (see e.g. [Sm], [A], [Si1], [B], [Si2], [W]). The invariant cones property is quite special, however. It is not enjoyed by general dynamical systems.

In the 1970 s and 80 s , an abstract nonuniform hyperbolic theory emerged. This theory is applicable to systems in which hyperbolicity is assumed only asymptotically in time and almost everywhere with respect to an invariant measure (see e.g. $[\mathrm{O}],[\mathrm{P}],[\mathrm{R}],[\mathrm{LY}]$ ). It is a very general theory with the potential for far-reaching consequences.

Yet using this abstract theory in concrete situations has proved to be difficult, in part because the assumptions on which this theory is based, such as the positivity of Lyapunov exponents or existence of SRB measures, are inherently difficult to verify. At the very least, the subject is in need of examples. To improve its utility, better techniques are needed to bridge the gap between theory and application. The project of which the present paper is a crucial component (see B and C below) is an attempt to address these needs.

We exhibit in this paper large numbers of nonuniformly hyperbolic attractors with controlled dynamics near every 1D map satisfying the well-known Misiurewicz condition. A detailed account of the mechanisms responsible for the hyperbolicity is given in Part II.

With a view toward applications, we sought to formulate conditions for the existence of SRB measures that are verifiable in concrete situations. These conditions cannot be placed on the map directly, for in the absence of invariant cones, to determine whether a map has this measure requires knowing it to infinite precision. We resolved this dilemma for the systems in question by identifying checkable conditions on 1-parameter families. These conditions guarantee the existence of SRB measures with positive probability, i.e. for positive measure sets of parameters. See Section 1.
B. In relation to one dimensional maps. In terms of techniques, this paper borrows heavily from the theory of iterated 1D maps, where much progress has been made in the last 25 years. Among the works that have influenced us the most are $[\mathrm{M}],[\mathrm{J}],[\mathrm{CE}],[\mathrm{BC} 1]$ and $[\mathrm{TTY}]$. The first breakthrough from 1D
to a family of strongly dissipative 2D maps is due to Benedicks and Carleson, whose paper $[\mathrm{BC} 2]$ is a tour de force analysis of the Hénon maps near the parameters $a=2, b=0$. Much of the local phase-space analysis in this paper is a generalization of their techniques, which in turn have their origins in 1D. Based on [BC2], SRB measures were constructed for the first time in [BY] for a (genuinely) nonuniformly hyperbolic attractor. The results in [BC2] were generalized in $[\mathrm{MV}]$ to small perturbations of the same maps. These papers form the core material referred to in the second box below.


All of the results in the second box depend on the formula of the Hénon maps. In going from the second to the third box, our aim is to take this mathematics to a more general setting, so that it can be leveraged in the analysis of attractors with similar characteristics (see below). Our treatment of the subject is necessarily more conceptual as we replace the equation of the Hénon maps by geometric conditions. A 2D version of these results was published in [WY1].

We believe the proper context for this set of ideas is $m$ dimensions, $m \geq 2$, where we retain the rank one character of the attractor but allow the number of stable directions to be arbitrary. We explain an important difference between this general setup and 2D: For strongly contractive maps $T$ with $T(X) \subset X$, by tracking $T^{n}(\partial X)$ for $n=1,2,3, \cdots$, one can obtain a great deal of information on the attractor $\cap_{n \geq 0} T^{n}(X)$. This is because the area or volume of $T^{n}(X)$ decreases to zero very quickly. Since the boundary of a 2 D domain consists of 1D curves, the study of planar attractors can be reduced to tracking a finite number of curves in the plane. This is what has been done in 2D, implicitly or explicitly. In $D>2$, both the analysis and the geometry become more complex; one is forced to deal directly with higher dimensional objects. The proofs in this paper work in all dimensions including $D=2$.
C. Further results and applications. We have a fairly complete dynamical description for the maps $T \in \mathcal{G}$ (see the beginning of this introduction), but in order to keep the length of the present paper reasonable, we have opted to publish these results separately. They include (1) a bound on the number of ergodic SRB measures, (2) conditions that imply ergodicity and mixing for SRB measures, (3) almost-everywhere behavior in the basin, (4) statistical properties of SRB measures such as correlation decay and CLT, and (5) coding of orbits on the attractor, growth of periodic points, etc. A 2D version of these results is published in [WY1]. Additional work is needed in higher dimensions due to the increased complexity in geometry.

We turn now to applications. First, by leveraging results of the type in this paper, we were able to recover and extend - by simply checking the conditions in Section 1 - previously known results on the Hénon maps and homoclinic bifurcations ([BC2],[MV],[V]).

The following new applications were found more recently: Forced oscillators are natural candidates for rank one attractors. We proved in [WY2],[WY3] that any limit cycle, when periodically kicked in a suitable way, can be turned into a strange attractor of the type studied here. It is also quite natural to associate systems with a single unstable direction with scenarios following a loss of stability. This is what led us to the result on the emergence of strange attractors from Hopf bifurcations in periodically kicked systems [WY3]. Finally, we mention some work in preparation in which we, together with K. Lu, bring some of the ideas discussed here including strange attractors and SRB measures to the arena of PDEs.

About this paper. This paper is self-contained, in part because relevant results from previously published works are inadequate for our purposes. The table of contents is self-explanatory. We have put all of the computational proofs in the Appendices so as not to obstruct the flow of ideas, and recommend that the reader omit some or all of the Appendices on first pass. This suggestion applies especially to Section 3, which, being a toolkit, is likely to acquire context only through subsequent sections. That having been said, we must emphasize also that the Appendices are an integral part of this paper; our proofs would not be complete without them.

## 1. Statement of results

We begin by introducing $\mathcal{M}$, the class of one-dimensional maps of which all maps studied in this paper are perturbations. In the definition below, $I$ denotes either a closed interval or a circle, $f: I \rightarrow I$ is a $C^{2}$ map, $C=\left\{f^{\prime}=0\right\}$ is the critical set of $f$, and $C_{\delta}$ is the $\delta$-neighborhood of $C$ in $I$. In the case of an interval, we assume $f(I) \subset \operatorname{int}(I)$, the interior of $I$. For $x \in I$, we let $d(x, C)=\min _{\hat{x} \in C}|x-\hat{x}|$.

Definition 1.1. We say $f \in \mathcal{M}$ if the following hold for some $\delta_{0}>0$ :
(a) Critical orbits: for all $\hat{x} \in C, d\left(f^{n}(\hat{x}), C\right)>2 \delta_{0}$ for all $n>0$.
(b) Outside of $C_{\delta_{0}}$ : there exist $\lambda_{0}>0, M_{0} \in \mathbb{Z}^{+}$and $0<c_{0} \leq 1$ such that
(i) for all $n \geq M_{0}$, if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta_{0}}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq e^{\lambda_{0} n}$;
(ii) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta_{0}}$ and $f^{n}(x) \in C_{\delta_{0}}$, any $n$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq c_{0} e^{\lambda_{0} n}$.
(c) Inside $C_{\delta_{0}}$ : there exists $K_{0}>1$ such that for all $x \in C_{\delta_{0}}$,
(i) $f^{\prime \prime}(x) \neq 0$;
(ii) $\exists p=p(x), K_{0}^{-1} \log \frac{1}{d(x, C)}<p(x)<K_{0} \log \frac{1}{d(x, C)}$, such that $f^{j}(x) \notin$ $C_{\delta_{0}} \forall j<p$ and $\left|\left(f^{p}\right)^{\prime}(x)\right| \geq c_{0}^{-1} e^{\frac{1}{3} \lambda_{0} p}$.

This definition may appear a little technical, but the properties are exactly those needed for our purposes. The class $\mathcal{M}$ is a slight generalization of the maps studied by Misiurewicz in $[\mathrm{M}]$.

Assume $f \in \mathcal{M}$ is a member of a one-parameter family $\left\{f_{a}\right\}$ with $f=f_{a^{*}}$. Certain orbits of $f$ have natural continuations to $a$ near $a^{*}$ : For $\hat{x} \in C, \hat{x}(a)$ denotes the corresponding critical point of $f_{a}$. For $q \in I$ with $\inf _{n \geq 0} d\left(f^{n}(q), C\right)$ $>0, q(a)$ is the unique point near $q$ whose symbolic itinerary under $f_{a}$ is identical to that of $q$ under $f$. For more detail, see Sections 2.1 and 2.4.

Let $X=I \times D_{m-1}$ where $I$ is as above and $D_{m-1}$ is the closed unit disk in $\mathbb{R}^{m-1}, m \geq 2$. Points in $X$ are denoted by $(x, y)$ where $x \in I$ and $y=\left(y^{1}, \cdots, y^{m-1}\right) \in D_{m-1}$. To $F: X \rightarrow I$ we associate two maps, $F^{\#}: X \rightarrow X$ where $F^{\#}(x, y)=(F(x, y), 0)$ and $f: I \rightarrow I$ where $f(x)=F(x, 0)$. Let $\|\cdot\|_{C^{r}}$ denote the $C^{r}$ norm of a map. A one-parameter family $F_{a}: X \rightarrow I$ (or $\left.T_{a}: X \rightarrow X\right)$ is said to be $C^{3}$ if the mapping $(x, y ; a) \mapsto F_{a}(x, y)$ (respectively $\left.(x, y ; a) \mapsto T_{a}(x, y)\right)$ is $C^{3}$.

Standing Hypotheses. We consider embeddings $T_{a}: X \rightarrow X, a \in\left[a_{0}, a_{1}\right]$, where $\left\|T_{a}-F_{a}^{\#}\right\|_{C^{3}}$ is small for some $F_{a}$ satisfying the following conditions:
(a) There exists $a^{*} \in\left[a_{0}, a_{1}\right]$ such that $f_{a^{*}} \in \mathcal{M}$.
(b) For every $\hat{x} \in C=C\left(f_{a^{*}}\right)$ and $q=f_{a^{*}}(\hat{x})$,

$$
\begin{equation*}
\frac{d}{d a} f_{a}(\hat{x}(a)) \neq \frac{d}{d a} q(a)^{1} \quad \text { at } \quad a=a^{*} \tag{1}
\end{equation*}
$$

(c) For every $\hat{x} \in C$, there exists $j \leq m-1$ such that

$$
\begin{equation*}
\frac{\partial F\left(\hat{x}, 0 ; a^{*}\right)}{\partial y^{j}} \neq 0 \tag{2}
\end{equation*}
$$

A $T$-invariant Borel probability measure $\nu$ is called an SRB measure if (i) $T$ has a positive Lyapunov exponent $\nu$-a.e.; (ii) the conditional measures of $\nu$ on unstable manifolds are absolutely continuous with respect to the Riemannian measures on these leaves.

Theorem. In addition to the Standing Hypotheses above, assume that $\left\|T_{a}-F_{a}^{\#}\right\|_{C^{3}}$ is sufficiently small depending on $\left\{F_{a}\right\}$. Then there is a positive measure set $\Delta \subset\left[a_{0}, a_{1}\right]$ such that for all $a \in \Delta, T=T_{a}$ admits an SRB measure.

[^1]Notation. For $z_{0} \in X$, let $z_{n}=T^{n}\left(z_{0}\right)$, and let $X_{z_{0}}$ be the tangent space at $z_{0}$. For $v_{0} \in X_{z_{0}}$, let $v_{n}=D T_{z_{0}}^{n}\left(v_{0}\right)$. We identify $X_{z}$ freely with $\mathbb{R}^{m}$, and work in $\mathbb{R}^{m}$ from time to time in local arguments. Distances between points in $X$ are denoted by $|\cdot-\cdot|$, and norms on $X_{z}$ by $|\cdot|$. The notation $\|\cdot\|$ is reserved for norms of maps (e.g. $\left\|T_{a}\right\|_{C^{3}}$ as above, $\|D T\|:=\sup _{z \in X}\left\|D T_{z}\right\|$ ).

For definiteness, our proofs are given for the case $I=S^{1}$. Small modifications are needed to deal with the case where $I$ is an interval. This is discussed in Section 3.9 at the end of Part I.

## PART I. PREPARATION

## 2. Relevant results from one dimension

The attractors studied in this paper have both an $m$-dimensional and a 1-dimensional character, the first having to do with how they are embedded in $m$-dimensional space, the second due the fact that the maps in question are perturbations of 1D maps. In this section, we present some results on 1D maps that are relevant for subsequent analysis. When specialized to the family $f_{a}(x)=1-a x^{2}$ with $a^{*}=2$, the material in Sections 2.2 and 2.3 is essentially contained in [BC2]; some of the ideas go back to [CE]. Part of Section 2.4 is a slight generalization of part of [TTY], which also contains an extension of [ BC 1$]$ and the 1 D part of $[\mathrm{BC} 2]$ to unimodal maps.

### 2.1. More on maps in $\mathcal{M}$

The maps in $\mathcal{M}$ are among the simplest maps with nonuniform expansion. The phase space is divided into two regions: $C_{\delta_{0}}$ and $I \backslash C_{\delta_{0}}$. Condition (b) in Definition 1.1 says that on $I \backslash C_{\delta_{0}}, f$ is essentially (uniformly) expanding. (c) says that every orbit from $C_{\delta_{0}}$, though contracted initially, is not allowed to return to $C_{\delta_{0}}$ until it has regained some amount of exponential growth.

An important feature of $f \in \mathcal{M}$ is that its Lyapunov exponents outside of $C_{\delta}$ are bounded below by a strictly positive number independent of $\delta$. Let $\delta_{0}$, $\lambda_{0}, M_{0}$ and $c_{0}$ be as in Definition 1.1.

Lemma 2.1. For $f \in \mathcal{M}, \exists c_{0}^{\prime}>0$ such that the following hold for all $\delta<\delta_{0}$ :
(a) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq c_{0}^{\prime} \delta e^{\frac{1}{3} \lambda_{0} n}$;
(b) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta}$ and $f^{n}(x) \in C_{\delta_{0}}$, any $n$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq$ $c_{0} e^{\frac{1}{3} \lambda_{0} n}$.

Obviously, as we perturb $f$, its critical orbits will not remain bounded away from $C$. The expanding properties of $f$ outside of $C_{\delta}$, however, will
persist in the manner to be described. Note the order in which $\varepsilon$ and $\delta$ are chosen in the next lemma.

Lemma 2.2. Let $f$ and $c_{0}^{\prime}$ be as in Lemma 2.1, and fix an arbitrary $\delta<\delta_{0}$. Then there exists $\varepsilon=\varepsilon(\delta)>0$ such that the following hold for all $g$ with $\|g-f\|_{C^{2}}<\varepsilon$ :
(a) if $x, g(x), \cdots, g^{n-1}(x) \notin C_{\delta}$, then $\left|\left(g^{n}\right)^{\prime}(x)\right| \geq \frac{1}{2} c_{0}^{\prime} \delta e^{\frac{1}{4} \lambda_{0} n}$;
(b) if $x, g(x), \cdots, g^{n-1}(x) \notin C_{\delta}$ and $g^{n}(x) \in C_{\delta_{0}}$, any $n$, then $\left|\left(g^{n}\right)^{\prime}(x)\right| \geq$ $\frac{1}{2} c_{0} e^{\frac{1}{4} \lambda_{0} n}$.

Lemmas 2.1 and 2.2 are proved in Appendix A.1.

### 2.2. A larger class of 1 D maps with good properties

We introduce next a class of maps more flexible than those in $\mathcal{M}$. These maps are located in small neighborhoods of $f_{0} \in \mathcal{M}$. They will be our model of controlled dynamical behavior in higher dimensions.

For the rest of this subsection, we fix $f_{0} \in \mathcal{M}$, and let $\delta_{0}, \lambda_{0}, M_{0}$ and $c_{0}$ be as in Definition 1.1. We fix also $\lambda<\frac{1}{5} \lambda_{0}$ and $\alpha \ll \min \{\lambda, 1\}$. The letter $K \geq 1$ is used as a generic constant that is allowed to depend only on $f_{0}$ and $\lambda$. By "generic", we mean $K$ may take on different values in different situations.

Let $\delta>0$, and consider $f$ with $\left\|f-f_{0}\right\|_{C^{2}} \ll \delta$. Let $C$ be the critical set of $f$. We assume that for all $\hat{x} \in C$, the following hold for all $n>0$ :
(G1) $d\left(f^{n}(\hat{x}), C\right)>\min \left\{\delta, e^{-\alpha n}\right\} ;{ }^{2}$
(G2) $\left|\left(f^{n}\right)^{\prime}(f(\hat{x}))\right| \geq \hat{c}_{1} e^{\lambda n}$ for some $\hat{c}_{1}>0$.
Proposition 2.1. Let $\delta>0$ be sufficiently small depending on $f_{0}$. Then there exists $\varepsilon=\varepsilon\left(f_{0}, \lambda, \alpha, \delta\right)>0$ such that if $\left\|f-f_{0}\right\|_{C^{2}}<\varepsilon$ and $f$ satisfies (G1) and (G2), then it has properties (P1)-(P3) below.
(P1) Outside of $C_{\delta}$. There exists $c_{1}>0$ such that the following hold:
(i) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta}$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq c_{1} \delta e^{\frac{1}{4} \lambda_{0} n}$;
(ii) if $x, f(x), \cdots, f^{n-1}(x) \notin C_{\delta}$ and $f^{n}(x) \in C_{\delta_{0}}$, any $n$, then $\left|\left(f^{n}\right)^{\prime}(x)\right| \geq$ $c_{1} e^{\frac{1}{4} \lambda_{0} n}$.

For $\hat{x} \in C$, let $C_{\delta}(\hat{x})=(\hat{x}-\delta, \hat{x}+\delta)$. We now introduce a partition $\mathcal{P}$ on $I$ : For each $\hat{x} \in C,\left.\mathcal{P}\right|_{C_{\delta}(\hat{x})}=\left\{I_{\mu j}^{\hat{x}}\right\}$ where $I_{\mu j}^{\hat{x}}$ are defined as follows: For

[^2]$\mu \geq \log \frac{1}{\delta}$ (which we may assume is an integer), let $I_{\mu}^{\hat{x}}=\left(\hat{x}+e^{-(\mu+1)}, \hat{x}+e^{-\mu}\right) ;$ for $\mu \leq \log \delta$, let $I_{\mu}^{\hat{x}}$ be the reflection of $I_{-\mu}^{\hat{x}}$ about $\hat{x}$. Each $I_{\mu}^{\hat{x}}$ is further subdivided into $\frac{1}{\mu^{2}}$ subintervals of equal length called $I_{\mu j}^{\hat{x}}$. We usually omit the superscript $\hat{x}$ in the notation above, with the understanding that $\hat{x}$ may vary from statement to statement. For example, " $x \in I_{\mu j}$ and $f^{n}(x) \in I_{\mu^{\prime} j^{\prime}}$ " may refer to $x \in I_{\mu j}^{\hat{x}}$ and $f^{n}(x) \in I_{\mu^{\prime} j^{\prime}}^{\hat{x}^{\prime}}$ for $\hat{x} \neq \hat{x}^{\prime}$. The rest of $I$, i.e. $I \backslash C_{\delta}$, is partitioned into intervals of length $\approx \delta$.
(P2) Partial derivative recovery for $x \in C_{\delta}(\hat{x}) \backslash\{\hat{x}\}$. For $x \in C_{\delta}(\hat{x}) \backslash\{\hat{x}\}$, let $p(x)$, the bound period of $x$, be the largest integer such that $\left|f^{i}(x)-f^{i}(\hat{x})\right| \leq$ $e^{-2 \alpha i} \forall j<p(x)$. Then
(i) $K^{-1} \log \frac{1}{|x-\hat{x}|} \leq p(x) \leq K \log \frac{1}{|x-\hat{x}|}$.
(ii) $\left|\left(f^{p(x)}\right)^{\prime}(x)\right|>e^{\frac{\lambda}{3} p(x)}$.
(iii) If $\omega=I_{\mu j}$, then $\left|f^{p(x)}\left(I_{\mu j}\right)\right|>e^{-K \alpha|\mu|}$ for all $x \in \omega$.

The idea behind (P1) and (P2) is as follows: By choosing $\varepsilon$ sufficiently small depending on $\delta$, we are assured that there is a neighborhood $\mathcal{N}$ of $f_{0}$ such that all $f \in \mathcal{N}$ are essentially expanding outside of $C_{\delta}$. Non-expanding behavior must, therefore, originate from inside $C_{\delta}$. We hope to control that by imposing conditions (G1) and (G2) on $C$, and to pass these properties on to other orbits starting from $C_{\delta}$ via ( P 2 ).
(P2) leads to the following view of an orbit:
Returns to $C_{\delta}$ and ensuing bound periods. For $x \in I$ such that $f^{i}(x) \notin C$ for all $i \geq 0$, we define (free) return times $\left\{t_{k}\right\}$ and bound periods $\left\{p_{k}\right\}$ with

$$
t_{1}<t_{1}+p_{1} \leq t_{2}<t_{2}+p_{2} \leq \cdots
$$

as follows: $t_{1}$ is the smallest $j \geq 0$ such that $f^{j}(x) \in C_{\delta}$. For $k \geq 1, p_{k}$ is the bound period of $f^{t_{k}}(x)$, and $t_{k+1}$ is the smallest $j \geq t_{k}+p_{k}$ such that $f^{j}(x) \in C_{\delta}$. Note that an orbit may return to $C_{\delta}$ during its bound periods, i.e. $t_{i}$ are not the only return times to $C_{\delta}$.

The following notation is used: If $P \in \mathcal{P}$, then $P^{+}$denotes the union of $P$ and the two elements of $\mathcal{P}$ adjacent to it. For an interval $Q \subset I$ and $P \in \mathcal{P}$, we say $Q \approx P$ if $P \subset Q \subset P^{+}$. For practical purposes, $P^{+}$containing boundary points of $C_{\delta}$ can be treated as "inside $C_{\delta}$ " or "outside $C_{\delta}$ ". ${ }^{3}$ For an interval $Q \subset I_{\mu j}^{+}$, we define the bound period of $Q$ to be $p(Q)=\min _{x \in Q}\{p(x)\}$.
(P3) is about comparisons of derivatives for nearby orbits. For $x, y \in I$, let $[x, y]$ denote the segment connecting $x$ and $y$. We say $x$ and $y$ have the same

[^3]itinerary (with respect to $\mathcal{P}$ ) through time $n-1$ if there exist $t_{1}<t_{1}+p_{1} \leq$ $t_{2}<t_{2}+p_{2} \leq \cdots \leq n$ such that for every $k, f^{t_{k}}([x, y]) \subset P^{+}$for some $P \subset C_{\delta}$, $p_{k}=p\left(f^{t_{k}}([x, y])\right)$, and for all $i \in[0, n) \backslash \cup_{k}\left[t_{k}, t_{k}+p_{k}\right), f^{i}([x, y]) \subset P^{+}$for some $P \cap C_{\delta}=\emptyset$.
(P3) Distortion estimate. There exists $K$ (independent of $\delta, x, y$ or $n$ ) such that if $x$ and $y$ have the same itinerary through time $n-1$, then
$$
\left|\frac{\left(f^{n}\right)^{\prime}(x)}{\left(f^{n}\right)^{\prime}(y)}\right| \leq K
$$

We remark that the partition of $I_{\mu}$ into $I_{\mu j}$-intervals is solely for purposes of this estimate. A proof of Proposition 2.1 is given in Appendix A.1.

### 2.3. Statistical properties of maps satisfying (P1)-(P3)

We assume in this subsection that $f$ satisfies the assumptions of Proposition 2.1, so that in particular (P1)-(P3) hold. Let $\omega \subset I$ be an interval. For reasons to become clear later, we write $\gamma_{i}=f^{i}$, i.e. we consider $\gamma_{i}: \omega \rightarrow I$, $i=0,1,2, \cdots$.

Lemma 2.3. For $\omega \approx I_{\mu_{0} j_{0}}$, let $n$ be the largest $j$ such that all $s \in \omega$ have the same itinerary up to time $j$. Then $n \leq K\left|\mu_{0}\right|$.

We call $n+1$ the extended bound period for $\omega$. The next result, the proof of which we leave as an exercise, is used only in Lemma 8.2.

Lemma 2.4. For $\omega \approx I_{\mu_{0} j_{0}}$, there exists $n \leq K\left|\mu_{0}\right|$ such that $\gamma_{n}(\omega) \supset$ $C_{\delta}(\hat{x})$ for some $\hat{x} \in C$.

The results in the rest of this subsection require that we track the evolution of $\gamma_{i}$ to infinite time. To maintain control of distortion, it is necessary to divide $\omega$ into shorter intervals. The increasing sequence of partitions $\mathcal{Q}_{0}<\mathcal{Q}_{1}<\mathcal{Q}_{2}<$ $\cdots$ defined below is referred to as a canonical subdivision by itinerary for the interval $\omega: \mathcal{Q}_{0}$ is equal to $\left.\mathcal{P}\right|_{\omega}$ except that the end intervals are attached to their neighbors if they are strictly shorter than the elements of $\mathcal{P}$ containing them. We assume inductively that all $\hat{\omega} \in \mathcal{Q}_{i}$ are intervals and all points in $\hat{\omega}$ have the same itinerary through time $i$. To go from $\mathcal{Q}_{i}$ to $\mathcal{Q}_{i+1}$, we consider one $\hat{\omega} \in \mathcal{Q}_{i}$ at a time.

- If $\gamma_{i+1}(\hat{\omega})$ is in a bound period, then $\hat{\omega}$ is automatically put into $\mathcal{Q}_{i+1}$. (Observe that if $\gamma_{i+1}(\hat{\omega}) \cap C_{\delta} \neq \emptyset$, then $\gamma_{i+1}(\hat{\omega}) \subset I_{\mu^{\prime} j^{\prime}}^{+}$for some $\mu^{\prime}, j^{\prime}$; i.e., no cutting is needed during bound periods. This is an easy exercise.)
- If $\gamma_{i+1}(\hat{\omega})$ is not in a bound period, but all points in $\hat{\omega}$ have the same itinerary through time $i+1$, we again put $\hat{\omega} \in \mathcal{Q}_{i+1}$.
- If neither of the last two cases hold, then we partition $\hat{\omega}$ into segments $\left\{\hat{\omega}^{\prime}\right\}$ that have the same itineraries through time $i+1$ and with $\gamma_{i+1}\left(\hat{\omega}^{\prime}\right) \approx$ $P$ for some $P \in \mathcal{P}$. (If, for example, a segment appears that is strictly shorter than the $I_{\mu j}$ containing it, then it is attached to a neighboring segment.) The resulting partition is $\left.\mathcal{Q}_{i+1}\right|_{\hat{\omega}}$.

For $s \in \omega$, let $\mathcal{Q}_{i}(s)$ be the element of $\mathcal{Q}_{i}$ to which $s$ belongs. We consider the stopping time $S$ on $\omega$ defined as follows: For $s \in \omega$, let $S(s)$ be the smallest $i$ such that $\gamma_{i}\left(\mathcal{Q}_{i-1}(s)\right)$ is not in a bound period and has length $>\delta$.

Lemma 2.5. Assume $\delta$ is sufficiently small, and let $\omega \approx I_{\mu_{0} j_{0}}$. Then

$$
|\{s \in \omega: S(s)>n\}|<e^{-\frac{1}{2} K^{-1} n}|\omega| \text { for all } n>K\left|\mu_{0}\right| .
$$

Here $K$ is the constant in the statement of Lemma 2.2.
Corollary 2.1. There exists $\hat{K}>0$ such that for any $\omega \subset I$ with $\delta<$ $|\omega|<3 \delta$,

$$
|\{s \in \omega: S(s)>n\}|<e^{-\hat{K}^{-1} n}|\omega| \quad \text { for } n>\hat{K} \log \delta^{-1} .
$$

For $\hat{\delta}<\delta, s \in \omega$ and $n \geq 0$, let $B_{n}(s)$ be the number of $i \leq n$ such that $\gamma_{i}(s)$ is in a bound period initiated from a visit to $C_{\hat{\delta}}$.

Proposition 2.2. Given any $\sigma>0$, there exists $\varepsilon_{1}>0$ such that for all $\hat{\delta}>0$ sufficiently small, the following holds for all $\omega \approx I_{\mu_{0} j_{0}}$ :

$$
\left|\left\{s \in \omega: B_{n}>\sigma n\right\}\right|<e^{-\varepsilon_{1} n}|\omega| \quad \text { for all } n \geq \sigma^{-1} K \mu_{0} .
$$

Proofs of all the results in this subsection are given in Appendix A. 2 except that of Lemma 2.4, which is left to the reader as an exercise.

Remark. The main use of Proposition 2.2 in this paper is in parameter estimates. When used in that context, it will be necessary for us to stop considering certain elements $\omega^{\prime}$ of $\mathcal{Q}_{i}$ corresponding to deletions. Without going further into parameter considerations, we introduce the following notation. Let $*$ be the "garbage symbol". At step $i$, we may, in principle, choose to set $\gamma_{i}=*$ on any collection of elements of $\mathcal{Q}_{i}$. Once we set $\left.\gamma_{i}\right|_{\omega^{\prime}}=*$, it follows automatically that $\left.\gamma_{j}\right|_{\omega^{\prime}}=*$ for all $j \geq i$, i.e. we do not iterate $\omega^{\prime}$ forward from time $i$ on. We leave it as an (easy) exercise to verify that Proposition 2.2 remains valid in this slightly more general setting if we count only those $i$ for which $\gamma_{i}(s) \neq *$ in the definition of $B_{n}(s)$.

### 2.4. Parameter transversality

We begin with a description of the structure of $f \in \mathcal{M}$ in terms of its symbolic dynamics. Let $\mathcal{J}=\left\{J_{1}, \cdots, J_{q}\right\}$ be the components of $I \backslash C$. For
$x \in I$ such that $f^{i}(x) \notin C$ for all $i \geq 0$, let $\phi(x)=\left(\iota_{i}\right)_{i=0,1, \ldots}$ be given by $\iota_{i}=k$ if $f^{i}(x) \in J_{k}$.

Lemma 2.6. For $f \in \mathcal{M}$, there exists an increasing sequence of compact sets $\Lambda^{(n)}$ with $\cup_{n} \Lambda^{(n)}$ dense in I such that the following hold:
(a) $\Lambda^{(n)} \cap C=\emptyset, f\left(\Lambda^{(n)}\right) \subset \Lambda^{(n)}$, and $\left.f\right|_{\Lambda^{(n)}}$ is conjugate to a shift of finite type;
(b) if $\inf _{i \geq 0} d\left(f^{i}(x), C\right)>0$, then $x \in \Lambda^{(n)}$ for some $n$.

Our next result, which is a corollary of Lemmas 2.2 and 2.6, guarantees that continuations of the type in Standing Hypothesis (b) are well defined.

Corollary 2.2. Let $f \in \mathcal{M}$, and let $q \in f(I)$ be such that $\delta_{1}:=$ $\inf _{n \geq 0} d\left(f^{n}(q), C\right)>0$. Then for all $g$ with $\|g-f\|_{C^{2}}<\varepsilon$ where $\varepsilon=\varepsilon\left(\delta_{1}\right)$ is as in Lemma 2.2, there is a unique point $q_{g} \in I$ with $\phi_{g}\left(q_{g}\right)=\phi_{f}(q)$.

Let $\left\{f_{a}\right\}$ be as in Section 1, with $f_{a^{*}} \in \mathcal{M}$. We fix $\hat{x} \in C\left(f_{a^{*}}\right)$, and let $q=f_{a^{*}}(\hat{x})$. Let $\omega$ be an interval containing $a^{*}$ on which $\hat{x}(a)$ and $q(a)$ (as given by Corollary 2.2) are well defined. We write $\hat{x}_{k}(a)=f_{a}^{k}(\hat{x}(a))$.

Proposition 2.3. (i) $a \mapsto q(a)$ is differentiable;
(ii) as $k \rightarrow \infty$,

$$
Q_{k}\left(a^{*}\right):=\frac{\frac{d \hat{x}_{k}}{d a}\left(a^{*}\right)}{\left(f_{a^{*}}^{k-1}\right)^{\prime}\left(\hat{x}_{1}\left(a^{*}\right)\right)} \rightarrow \frac{d \hat{x}_{1}}{d a}\left(a^{*}\right)-\frac{d q}{d a}\left(a^{*}\right)=\sum_{i=0}^{\infty} \frac{\left.\partial_{a} f_{a}\left(\hat{x}_{i}\left(a^{*}\right)\right)\right|_{a=a^{*}}}{\left(f_{a^{*}}^{*}\right)^{\prime}\left(\hat{x}_{1}\left(a^{*}\right)\right)} .
$$

A proof of this proposition, which is a slight adaptation of a result in [TTY], is given in Appendix A.3. Hypothesis (b) states that the expression on the right is nonzero. This condition, which can be viewed as a transversality condition for one-parameter families in the space of $C^{2}$ maps, is open and dense among the set of all 1-parameter families $f_{a}$ passing through a given $f \in \mathcal{M}$. The proof in [TTY] is easily adapted to the present setting.

## 3. Tools for analyzing rank one maps

This section is a toolkit for the analysis of maps $T: X \rightarrow X$ that are small perturbation of maps from $X$ to $I \times\{0\}$. More conditions are assumed as needed, but detailed structures of the maps in question are largely unimportant. The purpose of this section is to develop basic techniques for use in the rest of the paper.

Notation. The following rules on the use of constants are observed throughout:

- Two constants, $K_{0} \geq 1$ and $0<b \ll 1$, are used to bound the sizes of the objects being studied; they appear in assumptions.
- $K$ is used as a generic constant; it appears in statements of results. In Sections 3.1-3.4, $K$ depends only on $K_{0}$ and $m$, the dimension of $X$; from Section 3.5 on, it depends on an additional object to be specified.
- $b$ is assumed to be as small as need be; it is shrunk a finite number of times as we go along. Under no conditions is $K$ allowed to depend on $b$.

For small angles, $\theta$ is often confused with $|\sin \theta|$.

### 3.1. Stability of most contracted directions

Most contracted directions on planes. Consider first $M \in \mathrm{~L}(2, \mathbb{R})$ and assume $M \neq c O$ where $O$ is orthogonal and $c \in \mathbb{R}$. Then there is a unit vector $e$, uniquely defined up to sign, that represents the most contracted direction of $M$, i.e. $|M e| \leq|M u|$ for all unit vectors $u$. From standard linear algebra, we know $e^{\perp}$ is the most expanded direction, meaning $\left|M e^{\perp}\right| \geq|M u|$ for all unit vectors $u$, and $M e \perp M e^{\perp}$. The numbers $|M e|$ and $\left|M e^{\perp}\right|$ are the singular values of $M$.

Next let $M \in \mathrm{~L}(m, \mathbb{R})$ for $m \geq 2$, and let $S \subset \mathbb{R}^{m}$ be a 2 D linear subspace. Then the ideas in the last paragraph clearly apply to $\left.M\right|_{S}$, and we say $e=e(S)$ is a most contracted direction of $M$ restricted to $S$ if $|M e| \geq|M u|$ for all unit vectors $u \in S$. We let $f$ denote one of the two unit vectors in $S$ orthogonal to $e$, i.e. $f$ represents the most expanded direction in $S$, and $|M f|=\left\|\left.M\right|_{S}\right\|$, the norm of $M$ restricted to $S$.

Two notions of stability for most contracted directions. For $M_{1}, M_{2}, \cdots \in$ $\mathrm{L}(m, \mathbb{R})$, we let $M^{(i)}$ denote the composition $M_{i} \cdots M_{2} M_{1}$.
(1) Let $S \subset \mathbb{R}^{m}$ be as above, and let $e_{i}(S)$ be the most contracted direction of $\left.M^{(i)}\right|_{S}$ assuming it is well defined. It is known that if $\left.M^{(i)}\right|_{S}, i=1,2, \cdots$, has two distinct Lyapunov exponents as $i \rightarrow \infty$, then $e_{i}(S)$ converges to some $e_{\infty}(S)$ as $i \rightarrow \infty$. We are interested in the speed of this convergence.
(2) For parametrized families of linear maps $M_{i}(s)$ and plane fields $S(s)$ where $s=\left(s_{1}, \cdots, s_{q}\right)$ is a $q$-tuple of numbers, control of $\partial^{k} e_{i}$ and $\partial^{k} M^{(n)} e_{i}$ represents another form of stability for $e_{i}$. Here $\partial^{k}$ denotes any one of the $k$ th partial derivatives in $s$.

Main results. The ideas above are used to study the relation between pairs of vectors under the action of $D T^{n}$. To accommodate the many situations in which this analysis will be applied, we formulate our next lemma in terms of abstract linear maps. For motivation, the reader should think of $M_{i}$ as $D T_{z_{i-1}}$ where $z_{0} \in X$ and $T: X \rightarrow X$ is as in Section 1.1. For (H2), consider $z_{0}(s) \in X, S(s) \subset X_{z_{0}(s)}$, and $M_{i}(s)=D T_{z_{i-1}(s)}$.
(H1) Let $M_{i}=\left(\hat{M}_{i}^{1}, \cdots, \hat{M}_{i}^{m}\right) \in \mathrm{L}(m, \mathbb{R})$; i.e., $\hat{M}_{i}^{j}: \mathbb{R}^{m} \rightarrow \mathbb{R}$. Then for all $i \geq 1$,
(i) $\left|\hat{M}_{i}^{1}\right|<K_{0}$;
(ii) $\left|\hat{M}_{i}^{j}\right|<b$ for $j=2, \cdots, m$.
(H2) Let $u(s)$ and $v(s) \in \mathbb{R}^{m}$ be linearly independent, and let $S(s)=$ $S(u(s), v(s))$ be the 2D subspace spanned by $u$ and $v$. Let $M_{i}(s) \in$ $\mathrm{L}(m, \mathbb{R})$. We assume the maps $s \mapsto u(s), v(s), M_{i}(s)$ are $C^{2}$ with
(i) $\|u\|_{C^{2}},\|v\|_{C^{2}}<K_{0}$;
(ii) $\left\|\hat{M}_{i}^{1}\right\|_{C^{2}}<K_{0}^{i}$;
(iii) $\left\|\hat{M}_{i}^{j}\right\|_{C^{2}}<K_{0}^{i} b$ for $j=2, \cdots, m$.

Lemma 3.1. (a) Let $M_{i}$ be as in (H1), let $S \subset \mathbb{R}^{m}$ be an arbitrary 2D subspace, and let $\kappa$ be such that $b^{\frac{1}{3}}<\kappa \leq 1$. If $\left\|\left.M^{(i)}\right|_{S}\right\|>K_{0}^{-1} \kappa^{i-1}$ for all $1 \leq i \leq n$, then

$$
\begin{aligned}
\left|e_{i+1}(S)-e_{i}(S)\right|<\left(K b \kappa^{-2}\right)^{i} & \text { for } i<n ; \\
\left|M^{(i)} e_{n}(S)\right|<\left(K b \kappa^{-2}\right)^{i} & \text { for } i \leq n .
\end{aligned}
$$

(b) Let $M_{i}(s)$ and $S(s)$ be as in (H2), and $b^{\frac{1}{5}} \leq \kappa \leq 1$. If for $1 \leq i \leq n$, $\left\|M^{(i)}{ }_{S}\right\|>K_{0}^{-1} \kappa^{i-1}$ for all $s$, then for $k=1,2$,

$$
\begin{aligned}
\left|\partial^{k} e_{1}(S)\right| & <K \\
\left|\partial^{k}\left(e_{i+1}(S)-e_{i}(S)\right)\right| & <\left(K b \kappa^{-(2+k)}\right)^{i} \\
\left|\partial^{k} M^{(i)} e_{n}(S)\right|<\left(K b \kappa^{-(2+k)}\right)^{i} & \text { for } i<n
\end{aligned}
$$

A proof of Lemma 3.1 is given in Appendix A.5, after some preliminary material in Appendix A.4.

Assumptions for the rest of Section 3. We consider $T: X \rightarrow X$ with the following properties: Let $T=\left(\hat{T}^{1}, \cdots, \hat{T}^{m}\right)$ be the coordinate maps of $T$. Then
(i) $\left\|\hat{T}^{1}\right\|_{C^{3}}<K_{0}$;
(ii) $\left\|\hat{T}^{j}\right\|_{C^{3}}<b$ for $j=2, \cdots, m$.

### 3.2. A perturbation lemma

The next lemma compares $w_{n}=D T_{z_{0}}^{n}\left(w_{0}\right)$ and $w_{n}^{\prime}=D T_{z_{0}^{\prime}}^{n}\left(w_{0}^{\prime}\right)$ where $z_{i}$ is near $z_{i}^{\prime}$ for $0 \leq i<n$ and $w_{0} \in X_{z_{0}}$ and $w_{0}^{\prime} \in X_{z_{0}^{\prime}}$ are unit vectors such that $w_{0} \approx w_{0}^{\prime}$.

Lemma 3.2. There exists $K_{1}$ depending on $K_{0}$ such that for $\kappa$ and $\eta$ satisfying $\kappa \leq 1$ and $b^{\frac{1}{2}}<\eta<K_{1}^{-1} \kappa^{8}$, the following hold: Let $\left(z_{0}, w_{0}\right)$ and $\left(z_{0}^{\prime}, w_{0}^{\prime}\right)$ be such that $\angle\left(w_{0}, w_{0}^{\prime}\right)<\eta^{\frac{1}{4}},\left|w_{i}\right|>K_{0}^{-1} \kappa^{i-1}$ and $\left|z_{i}-z_{i}^{\prime}\right|<\eta^{i+1}$ for $1 \leq i<n$. Then
(a) $\left|w_{n}^{\prime}\right|>\frac{1}{2} K_{0}^{-1} \kappa^{n-1}$;
(b) $\angle\left(w_{n}, w_{n}^{\prime}\right)<\eta^{\frac{n+1}{4}}$.

Lemma 3.2 is proved in Appendix A. 6.

### 3.3. Temporary stable curves and manifolds

One dimensional strong stable curves - temporary or infinite-time - can be obtained by integrating vector fields of most contracted directions. In the proposition below, a neighborhood of 0 in $X_{z_{0}}$ is identified with a neighborhood of $z_{0}$ in $X$, which in turn is identified with an open set of $\mathbb{R}^{m}$.

Proposition 3.1. Let $\kappa$ and $\eta$ be as in Lemma 3.2, and let $z_{0} \in X$ and $w_{0} \in X_{z_{0}}$ be such that $\left|w_{i}\right| \geq K_{0}^{-1} \kappa^{i-1}\left|w_{0}\right|$ for $i=1, \cdots, n$. Let $S$ be a 2D plane in $X$ containing $z_{0}$ and $z_{0}+w_{0}$. For any $n \geq 1$, we view $e_{n}(S)$ as a vector field on $S$, defined where it makes sense, and let $\gamma_{n}=\gamma_{n}\left(z_{0}, S\right)$ be the integral curve to $e_{n}(S)$ with $\gamma_{n}(0)=z_{0}$. Then
(a) $\gamma_{n}$ is defined on $[-\eta, \eta]$ or until it runs out of $X$;
(b) for all $z \in \gamma_{n},\left|T^{i} z_{0}-T^{i} z\right|<\left(\frac{K b}{\kappa^{2}}\right)^{i} \eta$ for all $i \leq n$.

Proposition 3.1 is proved in Appendix A.7.
We call $\gamma_{n}$ a temporary stable curve or stable curve of order $n$ through $z_{0}$. To obtain the full temporary stable manifold through $z_{0}$, we let $S$ vary over all 2D planes containing $z_{0}$ and $z_{0}+w_{0}$, obtaining

$$
W_{n}^{s}\left(z_{0}\right):=\cup_{S} \gamma_{n}\left(z_{0}, S\right)
$$

which we call a temporary stable manifold of order $n$ through $z_{0}$. Observe that $W_{n}^{s}\left(z_{0}\right)$ is a $C^{1}$-embedded disk of co-dimension one. (The fact that $W_{n}^{s}\left(z_{0}\right)$ is $C^{1}$ away from $z_{0}$ follows from Lemma 3.1; at $z_{0}$ it has continuous partial derivatives.)

### 3.4. A curvature estimate

Let $\gamma_{0}:\left[c_{1}, c_{2}\right] \rightarrow X$ be a $C^{2}$ curve, and let $\gamma_{i}(s)=T^{i}\left(\gamma_{0}(s)\right)$. We denote the curvature of $\gamma_{i}$ at $\gamma_{i}(s)$ by $k_{i}(s)$. Here $\gamma_{i}^{\prime}(s)$ is the tangent vector to $\gamma_{i}(s)$.

Lemma 3.3. Let $\kappa>b^{\frac{1}{3}}$, and let $\gamma_{0}$ be such that $k_{0}(s) \leq 1$ for all $s$. Then the following hold for every $n>0$ : If

$$
\left|D T_{\gamma_{j}(s)}^{n-j}\left(\gamma_{j}^{\prime}(s)\right)\right| \geq \kappa^{n-j}\left|\gamma_{j}^{\prime}(s)\right|
$$

for every $j<n$, then

$$
k_{n}(s) \leq \frac{K b}{\kappa^{3}}
$$

Lemma 3.3 is proved in Appendix A.8.
Additional assumptions for Sections 3.5-3.8. Let $\delta>0$ be a small number.
(1) The following is assumed about $\hat{T}^{1}: X \rightarrow I$ and $f:=\left.\hat{T}^{1}\right|_{I \times\{0\}}$. Let $C=\left\{f^{\prime}=0\right\}$. Then
(i) outside of $C_{\delta}, f$ satisfies (P1) in Section 2.2;
(ii) inside $C_{\delta},\left|f^{\prime \prime}\right|>K_{0}^{-1}$;
(iii) for all $\hat{x} \in C$, there exists $i$ such that $\left|\partial_{y^{i}} \hat{T}^{1}(x, 0)\right|>K_{0}^{-1}$ for all $x \in C_{\delta}(\hat{x})$.
(2) From here on we restrict $T$ to $R_{1}:=I \times\left\{|y| \leq(m-1)^{\frac{1}{2}} b\right\}$. Note that $T\left(R_{1}\right) \subset R_{1}$ (see assumption (ii) at the end of Section 3.1).

From here on in this section the generic constant $K$ depends on the map $\hat{T}^{1}$ as well as $K_{0}$ and $m$. We introduce the following notation used in the rest of the paper:

- The first critical region $\mathcal{C}^{(1)}$ is defined to be

$$
\mathcal{C}^{(1)}=\left\{(x, y) \in R_{1}:|x-\hat{x}|<\delta, \hat{x} \in C(f)\right\} .
$$

- $\mathbf{v} \in \mathbb{R}^{m}$ (identified with $X_{z}$, any $z$ ) is a fixed unit vector with zero $x$ component such that $\left|D \hat{T}_{(x, 0)}^{1} \mathbf{v}\right|>K_{0}^{-1}$ for all $x \in C_{\delta}$. The existence of $\mathbf{v}$ is guaranteed by assumption (1)(iii) above. (We may take it to be orthogonal to the kernel of $D \hat{T}_{(\hat{x}, 0)}^{1}$ for $\hat{x} \in C$ but that is not necessary.) In general, $\mathbf{v}$ will be thought of as a reference vector in the "vertical" direction.


### 3.5. Dynamics outside of $\mathcal{C}^{(1)}$

For $u \in \mathbb{R}^{m}$, let $\left(u_{x}, u_{y}\right)$ denote its $x$ and $y$ (or first and last $m-1$ ) components, and let $s(u)=\frac{\left|u_{y}\right|}{\left|u_{x}\right|}$. Curvature continues to be denoted by $k$.

Definition 3.1. Assuming $\left|f^{\prime}\right|>K_{0}^{-1} \delta$ outside of $\mathcal{C}^{(1)}$, we say $u \in \mathbb{R}^{m}$ is $b$-horizontal if $s(u)<\frac{3 K_{0}}{\delta} b$. A curve $\gamma$ in $R_{1}$ is called a $C^{2}(b)$-curve if $\gamma^{\prime}(s)$ is $b$-horizontal and $k(s)$ is $<\frac{K_{1} b}{\delta^{3}}$ for all $s$ where $K_{1}$ is as defined explicitly in the proof of Lemma 3.4. ${ }^{4}$

[^4]Lemma 3.4. (a) For $z \notin \mathcal{C}^{(1)}$, if $u \in X_{z}$ is b-horizontal, then so is $D T_{z}(u)$; in fact, $s\left(D T_{z}(u)\right)<\frac{3 K_{0}}{2 \delta} b$. Also, for $z \in \mathcal{C}^{(1)}, D T_{z}(\mathbf{v})$ is b-horizontal.
(b) If $\gamma$ is a $C^{2}(b)$-curve outside of $\mathcal{C}^{(1)}$, then $T(\gamma)$ is again a $C^{2}(b)$-curve.

Proof. The first assertion in (a) follows from the following invariant cones condition: Let $u$ be such that $\left|u_{x}\right|=1$ and $\left|u_{y}\right|<\frac{3 K_{0}}{\delta} b$. Then

$$
s\left(D T_{z}(u)\right)<\frac{b\left(1+\frac{3 K_{0}}{\delta} b\right)}{K_{0}^{-1} \delta-K_{0} \frac{3 K_{0}}{\delta} b}<\frac{3 K_{0}}{2 \delta} b
$$

provided $b$ is sufficiently small. For $z \in \mathcal{C}^{(1)}, s\left(D T_{z}(\mathbf{v})\right)<2 K_{0} b$. For (b) we apply Lemma 3.3 to one iteration of $T$ : Since $T$ is a small perturbation of $f$, we have $\left|D T\left(\gamma^{\prime}\right)\right|>\frac{1}{2} c_{1} \delta\left|\gamma^{\prime}\right|$ where $c_{1}$ is as in (P1). This together with Lemma 3.3 gives $k<\frac{K_{1}}{\delta^{3}} b$ where $K_{1}=8 c_{1}^{-3} K$ and $K$ is as in Lemma 3.3.

The next lemma says that outside of $\mathcal{C}^{(1)}$, iterates of $b$-horizontal vectors behave in a way very similar to that in 1D. Its proof is an easy adaption of the arguments in Sections 2.1 and 2.2 made possible by part (a) of the last lemma.

LEMMA 3.5. There exists $c_{2}>0$ independent of $\delta$ such that the following hold: Let $z_{0} \in R_{1}$ be such that $z_{i} \in R_{1} \backslash \mathcal{C}^{(1)}$ for $i=0,1, \cdots, n-1$, and let $w_{0} \in X_{z_{0}}$ be b-horizontal. Then
(i) $\left|w_{n}\right|>c_{2} \delta e^{\frac{1}{4} \lambda_{0} n}\left|w_{0}\right|$;
(ii) if, in addition, $z_{n} \in \mathcal{C}^{(1)}$, then $\left|w_{n}\right| \geq c_{2} e^{\frac{1}{4} \lambda_{0} n}\left|w_{0}\right|$.

### 3.6. Properties of $e_{1}(S)$ for suitable $S$

We consider in this subsection $e_{1}$ of $D T$ restricted to suitable choices of $S$.
Lemma 3.6. For $z_{0} \notin \mathcal{C}^{(1)}$, let $w \in X_{z_{0}}$ be b-horizontal, and let $S \subset X_{z_{0}}$ be any 2 D plane containing $w$. Then $\angle\left(e_{1}(S), w\right)>K^{-1} \delta$.

Proof. Assuming $|w|=1$, write $e_{1}=a_{1} w+a_{2} v$ where $v \in S$ is a unit vector $\perp w$. Then $K b>\left|D T\left(e_{1}\right)\right|=\left|a_{1} D T(w)+a_{2} D T(v)\right|$. Since $|D T(w)|>K^{-1} \delta$, it follows that $\left|a_{2}\right|>K^{-1} \delta$.

Let $\gamma$ be a $C^{2}(b)$ curve in $\mathcal{C}^{(1)}$ parametrized by arclength. At each point $\gamma(s)$, we let $S(s)=S\left(\gamma^{\prime}(s), \mathbf{v}\right)$. Let $u(s)=\gamma^{\prime}(s), v(s)=\frac{\mathbf{v}-\langle u, \mathbf{v}\rangle u}{|\mathbf{v}-\langle u, \mathbf{v}\rangle u|}$ (i.e. $v(s)$ is a unit vector in $S(s)$ perpendicular to $u(s))$ and let $\eta(s)=\left\langle e_{1}(S(s)), v(s)\right\rangle$.

Lemma 3.7. Let $\gamma(s), S(s)$ and $\eta(s)$ be as above. Then $e_{1}(S(s))$ is welldefined on all of $\gamma$, and

$$
\begin{equation*}
\left|\frac{d \eta(s)}{d s}\right|>K_{1}^{-1} \tag{3}
\end{equation*}
$$

for some $K_{1}$ independent of $\gamma$.

Lemma 3.7 is a direct consequence of our assumptions that $f^{\prime \prime}(\hat{x}) \neq 0$ and $\partial_{y^{i}} \hat{T}_{(\hat{x}, 0)}^{1} \neq 0$ for $\hat{x} \in C$. A proof is given in Appendix A.9.

### 3.7. Critical points on $C^{2}(b)$ curves in $\mathcal{C}^{(1)}$

We fix $\hat{K}_{0}>10 K_{0}$ where $K_{0}$ satisfies $\left|D \hat{T}_{(x, 0)}^{1} \mathbf{v}\right|>K_{0}^{-1}$.
Definition 3.2. Let $\gamma$ be a $C^{2}(b)$-curve in $\mathcal{C}^{(1)}$. We say that $z_{0}$ is a critical point of order $n$ on $\gamma$ if
(a) $\left|D T_{z_{0}}^{i}(\mathbf{v})\right| \geq \hat{K}_{0}^{-1}$ for $i=1,2, \cdots, n$;
(b) at $z_{0}, \angle\left(e_{n}(S), \gamma^{\prime}\right)=0$ with $S=S\left(\gamma^{\prime}, \mathbf{v}\right)$.

Corollary 3.1 (Corollary to Lemma 3.7). On any $C^{2}(b)$-curve traversing the full length of a component of $\mathcal{C}^{(1)}$, there exists a unique critical point of order 1 .

We now turn to the problem of inducing new critical points on nearby curves starting from a known critical point on a $C^{2}(b)$-curve. We begin with two lemmas the exact form of which will be used.

Lemma 3.8. Let $\gamma$ and $\hat{\gamma}$ be $C^{2}(b)$-curves parametrized by arclength in $\mathcal{C}^{(1)}$. Assume
(a) $\gamma(0)$ is a critical point of order $n$ on $\gamma$ with $\left|D T_{\gamma(0)}^{i}(\mathbf{v})\right| \geq 2 \hat{K}_{0}^{-1}$ for $i \leq n$;
(b) $|\gamma(0)-\hat{\gamma}(0)|,\left|\gamma^{\prime}(0)-\hat{\gamma}^{\prime}(0)\right|<b^{\frac{n}{4}}$; and
(c) $\hat{\gamma}(s)$ is defined for all $s \in\left[-b^{\frac{n}{5}}, b^{\frac{n}{5}}\right]$.

Then there exists a unique $s,|s|<K b^{\frac{n}{4}}$, such that $\hat{\gamma}(s)$ is a critical point of order $n$ on $\hat{\gamma}$.

Lemma 3.9. There exists $K_{2}$ for which the following holds: Let $\gamma$ be a $C^{2}(b)$-curve parametrized by arclength in $\mathcal{C}^{(1)}$, and let $z=\gamma(0)$ be a critical point of order $n$. If
(a) $\left|D T_{z}^{i}(\mathbf{v})\right| \geq 2 \hat{K}_{0}^{-1}$ for $i=1,2, \cdots, n+m$, and
(b) $\gamma(s)$ is defined for $s \in\left[-K_{2}(K b)^{n}, K_{2}(K b)^{n}\right]$,
then there exists a unique critical point $\hat{z}$ of order $n+m$ on $\gamma$, and $|\hat{z}-z|<$ $K_{2}(K b)^{n}$.

Proofs of Corollary 3.1 and Lemmas 3.8 and 3.9 are given in Appendix A. 10 .

### 3.8. Tracking $w_{n}=D T_{z_{0}}^{n}\left(w_{0}\right)$ : a splitting algorithm

Let $z_{0} \in R_{1}$, and let $w_{0} \in X_{z_{0}}$ be a $b$-horizontal unit vector. In the case where $z_{i} \notin \mathcal{C}^{(1)}$ for all $i$, the resemblance to 1D dynamics is made clear in

Lemmas 3.4 and 3.5. Consider next an orbit $z_{0}, z_{1}, \cdots$ that visits $\mathcal{C}^{(1)}$ exactly once, say at time $t>0$. Assume:
(i) There exists $\ell>1$ such that $\left|D T_{z_{t}}^{i}(\mathbf{v})\right| \geq K_{0}^{-1}$ for all $i \leq \ell$, so that in particular $e_{\ell}(S)$ is defined at $z_{t}$ with $S=S\left(\mathbf{v}, w_{t}\right)$.
(ii) $\angle\left(w_{t}, e_{\ell}(S)\right) \geq b^{\frac{\ell}{2}}$.

Then $D T_{z_{0}}^{i}\left(w_{0}\right)$ can be analyzed as follows. We split $w_{t}$ into $w_{t}=\hat{w}_{t}+\hat{E}$ where $\hat{w}_{t}$ is a scalar multiple of $\mathbf{v}$ and $\hat{E}$ is a scalar multiple of $e_{\ell}(S)$. For $i \leq t$ and $i \geq t+\ell$, let $w_{i}^{*}=w_{i}$. For $i$ with $t<i<t+\ell$, let $w_{i}^{*}=D T_{z_{t}}^{i-t}\left(\hat{w}_{t}\right)$. We claim that all the $w_{i}^{*}$ are $b$-horizontal vectors, and that $\left\{\left|w_{i+1}^{*}\right| /\left|w_{i}^{*}\right|\right\}_{i=0,1,2, \cdots}$ resembles a sequence of 1D derivatives, with $\left|w_{t+1}^{*}\right| /\left|w_{t}^{*}\right|$ simulating a drop in the derivative when an orbit comes near a critical point in 1D.

In light of Lemma 3.4, to show that $w_{i}^{*}$ is $b$-horizontal, it suffices to consider $w_{t+\ell}^{*}$. Observe from assumption (ii) above that $\left|\hat{w}_{t}\right|>b^{\frac{\ell}{2}}|\hat{E}|$. (Note that $e_{\ell}$ is close to $e_{1}$ from Lemma 3.1, and $s\left(e_{1}\right)<K \delta$ for $z \in \mathcal{C}^{(1)}$.) This together with assumption (i) implies that

$$
\left|D T_{z_{t}}^{\ell}(\hat{E})\right| \leq(K b)^{\ell}|\hat{E}| \leq K^{\ell} b^{\frac{\ell}{2}}\left|\hat{w}_{t}\right| \leq K_{0} K^{\ell} b^{\frac{\ell}{2}}\left|D T_{z_{t}}^{\ell}\left(\hat{w}_{t}\right)\right| .
$$

Since $s\left(D T_{z_{t}}^{\ell}\left(\hat{w}_{t}\right)\right)<\frac{3 K_{0}}{2 \delta} b$ (see Lemma 3.4), $w_{t+\ell}^{*}=D T_{z_{t}}^{\ell}\left(\hat{w}_{t}\right)+D T_{z_{t}}^{\ell}(\hat{E})$ is $b$-horizontal.

The discussion above motivates the following:
Splitting algorithm. We give this algorithm only for $z_{0} \in \mathcal{C}^{(1)}$ and $w_{0}=\mathbf{v}$ since this is mostly how it will be used. Let $t_{1}<t_{2}<\cdots$ be the times $>0$ when $z_{i} \in \mathcal{C}^{(1)}$. For each $t_{j}$, fix $\ell_{t_{j}} \geq 2$ with the property that $\left|D T_{z_{t_{j}}}^{i}(\mathbf{v})\right|>K_{0}^{-1}$ for $i=1, \cdots, \ell_{t_{j}}$ (such $\ell_{t_{j}}$ always exist). The following algorithm generates two sequences of vectors $w_{i}^{*}$ and $\hat{w}_{i}$ :

1. For $0 \leq i<t_{1}$, let $w_{i}^{*}=\hat{w}_{i}=w_{i}$.
2. At $i=t_{1}$, set $w_{i}^{*}=w_{i}$, and define $\hat{w}_{i}$ as follows: If $w_{i}^{*}$ is a scalar multiple of $\mathbf{v}$, let $\hat{w}_{i}=w_{i}^{*}$. If not, let $S=S\left(w_{i}^{*}, \mathbf{v}\right)$. Then split $w_{i}^{*}$ into

$$
w_{i}^{*}=\hat{w}_{i}+\hat{E}_{i}
$$

where $\hat{w}_{i}$ is a scalar multiple of $\mathbf{v}$ and $\hat{E}_{i}$ is a scalar times $e_{\ell_{i}}(S)$.
3. For $i>t_{1}$, we let

$$
\begin{equation*}
w_{i}^{*}=D T_{z_{i-1}}\left(\hat{w}_{i-1}\right)+\sum_{j: t_{j}+\ell_{t_{j}}=i} D T_{z_{t_{j}}}^{\ell_{t_{j}}}\left(\hat{E}_{t_{j}}\right), \tag{4}
\end{equation*}
$$

and define $\hat{w}_{i}$ as follows: if $i=t_{j}$, split $w_{i}^{*}$ into $w_{i}^{*}=\hat{w}_{i}+\hat{E}_{i}$ as in item 2 ; if $i \neq t_{j}$ for any $j$, set $\hat{w}_{i}=w_{i}^{*}$.

This algorithm is of interest when the contributions from the $\hat{E}_{i}$-terms as they rejoin $w_{i}^{*}$ are negligible; the meaning of $w_{i}^{*}$ and $\hat{w}_{i}$ are unclear otherwise.

The next lemma contains a set of technical conditions describing a "good" situation:

Lemma 3.10. Let $z_{0}, \ell_{t_{j}}, w_{i}$ and $w_{i}^{*}$ be as above, and let $I_{j}:=\left[t_{j}, t_{j}+\ell_{t_{j}}\right)$. Assume
(a) for each $i=t_{j},\left|\hat{w}_{i}\right|>b^{\frac{\ell_{i}}{2}}\left|\hat{E}_{i}\right|$;
(b) the $I_{j}$ are nested, i.e. for $j<j^{\prime}$, either $I_{j} \cap I_{j^{\prime}}=\emptyset$ or $I_{j^{\prime}} \subset I_{j}$.

Then the $w_{i}^{*}$ are b-horizontal.
A proof of Lemma 3.10 is given in Appendix A.11.

### 3.9. Attractors arising from interval maps

We explain how to deal with the endpoints of $I$ in the case where $I$ is an interval.

Let $f \in \mathcal{M}$. By assumption, $f(I) \subset \operatorname{int}(I)$. We let $\Lambda=\Lambda^{(n)}$ be as in Lemma 2.6 where $n$ is large enough that $f(I)$ is well inside $\left[x_{1}, x_{2}\right]$, the shortest interval containing $\Lambda$. It is a standard fact that periodic points are dense in topologically transitive shifts of finite type. From this, one deduces easily that pre-periodic points are dense in all shifts of finite type, transitive or not. Let $y_{1}$ and $y_{2}$ be pre-periodic points so that $f(I)$ is well inside $\left[y_{1}, y_{2}\right]$. For $i=1,2$, let $k_{i}$ and $n_{i}$ be such that $f^{k_{i}+n_{i}}\left(y_{i}\right)=f^{n_{i}}\left(y_{i}\right)$. Our plan is to prove the following for $T$ when $b$ is sufficiently small:
(i) Near $\left(f^{n_{i}}\left(y_{i}\right), 0\right), i=1,2, T$ has a periodic point $z_{i}$ of period $k_{i}$.
(ii) $z_{i}$ is hyperbolic; it therefore has a codimension one stable manifold $W^{s}\left(z_{i}\right)$. We claim that $W_{i}$, the connected component of $W^{s}\left(z_{i}\right)$ containing $z_{i}$, spans $R_{1}$ in the sense that it is the graph of a function from $\left\{|y| \leq(m-1)^{\frac{1}{2}} b\right\}$ to $I$.
(iii) Near $\left(y_{i}, 0\right)$ there is a connected component $V_{i}$ of $W^{s}\left(z_{i}\right) ; V_{i}$ also spans $R_{1}$.
(iv) If $\hat{R}_{1}$ is the part of $R_{1}$ between $V_{1}$ and $V_{2}$, then $T\left(\hat{R}_{1}\right) \subset \hat{R}_{1}$.

The existence and hyperbolicity of $z_{i}$ follows from the fact that $\left|\left(f^{k_{i}}\right)^{\prime}\left(f^{n_{i}}\left(y_{i}\right)\right)\right|>1$ (Lemma 2.1). That $W_{i}$ spans the cross-section of $R_{1}$ follows from Lemma 3.1 and the construction in Section 3.3 with $n \rightarrow \infty$. Moving on to (iii), the existence of a component of $T^{-n_{i}} W_{i}$ near $\left(y_{i}, 0\right)$ follows by continuity. Repeating the arguments at $z_{i}$ on a (any) point in $V_{i}$, we see that not only does $V_{i}$ span $R_{1}$ but its tangent vectors make angles $>K^{-1} \delta$ with the $x$-axis. Thus the diameter of $V_{i}$ is arbitrarily small as $b \rightarrow 0$, and (iv) follows from $f(I) \subset\left(y_{1}, y_{2}\right)$.

In Part II, we restrict the domain of $T$ to $\hat{R}_{1}$. The two ends of $\hat{R}_{1}$, namely $V_{1} \cup V_{2}$, are asymptotic to the periodic orbits of $z_{1}$ and $z_{2}$. In particular, they stay away from $\mathcal{C}^{(1)}$. This part of $\partial \hat{R}_{1}$ is not visible in local arguments. In Sections 7 and 8, in the treatment of monotone branches, there will be some special branches that end in $T^{j}\left(V_{i}\right)$. Modifications in the arguments are straightforward.

In Part III, we take $z_{i}(a)$ to be continuations of the same periodic orbits, so that $\hat{R}_{1}(a)$ varies continuously with $a$.

Notation for the rest of the paper.

- We assume $T=\left(\hat{T}^{1}, \cdots, \hat{T}^{m}\right): X \rightarrow X$ is such that $\left\|\hat{T}^{j}\right\|_{C^{3}}<b$ for $j=2, \cdots, m$.
- $R_{1}:=I \times\left\{y \in \mathbb{R}^{m-1}:|y|<(m-1)^{\frac{1}{2}} b\right\} ; R_{k}:=T^{k-1} R_{1}$ for $k=2,3, \cdots$.
- For definiteness, we let $\mathcal{F}_{1}$ be the foliation on $R_{1}$ given by $\{y=$ constant $\}$ (this can be replaced by any foliation whose leaves are $C^{2}(b)$ curves); for $k>1, \mathcal{F}_{k}:=T_{*}^{k-1}\left(\mathcal{F}_{1}\right)$; i.e., the leaves of $\mathcal{F}_{k}$ are the $T^{k-1}$-images of those of $\mathcal{F}_{1}$.
- A subset $H \subset R_{j}$ is called a section of $R_{j}$ if it is the diffeomorphic image of $\Phi:[-1,1] \times D_{m-1} \rightarrow R_{j}$ with $\Phi^{-1}\left(\partial R_{j}\right)=[-1,1] \times \partial D_{m-1}$. A section $H$ of $R_{j}$ is called horizontal if each component of $\Phi\left(\{ \pm 1\} \times D_{m-1}\right)$ is contained in a hyperplane $\{x=$ const $\}$ and all the leaves of $\left.\mathcal{F}_{j}\right|_{H}$ are $C^{2}(b)$-curves. The cross-sectional diameter of a horizontal section $H$ is defined to be the supremum of $\operatorname{diam}(V \cap H)$ as $V$ varies over all hyperplanes perpendicular to $S^{1}$.
- The distance from $z$ to $z^{\prime}$ in $R_{1}$ is denoted by $\left|z-z^{\prime}\right|$, and their horizontal distance, i.e. difference in $x$-coordinates, is denoted by $\left|z-z^{\prime}\right|_{h}$.


## PART II. PHASE-SPACE DYNAMICS

The goal of Part II is to identify, among all maps $T: X \rightarrow X$ that are near small perturbations of 1D maps, a class $\mathcal{G}$ with certain desirable features. To explain what we have in mind, consider the situation in 1D. In Section 2.2, we show that for maps sufficiently near $f_{0} \in \mathcal{M}$, two relatively simple conditions, (G1) and (G2), imply dynamical properties (P1)-(P3), which in turn lead to other desirable characteristics. Our class $\mathcal{G}$ will be modelled after these maps.

The first major hurdle we encounter as we attempt to formulate higher dimensional analogs of (G1) and (G2) is the absence of a well defined critical set. As we will show, the concept of a critical set can be defined, but only inductively and only for certain maps. This implies that our "good maps" can
only be identified inductively. The task before us, therefore, is the inductive construction of $\mathcal{G}_{n}, n=1,2, \cdots$, consisting of maps that are "good" in their first $n$ iterates, and $\mathcal{G}$ is taken to be $\cap_{n \geq 0} \mathcal{G}_{n}$.

We do not claim in Part II that $\mathcal{G}$ is nonempty, and we consider one map at a time to determine if it is in $\mathcal{G}$; no parameters are involved. The existence (and abundance) of maps in $\mathcal{G}$ is proved in Part III.

Organization. Sections 4-9, which comprise Part II, are organized as follows: Section 4.1 contains five statements describing five aspects of dynamical behavior. Together, these statements give a snapshot of the maps in $\mathcal{G}_{n}$ for certain $n$. The rest of Section 4 is devoted to the elucidation of the ideas introduced.

Implications of these ideas are developed in Section 5, and a formal inductive construction of $\mathcal{G}_{n}$ for $n \leq N_{0} \sim\left(\log \frac{1}{b}\right)^{2}$ is given in Section 6 .

After $N_{0}$ iterates, a fundamental, qualitative change in geometry occurs. The new complexities that arise are dealt with in Sections 7 and 8.

The existence of SRB measures for $T \in \mathcal{G}$ is proved in Section 9 .
The notation is as in Section 1, namely that $f: S^{1} \rightarrow S^{1}, F: R_{1} \rightarrow S^{1}$ and $F^{\#}: R_{1} \rightarrow R_{1}$ are related by $F(x, 0)=f(x)$ and $F^{\#}(x, y)=(F(x, y), 0)$, and $T: R_{1} \rightarrow R_{1}$ is a $C^{3}$ embedding.

Standing hypotheses. Throughout Part II, we fix $f_{0} \in \mathcal{M}$ and $K_{0}>1$, and consider

- $f: S^{1} \rightarrow S^{1}$ with $\left\|f-f_{0}\right\|_{C^{2}}<a$,
- $F: R_{1} \rightarrow S^{1}$ with $\|F\|_{C^{3}}<K_{0}$ and $\left|D F_{(\hat{x}, 0)}(\mathbf{v})\right|>K_{0}^{-1}$ for $\hat{x} \in C\left(f_{0}\right)$, and
- $T: R_{1} \rightarrow R_{1}$ with $\left\|T-F^{\#}\right\|_{C^{3}}<b$
where $a, b>0$ are as small as need be. The letter $K$ is used as a generic constant which, in Part II, is allowed to depend only on $f_{0}, K_{0}$ and our choice of $\lambda$.


## 4. Critical structure and orbits

### 4.1. Formal assumptions

We describe in this subsection several aspects of geometric and dynamical behaviors to be viewed as desirable. These assumptions, labelled (A1)- (A5), will eventually be part of the inductive cycle up to a certain time. For the moment they are only formal statements.

For purposes of the present discussion, $\lambda>0$ can be any number $<\frac{1}{5} \lambda_{0}$ (see $\S 2.2$ ). We choose $\alpha$ so that $b \ll \alpha \ll \min (\lambda, 1)$, and let $\alpha^{*}=\frac{6}{\lambda} \alpha$. Let $\theta=\frac{K}{\log \frac{1}{b}}$ where $K$ is chosen so that $b^{\theta}<\|D T\|^{-20}$. Let $N$ be a positive integer $\gg 1$. For simplicity of notation, we assume $\theta N, \theta^{-1}, \frac{1}{\alpha^{*}} \in \mathbb{Z}^{+}$(otherwise write $\left.[\theta N],\left[\theta^{-1}\right],\left[\frac{1}{\alpha^{*}}\right]\right)$.
(A1) Geometry of critical regions. There are sets $\mathcal{C}^{(1)} \supset \mathcal{C}^{(2)} \supset \cdots \supset$ $\mathcal{C}^{(\theta N)}$ called critical regions with the following properties:
(i) $\mathcal{C}^{(1)}$ is as introduced in Section 3.4. For $1<k \leq \theta N, \mathcal{C}^{(k)}$ is the union of a finite number of connected components $\left\{Q^{(k)}\right\}$ each one of which is a horizontal section of $R_{k}$ of length $\min \left(2 \delta, 2 e^{-\lambda k}\right)$ and cross-sectional diameter $<b^{\frac{k}{2}}$.
(ii) $\mathcal{C}^{(k)}$ is related to $\mathcal{C}^{(k-1)}$ as follows: For each $Q^{(k-1)}$, either $R_{k} \cap Q^{(k-1)}=\emptyset$ or it meets $Q^{(k-1)}$ in a finite number of horizontal sections $\{H\}$ each one of which extends $>\frac{1}{2} e^{-\alpha k}$ beyond the two ends of $Q^{(k-1)}$. Each $H \cap Q^{(k-1)}$ contains exactly one component of $\mathcal{C}^{(k)}$ located roughly in the middle. (See Figure 1.)
(iii) Inside each $Q^{(k)}$, a point $z_{0}=z_{0}^{*}\left(Q^{(k)}\right)$ whose $x$-coordinate is exactly half-way between those of the two ends of $Q^{(k)}$ is singled out; $z_{0}$ is a critical point of order $k$ in the sense of Definition 3.2 with respect to the leaf of the foliation $\mathcal{F}_{k}$ containing it.


Figure 1. Structure of critical regions
We call $z_{0}^{*}\left(Q^{(k)}\right)$ a critical point of generation $k$, and let $\Gamma_{k}$ denote the set of all critical points of generation $\leq k$. Let $Q^{(k)}\left(z_{0}\right)$ denote the component of $\mathcal{C}^{(k)}$ containing $z_{0}$.

The next three assumptions prescribe certain behaviors on the orbits of $z_{0} \in \Gamma_{\theta N}$. To state them, we need the following definitions:

First, we define a notion of distance to critical set for $z_{i}$, denoted $d_{\mathcal{C}}\left(z_{i}\right)$. If $z_{i} \notin \mathcal{C}^{(1)}$, let $d_{\mathcal{C}}\left(z_{i}\right)=\delta+d\left(z_{i}, \mathcal{C}^{(1)}\right)$. If $z_{i} \in \mathcal{C}^{(1)}$, we let $d_{\mathcal{C}}\left(z_{i}\right)=\left|z_{i}-\phi\left(z_{i}\right)\right|$ where $\phi\left(z_{i}\right)$ is defined as follows. Let $j$ be the largest integer $\leq \alpha^{*} \theta i$ with the property that $z_{i} \in \mathcal{C}^{(j)}$. Then $\phi\left(z_{i}\right):=z_{0}^{*}\left(Q^{(j)}\left(z_{i}\right)\right)$ is called the guiding critical point for $z_{i}$. As the name suggests, the orbit of $\phi\left(z_{i}\right)$ will be thought of as guiding that of $z_{i}$ through its derivative recovery. Suppose $z_{i} \in \mathcal{C}^{(1)}$ and $\phi\left(z_{i}\right)$ is of generation $j$. We say $w \in X_{z_{i}}$ is correctly aligned, or correctly aligned
with respect to the leaves of the $\mathcal{F}_{j}$-foliation, if $\angle\left(\tau_{j}\left(z_{i}\right), w\right) \ll K_{1}^{-1} d_{\mathcal{C}}\left(z_{i}\right)$ where $K_{1}^{-1}$ is the lower bound on $\left|\frac{d}{d s} \eta\right|$ along $C^{2}(b)$-curves in $\mathcal{C}^{(1)}$ in Lemma 3.7 and $\tau_{j}\left(z_{i}\right)$ is tangent to the leaf of $\mathcal{F}_{j}$ through $z_{i}$. We say $w$ is correctly aligned with $\varepsilon$-error if $\varepsilon \ll K_{1}^{-1}$ and $\angle\left(\tau_{j}\left(z_{i}\right), w\right)<\varepsilon d_{\mathcal{C}}\left(z_{i}\right)$.

For $z_{0} \in \Gamma_{\theta N}$, we let $w_{0}=\mathbf{v}$, and for a chosen family of $\ell_{i}$ corresponding to $z_{i} \in \mathcal{C}^{(1)}$, let $w_{i}^{*}, i=0,1,2, \cdots$, be given by the splitting algorithm in Section 3.8. The numbers $\left\{\ell_{i}\right\}$ are called the splitting periods for $z_{0}$. Let $\varepsilon_{0} \ll K_{1}^{-1}$ be fixed. We shrink $\delta$ if necessary so that it is $\ll \varepsilon_{0}$.
(A2)-(A4) Properties of critical orbits. For $z_{0} \in \Gamma_{\theta N}$ of generation $k$, the following hold for all $i \leq k \theta^{-1}$ :
(A2) $d_{\mathcal{C}}\left(z_{i}\right)>\min \left(\delta, e^{-\alpha i}\right)$.
(A3) There exist $\left\{\ell_{j}\right\}$ (to be specified in $\S 4.4$ ) so that $w_{i}^{*}$ is correctly aligned with $\varepsilon_{0}$-error when $z_{i} \in \mathcal{C}^{(1)}$.
(A4) $\left|w_{i}^{*}\right|>\frac{1}{2} c_{2} e^{\lambda i}$ where $c_{2}$ is as in Lemma 3.5.
Our next assumption gives the relation between $z_{i}$ and $\phi\left(z_{i}\right)$. Let $\hat{\beta}$ be such that $\alpha \ll \hat{\beta} \ll 1$. For $z_{0}, \xi_{0} \in R_{1}$, let $\hat{p}\left(z_{0}, \xi_{0}\right)$ be the smallest $j>0$ such that $\left|z_{j}-\xi_{j}\right| \geq e^{-\hat{\beta} j}$. For reasons to be explained in Section 4.3B, we will be interested in a range of $p$ near $\hat{p}\left(z_{0}, \xi_{0}\right)$. Inside each $Q^{(k)}$, let

$$
B^{(k)}=\left\{z \in Q^{(k)}:\left|z-z_{0}^{*}\left(Q^{(k)}\right)\right|_{h}<b^{\frac{1}{5}} k .\right.
$$

(A5) How critical orbits influence nearby orbits. For $z_{0}=z_{0}^{*}\left(Q^{(k)}\right)$ and $\xi_{0} \in Q^{(k)} \backslash B^{(k)}, k \leq \theta N$, the following hold for all $p \in\left[\hat{p}\left(z_{0}, \xi_{0}\right)\right.$, $\left.\left(1+\frac{9}{\lambda} \alpha\right) \hat{p}\left(z_{0}, \xi_{0}\right)\right]:$
(i) (Length of bound period). Suppose $\left|z_{0}-\xi_{0}\right|=e^{-h}$. Then

$$
\frac{1}{3 \ln \|D T\|} h \leq p \leq \frac{3}{\lambda} h
$$

the first inequality being valid if $\frac{1}{3 \ln \|D T\|} h \leq k \theta^{-1}$ and the second if $\frac{3}{\lambda} h \leq k \theta^{-1}$.
(ii) (Partial derivative recovery). If $p \leq k \theta^{-1}$, then $\left|w_{p}\left(z_{0}\right)\right|\left|\xi_{0}-z_{0}\right| \geq e^{\frac{1}{3} \lambda p}$.
(iii) (Quadratic nature of turns). Let $\gamma$ be the $\mathcal{F}_{k}$-leaf segment joining $\xi_{0}$ to $B^{(k)}$. Then for all $\eta_{0} \in \gamma$ and $\ell\left(\eta_{0}\right)<i \leq \min \left\{p, k \theta^{-1}\right\}$,

$$
\left|\eta_{i}-z_{i}\right|=\frac{1}{2}\left(\left|\frac{d e_{1}}{d s}\left(z_{0}\right)\right| \pm \mathcal{O}(b)\right) \cdot\left(\left|w_{i}\left(z_{0}\right)\right| \pm \mathcal{O}\left(\left|\eta_{0}-z_{0}\right|^{\frac{1}{2}}\right)\right) \cdot\left|\eta_{0}-z_{0}\right|^{2}
$$

Here $\ell\left(\eta_{0}\right)$ is defined by $b^{\frac{\ell\left(\eta_{0}\right)}{2}}=\left|\eta_{0}-z_{0}\right|$, and $e_{1}=e_{1}(S)$ where $S=$ $S\left(\mathbf{v}, \tau_{k}\right), \tau_{k}$ being the tangent to the $\mathcal{F}_{k}$-leaf through $z_{0}$.

This completes the formulation of the five statements (A1)-(A5). We also write (A1)(N)-(A5)(N) when more than one time frame is involved. The rest of this section contains some immediate clarifications.

Three important time scales. We point out that in the dynamical picture described by (A1)-(A5), there are three distinct time scales: $\theta N \ll \alpha N \ll N$. The fastest time scale, $N$, gives the number of times the map is iterated. The slowest, $\theta N$, is the number of generations of critical regions and critical points constructed. The middle time scale, which is on the order of $\alpha N\left(\alpha^{*} N\right.$ to be precise), is an upper bound for the lengths of the bound periods initiated by critical orbits returning to $\mathcal{C}^{(1)}$ at times $\leq N$ (this follows from (A2) and (A5)(i) combined).

We assume (A1)-(A5) for the rest of Section 4.

### 4.2. Clustering of critical orbits

In Section 4.1, we presented a viewpoint - convenient for some practical purposes - in which a critical point $z_{0}^{*}\left(Q^{(k)}\right)$ in each component $Q^{(k)}$ of $\mathcal{C}^{(k)}$ is singled out for special consideration. To understand the relation among the points in $\Gamma_{\theta N}$, it is more fruitful to group them into clusters. We propose here to view these clusters as represented by $B^{(k)}$. To justify this view, we prove

Lemma 4.1. For all $k<\hat{k}<\theta N$, if $Q^{(\hat{k})} \subset Q^{(k)}$, then

$$
\left|z_{0}^{*}\left(Q^{(k)}\right)-z_{0}^{*}\left(Q^{(\hat{k})}\right)\right|<K b^{\frac{k}{4}}
$$

and $B^{(\hat{k})} \subset B^{(k)}$.
The proof of this lemma uses the technical estimate below. Both results rely on the geometric information on $Q^{(k)}$ in (A1). Proofs are given in Appendix A.12.

Lemma 4.2. Let $k<\hat{k}, Q^{(\hat{k})} \subset Q^{(k)}, z \in Q^{(k)}, \hat{z} \in Q^{(\hat{k})}$, and let $\gamma$ and $\hat{\gamma}$ be the $\mathcal{F}_{k^{-}}$and $\mathcal{F}_{\hat{k}}$-leaves containing $z$ and $\hat{z}$ respectively. Let $\tau$ and $\hat{\tau}$ be the tangent vectors to $\gamma$ and $\hat{\gamma}$ at $z$ and $\hat{z}$. Then

$$
\angle(\tau, \hat{\tau}) \leq b^{\frac{k}{4}}+K \delta^{-3} b \cdot|z-\hat{z}|_{h} .
$$

Evolution of critical blobs. A theme that runs through our discussion is that orbits emanating from the same $B^{(k)}$ are viewed as essentially indistinguishable for $k \theta^{-1}$ iterates. Informally, we call these finite orbits of $B^{(k)}$ critical blobs.

Recall that $\theta$ is assumed so that $b^{\theta}<\|D T\|^{-20}$. This implies that for all $i \leq k \theta^{-1}, \operatorname{diam}\left(T^{i} B^{(k)}\right)<b^{\frac{1}{5} k}\|D T\|^{i}<\left(b^{\theta}\right)^{\frac{1}{5} i}\|D T\|^{i}$. This is $\ll e^{-\alpha i}$, the minimum allowed distance to the critical set (see (A2)).

Obviously, we cannot iterate indefinitely and hope that $T^{i} B^{(k)}$ remains small; that is why we regard $z_{0}^{*}\left(Q^{(k)}\right)$ as active for only $k \theta^{-1}$ iterates. The word "active" here refers to both (i) prescribed behavior for $z_{i}$ (as in (A2)(A4)) and (ii) the use of $z_{i}$ as guiding critical orbit or in the sense of (A5).

It is useful to keep in mind the following dynamical picture: At time $i=0$, $T$ has a set $B^{(1)}$ corresponding to each critical point of $f$. For $i \leq \theta^{-1}$, the $T^{i}$-images of $B^{(1)}$ are relatively small, so that $\left\{T^{i} B^{(1)}\right\}_{i=0,1, \cdots, \theta^{-1}}$ for each $B^{(1)}$ can be treated as a single orbit.

As $i$ increases, the sizes of $T^{i} B^{(1)}$ become larger, eventually becoming too large for $\left\{T^{i} B^{(1)}\right\}_{i=0,1, \ldots}$ to be treated as a single orbit. We stop considering these critical blobs long before that time, however. At time $i=\theta^{-1}$, we replace each $T^{\theta^{-1}} B^{(1)}$ by the collection of $T^{\theta^{-1}} B^{(2)}$ contained in it. For $\theta^{-1}<i \leq 2 \theta^{-1}$, $T^{i} B^{(2)}$ are again relatively small, and so can be viewed as a finite collection of orbits. At time $i=2 \theta^{-1}$, each $T^{2 \theta^{-1}} B^{(2)}$ is replaced by the collection of $T^{2 \theta^{-1}} B^{(3)}$ inside it, and so on.

As $i$ increases, the number of relevant critical blobs increases, each becoming smaller in size. Blobs that have separated move about "independently". By virtue of (A2), they are allowed to come closer to the critical set with the passage of time.

We finish by recording a technical fact that will be used in conjunction with Lemma 3.8.

Lemma 4.3. For any $C^{2}(b)$-curve $s \mapsto l(s)$ traversing a given $B^{(k)} \subset Q^{(k)}$, there exists a point in l, denoted by l(0), such that

$$
\angle\left(l^{\prime}(0), \tau\left(z_{0}\right)\right)<b^{\frac{k}{4}}
$$

where $z_{0}=z_{0}^{*}\left(Q^{(k)}\right)$ and $\tau\left(z_{0}\right)$ is tangent to the leaf of $\mathcal{F}_{k}$ at $z_{0}$.
As with Lemma 4.2, Lemma 4.3 is proved by a straightforward application of Sublemma A.12.1 in Appendix A.12. We leave it as an exercise.

### 4.3. Bound periods

Let $z_{0} \in \Gamma_{\theta N}$ be of generation $k$, and let $z_{i} \in \mathcal{C}^{(1)}, i \leq k \theta^{-1}$. In Section 4.1, we assigned to $z_{i}$ a guiding critical point $\phi\left(z_{i}\right) \in \Gamma_{\theta N}$. (A5)(i)-(iii) hold for all $p \in\left[\hat{p},\left(1+\frac{9}{\lambda} \alpha\right) \hat{p}\right]$ where $\hat{p}=\hat{p}\left(z_{i}, \phi\left(z_{i}\right)\right)$. We now choose a specific number $p=p\left(z_{i}\right)$ in this range with certain desirable properties. This number will be called the bound period of $z_{i}$.
A. Remarks on $\phi(\cdot)$ and $d_{\mathcal{C}}(\cdot)$. In general, when $z_{i} \in \mathcal{C}^{(1)}$, it is in many $Q^{(j)}$. Since $\mathcal{C}^{(j)}$ for larger $j$ give better approximations of the eventual critical set, it is natural to want to define $d_{\mathcal{C}}\left(z_{i}\right)$ using the largest $j$ possible. We do not do exactly that; instead, we take $\phi\left(z_{i}\right)$ to be $z_{0}^{*}\left(Q^{(\hat{j})}\left(z_{i}\right)\right)$ where $\hat{j}$ is the
largest $j \leq \alpha^{*} \theta i$ such that $z_{i} \in Q^{(j)}$. The significance of this upper bound on $j$ will become clear in Section 6. For now we observe

Lemma 4.4. (i) $\left|z_{i}-\phi\left(z_{i}\right)\right| \gg b^{\frac{\hat{2}}{5}}$; in particular, $z_{i} \in Q^{(\hat{j})} \backslash B^{(\hat{j})}$, so that (A5) applies.
(ii) Let $p \in\left[\hat{p},\left(1+\frac{9}{\lambda} \alpha\right) \hat{p}\right]$ be as in (A5). Then $p \leq \hat{j} \theta^{-1}$.

Proof: Case 1. $\hat{j}+1 \leq \alpha^{*} \theta i$. This implies $z_{i} \in\left(Q^{(\hat{j})} \cap R_{\hat{j}+1}\right) \backslash \mathcal{C}^{(\hat{j}+1)}$, i.e. $d_{\mathcal{C}}\left(z_{i}\right)>e^{-\lambda(\hat{j}+1)}$. Hence $b^{\frac{\hat{3}}{5}} \ll d_{\mathcal{C}}\left(z_{i}\right)$ and $p \ll \hat{j} \theta^{-1}$ by (A5)(i).

Case 2. $\hat{j}+1>\alpha^{*} \theta i$. Using this relation between $i$ and $\hat{j}$, we see that $d_{\mathcal{C}}\left(z_{i}\right)>e^{-\alpha i}>e^{-\frac{\alpha}{\alpha^{*}} \theta^{-1}(\hat{j}+1)}$, which we check is $\gg b^{\frac{\hat{y}}{5}}$ by the definition of $b^{\theta}$ and the facts that $\frac{\alpha}{\alpha^{*}}=\frac{\lambda}{6}$ and $e^{\lambda}<\|D T\|$. Also, $p \leq \frac{3}{\lambda} \alpha i$ by (A2) and (A5)(i). This upper bound is $=\frac{1}{2} \alpha^{*} i \leq \frac{1}{2}(\hat{j}+1) \theta^{-1} \leq \hat{j} \theta^{-1}$.

We use $\phi\left(z_{i}\right)$ to define $d_{\mathcal{C}}\left(z_{i}\right)$. One may ask if it makes a significant difference if some other critical point is used. The answer is that when $d_{\mathcal{C}}\left(z_{i}\right)$ is relatively large, for example when $d_{\mathcal{C}}\left(z_{i}\right)>b^{\frac{1}{5}}$, it does not matter much, but when $d_{\mathcal{C}}\left(z_{i}\right)$ is small, the values of $\left|\hat{z}-z_{i}\right|$ or even $\left|\hat{z}-z_{i}\right|_{h}$ can vary nontrivially as $\hat{z}$ varies over $\Gamma_{\theta N}$. For the same reason, for $z_{i}, z_{j}^{\prime} \in \mathcal{C}^{(1)}$, we cannot conclude - without further information - that $\left|d_{\mathcal{C}}\left(z_{i}\right)-d_{\mathcal{C}}\left(z_{j}^{\prime}\right)\right| \approx\left|z_{i}-z_{j}^{\prime}\right|$, for $z_{i}$ and $z_{j}^{\prime}$ can be in very different "layers" of the critical structure, resulting in $\phi\left(z_{i}\right)$ and $\phi\left(z_{j}^{\prime}\right)$ being relatively far apart.

We do have the following:
Lemma 4.5. (i) Let $z \in Q^{(k)} \backslash B^{(k)}$. Then for all $\hat{z}, \tilde{z} \in \Gamma_{\theta N} \cap B^{(k)}$, we have $|z-\hat{z}|=\left(1 \pm \mathcal{O}\left(b^{\frac{k}{20}}\right)\right)|z-\tilde{z}|$.
(ii) Suppose $\hat{z}_{0}=\phi\left(z_{i}\right)$, and $\hat{z}_{j} \in \mathcal{C}^{(1)}$ for some $0<j<\hat{p}\left(\hat{z}_{0}, z_{i}\right)$. Then $d_{\mathcal{C}}\left(z_{i+j}\right)=\left(1 \pm \mathcal{O}\left(e^{-\frac{1}{2} \hat{\beta} j}\right)\right) d_{\mathcal{C}}\left(\hat{z}_{j}\right)$.

Proof. (i) By Lemma 4.1, $|\hat{z}-\tilde{z}|<K b^{\frac{k}{4}}$, and by assumption, $z$ is $>b^{\frac{k}{5}}$ from the center of $Q^{(k)}$. This proves $|z-\hat{z}|=\left(1 \pm \mathcal{O}\left(b^{\frac{k}{20}}\right)\right)|z-\tilde{z}|$.
(ii) By definition, $\left|z_{i+j}-\hat{z}_{j}\right|<e^{-\hat{\beta} j} \ll e^{-\alpha j}$, which is $<d_{\mathcal{C}}\left(\hat{z}_{j}\right)$ by (A2). As explained above, this in itself is insufficient for guaranteeing the asserted relationship between $d_{\mathcal{C}}\left(z_{i+j}\right)$ and $d_{\mathcal{C}}\left(\hat{z}_{j}\right)$. We have, however, the following additional information: $\operatorname{By}(\mathrm{A} 5)($ iii $)$, there is a curve $\omega$ joining $z_{i}$ to $\hat{z}_{0}$ such that $\operatorname{diam}\left(T^{j}(\omega)\right) \ll e^{-\alpha j}$. Now suppose $\hat{z}_{j} \in \mathcal{C}^{(1)}$ is such that $\phi\left(\hat{z}_{j}\right)=z_{0}^{*}\left(Q^{(\hat{k})}\right)$. Since $\hat{k} \ll j, T^{j}(\omega)$ is contained, or nearly contained, in $Q^{(\hat{k})}\left(\hat{z}_{j}\right)$. Part (i) now enables us to make the desired comparison.
B. Definition of bound periods. Consider $z_{0} \in \Gamma_{\theta N}$. For each $i$ such that $z_{i} \in \mathcal{C}^{(1)}$, let $p\left(z_{i}\right)$ be the bound period of $z_{i}$ to be defined. We say $\left\{p\left(z_{i}\right)\right\}$ has
a nested structure if whenever $i<j$ are such that $z_{i}, z_{j} \in \mathcal{C}^{(1)}$ and $j<i+p\left(z_{i}\right)$, we have $j+p\left(z_{j}\right) \leq i+p\left(z_{i}\right)$.

To define $p\left(z_{i}\right)$, we start with $\hat{p}_{i}:=\hat{p}\left(\phi\left(z_{i}\right), z_{i}\right)$ where $\hat{p}\left(\phi\left(z_{i}\right), z_{i}\right)$ is as defined in Section 4.1. There is no reason why $\left\{\hat{p}_{i}\right\}$ should have a nested structure. We call $j_{0}<j_{1}<\cdots<j_{n}$ a chain of overlapping bound intervals if $z_{j_{k}} \in \mathcal{C}^{(1)}$ and $j_{k} \in\left(j_{k-1}, j_{k-1}+\hat{p}_{j_{k-1}}\right)$ for every $k \leq n$. Let $\Lambda_{i}$ be the set of all integers $k>i$ such that there is a chain of overlapping bound intervals $j_{0}<j_{1}<\cdots<j_{n}$ with $j_{0}=i$ and $j_{n}+\hat{p}_{n} \geq k$. We define $p\left(z_{i}\right):=i^{\prime}-i$ where $i^{\prime}$ is the supremum of the set $\Lambda_{i}$. A priori, $p\left(z_{i}\right)$ can be $\gg \hat{p}_{i}$; it can even be infinite. We prove in Lemma 4.6 below that this is not the case.

Lemma 4.6. For all $z_{0} \in \Gamma_{\theta N}$ and all $z_{i} \in \mathcal{C}^{(1)}$,
(a) $p\left(z_{i}\right)<\left(1+\frac{6}{\lambda} \alpha\right) \hat{p}_{i}$.
(b) $\left\{p\left(z_{i}\right)\right\}$ has a nested structure.

Proof. (a) For $z_{i} \in \mathcal{C}^{(1)}$, let $j$ be such that $i<j<i+\hat{p}_{i}$. Then $d_{\mathcal{C}}\left(z_{j}\right) \approx$ $\left.d_{\mathcal{C}}\left(\left(\phi\left(z_{i}\right)\right)_{j-i}\right)\right)$ by Lemma 4.5(ii). Applying (A2) to $\phi\left(z_{i}\right)$ and then (A5)(i) to $z_{j}$, we obtain $\hat{p}_{j} \leq \frac{3}{\lambda} \alpha(j-i) \leq \frac{3}{\lambda} \alpha \hat{p}_{i}$. If $j_{0}<j_{1}<\cdots<j_{n}$ is a chain of overlapping bound intervals with $j_{0}=i$, then similar reasoning gives $\hat{p}_{j_{k}} \leq$ $\frac{3}{\lambda} \alpha \hat{p}_{j_{k-1}}$, so that

$$
\hat{p}_{j_{0}}+\hat{p}_{j_{1}}+\cdots+\hat{p}_{j_{n}}<\left(1+\frac{3}{\lambda} \alpha+\left(\frac{3}{\lambda} \alpha\right)^{2}+\cdots\right) \hat{p}_{i}<\left(1+\frac{6}{\lambda} \alpha\right) \hat{p}_{i} .
$$

Since this bound is valid for all chains, we have $p\left(z_{i}\right)<\left(1+\frac{6}{\lambda} \alpha\right) \hat{p}_{i}$.
(b) We need to show that if $j \in\left(i, i+p\left(z_{i}\right)\right)$, then $j+p\left(z_{j}\right) \leq i+p\left(z_{i}\right)$. Note that since $p(\cdot)$ is finite, there exists a chain of overlapping intervals $i=$ $j_{0}<\cdots<j_{n}$ such that $j_{n}+\hat{p}_{j_{n}}=i+p\left(z_{i}\right)$. If $j+p\left(z_{j}\right)>i+p\left(z_{i}\right)$, then the chain that goes from $i$ to $i+p\left(z_{i}\right)$ combined with the one that goes from $j$ to $j+p\left(z_{j}\right)$ forms a new chain starting from $i$ and extending beyond $i+p\left(z_{i}\right)$. This contradicts the definition of $p\left(z_{i}\right)$.

Let $\beta=\hat{\beta}-\frac{9}{\lambda} \ln \|D T\| \alpha$, and let $p\left(z_{0}, \xi_{0}\right)$ be the smallest $j$ such that $\left|z_{j}-\xi_{j}\right| \geq e^{-\beta j}$. An easy calculation gives $\hat{p}\left(z_{0}, \xi_{0}\right)\left(1+\frac{9}{\lambda} \alpha\right) \leq p\left(z_{0}, \xi_{0}\right)$.

Clarification: Relation between $\hat{p}(\cdot, \cdot), p(\cdot, \cdot)$ and $p\left(z_{i}\right)$ for $z_{0} \in \Gamma_{\theta N}$.

1. These definitions are brought about by the tension between our desire to define "bound periods" in terms of the distances separating two orbits, and the advantages of having a nested structure for bound periods along individual orbits. We showed in Lemma 4.6 that a nested structure can be arranged if we allow some flexibility in scale when measuring distances, so that for $z_{0} \in \Gamma_{\theta N}$, there exist $\left\{p\left(z_{i}\right)\right\}$ with a nested structure and satisfying $\hat{p}\left(z_{i}, \phi\left(z_{i}\right)\right) \leq p\left(z_{i}\right) \leq$ $p\left(z_{i}, \phi\left(z_{i}\right)\right)$.
2. In general, in results pertaining to a single bound period (e.g. Propositions 5.1), we use $p(\cdot, \cdot)$, so that the result is valid for as long a duration as possible. In situations in which we follow the long range evolution of single orbits (e.g. Section 5.2), a nested structure arranged as above is used.
C. Bound and free states. For $z_{0} \in \Gamma_{\theta N}$ of generation $k$, we now have a decomposition of the orbit $z_{0}, z_{1}, \cdots, z_{k \theta^{-1}}$ into intervals of bound and free periods, i.e. we say $z_{i}$ is free if and only if it is not in a bound period. If we call the maximal bound intervals primary bound periods, the nested structure above allows us to speak of secondary bound periods, tertiary bound periods, and so on. Returns to $\mathcal{C}^{(1)}$ at the beginning of primary bound periods are called free returns, while returns at the start of seconding or higher order bound periods are called bound returns.

### 4.4. The splitting algorithm applied to $D T_{z_{0}}^{i}(\mathbf{v}), z_{0} \in \Gamma_{\theta N}$

The considerations below are motivated by the discussion in Section 3.8 and by Lemma 3.10 in particular. We continue to use the notation here.
A. Splitting periods. Fix $z_{0} \in \Gamma_{\theta N}$. We explain how the $\ell_{i}$ at return times $i$ in (A3) are chosen. From Section 3.8, we see that the following properties are desirable:
(i) $\ell_{i} \geq 2$;
(ii) $\left|D T_{z_{i}}^{j}(\mathbf{v})\right|>K^{-1}$ for $j=1,2, \cdots, \ell_{i}$;
(iii) the intervals $I_{i}=\left[i, i+\ell_{i}\right)$ have the nested property.

We explain why these properties can, in principle, be arranged. As a first approximation, let $\hat{\ell}$ be such that $b^{\frac{\hat{\imath}}{3}}=d_{\mathcal{C}}\left(z_{i}\right)$. Then (ii) holds for all $\ell \leq$ $\frac{5}{3} \hat{\ell}$. To justify this claim, we need to check that $\frac{5}{3} \hat{\ell} \leq$ the order of $\phi\left(z_{i}\right)$ as a critical point (this follows from Lemma 4.4(i)), and that the expanding property $\left|D T_{\phi\left(z_{i}\right)}^{j}(\mathbf{v})\right|>K_{0}^{-1}$ passes to a disk of radius $>d_{\mathcal{C}}\left(z_{i}\right)$ (Lemma 3.2). To achieve (iii), we need to show that if $z_{j}$ is a return for $i<j<i+\hat{\ell}_{i}$, then $\hat{\ell}_{j}<K \alpha\left(\log \frac{1}{b}\right)^{-1} \hat{\ell}_{i}$ (for which we follow the proof of Lemma 4.6).

Algorithm for choosing $\ell_{i}$ in (A3). Let $\hat{\ell}_{i}$ be as above. First we set $\ell_{i}^{\prime}=$ $\max \left\{2, \hat{\ell}_{i}\right\}$, then increase $\ell_{i}^{\prime}$ to $\ell_{i}^{*}$ if necessary so that the intervals $I_{i}=\left[i, i+\ell_{i}^{*}\right)$ are nested, and finally, for convenience, let $\ell_{i}=\ell_{i}^{*}+1$ or 2 to ensure that no splitting period ends at a return or at the step immediately after a return.
B. Correct alignment implies correct splitting. For $z_{0} \in \Gamma_{\theta N}$, we let $w_{i}^{*}, i=1,2, \cdots$ be generated by the splitting algorithm in Section 3.8 using the $\ell_{i}$ above. Our next lemma connects the "correct alignment" assumption in (A3) to hypothesis (a) in Lemma 3.10. Suppose $z_{i} \in \mathcal{C}^{(1)}$ and write $w_{i}^{*}=$ $A_{i} e_{\ell_{i}}(S)+B_{i} \mathbf{v}$ where $S=S\left(\mathbf{v}, w_{i}^{*}\right)$.

LEMMA 4.7. If $w_{i}^{*}$ is correctly aligned with $\varepsilon$-error where $\varepsilon \ll K_{1}^{-1}$, then

$$
\frac{\left|B_{i}\right|}{\left|A_{i}\right|}>\frac{1}{2} K_{1}^{-1} d_{\mathcal{C}}\left(z_{i}\right)
$$

where $K_{1}^{-1}$ is the lower bound of $\left|\frac{d}{d s} \eta\right|$ as in Lemma 3.7.
When the conclusion of Lemma 4.7 holds, we say $w_{i}^{*}$ splits correctly. We caution that when $d_{\mathcal{C}}(\cdot)$ is very small, $b$-horizontal vectors do not necessarily split correctly.

Corollary 4.1. If at all returns, $w_{i}^{*}$ is correctly aligned with $\varepsilon$-error where $\varepsilon \ll K_{1}^{-1}$, then $w_{i}^{*}$ is b-horizontal and splits correctly.

A proof of Lemma 4.7 is given in Appendix A.13. Corollary 4.1 follows from a direct application of Lemma 3.10 once we note that $\frac{1}{2} K_{1}^{-1} d_{\mathcal{C}}\left(z_{i}\right) \gg b^{\frac{\ell_{i}}{2}}$.

## 5. Properties of orbits controlled by the critical set

We continue to assume (A1)-(A5). This section contains a general discussion of the extent to which the orbits of $z_{0} \in \Gamma_{\theta N}$ can be used to guide other (noncritical) orbits, or, put differently, the extent to which $\left(\xi_{0}, w_{0}\right)$ for arbitrary $\xi_{0} \in R_{1}$ and $w_{0} \in X_{\xi_{0}}$ can be controlled by $\Gamma_{\theta N}$. The word control is given a formal definition in Section 5.2.

### 5.1. Copying segments of critical orbits

For $z, \xi$ in the same component of $\mathcal{C}^{(1)}$, let $p(z, \xi)$ be as defined as in Section 4.3B; i.e., it is the smallest $j$ such that $\left|T^{j} z-T^{j} \xi\right|>e^{-\beta j}$. For $z_{0}=z_{0}^{*}\left(Q^{(k)}\right) \in$ $\Gamma_{\theta N}$ and $\xi_{0}, \xi_{0}^{\prime} \in Q^{(1)}\left(z_{0}\right)$, we let $p\left(z_{0} ; \xi_{0}, \xi_{0}^{\prime}\right):=\min \left\{p\left(z_{0}, \xi_{0}\right), p\left(z_{0}, \xi_{0}^{\prime}\right), k \theta^{-1}\right\}$. Unlike (A5), we do not presuppose here any geometric relationship between $\xi_{0}$, $\xi_{0}^{\prime}$ and $z_{0}$. In particular, $p\left(z_{0} ; \xi_{0}, \xi_{0}^{\prime}\right)$ may not be in the time range for which (A5) is applicable.

Let $w_{0}\left(\xi_{0}\right)=w_{0}\left(\xi_{0}^{\prime}\right)=w_{0}\left(z_{0}\right)=\mathbf{v}$. We apply the splitting algorithm to $z_{0}, \xi_{0}$ and $\xi_{0}^{\prime}$ for $i \leq p\left(z_{0} ; \xi_{0}, \xi_{0}^{\prime}\right)$ using for all three points the splitting periods for $z_{0}$ as specified in Section 4.4. Our next proposition compares $w_{i}^{*}\left(\xi_{0}\right)$ and $w_{i}^{*}\left(\xi_{0}^{\prime}\right)$. Let

$$
\begin{equation*}
\Delta_{n}\left(\xi_{0}, \xi_{0}^{\prime}\right):=\sum_{s=0}^{n} b^{\frac{s}{4}} 2^{\ell_{n-s}}\left|\xi_{n-s}-\xi_{n-s}^{\prime}\right| \tag{5}
\end{equation*}
$$

where $\ell_{n-s}$ is the length of the longest splitting period $z_{n-s}$ find itself in, 0 if $z_{n-s}$ is out of all splitting periods.

Proposition 5.1. There is a constant $K_{1}$ such that for all $\xi_{0}, \xi_{0}^{\prime}$ and $z_{0}$ as above and $i<p\left(z_{0} ; \xi_{0}, \xi_{0}^{\prime}\right)$,

$$
\begin{equation*}
\frac{\left|w_{i}^{*}\left(\xi_{0}\right)\right|}{\left|w_{i}^{*}\left(\xi_{0}^{\prime}\right)\right|}, \quad \frac{\left|w_{i}^{*}\left(\xi_{0}^{\prime}\right)\right|}{\left|w_{i}^{*}\left(\xi_{0}\right)\right|} \leq \exp \left\{K_{1} \sum_{n=1}^{i-1} \frac{\Delta_{n}\left(\xi_{0}, \xi_{0}^{\prime}\right)}{d_{\mathcal{C}}\left(z_{n}\right)}\right\} \tag{6}
\end{equation*}
$$

and

$$
\begin{equation*}
\left.\angle\left(w_{i}^{*}\left(\xi_{0}\right), w_{i}^{*}\left(\xi_{0}^{\prime}\right)\right)\right) \leq b^{\frac{1}{2}} \Delta_{i-1}\left(\xi_{0}, \xi_{0}^{\prime}\right) \tag{7}
\end{equation*}
$$

This proposition would not be very useful without a priori bounds for the quantities involved. We explain how a bound for the right side of equations (6) and (7) can be arranged.

Lemma 5.1. Assume that (i) $\beta$ is sufficiently large compared to $\alpha$, (ii) $\delta$ is sufficiently small depending on $\alpha$ and $\beta$, and (iii) $b$ is small enough. Then for all $z_{0}, \xi_{0}, \xi_{0}^{\prime}, i$ and $n$ as above,

$$
\Delta_{n}<2 e^{-\frac{1}{2} \beta n} \ll \varepsilon_{0} d_{\mathcal{C}}\left(z_{n}\right)
$$

and

$$
K_{1} \sum_{n=1}^{i-1} \frac{\Delta_{n}\left(\xi_{0}, \xi_{0}^{\prime}\right)}{d_{\mathcal{C}}\left(z_{n}\right)} \ll 1
$$

Proposition 5.1 and Lemma 5.1 are proved in Appendix A.14. Our first application of Proposition 5.1 is to the case where $\xi_{0}^{\prime}=z_{0}$. We assume $\alpha, \beta, \delta$ and $b$ are chosen so that the following is an immediate corollary of Proposition 5.1 and Lemma 5.1.

Corollary 5.1. Let $z_{0} \in \Gamma_{\theta N}$ be of generation $k$. Then for $\xi_{0} \in Q^{(1)}\left(z_{0}\right)$ and $i<\min \left\{k \theta^{-1}, p\left(z_{0}, \xi_{0}\right)\right\}$,
(i) $\left|w_{i}^{*}\left(\xi_{0}\right)\right|>\frac{1}{4} c_{2} e^{\lambda i}$;
(ii) at return to $\mathcal{C}^{(1)}$, $w_{i}^{*}\left(\xi_{0}\right)$ is correctly aligned with $2 \varepsilon_{0}$-error.

### 5.2. A formal notion of "control"

Very roughly, a controlled orbit is one obtained by splicing together a finite number of orbit segments each one of which is either free or bound to a critical orbit. The goal of this subsection is to identify sufficient conditions at the joints that will guarantee that the resulting orbit has desirable properties.

Let $\xi_{0} \in R_{1}$ be an arbitrary point.

Definition 5.1. We say $\xi_{0}$ is controlled by $\Gamma_{\theta N}$ for $M$ iterates, or equivalently, the orbit segment $\xi_{0}, \xi_{1}, \cdots, \xi_{M-1}$ is controlled by $\Gamma_{\theta N}$, if the following hold: whenever $\xi_{i} \in \mathcal{C}^{(1)}, 0 \leq i<M$, there exists $Q^{(k)}, k \leq \theta N$, such that
(i) $\xi_{i} \in Q^{(k)} \backslash B^{(k)}$, and
(ii) $\min \left(\hat{p}\left(z_{0}, \xi_{i}\right), M-i\right) \leq k \theta^{-1}$ where $z_{0}=z_{0}^{*}\left(Q^{(k)}\right)$.

Condition (i) guarantees that (A5) applies to $\xi_{i}$. Condition (ii) guarantees that the guiding orbit $z_{0}$ remains active until either the bound period or the period of control expires.

Orbits controlled by $\Gamma_{\theta N}$ can be seen as follows: Let $n_{1} \geq 0$ be the first time $\xi_{i} \in \mathcal{C}^{(1)}$. For $i \leq n_{1}$, we regard $\xi_{i}$ as free. At time $n_{1}$, we assume there exists $z_{0} \in \Gamma_{\theta N}$ satisfying the conditions in Definition 5.1. Such a critical point is usually not unique. We make an arbitrary choice, call it $\tilde{\phi}\left(\xi_{n_{1}}\right)$, and define $\tilde{d}_{\mathcal{C}}\left(\xi_{n_{1}}\right):=\left|\tilde{\phi}\left(\xi_{n_{1}}\right)-\xi_{n_{1}}\right|$. From Lemma 4.1 we see that among the admissible choices $\tilde{\phi}\left(\xi_{n_{1}}\right), \tilde{d}_{\mathcal{C}}\left(\xi_{n_{1}}\right)$ do not differ substantially. Instead of $\tilde{\phi}(\cdot)$ and $\tilde{d}_{\mathcal{C}}(\cdot)$, we write $\phi(\cdot)$ and $d_{\mathcal{C}}(\cdot)$ for notational simplicity.

For the next $\hat{p}\left(\xi_{n_{1}}, \phi\left(\xi_{n_{1}}\right)\right)$ iterates, we think of $\xi_{n_{1}}$ as bound to $\phi\left(\xi_{n_{1}}\right)$ as in Section 5.1, inheriting from the orbit of $\phi\left(\xi_{n_{1}}\right)$ bound and splitting periods. At the end of the $\hat{p}\left(\xi_{n_{1}}, \phi\left(\xi_{n_{1}}\right)\right)$ iterates, there may be some bound periods that have not expired. In the interest of a nested structure for bound periods, we extend $\hat{p}\left(\xi_{n_{1}}, \phi\left(\xi_{n_{1}}\right)\right)$ to $p_{1}$, so that $n_{1}+p_{1}$ is the first moment when all bound periods initiated before $n_{1}+p_{1}$ have expired. For the same reason as in the proof of Lemma 4.6, we have $p_{1}<\left(1+\frac{6}{\lambda} \alpha\right) \hat{p}\left(\xi_{n_{1}}, \phi\left(\xi_{n_{1}}\right)\right)$. (This uses condition (i) in Definition 5.1.)

We regard $\xi_{n_{1}+p_{1}}$ as "free", and think of its orbit as remaining free until $n_{2}$, the first time $\geq n_{1}+p_{1}$ when $\xi_{n_{2}} \in \mathcal{C}^{(1)}$. For a controlled orbit, we are guaranteed the existence of at least one critical point satisfying the conditions of Definition 5.1 with respect to $\xi_{n_{2}}$. We think of $\xi_{n_{2}}$ as bound to $\phi\left(\xi_{n_{2}}\right)$ for $p_{2}$ iterates, and so on.

The process continues until the period of control expires. Splitting periods with a nested structure are defined similarly.

Next we discuss what it means for a ( $\xi_{0}, w_{0}$ )-pair to be controlled. Let $\varepsilon_{1}$ be such that $4 \varepsilon_{0}<\varepsilon_{1} \ll K_{1}^{-1}$ where $\varepsilon_{0}$ and $K_{1}$ are as in Section 4.1. Let $\xi_{0}$ be a controlled orbit, and let $w_{0} \in X_{\xi_{0}}$ be an arbitrary unit vector. The vectors $w_{i}^{*}\left(\xi_{0}\right)$ are obtained by using the splitting periods defined above.

Definition 5.2. We say $\left(\xi_{0}, w_{0}\right)$ is controlled by $\Gamma_{\theta N}$ for $M$ iterates, or equivalently, the sequence $\left(\xi_{0}, w_{0}\right), \cdots,\left(\xi_{M-1}, w_{M-1}\right)$ is controlled by $\Gamma_{\theta N}$, if $\xi_{0}$ is controlled for $M$ iterates and the following holds: whenever $\xi_{i} \in \mathcal{C}^{(1)}$, $0 \leq i<M, w_{i}^{*}$ is correctly aligned with $\varepsilon_{1}$-error, i.e. if $\phi\left(\xi_{i}\right)$ is of generation $j$ and $d_{\mathcal{C}}\left(\xi_{i}\right)$ is as above, then $\angle\left(w_{i}^{*}\left(\xi_{0}\right), \tau\right)<\varepsilon_{1} d_{\mathcal{C}}\left(\xi_{i}\right)$ where $\tau$ is the tangent to the leaf of $\mathcal{F}_{j}$ through $\xi_{i}$.

A slightly expanded definition. It is convenient to expand the definition of control to allow the following initial condition: If $\xi_{0} \in \mathcal{C}^{(1)}$ and $w_{0}=\mathbf{v}$, then the conditions in Definitions 5.1 and 5.2 are waived at time 0. (The rationale for this inclusion is that since no splitting occurs at time 0 , derivative recovery is automatic.)

The properties of a controlled $\left(\xi_{0}, w_{0}\right)$-pair can be summarized as follows:
Proposition 5.2. Assume that $\left(\xi_{0}, w_{0}\right)$ is controlled by $\Gamma_{\theta N}$ for $M$ iterates. Then
(1) there exist $0 \leq n_{1}<n_{1}+p_{1} \leq n_{2}<n_{2}+p_{2} \leq n_{3} \cdots<M$ such that for each $i$,
(i) there is $\phi\left(\xi_{n_{i}}\right) \in \Gamma_{\theta N}$ to which $\xi_{n_{i}}$ is bound for $p_{i}$ iterates, $p_{i} \sim$ $\log \frac{1}{d_{\mathcal{C}}\left(\xi_{n_{i}}\right)}$;
(ii) $\xi_{j} \notin \mathcal{C}^{(1)}$ for $n_{i}+p_{i} \leq j<n_{i+1}$;
(2) $w_{i}^{*}$ has the following growth properties:

$$
\frac{\left|w_{n_{i}+p_{i}}^{*}\right|}{\left|w_{n_{i}}^{*}\right|}>K^{-1} e^{\frac{1}{3} \lambda p_{i}} ; \quad \frac{\left|w_{n_{i+1}}^{*}\right|}{\left|w_{n_{i}+p_{i}}^{*}\right|}>\frac{1}{2} c_{2} e^{\frac{1}{4} \lambda_{0}\left(n_{i+1}-\left(n_{i}+p_{i}\right)\right)} .
$$

Proof. (1) is a summary of the discussion following Definition 5.1; the estimate for $p_{i}$ uses (A5)(i). The second inequality in (2) follows immediately from Lemma 3.5. The first is proved as follows: By Proposition 5.1, we have $\left|D T_{\xi_{n_{i}}}^{p_{i}}(\mathbf{v})\right|>\frac{1}{2}\left|D T_{\phi\left(\xi_{n_{i}}\right)}^{p_{i}}(\mathbf{v})\right|$. For purposes of this proof, it is simplest to split off a vector from $w_{n_{i}}^{*}$ that is known to contract for $p_{i}$ iterates. Let $e=e_{p_{i}}$ be the most contracted direction for $D T_{\xi_{n_{i}}}^{p_{i}}$ on $S=S\left(w_{n_{i}}, \mathbf{v}\right)$. We claim that if $w_{n_{i}}^{*}=$ $A e+B \mathbf{v}$, then $|B|>K^{-1} d_{\mathcal{C}}\left(\xi_{n_{i}}\right)$. (Reason: correct splitting is assumed at time $n_{i}$; the (normal) splitting period, $\ell$, is $\ll p_{i}$; and so $\angle\left(e, e_{\ell}\right)<(K b)^{\ell}$, which is $\ll$ $b^{\frac{\ell}{3}} \approx d_{\mathcal{C}}\left(\xi_{n_{1}}\right)$.) (A5)(ii) then gives $\left|D T_{\xi_{n_{i}}}^{p_{i}}(B \mathbf{v})\right|>K^{-1}\left|D T_{\phi\left(\xi_{n_{i}}\right)}^{p_{i}}(\mathbf{v})\right| d_{\mathcal{C}}\left(\xi_{n_{i}}\right)>$ $K^{-1} e^{\frac{1}{3} \lambda p_{i}}$. The addition of $D T_{\xi_{n_{i}}}^{p_{i}}(A e)$ has negligible effect.

In Section 2.2, we proved that for a class of "good" 1D maps, every orbit not passing through the critical set has the properties in Proposition 5.2. A consequence of the definition of control, therefore, is that $\left(\xi_{0}, w_{0}\right)$-pairs have 1D behavior.

### 5.3. A collection of useful facts

We record in this subsection a miscellaneous collection of facts related to controlled orbits that are used in the future. Lemmas 5.2-5.6 are proved in Appendices A.15-A.17. Proposition 5.3 is proved in Appendix A.18.
A. Relation between $\left|w_{i}\right|$ and $\left|w_{i}^{*}\right|$.

Lemma 5.2. Assume that $\left(\xi_{0}, w_{0}\right)$ is controlled by $\Gamma_{\theta N}$ for $M$ iterates. Under the additional assumption that $d_{\mathcal{C}}\left(\xi_{i}\right)>\min \left(\delta, e^{-\alpha i}\right)$ for all $i<M$,

$$
\begin{equation*}
K^{-\varepsilon i}\left|w_{i}^{*}\left(\xi_{0}\right)\right| \leq w_{i}\left(\xi_{0}\right) \leq K^{\varepsilon i} e^{2 \alpha i}\left|w_{i}^{*}\left(\xi_{0}\right)\right|, \quad \varepsilon=K \alpha \theta \tag{8}
\end{equation*}
$$

## B. Angles at bound returns.

Lemma 5.3. Let $\xi_{0}$ be controlled by $\Gamma_{\theta N}$ for $M$ iterates, and assume that at all free returns, $w_{i}^{*}$ is correctly aligned with $<\varepsilon_{1}$-error. Then at all bound returns, $w_{i}^{*}$ is correctly aligned with $<3 \varepsilon_{0}$-error.

Since $\varepsilon_{1}$, the error in alignment of $w_{i}^{*}$ at free returns, can be $\gg 3 \varepsilon_{0}$, Lemma 5.3 implies that the magnitudes of the errors at free returns are not reflected in the angles at returns during ensuing bound periods provided they are within an acceptable range.
C. Growth of $\left|w_{i}\right|,\left|w_{i}^{*}\right|$ and $\left\|D T^{i}\right\|$. The next three results provide more detailed information on derivative growth than Proposition 5.2.

Lemma 5.4. There exists $\lambda^{\prime}$ slightly smaller than $\frac{1}{3} \lambda$ such that if $\left(\xi_{0}, w_{0}\right)$ is controlled by $\Gamma_{\theta N}$ for $M$ iterates, then for every $0 \leq k<n<M$,

$$
\left|w_{n}^{*}\right| \geq K^{-1} d_{\mathcal{C}}\left(\xi_{j}\right) e^{\lambda^{\prime}(n-k)}\left|w_{k}^{*}\right|
$$

where $j$ is the first time $\geq k$ when a bound period extending beyond time $n$ is initiated. If no such $j$ exists, set $d_{\mathcal{C}}\left(\xi_{j}\right)=\delta$.

Lemma 5.5. The setting and notation are as in Lemma 5.4. If in addition $\xi_{n}$ is free, then

$$
\left|w_{n}\right|>\delta K^{-K \theta(n-k)} e^{\lambda^{\prime}(n-k)}\left|w_{k}\right| .
$$

If $\xi_{n}$ is a free return, then $\delta$ on the right side can be omitted.
We finish by recording a technical lemma that will be used in Part III.
Lemma 5.6. Suppose $\left(\xi_{0}, w_{0}\right)$ is controlled for $M$ iterates by $\Gamma_{\theta N}$, and that $d_{\mathcal{C}}\left(\xi_{i}\right)>e^{-\alpha i}$ for all $i \leq M$. Then there exist constants $K$ and $\hat{\lambda}>0$ slightly smaller than $\frac{1}{2} \lambda^{\prime}$ such that for every $0 \leq s<i<M$,

$$
\left\|D T_{\xi_{s}}^{i-s}\right\| \leq K e^{-\hat{\lambda} s}\left|w_{i}\right|
$$

D. Quadratic properties of turns. We consider in this paragraph the special situation where the critical point on a $C^{2}(b)$-curve is controlled. The quadratic distance formula in Proposition 5.3 is used to prove estimates of the kind in (A5).

The precise setting is as follows: Let $\gamma \subset \mathcal{C}^{(1)}$ be a $C^{2}(b)$-curve, and let $z_{0} \in \gamma$ be a critical point of order $M$ on $\gamma$ in the sense of Definition 3.2. (There is no restriction on the size of $M$; it can be $>N$.) We assume that
(1) $\left(z_{0}, \mathbf{v}\right)$ is controlled by $\Gamma_{\theta N}$ for $M$ iterates; and
(2) $d_{\mathcal{C}}\left(z_{i}\right)>\min \left(\delta, e^{-\alpha i}\right)$ for all $0<i \leq M$.

Let $s \mapsto \xi_{0}(s)$ be the parametrization of $\gamma$ by arclength with $\xi_{0}(0)=z_{0}$.

Proposition 5.3. For given $s_{1}>0$, let $p\left(s_{1}\right)=\min \left\{p\left(\xi_{0}\left(s_{1}\right), z_{0}\right), M\right\}$. Then for all $0<s \leq s_{1}$ and $i \in\left[\ell(s), p\left(s_{1}\right)\right]$ with $\ell(s)=2 \frac{\log s}{\log b}$,

$$
\left|\xi_{i}(s)-z_{i}\right| \approx \frac{1}{2}\left|\frac{d}{d s} e_{1}(0)\right|\left|w_{i}(0)\right| s^{2}
$$

where $e_{1}=e_{1}(S)$ and $S=S\left(\gamma^{\prime}, \mathbf{v}\right)$.

## 6. Identification of hyperbolic behavior: Formal inductive procedure

### 6.1. Global constants (mostly review)

For $N=1,2, \cdots$, we define below a set of "good" maps $T: X \rightarrow X$ denoted by

$$
\mathcal{G}_{N}=\mathcal{G}_{N}\left(f_{0}, K_{0}, a, b ; \lambda, \alpha ; \delta, \beta, \varepsilon_{0}, \theta\right) .
$$

The arguments on the right side can be understood conceptually as follows:

1. The first group consists of $f_{0} \in \mathcal{M}$ and three constants, $K_{0}, a$ and $b$. These items appear in the Standing Hypotheses at the beginning of Part II; they define an open set in the space of $C^{3}$ embeddings of $X$ into itself.
2. In the next group are $\lambda$ and $\alpha$, two constants that appear in (A2) and (A4). As we will see, (A2) and (A4) play a special role in determining if $T$ in the open set above is in $\mathcal{G}_{N}$; they are analogous to (G1) and (G2) for 1D maps (see §2.2).
3. Unlike the situation in 1D, auxiliary constructions are needed before we are able to properly formulate (A2) and (A4). The constants in the last group, namely $\delta, \beta, \varepsilon_{0}$ and $\theta$, appear in these auxiliary constructions. They do not directly impact whether a map is in $\mathcal{G}_{N}$, but help maintain uniform estimates in the constructions.

In the definition of $\mathcal{G}_{N}, f_{0}$ is chosen first; it can be any element of $\mathcal{M}$. We then fix $K_{0}$, which can be any number $>\left\|f_{0}\right\|_{C^{3}}$. Precise conditions imposed on the rest of the constants are given in the text. We review below their (rough) meanings and give the order in which they are chosen. To ensure consistency in our choices, it is important that (i) only upper bounds are imposed on each constant, and (ii) these bounds are allowed to depend only on the constants higher up on the list (in addition to $f_{0}, K_{0}$, and $m$, the dimension of $X$ ). Except for $\lambda$, all the constants listed below are $\ll 1$ and must be taken to be as small as necessary.

Important constants, their meanings, and the order in which they are chosen.
$-\lambda$ is our targeted Lyapunov exponent; it can be anything $<\frac{1}{5} \lambda_{0}$ where $\lambda_{0}$ is a growth rate of $\left|f_{0}^{\prime}\right|$ (see Definition 1.1). Once chosen, it is fixed throughout.

- Next we fix $\alpha$ and $\beta$ and think of $e^{-\alpha n}$ and $e^{-\beta n}$ as representing two small scales. The requirements are that $0<\alpha, \beta \ll \min \{\lambda, 1\}$ and $\beta>K \alpha$ for some $K$ depending on $f_{0}$ and $K_{0}$. The meaning of $\alpha$ is that critical orbits are not allowed to approach the critical set at speeds faster than $e^{-\alpha n}$. Two orbits $\left\{z_{i}\right\}$ and $\left\{z_{i}^{\prime}\right\}$ with $\left|z_{i}-z_{i}^{\prime}\right|<e^{-\beta i}$ are to be thought of as "bound together".
- $\varepsilon_{0}$, which depends only on $f_{0}, K_{0}$ and $m$, has the following meaning: For $z \in \mathcal{C}^{(1)}$, vectors $v \in X_{z}$ that make angles $<\varepsilon_{0} d_{\mathcal{C}}(z)$ with certain $\mathcal{F}_{k}$-leaves are viewed as "correctly aligned".
- The size of $\delta$ is limited by many factors: examples of which include $\delta<\delta_{0}$ where $\delta_{0}$ is as in Definition 1.1, a bound used in distortion (Lemma 5.1), the Taylor formula estimate at "turns" (Proposition 5.3), $\delta \ll \varepsilon_{0}$, and some purely numerical inequalities (e.g. if $\delta=e^{-\mu}$, then $\frac{1}{\mu^{2}} \ll e^{-\mu}$ ).
- Chosen last are $a$ and $b$. It is best to think of $a$ and $b$ as very small numbers that we may need to decrease a finite number of times as we go along.
- The smaller $a$ is, the longer $f^{n}(\hat{x}), \hat{x} \in C$, can be kept out of $C_{\delta_{0}}$.
- The smaller $b$ is, the more closely $T$ mimics $F^{\#}$.
- There is an important constant defined at the same time as $b$, namely $\theta:=\frac{K}{\log \frac{1}{b}}$ where $K$ is chosen so that $b^{\theta}=\|D T\|^{-20}$. With this choice of $\theta$, critical orbits emanating from the same $B^{(k)}$ can be viewed as a single orbit for $k \theta^{-1}$ iterates. We may, therefore, regard the number of critical orbits (or "critical blobs") present at time $N$ as $\leq K^{\theta N}$ for some $K$ depending on $f_{0}$.

When referring to $\mathcal{G}_{N}$ in the future, it will be understood that the arguments above are implicit. In particular, $\mathcal{G}_{0}$ is the set of maps $T$ satisfying the conditions at the beginning of Part II with regard to some fixed $f_{0}, K_{0}, a$ and $b$. Constants (such as $K_{1}$ ) not on this list are regarded as local in context; they must be specified each time they are used. Finally, we emphasize that the generic constant $K$ that appears in many of our results is allowed only to depend on $f_{0}, K_{0}, m$ and $\lambda$ provided that the other constants are appropriately small.

### 6.2. Three stages of evolution

Our construction of $\mathcal{G}_{N}$ comes in three distinct stages: For $N \leq \theta^{-1}$, the situation is, in many ways, not far from that in 1D. This part is simple and is disposed of immediately in the next paragraph. At time $N=\theta^{-1}$, certain local complexities of higher dimensional maps begin to develop, "turns" play a
more prominent role, and the definition of $\mathcal{G}_{N}$ becomes necessarily inductive. We have been building up the dynamical picture for this part in Sections 4 and 5 and will complete its construction in Sections 6.3 and 6.4. At $N=\theta^{-2}$, the global structure of $T$ begins to depart from those of 1D maps. New ideas are needed; they are discussed in Sections 7 and 8.

Getting started: the first $\theta^{-1}$ steps. Let $T \in \mathcal{G}_{0}$, and assume the leaves of $\mathcal{F}_{1}$ are parallel to the $x$-axis. Let $C=\left\{\hat{x}_{1}, \cdots, \hat{x}_{q}\right\}$ be the critical set of $f$. Then near each $\hat{x}_{i}$, there exists $\tilde{x}_{i}$ such that $e_{1}(S)=\partial_{x}$ at $\left(\tilde{x}_{i}, 0\right)$ where $S=S\left(\partial_{x}, \mathbf{v}\right)$. A simple computation gives $\left|\tilde{x}_{i}-\hat{x}_{i}\right|<K b$. Let $\Gamma_{1}=\left\{\left(\tilde{x}_{1}, 0\right), \cdots,\left(\tilde{x}_{q}, 0\right)\right\}$. These are the only critical points for the first $\theta^{-1}$ iterates.

Before proceeding further, we observe that if $\gamma_{0}$ is a $C^{2}(b)$-segment with the property that $\gamma_{i}:=T^{i}\left(\gamma_{0}\right)$ does not meet $B^{(1)}$ for all $i<n$, then the curves $\gamma_{i}$ are roughly horizontal for all $i \leq n$. This follows immediately from the fact that for $z$ with $d_{\mathcal{C}}(z)>b^{\frac{1}{5}}$ and $u \in X_{z}$ with $s(u)<b^{\frac{3}{4}}, s\left(D T_{z}(u)\right)<b^{\frac{3}{4}}$ (see §3.5).

For $N=1,2, \cdots, \theta^{-1}$, let

$$
\mathcal{G}_{N}=\left\{T \in \mathcal{G}_{0} \mid(\mathrm{A} 2) \text { and (A4) hold for all } z_{0} \in \Gamma_{1} \text { and } i \leq N\right\} .
$$

(A2) and (A4) are, as noted earlier, analogs of (G1) and (G2) in Section 2.2.
We claim that for $T \in \mathcal{G}_{N}$, (A3) and (A5) are satisfied automatically. (A3) is easily verified since all $b$-horizontal vectors are correctly aligned at $d_{\mathcal{C}}>e^{-\alpha \theta^{-1}}>b^{\frac{1}{5}}\left(b^{\theta}<e^{-5 \alpha}\right.$ by definition, so $\left.b^{\frac{1}{5}} \ll e^{-\alpha \theta^{-1}}\right)$. (A5) follows from 1D estimates: Let $\gamma_{0}$ be the curve joining $\xi_{0} \in Q^{(1)} \backslash B^{(1)}$ to $B^{(1)}$ in (A5). Then during its bound period, all tangent vectors to $\gamma_{i}$ are roughly horizontal as explained above and curvatures are $<K b^{\frac{2}{5}}$. An argument entirely parallel to that in Appendix A. 1 proves (A5)(i)-(iii).

Inductive scheme for going from $N=\theta^{-1}$ to $N=\theta^{-2}$. Beyond $N=\theta^{-1}$, more critical points are needed as orbits emanating from $B^{(1)}$ begin to diverge. To help describe the structures needed for the identification of new critical points, we have introduced a set of assumptions, namely (A1)-(A5). In Section 6.3, we will add another one, (A6), which is also trivially satisfied up to time $\theta^{-1}$. Let

$$
\mathcal{G}_{N}:=\left\{T \in \mathcal{G}_{0} \mid(\mathrm{A} 1)(N)-(\mathrm{A} 6)(N) \text { hold, } \quad N \leq \theta^{-2} .\right.
$$

Observe that this definition is consistent with the one defined earlier for $N \leq$ $\theta^{-1}$. The goal of Sections 6.3 and 6.4 is to prove the following:
(*) Let $\theta^{-1}<N<\frac{1}{\alpha^{*}} N \leq \theta^{-2}$. Assume $T \in \mathcal{G}_{N}$, and prove that if $T$ satisfies (A2) and (A4) up to time $\frac{1}{\alpha^{*}} N$, then it is in $\mathcal{G}_{\frac{1}{\alpha^{*}} N}$.
The time step of the construction above is determined by the fact that at times $\leq \frac{1}{\alpha^{*}} N$, the lengths of the bound periods are $\leq N$. This ratio is noted in the paragraph on "three important time scales" in Section 4.1.

Why stop at $N=\theta^{-2}$ ? We emphasize that the material in this section is for iterates $N \leq \theta^{-2}$. The reason for this time restriction is that as mentioned above, the sets $T^{k} B^{(1)}$ begin to get "large" at $k=\theta^{-1}$, affecting the geometry of the critical regions. (A1) and (A6), which we will introduce shortly, cannot be sustained as formulated for $N>\theta^{-2}$.

Notation. In this section and the next, we will be working with the foliations $\mathcal{F}_{k}$. Given that we have defined $\mathcal{F}_{1}$ to be the initial foliation on $R_{1}$, it is advantageous in discussions involving $\mathcal{F}_{k}$ to let $\xi_{1}$ denote arbitrary points in $R_{1}$ and $\tau_{1}$ unit tangent vectors to the leaves of $\mathcal{F}_{1}$. This convention (instead of the usual $\left.\left(\xi_{0}, \tau_{0}\right)\right)$ leads to more pleasing notation such as $\xi_{k} \in R_{k}$ and $\tau_{k}$ as tangent vectors to the leaves of $\mathcal{F}_{k}$.

### 6.3. Controlling $\mathcal{F}_{k}$, and pushing forward (A1) and (A6)

One way to gain a better grip on the geometry of $R_{k}$ is to control $\left(\xi_{i}, \tau_{i}\right)$ for $\xi_{1} \in R_{1}$.

Rules for setting control. (1) We stop controlling ( $\xi_{i}, \tau_{i}$ ) once $\xi_{i}$ enters $B^{(i)}$; this is compatible with the idea that $T^{k} B^{(i)}, k=1,2, \cdots, i \theta^{-1}$, is to be seen as the orbit of a single point.
(2) In our inductive scheme to be detailed shortly, the control of $\left(\xi_{i}, \tau_{i}\right)$ proceeds in parallel with the construction of $\Gamma_{i}$. For this reason, we will take $\phi\left(\xi_{i}\right) \in \Gamma_{i}$.
(3) As explained in Section 5.2, it suffices to set control at free returns. Let $i$ be a free return. Then $\phi\left(\xi_{i}\right)$ is chosen as follows: If there exists $j<i$ such that $\xi_{i} \in \mathcal{C}^{(j)} \backslash \mathcal{C}^{(j+1)}$, then we let $\phi\left(\xi_{i}\right)=z_{0}^{*}\left(Q^{(j)}\right)$ where $Q^{(j)}=Q^{(j)}\left(\xi_{i}\right)$. If $\xi_{i} \in Q^{(i)}$, we have no choice (in view of (2)) but to let $\phi\left(\xi_{i}\right)=z_{0}^{*}\left(Q^{(i)}\right)$.

To the five assumptions (A1)-(A5) in Section 4.1, we now add another one. We say the foliation $\mathcal{F}_{k+1}$ is controlled on $R_{k+1}$ by $\Gamma_{k}$ if for all $\xi_{1} \in R_{1}$ and $i \leq k,\left(\xi_{i}, \tau_{i}\right)$ is controlled by $\Gamma_{k}$ provided $\xi_{i} \notin B^{(i)}$ for all $i \leq k$. (The indices in the last sentence are intended as written: we say $\mathcal{F}_{k+1}$ is controlled because control of ( $\xi_{i}, \tau_{i}$ ) for $i \leq k$ leads to geometric knowledge of the leaves of $\mathcal{F}_{k+1}$.)
(A6) ( $N$ ) For all $k \leq \theta N, \mathcal{F}_{k+1}$ is controlled on $R_{k+1}$ by $\Gamma_{k}$.
At this point we would like to assert that (A1) $\left(\frac{1}{\alpha^{*}} N\right)$ and (A6) $\left(\frac{1}{\alpha^{*}} N\right)$ hold for $T \in \mathcal{G}_{N}$. A proof would involve simultaneously constructing $\mathcal{C}^{(k)}$ and $\Gamma_{k}$, and using $\Gamma_{k}$ to control $\mathcal{F}_{k+1}$. What prevents us from making a clean statement to this effect at this time is that without having first assumed or proved (A2) $\left(k \theta^{-1}\right)-(\mathrm{A} 5)\left(k \theta^{-1}\right)$ for $k>\theta N$, we cannot, in principle, conclude that orbits controlled by $\Gamma_{k}$ have the properties in Section 5.3.

We examine the situation more closely: Assume $T$ satisfies (A1)( $N$ )$(\mathrm{A} 6)(N)$. Fix $\theta N<i \leq \frac{1}{\alpha^{*}} \theta N$, and assume (A1) $\left(i \theta^{-1}\right)$ holds. Let $\xi_{i} \in \mathcal{C}^{(1)}$ be an arbitrary point. We define $\phi\left(\xi_{i}\right)$ as in (3) above and assume that $\tau_{i}$ is correctly aligned (with respect to $\mathcal{F}_{j}$ where $j$ is the generation of $\phi\left(\xi_{i}\right)$ ). The discussion below pertains only to time $\leq \frac{1}{\alpha^{*}} \theta N$.

Case 1. $j \leq \theta N$. In this case, $\left(\xi_{i}, \tau_{i}\right)$ is controlled by $\Gamma_{\theta N}$ for the next $\min \left(p, \frac{1}{\alpha^{*}} \theta N-i\right)$ iterates where $p$ is the bound period between $\xi_{i}$ and $\phi\left(\xi_{i}\right)$.

Case 2. $j>\theta N$, and $\xi_{i} \notin B^{(\theta N)}$. The conclusion is as in Case 1. The orbit of $\hat{z}_{0}:=\phi\left(\xi_{i}\right)$ and that of $z_{0}=z_{0}^{*}\left(B^{(\theta N)}\left(\phi\left(\xi_{i}\right)\right)\right)$ remain extremely close during the period in question (more precisely, $\left|\hat{z}_{k}-z_{k}\right|<\|D T\|^{k} b^{\frac{\theta N}{5}} \ll e^{-\beta k}$ ), and it makes no difference whether we view $\xi_{i}$ as bound to $\hat{z}_{0}$ or to $z_{0}$.

Case 3. $j>\theta N$ and $\xi_{i} \in B^{(\theta N)}$. The estimate in the last paragraph shows that $\xi_{i}$ is bound to $\phi\left(\xi_{i}\right)$ - and to $z_{0}^{*}\left(B^{(\theta N)}\left(\phi\left(\xi_{i}\right)\right)\right.$ - through time $\frac{1}{\alpha^{*}} \theta N$. From Proposition 5.1, we know that $e_{\ell}$ is well defined on all of $B^{(\theta N)}$ for all $\ell \leq N$, and by our correct alignment assumption together with Lemma 4.7, $\tau_{i}$ splits correctly. The evolution of the $\mathbf{v}$-component at $\xi_{i}$ can then be compared to that at $z_{0}^{*}\left(B^{(\theta N)}\right)$ by Proposition 5.1 and Lemma 5.3.

The discussion above tells us that in the control $\left(\xi_{1}, \tau_{1}\right),\left(\xi_{2}, \tau_{2}\right), \cdots$ up to time $\frac{1}{\alpha^{*}} \theta N$, the only role played by $\Gamma_{k} \backslash \Gamma_{\theta N}$ and $\mathcal{F}_{k}$ for $k>\theta N$ is to determine the correctness of alignment and subsequent splitting period at free returns. The rest of the control is really provided by $\Gamma_{\theta N}$. We have argued that Lemmas 5.2-5.6 apply up to time $\frac{1}{\alpha^{*}} \theta N$. Nevertheless, to distinguish between the present situation and that after we have conferred (A2)-(A5) upon $\Gamma_{k}$, we will say, if correct alignment holds for all $i \leq k$, that the sequence $\left(\xi_{1}, \tau_{1}\right), \cdots,\left(\xi_{k}, \tau_{k}\right)$ is provisionally controlled by $\Gamma_{k}$.

We now state the main result of this subsection.
Proposition 6.1. Let $\theta^{-1} \leq N<\frac{1}{\alpha^{*}} N \leq \theta^{-2}$, and assume $T$ satisfies (A1) $(N)-(\operatorname{A} 6)(N)$. Then for $\theta N<k \leq \frac{1}{\alpha^{*}} \theta N$ :
$(\mathrm{a})_{k} \mathcal{C}^{(k)}$ and $\Gamma_{k}$ with the properties in (A1) can be constructed;
(b) ${ }_{k}$ if $\xi_{1} \in R_{1}$ is such that $\xi_{i} \notin B^{(i)}$ for all $i \leq k$, the sequence $\left(\xi_{1}, \tau_{1}\right), \cdots\left(\xi_{k}, \tau_{k}\right)$ is provisionally controlled by $\Gamma_{k}$.

Proof. We assume (a) ${ }_{i}$ and (b) ${ }_{i}$ for all $i<k$.
Proof of (a) $)_{k}$. Noting that it makes sense to speak about those segments of $\mathcal{F}_{k}$-leaves that are provisionally controlled as being in a bound or free state, we begin with the following result of independent interest:

Lemma 6.1. Let $\gamma$ be a leaf of $\mathcal{F}_{k}$. If every $\xi_{k} \in \gamma$ is free, then $\gamma$ is a $C^{2}(b)$-curve.

Proof. That $\tau_{k}$ is $b$-horizontal follows from Corollary 4.1. As for curvature, we appeal to Lemma 3.3 after using Lemma 5.5 to establish that $\left|\tau_{k}\right| \geq \delta K^{-K \theta(k-i)}\left|\tau_{i}\right|$ for all $i<k$.

Let $\gamma$ be a leaf segment of $\mathcal{F}_{k}$ meeting $\mathcal{C}^{(k-1)}$. We claim that it is contained in a maximal free segment that traverses the entire length of $Q^{(k-1)}$, extending as a $C^{2}(b)$-curve by $>\frac{1}{2} e^{-\alpha k}$ on both sides. To see this, let $\xi_{k} \in R_{k}$ be such that $d_{\mathcal{C}}\left(\xi_{k}\right)<\frac{1}{2} e^{-\alpha k}$, and suppose it is not free. Then there are only two possibilities: (1) For some $i<k, \xi_{i} \in B^{(i)}$ and we stopped controlling its orbit, or (2) $\xi_{i}$ is controlled for all $i<k$, and $\xi_{k}$ is in a bound period initiated at some time $i<k$. (1) is not feasible, for if we let $z_{0}=z_{0}^{*}\left(B^{(i)}\right)$, then $d_{\mathcal{C}}\left(z_{k-i}\right)>e^{-\alpha(k-i)}$, and $\operatorname{diam}\left(T^{k-i} B^{(i)}\right) \ll e^{-\alpha(k-i)}$ (here we need $k<\theta^{-1}$ ), contradicting $d_{\mathcal{C}}\left(\xi_{k}\right)<\frac{1}{2} e^{-\alpha k}$. (2) is also impossible, for if we let $z_{0}=\phi\left(\xi_{i}\right)$, then $d_{\mathcal{C}}\left(z_{k-i}\right)>e^{-\alpha(k-i)}$ while $\left|\xi_{k}-z_{k-i}\right|<e^{-\beta(k-i)} \ll e^{-\alpha(k-i)}$.

We have proved that $R_{k} \cap Q^{(k-1)}$, if non-empty, is the union of a collection of horizontal sections $\{H\}$. In each $H$, we arbitrarily pick an $\mathcal{F}_{k}$-leaf $\gamma$. The critical point $z_{0}^{*}\left(Q^{(k-1)}\right)$ constructed in step $(\text { a })_{k-1}$ induces a critical point of order $k-1$ on $\gamma$ (Lemma 4.3 and 3.8). By Lemma 3.9, this critical point can be upgraded to one of order $k$. We make it an element of $\Gamma_{k}$, and construct a $Q^{(k)}$ of length $\min \left(2 \delta, e^{-\lambda k}\right)$ centered at it. Doing this for every horizontal section $H$ that passes through every $Q^{(k-1)}$ completes the construction of $\mathcal{C}^{(k)}$ and $\Gamma_{k}$.

It follows directly from the next lemma that the sectional diameter of $Q^{(k)}$ is $<b^{\frac{k}{2}}$.

Lemma 6.2. Every $\xi_{k} \in Q^{(k)}$ is contained in a codimension one manifold $W$ with the property that
(i) $W$ meets every connected component of $\mathcal{F}_{k}$-leaf in $Q^{(k)}$ in exactly one point;
(ii) for all $\xi_{1}, \xi_{1}^{\prime} \in T^{-(k-1)} W,\left|\xi_{i}-\xi_{i}^{\prime}\right|<b^{\frac{i}{2}}$ for all $i \leq k$.

Lemma 6.2 is proved in Appendix A. 19.
Proof of $(\mathrm{b})_{k}$. As noted earlier, it suffices to consider the case where $\xi_{k}$ is a free return, and it suffices to show correct alignment of $\tau_{k}$ at $\xi_{k}$. Let $\gamma$ be the maximal free segment of $\mathcal{F}_{k}$-leaf containing $\xi_{k}$. Then the endpoints of $\gamma$ are in a bound state, and so are outside of $\mathcal{C}^{(k)}$. This leaves two possibilities for the relation between $\gamma$ and $\mathcal{C}^{(k)}$.

Case 1. $\gamma$ passes through the entire length of some $Q^{(k)}$. We consider $Q^{(k)} \subset Q^{(k-1)} \subset Q^{(k-2)} \subset \cdots$ until we reach the first $Q^{(j)}$ that contains $\xi_{k}$. Since $\xi_{k} \in\left(Q^{(j)} \cap R_{j+1}\right) \backslash \mathcal{C}^{(j+1)}$, $d_{\mathcal{C}}\left(\xi_{k}\right)>e^{-\lambda(j+1)}$. We let $\hat{\gamma}$ be the $\mathcal{F}_{j}$ leaf
through $\xi_{k}$, and apply Lemma 4.3 to obtain two points $\gamma(0)$ and $\hat{\gamma}(0)$ in $\gamma$ and $\hat{\gamma}$ respectively with
(i) $\left|\xi_{k}-\gamma(0)\right| \approx\left|\xi_{k}-\hat{\gamma}(0)\right| \approx d_{\mathcal{C}}\left(\xi_{k}\right)$, and
(ii) $\angle\left(\hat{\gamma}^{\prime}(0), \gamma^{\prime}(0)\right)<K b^{\frac{j}{4}}$.

Letting $\hat{\tau}_{j}$ and $\tau_{k}$ denote the tangents to $\hat{\gamma}$ and $\gamma$ respectively at $\xi_{k}$ and using the fact that $\gamma$ and $\hat{\gamma}$ are $C^{2}(b)$-curves between the points in question, we have

$$
\begin{aligned}
\angle\left(\hat{\tau}_{j}, \tau_{k}\right) & \leq \angle\left(\hat{\tau}_{j}, \hat{\gamma}^{\prime}(0)\right)+\angle\left(\hat{\gamma}^{\prime}(0), \gamma^{\prime}(0)\right)+\angle\left(\gamma^{\prime}(0), \tau_{k}\right) \\
& <\frac{K b}{\delta^{3}} d_{\mathcal{C}}\left(\xi_{k}\right)+K b^{\frac{j}{4}}+\frac{K b}{\delta^{3}} d_{\mathcal{C}}\left(\xi_{k}\right) \\
& <b^{\frac{1}{5}} d_{\mathcal{C}}\left(\xi_{k}\right) .
\end{aligned}
$$

Case 2. $\gamma$ does not meet $\mathcal{C}^{(k)}$. We first formally treat the geometry before making the required angle estimates.

Geometry: (i) Let $j$ be the largest integer such that $\xi_{k} \in Q^{(j)}$, so that $\xi_{k} \in H \backslash Q^{(j+1)}$ where $H$ is a component of $R_{j+1} \cap Q^{(j)}$. Suppose for definiteness that $\xi_{k}$ lies in the right chamber of $H \backslash Q^{(j+1)}$. We move along $\gamma$ to the left until we reach either $\xi$, the left endpoint of $\gamma$, or the right boundary of $Q^{(j+1)}$, whichever happens first. Once $\xi$ is reached, we stop. Otherwise we continue moving through $Q^{(j+1)}$ until we reach either $\xi$ or the right boundary of $Q^{(j+2)}$. (We have used implicitly the fact that $\gamma$, which is a leaf of $\mathcal{F}_{k}$, does not meet $\partial R_{i}$ for $i<k$.) By assumption, $\xi$ is reached before we arrive at $Q^{(k)}$, so that $\xi \in Q^{\left(j^{\prime}\right)} \backslash \mathcal{C}^{\left(j^{\prime}+1\right)}$ for some $j^{\prime}$ with $k>j^{\prime} \geq j$.
(ii) We note that $\xi$ can also be regarded as in bound state, and argue now that $\phi(\xi)$ is to the left of $\xi$. More precisely, we write $\xi=\eta_{k}$, let $\eta_{i}, i<k$, be the last free return, and let $\phi\left(\eta_{i}\right)=\hat{z}_{0}$. Recalling the definition of $\phi(\cdot)$ for critical orbits ( $\S \S 4.1$ and 4.3 A ), we deduce that $\phi(\xi)=\phi\left(\hat{z}_{k-i}\right)$ is of generation $j^{\prime \prime}$ for some $j^{\prime \prime} \leq j^{\prime}$ (it can be considerably smaller), and that both $\xi$ and $\hat{z}_{k-i}$ are in $Q^{\left(j^{\prime \prime}\right)} \backslash B^{\left(j^{\prime \prime}\right)}$. To see that $\xi$ is in the right chamber of $Q^{\left(j^{\prime \prime}\right)} \backslash B^{\left(j^{\prime \prime}\right)}$, we interpolate between $Q^{\left(j^{\prime}\right)} \subset \cdots \subset Q^{\left(j^{\prime \prime}\right)}$, noting that the right chamber of each $Q^{(i)}$ does not meet the left chamber of $Q^{(i-1)}$.

Angles: Let $\hat{\gamma}$ be the leaf of $\mathcal{F}_{j}$ through $\xi_{k}$ and $\tilde{\gamma}$ the leaf of $\mathcal{F}_{j^{\prime \prime}}$ through $\xi$. We will use the following notation: $\tau_{k, \xi_{k}}$ and $\tau_{k, \xi}$ are tangents to $\gamma$ at $\xi_{k}$ and $\xi$ respectively; $\hat{\tau}_{j, \xi_{k}}$ is tangent to $\hat{\gamma}$ at $\xi_{k}$, and $\tilde{\tau}_{j^{\prime \prime}, \xi}$ is tangent to $\tilde{\gamma}$ at $\xi$. Then

$$
\angle\left(\tau_{k, \xi_{k}}, \hat{\tau}_{j, \xi_{k}}\right) \leq \angle\left(\tau_{k, \xi_{k}}, \tau_{k, \xi}\right)+\angle\left(\tau_{k, \xi}, \tilde{\tau}_{j^{\prime \prime}, \xi}\right)+\angle\left(\tilde{\tau}_{j^{\prime \prime}, \xi}, \hat{\tau}_{j, \xi_{k}}\right) .
$$

The terms above are estimated by
(i) $\angle\left(\tau_{k, \xi_{k}}, \tau_{k, \xi}\right)<\frac{K b}{\delta^{3}}\left|\xi-\xi_{k}\right|$ since $\gamma$ is free and hence $C^{2}(b)$;
(ii) $\angle\left(\tau_{k, \xi}, \tilde{\tau}_{j^{\prime \prime}, \xi}\right)<3 \varepsilon_{0}\left|\xi-\phi\left(\hat{z}_{k-i}\right)\right|$ since $\xi=\eta_{k}$ is a bound return (Lemma 5.3);
(iii) $\angle\left(\tilde{\tau}_{j^{\prime \prime}, \xi}, \hat{\tau}_{j, \xi_{k}}\right)<b^{\frac{1}{4}} \min \left(j, j^{\prime \prime}\right)+\frac{K b}{\delta^{3}}\left|\xi-\xi_{k}\right|$ from Lemma 4.2.

Also, we have argued that $\phi\left(\hat{z}_{k-i}\right)$ is to the left of $\xi$, so that $\left|\xi-\xi_{k}\right|$, $\left|\xi-\phi\left(\hat{z}_{k-i}\right)\right|<d_{\mathcal{C}}\left(\xi_{k}\right)$. These inequalities together with $d_{\mathcal{C}}\left(\xi_{k}\right) \gg b^{\frac{j}{4}}$ and $d_{\mathcal{C}}\left(\xi_{k}\right)>d_{\mathcal{C}}(\xi)>b^{\frac{j^{\prime \prime}}{5}}$ give $\angle\left(\tau_{k, \xi_{k}}, \hat{\tau}_{j, \xi_{k}}\right)<4 \varepsilon_{0} d_{\mathcal{C}}\left(\xi_{k}\right)$.

### 6.4. What it takes to go from $\mathcal{G}_{N}$ to $\mathcal{G}_{\frac{1}{\alpha^{*}} N}$

In this subsection we fix $N$ with $\theta^{-1} \leq N<\frac{1}{\alpha^{*}} N \leq \theta^{-2}$, and assume that $T \in \mathcal{G}_{N}$, i.e. $(\mathrm{A} 1)(N)-(\mathrm{A} 6)(N)$ hold for $T$. It is proved in Proposition 6.1 that without further assumptions, (A1) $\left(\frac{1}{\alpha^{*}} N\right)$ holds automatically. The purpose of this subsection is to determine what constraints we need to impose on $T$ to put it in $\mathcal{G}_{\frac{1}{\alpha^{*}} N}$.
(A2): Rate of approach to critical set. By assumption, all $z_{0} \in \Gamma_{\theta N}$ obey (A2). At issue is whether or not $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N} \backslash \Gamma_{\theta N}$ obeys (A2) as stated in Section 4.1. We distinguish between the two time intervals $[1, N]$ and $\left[N+1, \frac{1}{\alpha^{*}} N\right]$ : On $[1, N]$, each $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N} \backslash \Gamma_{\theta N}$ follows closely a critical point of generation $\theta N$ (Corollary 5.1). The two orbits do differ by a little, however, so it is possible for $z_{i}$ to violate slightly the condition in (A2). On $\left[N+1, \frac{1}{\alpha^{*}} N\right]$, there is no reason why (A2) is respected by $z_{0} \in \Gamma_{\frac{1}{\alpha^{*} \theta N}} \backslash \Gamma_{\theta N}$. We conclude that (A2) $\left(\frac{1}{\alpha^{*}} N\right)$ is a new condition that must be imposed on $T$ if it is to belong in $\mathcal{G}_{\frac{1}{\alpha^{*}} N}$.
(A3): Correct alignment. We will prove that (A3), in fact, comes for free. The mechanisms for ensuring correct alignment of $w_{i}^{*}$ at free returns and at bound returns are entirely different. At free returns, this comes from geometry, from the "rank one" character of $T$ in particular. At bound returns, it comes from copying. We emphasize that we do not deduce (A3) ( $\left.\frac{1}{\alpha^{*}} N\right)$ directly from $(\mathrm{A} 3)(N)$. We prove it from scratch, in a sense, keeping track of the increase in error each time the picture is copied.

Proposition 6.2. Let $T \in \mathcal{G}_{N}$. We assume $\Gamma_{\frac{1}{\alpha^{*}} \theta N}$ is constructed, and fix $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N}$. Also, assume the condition in (A2) is imposed on $z_{0}$ up to time $\frac{1}{\alpha^{*}} N .{ }^{5}$ Then $\left(z_{0}, \mathbf{v}\right)$ is controlled by $\Gamma_{\theta N}$ up to time $\frac{1}{\alpha^{*}} N$. In fact, we have the following stronger results:
(i) If $z_{i}$ is a free return, then $w_{i}$ is aligned correctly with error $<d_{\mathcal{C}}\left(z_{i}\right) \ll \varepsilon_{0}$.
(ii) If $z_{i}$ is a bound return, then $w_{i}^{*}$ is aligned correctly with $<\varepsilon_{0}$-error.

Proof. To establish control for the orbit of $\left(z_{0}, w_{0}\right)$ with $w_{0}=\mathbf{v}$, we define $\phi\left(z_{i}\right)$ for $i \leq \frac{1}{\alpha^{*}} N$ according to the rule in Section 4.1, i.e. $\phi\left(z_{i}\right)$ is the critical

[^5]point of highest generation $j \leq \alpha^{*} \theta i \leq \theta N$ with the property that $z_{i} \in \mathcal{C}^{(j)}$. Note that this implies $\phi\left(z_{i}\right) \in \Gamma_{\theta N} .{ }^{6}$ With (A2) imposed on $z_{0}$, Lemma 4.4 shows that the conditions in Definition 5.1 are met. Correct splitting is proved inductively as follows. Assume that $\left(z_{0}, w_{0}\right)$ is controlled for $k-1$ iterates by $\Gamma_{\theta N}$ for some $k \leq \frac{1}{\alpha^{*}} N$.
(i) $z_{k}$ is a free return. Let $j+1$ be the generation of $\phi\left(z_{k}\right)$, and let $\tau_{z_{k}, j+1}$ be tangent to $\mathcal{F}_{j+1}$ at $z_{k}$. We need to show $\angle\left(w_{k}, \tau_{z_{k}, j+1}\right) \ll \varepsilon_{0} d_{\mathcal{C}}\left(z_{k}\right)$. Let $\xi_{1}=T^{-j} z_{k}$, i.e. $z_{k}=\xi_{j+1}$. From $(\operatorname{A} 6)(N)$, we know that $\left(\xi_{1}, \tau_{1}\right)$ is controlled for $j$ iterates unless $\xi_{n} \in B^{(n)}$ for some $n \leq j$, which is impossible because that would contradict our assumption that $z_{k}$ is free $\left(j+1<\theta^{-1}\right.$ is again used here). Observe that
\[

$$
\begin{aligned}
\angle\left(w_{k}, \tau_{z_{k}, j+1}\right) & =\angle\left(D T_{\xi_{1}}^{j} w_{k-j}, D T_{\xi_{1}}^{j} \tau_{1}\right) \leq \frac{\left|D T_{\xi_{1}}^{j} w_{k-j} \wedge D T_{\xi_{1}}^{j} \tau_{1}\right|}{\left|w_{k}\right|\left|\tau_{j+1}\right|} \\
& \leq(K b)^{j} \frac{\left|w_{k-j} \wedge \tau_{1}\right|}{\left|w_{k}\right|\left|\tau_{j+1}\right|} \leq(K b)^{j} \frac{\left|w_{k-j}\right|}{\left|w_{k}\right|\left|\tau_{j+1}\right|}
\end{aligned}
$$
\]

The second inequality above comes from Sublemma A.4.2 in Appendix A.4; the rest are straightforward. To estimate the quantity in the final bound, we claim that $\frac{\left|w_{k-j}\right|}{\left|w_{k}\right|}<K e^{-\lambda^{\prime \prime} j}$ where $\lambda^{\prime \prime}$ is slightly smaller than $\frac{1}{3} \lambda$. This is because $\left(z_{0}, w_{0}\right)$ is, by inductive assumption, controlled by $\Gamma_{\theta N}$ for $k$ iterates, and $z_{k}$ being a free return, Lemma 5.5 applies. Next we claim that $\frac{1}{\left|\tau_{j+1}\right|}<K e^{-\lambda^{\prime \prime} j}$. This is again a consequence of Lemma 5.5, after we establish the following: With $j \leq \theta N,\left(\xi_{1}, \tau_{1}\right)$ is a controlled pair by $(\mathrm{A} 6)(N)$, and $\xi_{j+1}$ is a free return because the bound-free structure on the orbit segment $\xi_{1}, \xi_{2}, \cdots, \xi_{j+1}$ can be taken to be identical to that of $z_{k-j}, z_{k-j+1}, \cdots, z_{k}$ except that bound periods for $z_{i}$ initiated before time $k-j$ do not count for $\xi_{1}$. Thus $\angle\left(w_{k}, \tau_{z_{k}, j+1}\right)<$ $(K b)^{j}$.

Correct alignment at $z_{k}$ is now straightforward: If $j+1<\theta N$, then $\angle\left(w_{k}, \tau_{z_{k}, j+1}\right)<(K b)^{j} \ll e^{-2 \lambda j} \approx\left(d_{\mathcal{C}}\left(z_{k}\right)\right)^{2}$. In general, $d_{\mathcal{C}}\left(z_{k}\right)>e^{-\alpha k} \geq$ $e^{-\alpha\left(\frac{1}{\alpha^{*}} N\right)}$, so that if $j+1=\theta N$, then

$$
\angle\left(w_{k}, \tau_{z_{k}, j+1}\right)<(K b)^{j}=(K b)^{\theta N-1}<e^{-\frac{\lambda}{3} N}=e^{-\alpha\left(\frac{2}{\alpha^{*}} N\right)} \leq\left(d_{\mathcal{C}}\left(z_{k}\right)\right)^{2}
$$

This completes the proof of the free return case.
(ii) $z_{k}$ is a bound return. Consider first the following scenario: Suppose $z_{j}, j<k$, is a free return, and the bound period initiated at that time extends beyond time $k$. Let $\phi\left(z_{j}\right)=\hat{z}_{0}$, and assume $\hat{z}_{k-j}$ is a free return.

We estimate the error in alignment at $z_{k}$ as follows: Let $g$ and $\hat{g}$ be the generations of $\phi\left(z_{k}\right)$ and $\phi\left(\hat{z}_{k-j}\right)$. We claim that $g \geq \hat{g}-1$. This is because if $\hat{z}_{k-j} \in Q^{(n)}$, then $z_{k}$ is also in $Q^{(n)}$ - or it is just outside, in which case it is

[^6]in the $Q^{(n-1)}$ containing $Q^{(n)}$ (see Lemma 4.5). Now in the definition of $\phi(\cdot)$, there is an upper bound on the generation of the guiding critical orbit. Since $k>k-j$, a more stringent upper bound is imposed on $\hat{z}_{k-j}$ than on $z_{k}$, hence the assertion.

In the discussion to follow, we let $\tau_{z, n}$ denote the tangent to the $\mathcal{F}_{n}$-leaf at $z$. The angle to be estimated, $\angle\left(w_{k}^{*}\left(z_{0}\right), \tau_{z_{k}, g}\right)$, is bounded above by

$$
\begin{aligned}
\angle\left(w_{k}^{*}\left(z_{0}\right), \tau_{z_{k}, g}\right) \leq & \angle\left(w_{k}^{*}\left(z_{0}\right), w_{k-j}\left(\hat{z}_{0}\right)\right) \\
& +\angle\left(w_{k-j}\left(\hat{z}_{0}\right), \tau_{\hat{z}_{k-j}, \hat{g}}\right)+\angle\left(\tau_{\hat{z}_{k-j}}, \hat{g}, \tau_{z_{k}, g}\right)
\end{aligned}
$$

From part (i), we have that the second term on the right is $<\left(d_{\mathcal{C}}\left(z_{k}\right)\right)^{2}$ : $\hat{z}_{k-j}$ is free, and $d_{\mathcal{C}}\left(z_{k}\right) \approx d_{\mathcal{C}}\left(\hat{z}_{k-j}\right)$ from Lemma 4.5. The third term is $<K b^{\frac{1}{4}(\hat{g}-1)}+K b \delta^{-3} e^{-\beta(k-j)}$ from Lemma 4.2. To estimate the first term we write $w_{j}\left(z_{0}\right)=A e_{k-j}+B \mathbf{v}$ to obtain

$$
\begin{aligned}
\angle\left(w_{k}^{*}\left(z_{0}\right), w_{k-j}\left(\hat{z}_{0}\right)\right) & <\angle\left(w_{k-j}\left(z_{j}\right), w_{k-j}\left(\hat{z}_{0}\right)\right)+(K b)^{k-j} \\
& <\frac{1}{2} e^{-\frac{1}{2} \beta(k-j)}+(K b)^{k-j}
\end{aligned}
$$

the second inequality is from Proposition 5.1 and Lemma 5.1. Plugging $d_{\mathcal{C}}\left(z_{k}\right)>$ $b^{\frac{1}{5}(\hat{g}-1)}$ (Lemma 4.4) and $d_{\mathcal{C}}\left(z_{k}\right)>e^{-\alpha(k-j)}$ into the three estimates above, we obtain

$$
\begin{aligned}
\angle\left(w_{k}^{*}\left(z_{0}\right), \tau_{z_{k}, g}\right) & <d_{\mathcal{C}}\left(z_{k}\right)\left\{d_{\mathcal{C}}\left(z_{k}\right)+K e^{\left(\frac{1}{2} \beta-\alpha\right)(k-j)}+(K b)^{k-j}+K b^{\frac{1}{20}(\hat{g}-1)}\right\} \\
& <\varepsilon_{0} d_{\mathcal{C}}\left(z_{k}\right)
\end{aligned}
$$

In general, there exist $j_{1}<j_{2}<\cdots<j_{n}<k$ and $\hat{z}_{0}^{(1)}, \cdots, \hat{z}_{0}^{(n)} \in \Gamma_{\theta N}$ such that
$-z_{j_{1}}$ is the last free return before time $k$, with $\phi\left(z_{j_{1}}\right)=\hat{z}_{0}^{(1)} ;$
$-\hat{z}_{k-j_{1}}^{(1)}$ is not a free return; its last free return is at time $j_{2}$, with $\phi\left(\hat{z}_{j_{2}-j_{1}}^{(1)}\right)$ $=\hat{z}_{0}^{(2)}$;
$-\hat{z}_{k-j_{2}}^{(2)}$ is not a free return; its last free return is at time $j_{3}$, and so on, until

- finally, $\hat{z}_{k-j_{n}}^{(n)}$ is a free return.

Considerations similar to those above show that as we go through the different layers of bindings, the errors in alignment form a geometric series which add up to $\leq e^{-\frac{1}{2} \beta\left(k-j_{n}\right)}+(K b)^{k-j_{n}}$.
(A4): Growth of $\left|w_{i}^{*}\right|$.
Corollary 6.1. Under the hypothesis of Proposition 6.2, for the critical orbit $z_{0}$ in question,

$$
\frac{\left|w_{i}^{*}\left(z_{0}\right)\right|}{\left|w_{N}^{*}\left(z_{0}\right)\right|}>K^{-1} e^{\left(\frac{1}{3} \lambda-2 \alpha\right)(i-N)}, \quad N<i \leq \frac{1}{\alpha^{*}} N
$$

This corollary, which follows from the control of $\left(z_{0}, \mathbf{v}\right)$ proved in Proposition 6.2, the condition in (A2), and Lemma 5.4, is in the direction of maintaining (A4) $\left(\frac{1}{\alpha^{*}} N\right)$ but the gain in Lyapunov exponent is only about $\frac{1}{3}$ of what is needed for that purpose. Indeed, (A4) is not a self-perpetuating condition: the exponent in Corollary 6.1, when copied again in future inductive steps, may lead to a downward spiral in the Lyapunov exponent along critical orbits.

We conclude that (A4) ( $\frac{1}{\alpha^{*}} N$ ) must be imposed (by external means) to ensure that $T \in \mathcal{G}_{\frac{1}{\alpha^{*}}} N$.
(A5): Quadratic turns, lengths of bound periods and derivative recovery.
Proposition 6.3. Let $T \in \mathcal{G}_{N}$ be such that (A2) ( $\left.\frac{1}{\alpha^{*}} N\right)$ and (A4) $\left(\frac{1}{\alpha^{*}} N\right)$ hold. Then (A5) ( $\left.\frac{1}{\alpha^{*}} N\right)$ holds automatically.

Proof. Let $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N}$ be fixed. Suppose $z_{0}=z_{0}^{*}\left(Q^{(k)}\right)$, and let $\xi_{0} \in$ $Q^{(k)} \backslash B^{(k)}$. For the definitions of and relation between $\hat{p}\left(z_{0}, \xi_{0}\right)$ and $p\left(z_{0}, \xi_{0}\right)$, see Section 4.3B.

Proof of (A5)(iii). Let $\gamma$ be the $\mathcal{F}_{k}$-leaf containing $\xi_{0}$ in $Q^{(k)}$. Then there exists $\tilde{z}_{0} \in \gamma \cap B^{(k)}$ such that $\tilde{z}_{0}$ is a critical point of order $p=$ $\min \left\{p\left(z_{0}, \xi_{0}\right), k \theta^{-1}\right\}$ on $\gamma$. For all practical purposes, $z_{0}$ and $\tilde{z}_{0}$ are indistinguishable for $p$ iterates, so we may regard $\tilde{z}_{0}$ as satisfying the hypotheses of Proposition 5.3 (which $z_{0}$ has been shown to satisfy). Proposition 5.3 then gives the desired result (with $\tilde{z}_{0}$ instead of $z_{0}$ ).

Proof of (A5)(i). For the lower bound, we have, for all $j \leq \frac{h}{3 \log \|D T\|}$,

$$
\left|\xi_{j}-z_{j}\right|<\|D T\|^{j}\left|\xi_{0}-z_{0}\right|<e^{-\frac{2 h}{3}} \ll e^{-\hat{\beta} \frac{h}{3 \log \|D T\|}},
$$

proving $\hat{p}\left(z_{0}, \xi_{0}\right) \geq \frac{h}{3 \log \|D T\|}$.
For the upper bound, by Proposition 5.3, $p=p\left(z_{0}, \xi_{0}\right)$ is the smallest integer $i$ such that $\left|w_{i}(0)\right| \cdot\left|z_{0}-\xi_{0}\right|^{2}>K_{1}^{-1} e^{-\beta i}$. We claim that $p \leq \frac{3 h}{\lambda}$. If not, then

$$
\left|w_{\frac{3 h}{\lambda}}\left(z_{0}\right)\right| \cdot\left|z_{0}-\xi_{0}\right|^{2}>K^{-\varepsilon \frac{3 h}{\lambda}}\left|w_{\frac{3 h}{\lambda}}^{*}\left(z_{0}\right)\right| \cdot\left|z_{0}-\xi_{0}\right|^{2}>K^{-\varepsilon \frac{3 h}{\lambda}} e^{\lambda \cdot \frac{3 h}{\lambda}} e^{-2 h}>1 .
$$

Lemma 5.2 is used in the first inequality above.
Proof of (A5)(ii). First we consider $\hat{p}=\hat{p}\left(z_{0}, \xi_{0}\right)$, for which we have $\left|w_{\hat{p}}\left(\tilde{z}_{0}\right)\right| \approx\left|w_{\hat{p}}\left(z_{0}\right)\right|($ Corollary 5.1), and

$$
\begin{aligned}
\left|z_{0}-\xi_{0}\right|\left|w_{\hat{p}}\left(z_{0}\right)\right| & =\left(\left|z_{0}-\xi_{0}\right|\left|w_{\hat{p}}\left(z_{0}\right)\right|^{\frac{1}{2}}\right) \cdot\left|w_{\hat{p}}\left(z_{0}\right)\right|^{\frac{1}{2}} \\
& >K^{-1}\left|z_{\hat{p}}-\xi_{\hat{p}}\right| \cdot\left|w_{\hat{p}}\left(z_{0}\right)\right|^{\frac{1}{2}} \quad \text { by (A5)(iii) } \\
& >K^{-1} e^{-\frac{\hat{\beta}}{2} \hat{p}} \cdot e^{\frac{\lambda \hat{p}}{2}} .
\end{aligned}
$$

Now let $p \in\left[\hat{p}, \hat{p}\left(1+\frac{9}{\lambda} \alpha\right)\right]$. From Lemmas 5.2 and 5.4, $\left|w_{p}\left(z_{0}\right)\right| \geq e^{-4 \alpha p}\left|w_{\hat{p}}\left(z_{0}\right)\right|$, and so

$$
\left|z_{0}-\xi_{0}\right|\left|w_{p}\left(z_{0}\right)\right|>K^{-1} e^{-4 \alpha p} e^{-\frac{\hat{\beta}}{2} \hat{p}} \cdot e^{\frac{\lambda \hat{p}}{2}}>e^{\frac{1}{3} \lambda p}
$$

(A6): Control of foliations. With (A2) and (A4) assumed and (A3) and (A5) proved up to time $\frac{1}{\alpha^{*}} N$, the provisional control proved in Proposition 6.1 is, by definition, upgraded to control in the usual sense.

Summary. (*) in Section 6.2 is proved.

## 7. Global geometry via monotone branches

The purpose of this section is to introduce the main geometric ideas needed to construct $T \in \mathcal{G}_{n}$ beyond $N=\theta^{-2}$, and to reformulate (A1) and (A6) to accommodate these new geometric structures.

### 7.1. Introduction

The idea of studying piecewise monotonic 1D maps via their monotone branches has been used many times. We attempt in this section to introduce a corresponding notion for $T$. For small $n$, it is easy to see that $R_{n}$ is the union of sets that are tubular neighborhoods of 1D monotone curve segments. These should be, by any definition, monotone branches of $T$. For $n \leq \theta^{-1}$, we have seen that $R_{n}$ is punctuated by (tiny) sections that are $T^{i}$-images of $B^{(n-i)}$, i.e. the "critical blobs" of Section 4.2. Intuitively, these sets are located at sharp "turns"; they divide $R_{n}$ into connected components that are comparatively "straight". Leaving precise definitions for later, we think of these components as monotone branches of generation $n$.

The picture in the last paragraph cannot be maintained indefinitely, however, for it relies on the fact that critical blobs are very small compared to their distances to the critical set. As $n$ increases, it is inevitable that the images of $B^{(k)}$ will grow large, making it impossible to keep them away from the critical regions. See Figure 2. As explained in Section 6.2, the significance of time $N=\theta^{-2}$ is that at time $\theta^{-2}$, the geometry of $R_{\theta^{-1}}$ becomes relevant, and $\theta^{-1}$ is the time beyond which we cannot guarantee the smallness of the images of $B^{(1)}$.


Figure 2. The images of $B^{(k)}$ cannot avoid critical regions forever

To avoid dealing with the situation depicted in Figure 2, we declared in Section 4.2 that $B^{(k)}$ ceases to be active after $k \theta^{-1}$ iterates. Once $B^{(k)}$ ceases to be active, we must "discontinue", i.e. stop considering, the monotone branches that end in $T^{k \theta^{-1}} B^{(k)}$. We do not wish to relinquish control completely of the region occupied by a discontinued branch, however, for it is likely to contain part of the attractor.

Central to our scheme is the idea of branch replacement. We will prove that all monotone branches that are discontinued can be systematically replaced by branches of higher generations, so that at every step $n$, there is a collection of "good" branches of generations $\sim n$ that together account for all parts of the attractor. This replacement procedure is discussed at the end of Section 8, after we make precise the notion of a monotone branch and integrate these new geometric ideas into the dynamical picture described in Sections 4, 5 and 6 .

A 2D version of monotone branches was introduced in [WY1]. They were not used, however, in the inductive construction of the dynamical picture.

### 7.2. Formal definitions and assumptions

The idea of monotone branches is inseparable from that of critical regions. Any definition necessarily assumes that certain relevant critical structures have been identified. Likewise, the identification of these critical structures relies on the idea of monotone branches. Definition 7.1 below is how we have elected to enter this inductive cycle.

We say $A_{1} \cup A_{2} \cup A_{3}$ are contiguous sections of $R_{k}$ if (i) each $A_{i}$ is a section of $R_{k}$, and (ii) if $\Phi_{i}:[-1,1] \times D_{m-1} \rightarrow R_{k}$ are the defining maps for $A_{i}$ (see the definitions immediately preceding Part II), then for $i=1,2$, $\Phi_{i}\left(\{1\} \times D_{m-1}\right)=\Phi_{i+1}\left(\{-1\} \times D_{m-1}\right)$.

Definition 7.1. Let $T \in \mathcal{G}_{0}$ and $n \geq 2$. Suppose that for each $k<n$, a collection of sections $\left\{B^{(k)}\right\}$ of $R_{k}$ has been identified. Then a monotone branch $M$ of generation $n$ for $T$ is a section of $R_{n}$ that is the union of three contiguous sections $E \cup M^{\circ} \cup E^{\prime}$ with the following properties:
(a) There exist $i, i^{\prime}<n$ and $\hat{B}^{(n-i)}, \hat{B}^{\left(n-i^{\prime}\right)}$ such that
(i) $T^{-i}(E)=\hat{B}^{(n-i)}$ and $i \leq(n-i) \theta^{-1}$;
(ii) $T^{-i^{\prime}}\left(E^{\prime}\right)=\hat{B}^{\left(n-i^{\prime}\right)}$ and $i^{\prime} \leq\left(n-i^{\prime}\right) \theta^{-1}$;
(b) for all $i<n, T^{-i}\left(M^{\circ}\right) \cap B^{(n-i)}=\emptyset$ for all $B^{(n-i)}$.
$E$ and $E^{\prime}$ are called the ends of the monotone branch $M$, and $M^{\circ}$ is its main body. We say the end $E$ has reached the end of its period of activity if $i=(n-i) \theta^{-1}$. For $T \in \mathcal{G}_{\theta^{-2}}$ and $k \leq \theta^{-1},\left\{B^{(k)}\right\}$ is as in Section 4.1. Thus we know from Sections 4-6 that monotone branches of generation $n \leq \theta^{-1}+1$ are well defined. In this time range, every $R_{n}$ is the union of a finite number of
monotone branches, with adjacent branches overlapping in $T^{i} B^{(n-i)}$ for some $B^{(n-i)}$.

Remark. In the case that $I$ is an interval, there are two special branches of generation $n$ for all $n \geq 2$ with the property that one of their two ends is $T^{n-1}\left(V_{i}\right)$. Here $V_{1}$ and $V_{2}$ are the two "vertical" components in the boundary of $\hat{R}_{1}$; see Section 3.9.

Tree of monotone branches. Associated with $T \in \mathcal{G}_{\theta^{-2}}$ is a combinatorial object $\cup_{1 \leq k \leq \theta^{-1}} \mathcal{T}_{k}$ defined as follows: We declare $R_{1}$ to be the unique monotone branch of generation 1 (even though it has no ends), and let $\mathcal{T}_{1}=\left\{R_{1}\right\}$. In general, $\mathcal{T}_{k}$ consists of a collection of monotone branches of generation $k$. Let $M \in \mathcal{T}_{k}$ for some $k<\theta^{-1}$. Here is how $M$ reproduces: By construction, $M$ either does not intersect any of the $B^{(k)}$, or it contains in its main body a finite number of them, say $B_{1}, B_{2}, \cdots, B_{s}$, in that order. In the first case, $T(M) \in \mathcal{T}_{k+1}$. In the second, $T(M)$ is the union of $s+1$ elements of $\mathcal{T}_{k+1}$, the main bodies of which connect $T(E)$ to $T\left(B_{1}\right), T\left(B_{i}\right)$ to $T\left(B_{i+1}\right)$, and $T\left(B_{s}\right)$ to $T\left(E^{\prime}\right)$. This construction defines a finite tree with $\theta^{-1}$ levels.

Suppose for $n \geq \theta^{-1}$ the tree $\cup_{1 \leq k \leq n} \mathcal{T}_{k}$ is defined, i.e. each $\mathcal{T}_{k}$ consists of a collection of monotone branches of generation $k$ and $\mathcal{T}_{k}$ and $\mathcal{T}_{k+1}$ are related as above. We assume further that the critical regions $Q^{(n)}$ and $B^{(n)}$ have been identified. Then we may extend the tree to level $n+1$ : Consider one $M \in \mathcal{T}_{n}$ at a time. First we check to see if either one of the two ends $E$ and $E^{\prime}$ of $M$ has reached the end of its period of activity. If so, the branch $M$ is "discontinued"; i.e., we do not iterate it further. If not, then $M$ reproduces as in the last paragraph, and $\mathcal{T}_{n+1}$ consists of all the offspring so obtained as $M$ ranges over $\mathcal{T}_{n}$. (Note that no branch is discontinued for $n<\theta^{-1}$.)

Provided that the relevant monotone branches and $\left\{B^{(k)}\right\}$ are well defined and are related in the manner described above, one can extend $\mathcal{T}_{n}$ indefinitely and obtain, as $n \rightarrow \infty$, an infinite tree $\mathcal{T}:=\cup_{k \geq 1} \mathcal{T}_{k}$.

We turn next to the inductive construction of $\mathcal{T}_{k}$ and $\mathcal{C}^{(k)}$. From the previous discussion, it is clear that the construction of these two objects must proceed hand in hand. Moreover, after time $N=\theta^{-2}$, due to the discontinuation of certain branches, the structure of critical regions becomes more complex, and (A1) and (A6) have to be modified accordingly. In (A1') and ( $\mathrm{A} 6^{\prime}$ ) below, we try to give as complete a geometric description of these structures as possible, without seeking to present a minimal set of conditions.
$\left(\mathrm{A1}^{\prime}\right)(N)$ Critical regions. For $1 \leq k \leq \theta N$, there are sets $\mathcal{C}^{(k)}$ called critical regions with the following properties:
(I) Geometric structure. $\mathcal{C}^{(1)}$ is as defined after Section 3.4. For $k \leq \theta N$, $\mathcal{C}^{(k)}$ has a finite number of connected components $\left\{Q^{(k)}\right\}$ each one of which is a horizontal section of $R_{k}$.
(a) Relations among different $Q^{(k)}$ :
(i) For $Q^{(k)}$ and $Q^{\left(k^{\prime}\right)}$ with $k<k^{\prime}$, either $Q^{\left(k^{\prime}\right)} \subset Q^{(k)}$ or $Q^{(k)} \cap Q^{\left(k^{\prime}\right)}=\emptyset$. This defines a partial ordering on the set $\left\{Q^{(k)}, 1 \leq k \leq \theta N\right\}$ with $Q<Q^{\prime}$ if and only if $Q^{\prime} \subset Q$.
(ii) If $Q^{\left(k_{1}\right)}<\cdots<Q^{\left(k_{n}\right)}$ is a maximal chain, ${ }^{7}$ then $k_{i+1} \leq k_{i}(1+2 \theta)$.
(b) Properties of individual $Q^{(k)}$ :
(i) $Q^{(k)}$ has length $\min \left(2 \delta, 2 e^{-\lambda k}\right)$ and cross-sectional diameter $<b^{\frac{k}{2}}$.
(ii) Exactly halfway between the two ends of $Q^{(k)}$ a point $z_{0}=z_{0}^{*}\left(Q^{(k)}\right)$ is singled out; $z_{0}$ is a critical point of order $k$ with respect to the leaf of the foliation $\mathcal{F}_{k}$ containing it, and $B^{(k)}:=\left\{z \in Q^{(k)}:\left|z-z_{0}\right|_{h}<b^{\frac{k}{5}}\right\}$.
(II) Construction and relation to $\mathcal{T}_{k}$. Let $M \in \mathcal{T}_{k}, k \leq \theta N$, and let $Q$ be a component of $\mathcal{C}^{(j)},(1+2 \theta)^{-1} k \leq j<k$. Then either $M \cap Q=\emptyset$ or it is the union of a finite number of horizontal sections each one of which extends $>\frac{1}{2} e^{-\alpha k}$ on both sides of $Q$. If $M \cap Q \neq \emptyset$, then each connected component $H$ of $M \cap Q$ contains a unique $Q^{(k)}$, which is located in roughly the middle of $H$ in terms of the $x$-coordinate. All $Q^{(k)}$ are constructed this way.
$\left(\mathrm{A}^{\prime}\right)(N)$ Monotone branches. $\cup_{0 \leq k \leq \theta N} \mathcal{T}_{k+1}$ is defined with the following properties:
(I) Construction and relation to $\mathcal{C}^{(k)}$. For each $M \in \mathcal{T}_{k}, 1 \leq k \leq \theta N$, either $M$ does not meet any $B^{(k)}$, or it contains in its main body a finite number of them, and it reproduces as described in the paragraph on "Tree of monotone branches".
(II) Dynamical control. For $M \in \cup_{0 \leq k \leq \theta N} \mathcal{T}_{k+1}, \mathcal{F}_{k+1}$ is controlled on $M^{\circ}$ by $\Gamma_{k}$.
(III) Relation to $R_{k}$. For $k \leq \theta N(1+2 \theta)^{-1}$,

$$
R_{k(1+2 \theta)} \subset \cup\left\{M, M \in \cup_{k \leq j<k(1+2 \theta)} \mathcal{T}_{j}\right\} .
$$

As before, we call $z_{0}^{*}\left(Q^{(k)}\right)$ a critical point of generation $k$, and let $\Gamma_{k}$ denote the set of critical points of generation $\leq k$.

We are finally in a position to give the definition of $\mathcal{G}_{N}$ that is valid for all $N \in \mathbb{Z}^{+}$:

$$
\mathcal{G}_{N}:=\left\{T \in \mathcal{G}_{0} \mid\left(\mathrm{A}^{\prime}\right)(N),(\mathrm{A} 2)(N)-(\mathrm{A} 5)(N), \text { and }\left(\mathrm{A}^{\prime}\right)(N) \text { hold }\right\} .
$$

[^7]The goal for the remainder of Part II, then, is to prove the following, which is a more general statement than $\left({ }^{*}\right)$ in Section 6.2:
(**) For all $N \geq \theta^{-1}$, if $T \in \mathcal{G}_{N}$ satisfies (A2) and (A4) up to time $\frac{1}{\alpha^{*}} N$, then it is in $\mathcal{G}_{\frac{1}{\alpha^{*}} N}$.

### 7.3. Clarification and implications

We have tried to capture in ( $\mathrm{A} 1^{\prime}$ ) and ( $\mathrm{A} 6^{\prime}$ ) a relatively concise summary of the geometric structures that appear after generation $\theta^{-1}$. Before embarking on a formal proof of $(* *)$, we would like to take this subsection to elaborate on the implications of these statements and to highlight those features that are new.

In the discussion below, $T$ is assumed to be in $\mathcal{G}_{N}$.
(1) Neighborhoods of attractors. Our attractor $\Omega$ is defined to be $\Omega=$ $\cap_{k \geq 1} R_{k}$. For $k \leq \theta^{-1}$, it is natural to see $R_{k}$ as an approximation of $\Omega$ with "finite geometry". Beyond $k=\theta^{-1}$, we are forced to choose between $R_{k}$, the geometry of which becomes increasingly complex, and something with simpler geometry. We opted for the latter. The set $\cup_{M \in \mathcal{T}_{k}} M$, is in general not a good approximation of $R_{k}$; in particular, $\cup_{M \in \mathcal{T}_{k}} M \not \supset \Omega$. On the other hand, $\left(\mathrm{A} 6^{\prime}\right)(\mathrm{III})$ tells us that for all $j, \cup\left\{M, M \in \cup_{j \leq k<j(1+2 \theta)} \mathcal{T}_{k}\right\} \supset \Omega$. That is to say, while any one level of the tree $\mathcal{T}$ may not be adequate, sets that are unions of branches from $\sim j \theta$ levels starting with level $j$ are bona fide neighborhoods of $\Omega$ with finite geometry.
(2) Structure of critical regions. (a) The structure of $\mathcal{C}^{(k)}$ in ( $\mathrm{Al}^{\prime}$ ) is not as complete as that in (A1). First, these regions are no longer nested; i.e. it is not necessarily the case that $\mathcal{C}^{(k+1)} \subset \mathcal{C}^{(k)}$. From its construction in ( $\mathrm{A}^{\prime}$ )(II), however, it follows that

$$
\mathcal{C}^{(k)} \subset \bigcup_{(1+2 \theta)^{-1} k \leq j<k} \mathcal{C}^{(j)}
$$

(b) The following structure inside each component $Q$ of $\mathcal{C}^{(k)}$ is used many times in the analysis to follow: ${ }^{8}$

- With regard to the partial order in $\left(\mathrm{A1}^{\prime}\right)(\mathrm{I})(\mathrm{a})$, there exist components of $\mathcal{C}^{(j)}, j>k$, which lie immediately below $Q$, i.e. if we call these components $Q_{i}$, then $Q_{i}>Q$ and there is no other $Q^{\prime}$ with $Q_{i}>Q^{\prime}>Q$.
- By $\left(\mathrm{A} 1^{\prime}\right)(\mathrm{I})(\mathrm{a})(\mathrm{ii}), Q_{i}$ is of generation $k_{i}$ with $k<k_{i} \leq k(1+2 \theta)$. We remark that for $k \leq \theta^{-1}, k_{i}=k+1$. For $k>\theta^{-1}$, this is not necessarily the case. The phenomenon described here will be referred to as the

[^8]"skipping of generations"; it introduces a number of technical problems (that will be addressed in §8.1).

- By $\left(\mathrm{A} 1^{\prime}\right)(\mathrm{II})$, each $Q_{i}$ is contained in a horizontal section $H_{i}$ that stretches across $Q$, extending considerably beyond, and
- if $k \leq \theta N(1+2 \theta)^{-1}$, then by $\left(\mathrm{A}^{\prime}\right)(\mathrm{III}),\left(Q \cap R_{k(1+2 \theta)}\right) \subset \cup_{i} H_{i}$.

We caution that there may be points $z \in Q^{(k)} \cap R_{k+1}$ that are not in $\cup_{i} H_{i}$.
(3) Eventual set of critical points. At the end of our inductive construction, there is a set $\mathcal{C}$ defined by

$$
\mathcal{C}=\lim _{k \rightarrow \infty} \Gamma_{k} \quad \text { or, equivalently, } \quad \mathcal{C}=\cap_{n>0} \cup_{k \geq n} \mathcal{C}^{(k)}
$$

$\mathcal{C}$ is the set of critical points for $T \in \mathcal{G}$. Note that all orbits of $z_{0} \in \mathcal{C}$ satisfy (A2) and (A4).
(4) Critical blobs and geometry of monotone branches. (a) In the notation of Definition 7.1, $E$ and $E^{\prime}$ are precisely what we called critical blobs in Section 4.2. The requirements that $n-i \leq i \theta^{-1}$ and $n-i^{\prime} \leq i^{\prime} \theta^{-1}$ are equivalent to discontinuing $M$ as soon as one of $z_{0}^{*}\left(B^{(i)}\right)$ or $z_{0}^{*}\left(B^{\left(i^{\prime}\right)}\right)$ ceases to be active.
(b) We state a result which together with Lemma 6.1 reinforces our mental picture of what a monotone branch should be: either $M$ is relatively small (such as when all or most of it is in bound state), or it consists of a relatively long horizontal section, namely the part that is free, connecting two relatively small pieces at the ends consisting of points that are in bound state.

Lemma 7.1. Let $T \in \mathcal{G}_{N}, N \geq \theta^{-1}$. Then for $M \in \mathcal{T}_{k}, 1 \leq k \leq \theta N$, the set

$$
\left\{\xi_{k+1} \in M: \xi_{k+1} \text { is free }\right\}
$$

if nonempty, is a connected set omitting a neighborhood of $E \cup E^{\prime}$ in $M$.
Corollary 7.1 below is a direct consequence of Lemma 7.1. A bound of this type is needed in the treatment of parameter issues in Part III.

Corollary 7.1. There exist $K_{1}, K_{2}$ depending only on $f_{0}$ such that for all $T \in \mathcal{G}_{N}$, the following hold for $k \leq \theta N$ :
(i) $M \in \mathcal{T}_{k}$ has at most $K_{1}$ children;
(ii) $\mathcal{C}^{(k)}$ has at most $K_{2}^{k}$ connected components.

Lemma 7.1 and Corollary 7.1 are proved in Appendix A.20.
We note again that (A1') and (A6 $6^{\prime}$ ) are consistent with (A1) and (A6) for $N \leq \theta^{-2}$. This is because no monotone branches are discontinued before time $N=\theta^{-2}$.

## 8. Completion of induction

### 8.1. Preparation: Sections 4, 5 and 6 revisited

Assume $T \in \mathcal{G}_{N}$. We begin by bringing to the foreground how the new geometry introduced in Section 7 affects the statements and/or proofs in Sections 4,5 and 6 .

1. Angles between leaves of $\mathcal{F}_{k}$ for different $k$ (Lemma 4.2): The statement of Lemma 4.2 is unchanged. Its proof, as stated, compares the leaves of $\mathcal{F}_{k}$ and $\mathcal{F}_{k^{\prime}}$ by going through the leaves of $\mathcal{F}_{j}$ for all intermediate $j$. Since not all of the $Q^{(j)}$ are present, the argument needs to be modified slightly: Replace $k+i$ by $k_{i}, i=0,1, \cdots, n$, where $Q^{(k)}=Q^{\left(k_{0}\right)} \supset Q^{\left(k_{1}\right)} \supset \cdots \supset Q^{\left(k_{n}\right)}=Q^{(\hat{k})}$ are the critical regions present. Since $k_{i+1} \leq k_{i}(1+2 \theta)$, the proof goes through as is.
2. Distances between critical points (Lemma 4.1): The statement of Lemma 4.1 is unchanged. In the proof, which estimates $\left|z_{0}^{*}\left(Q^{(k)}\right)-z_{0}^{*}\left(Q^{(k+1)}\right)\right|$, replace $k+1$ by $k^{\prime}$ where $k^{\prime}$ is the generation of the next $Q^{(i)}$ inside $Q^{(k)}$. To use Lemma 3.8 to induce a new critical point it suffices to have $\hat{\gamma}$ traverse $B^{(k)}$. This holds easily because $k^{\prime} \leq k(1+2 \theta)$. The order of the newly induced critical point is then updated as before.
3. "Reproduction" of critical blobs (last paragraph of Section 4.2): Let $Q$ and $Q_{i}$ be as in paragraph (2)(b) in Section 7.3, and let $B^{(k)}$ and $B^{\left(k_{i}\right)}$ be associated with $Q$ and $Q_{i}$ respectively. Then at time $k \theta^{-1}, T^{k \theta^{-1}}\left(B^{(k)}\right)$ is replaced by $\left\{T^{k \theta^{-1}}\left(B^{\left(k_{1}\right)}\right), \cdots, T^{k \theta^{-1}}\left(B^{\left(k_{s}\right)}\right)\right\}$; i.e., in the absence of in-between generations, some of the critical blobs $T^{k \theta^{-1}}\left(B^{\left(k_{j}\right)}\right)$ may be born a little earlier than before.
4. Existence of suitable $\phi(\cdot)$ for critical orbits (Lemma 4.4): This is a genuine concern, since fewer critical regions and therefore fewer critical points are available. Both the definition of $\phi\left(z_{i}\right)$ and the statement of Lemma 4.4 are unchanged. Its proof is modified as follows:

Case 1. $\hat{j}(1+2 \theta) \leq \alpha^{*} \theta i$. This implies $z_{i} \in Q^{(\hat{j})} \cap\left(H \backslash Q^{(j)}\right)$ for some horizontal section $H$ of generation $j$ crossing $Q^{(\hat{j})}, \hat{j}<j \leq \hat{j}(1+2 \theta)$.

Case 2. $\hat{j} \leq \alpha^{*} \theta i<\hat{j}(1+2 \theta)$. Here all $\hat{j}+1$ are replaced by $\hat{j}(1+2 \theta)$, and $d_{\mathcal{C}}\left(z_{i}\right)>e^{-\alpha i}$ is used as before.

In both cases, there is enough room for the new estimates to go through.
With regard to Section 5, unlike the situation in the last paragraph, the results in this section assume the existence of guiding critical orbits and so are unaffected. We go directly to Section 6.3.
5. Setting control for $\left(\xi_{k}, \tau_{k}\right)$ (§6.3): Here we need to set control for all $\xi_{i} \in M^{\circ} \cap \mathcal{C}^{(1)}, M \in \mathcal{T}_{i}$. The rules on the selection of $\phi\left(\xi_{i}\right)$ at free returns are essentially the same as (1)-(3) in the beginning of Section 6.3, with (3)
modified to read as follows: If there exists $j<i$ such that $\xi_{i} \in \mathcal{C}^{(j)} \backslash \cup_{k>j} \mathcal{C}^{(k)}$, then we let $\phi\left(\xi_{i}\right)=z_{0}^{*}\left(Q^{(j)}\left(\xi_{i}\right)\right)$; otherwise $\phi\left(\xi_{i}\right)$ is of generation $i$. At issue is the suitability of this choice of $\phi\left(\xi_{i}\right)$. Given that $\phi\left(\xi_{i}\right)$ may be of a lower generation than it would have been (due to the skipping of generations) and that certain monotone branches are discontinued, two questions are: (i) Do we know that $\phi\left(\xi_{i}\right)$ will remain active for as long as it is needed? (ii) If $\phi\left(\xi_{i}\right)$ is of generation $j<i$, is $\xi_{i} \notin B^{(j)}$ ?

Let $j$ be the generation of $\phi\left(z_{i}\right)$. There are two cases to consider.
Case 1. $j \leq i(1+2 \theta)^{-1}$. Here $\xi_{i}$, which is in $R_{i}$, lies in a horizontal section of generation $j^{\prime}$ crossing $Q^{(j)}\left(\xi_{i}\right)$. We may assume $j<j^{\prime} \leq j(1+2 \theta)$; see (2)(b) in Section 7.3. Observe that since $j^{\prime}>j, \xi_{i}$ is not in $\mathcal{C}^{\left(j^{\prime}\right)}$ by assumption. It follows that $d_{\mathcal{C}}\left(\xi_{i}\right)>e^{-\lambda(j+1)(1+2 \theta)}$, so that $\xi_{i} \notin B^{(j)}$ and $p\left(\xi_{i}\right) \ll \theta^{-1} j$.

Case 2. $j>i(1+2 \theta)^{-1}$. By assumption, $\xi_{i} \in M_{i} \cap Q^{(j)}$ for some $M_{i} \in \mathcal{T}_{i}$. It follows, by $\left(\mathrm{A} 1^{\prime}\right)(\mathrm{II})$, that $M_{i}$ extends all the way across $Q^{(j)}$ and $M_{i} \cap Q^{(j)}$ contains a component $Q^{(i)}$ of $\mathcal{C}^{(i)}$. If $\xi_{i} \notin Q^{(i)}$, the situation is as in Case 1. If $\xi_{i} \in Q^{(i)}$, then $\phi\left(\xi_{i}\right)=z_{0}^{*}\left(Q^{(i)}\left(\xi_{i}\right)\right)$. It is easy to show that for $k=1,2, \cdots$, $T^{k} \xi_{i}$ and $T^{k} \phi\left(\xi_{i}\right)$ lie in the same element of $\mathcal{T}_{i+k}$ until either this branch is discontinued or the bound period of $\xi_{i}$ expires. Recall that when the orbit of $\phi\left(\xi_{i}\right)$ ceases to be active, the monotone branch containing it is automatically discontinued.
6. Proposition 6.1 (§6.3): With the idea of provisional control as before, we reformulate Proposition 6.1 as follows:

Proposition 8.1. For $T \in \mathcal{G}_{N}$, and $\theta N<k \leq \frac{1}{\alpha^{*}} \theta N$ :
$(\mathrm{a})_{k} \mathcal{C}^{(k)}$ and $\Gamma_{k}$ with the properties in $\left(\mathrm{A}^{\prime}\right)$ can be constructed;
$(\mathrm{b})_{k} \mathcal{T}_{k+1}$ with the properties in $\left(\mathrm{A} 6^{\prime}\right)$ - except for the provisional nature of the control in ( $\mathrm{A} 6^{\prime}$ )(II) - can be constructed.

The proof of this proposition is postponed to Section 8.2. Assuming it for now, we continue with our list of modifications.
7. Alignment of vectors at free returns (§6.4): In the proof of Proposition 6.2(i), we view $z_{k-j}$ as a point $\xi_{1} \in R_{1}$ and show that $\left(\xi_{1}, \tau_{1}\right)$ is controlled for $j$ iterates. In this type of argument, one needs to verify that there exists $M \in \mathcal{T}_{j+1}$ such that $\xi_{j+1} \in M$; i.e., the ancestors of this branch were not discontinued. Here we know $M$ exists because our choice of $\phi\left(z_{k}\right)$ implies the existence of $Q^{(j+1)}$ with $\phi\left(z_{k}\right)=z_{0}^{*}\left(Q^{(j+1)}\right)$ and $\xi_{j+1}=z_{k} \in Q^{(j+1)}$, and, by definition, every $Q^{(j+1)}$ is contained in some $M \in \mathcal{T}_{j+1}$. The rest of the proof of Proposition 6.2 is not affected.

The quadratic estimate in (A5) relies on the behavior of the guiding critical orbits and not on global geometry; it is therefore not affected.

This completes our list of modifications for Sections 4, 5 and 6 and hence the proof of $(* *)$ stated at the end of Section 7.2 - modulo the proof of Proposition 8.1.

### 8.2. Construction of critical regions and monotone branches

Proof of Proposition 8.1. We follow in outline the proof of Proposition 6.1, focusing on those aspects of the situation that are new.

Assume $(\mathrm{a})_{i}$ and $(\mathrm{b})_{i}$ for all $i<k$. For $M \in \mathcal{T}_{k}$, if $\xi_{1} \in T^{-(k-1)} M^{\circ}$, then the sequence $\left(\xi_{1}, \tau_{1}\right), \cdots\left(\xi_{k-1}, \tau_{k-1}\right)$ is provisionally controlled by $\Gamma_{k-1}$. Thus it makes sense to speak about the leaf segments of $\mathcal{F}_{k}$ on $M^{\circ}$ as being in bound or free states. We divide the proof of $(\mathrm{a})_{k}$ and $(\mathrm{b})_{k}$ into the following steps:

1. Construction of $\mathcal{C}^{(k)}$ and $\Gamma_{k}$. Let $M \in \mathcal{T}_{k}$, and let $Q=Q^{(j)}$, $k(1+2 \theta)^{-1} \leq j<k$, be such that $M \cap Q \neq \emptyset$. We observe that $M \cap Q$ is the union of horizontal sections each extending $>\frac{1}{2} e^{-\alpha k}$ on both sides of $Q$. This is true because by (A2) for a much earlier time, $\left(E \cup E^{\prime}\right) \cap Q=\emptyset$ where $E$ and $E^{\prime}$ are the ends of $M$. We then consider one $\mathcal{F}_{k}$-leaf segment in $M$ at a time and argue as in the proof of Proposition 6.1.

Next we explain where $Q^{(k)}$ is constructed (postponing how it is done to the next paragraph). For each $M \in \mathcal{T}_{k}$, let $\mathcal{I}(M)$ be the set of all connected components of $M \cap Q^{(j)}, k(1+2 \theta)^{-1} \leq j<k$. We define a partial order on $\mathcal{I}(M)$ by set inclusion, i.e. $S<S^{\prime}$ if $S \supset S^{\prime}$. Let $\mathcal{H}(M)$ be the set of maximal elements. A critical region $Q^{(k)}$ is constructed in each $H \in \mathcal{H}(M)$.

As to how to construct $Q^{(k)}$, let $H \in \mathcal{H}(M)$ be a component of $M \cap Q^{(j)}$. We fix an arbitrary $\mathcal{F}_{k^{-}}$-leaf $\gamma$ in $H$, and use Lemma 3.8 and $z_{0}^{*}\left(Q^{(j)}\right)$ to induce a (unique) critical point $z_{0}^{\prime}$ of order $j$ on $\gamma$. In the $x$-coordinate, we know from Lemma 3.8 that $z_{0}^{\prime}$ is $<K b^{\frac{j}{4}}$ away from the center of $Q^{(j)}$. Lemma 3.9 then tells us that near $z_{0}^{\prime}$ there is a unique critical point $z_{0}$ of order $k$ on $\gamma$. We put $z_{0} \in \Gamma_{k}$, and construct a $Q^{(k)}$, i.e. a section of length $2 \min \left(\delta, e^{-\lambda k}\right)$ centered at it.
$\mathcal{C}^{(k)}$ is defined to be the union of all the $Q^{(k)}$ constructed as we let $M$ vary over $\mathcal{T}_{k}$.
2. Verification of $\left(\mathrm{A1}^{\prime}\right)$. The procedure above immediately gives the following: (1) The $Q^{(k)}$ constructed are disjoint sections of $R_{k}$ and hence are genuine connected components of $\mathcal{C}^{(k)}$. (2) The partial order in $\left(\mathrm{A}^{\prime}\right)(\mathrm{I})(\mathrm{a})$ is extended to $\left\{Q^{\left(k^{\prime}\right)}, k^{\prime} \leq k\right\}$; this is because each $Q^{(k)}$ constructed lies immediately below a unique $Q^{(j)}, k(1+2 \theta)^{-1} \leq j<k$, in this partial order. (2) implies (3), namely that the jumps in generation in (A1')(I)(ii) are as claimed. Observe also that $\left(\mathrm{A1}^{\prime}\right)(\mathrm{II})$ is fulfilled. As for $\left(\mathrm{Al}^{\prime}\right)(\mathrm{I})(\mathrm{b})$, all statements are true by construction except the one regarding sectional diameter, which follows directly from the next lemma.

Lemma 8.1. Let $M \in \mathcal{T}_{k}$. Then for all $\xi_{k} \in M^{\circ}$, there exists a codimension one manifold $W$ with $\xi_{k} \in W$ such that

- $W$ meets every connected component of $\mathcal{F}_{k}$-leaf in $M$ in exactly one point;
- for all $\xi_{1}, \xi_{1}^{\prime} \in T^{-k+1} W,\left|\xi_{i}-\xi_{i}^{\prime}\right|<(K b)^{\frac{i}{2}}$ for $i \leq k$.

The proof of this lemma is a small modification of that of Lemma 6.2. Details are left to the reader.
3. Construction of $\mathcal{T}_{k+1}$ and verification of $\left(\mathrm{A} 6^{\prime}\right)(\mathrm{I})$ and (II). The relation between $M \in \mathcal{T}_{k}$ and $B^{(k)}$ follows immediately from our construction in Step 1. We construct $\mathcal{T}_{k+1}$ as described in the paragraph on "Tree of monotone branches" in Section 7.2, proving ( $\mathrm{A} 6^{\prime}$ )(I). To prove ( $\mathrm{A} 6^{\prime}$ )(II), it suffices, as in the proof of Proposition 6.1, to prove correct alignment of the $\tau$-vectors at free returns, and we consider the two cases as before. To make transparent the effect of the "missing generations", we describe in some detail the geometry in Case 2 , the case where $\gamma$, our maximal free $\mathcal{F}_{k}$-segment, does not meet $\mathcal{C}^{(k)}$.

Let $\xi_{k}$ be fixed. We let $j=j_{0}$ be the largest integer such that $\xi_{k} \in \mathcal{C}^{(j)}$, and let $Q^{\left(j_{0}\right)}=Q^{\left(j_{0}\right)}\left(\xi_{k}\right)$. Then $j_{0}(1+2 \theta)<k$; otherwise $\gamma$ would lie in a monotone branch crossing $Q^{\left(j_{0}\right)}$, contradicting our assumption that it does not meet $\mathcal{C}^{(k)}$. By $\left(\mathrm{A}^{\prime}\right)(\mathrm{III}), \xi_{k} \in R_{j_{0}(1+2 \theta)+1} \subset \cup\left\{M, M \in \cup_{j_{0}<\ell \leq j_{0}(1+2 \theta)} \mathcal{T}_{\ell}\right\}$. Thus $\xi_{k} \in H$ where $H$ is a horizontal section of generation $j_{1}, j_{0}<j_{1} \leq j_{0}(1+2 \theta)$, that crosses the entire length of $Q^{\left(j_{0}\right)}\left(\xi_{k}\right)$. We may assume $Q^{\left(j_{1}\right)}$ is immediately below $Q^{\left(j_{0}\right)}$ in the partial order. Suppose for definiteness that $\xi_{k}$ lies in the right chamber of $H \backslash Q^{\left(j_{1}\right)}$. We move left along $\gamma$ until we reach either $\xi$, the left endpoint of $\gamma$, or the right boundary of $Q^{\left(j_{1}\right)}$, noting that since $j_{1}<k, \gamma$ cannot meet $\partial R_{j_{1}}$. If we reach the boundary of $Q^{\left(j_{1}\right)}$ before reaching $\xi$, then we continue to move left along $\gamma$, going into $Q^{\left(j_{1}\right)}$. For the same reason as above, $j_{1}(1+2 \theta)<k$, so that as we enter $Q^{\left(j_{1}\right)}$, we have entered a horizontal section of generation $j_{2}, j_{1}<j_{2} \leq(1+2 \theta) j_{1}$, that crosses the entire length of $Q^{\left(j_{1}\right)}$, and so on. After going through a finite number of $Q^{\left(j_{i}\right)}$, we must arrive at $\xi$, for the $j_{i}$ are strictly increasing with $i$ and $<k(1+2 \theta)^{-1}$.

The argument showing $\phi(\xi)$ is to the left of $\xi$ is essentially the same, except for the fact that the interpolating chain $Q^{\left(j^{\prime}\right)} \subset \cdots \subset Q^{\left(j^{\prime \prime}\right)}$ also involves skips in generation. One way to see that such a chain exists is to start from $Q^{\left(j^{\prime \prime}\right)}$.

After these preparations, the angle estimates are unchanged. This completes the verification of $\left(\mathrm{A} 6^{\prime}\right)(\mathrm{II})$.
4. Proof of $\left(\mathrm{A}^{\prime}\right)(\mathrm{III})$. This step involves a very different set of ideas. We formulate the result as Proposition 8.2 and give the proof in the next subsection.

Proposition 8.2. Consider $T \in \mathcal{G}_{N}$, and let $n \leq \frac{1}{\alpha^{*}} \theta N$. Assume (a) ${ }_{k}$ and $(\mathrm{b})_{k}$ in Proposition 8.1 for all $k \leq n-1$. Then for all $k \leq n(1+2 \theta)^{-1}$,

$$
\begin{equation*}
R_{k(1+2 \theta)} \subset \cup\left\{M, M \in \cup_{k \leq j<k(1+2 \theta)} \mathcal{T}_{j}\right\} \tag{9}
\end{equation*}
$$

Modulo this result, the proof of Proposition 8.1 is now complete.

### 8.3. Branch replacement: Proof of Proposition 8.2

As explained in Section 8.1, a monotone branch is discontinued before either one of its ends becomes too large. The problem of "branch replacement", roughly speaking, is one of finding a collection of branches of higher generations that together cover the part of the attractor "exposed" by the removal of the discontinued branch. Proposition 8.2 tells us explicitly what neighborhoods are covered by which collections of branches.

We begin with some preliminary definitions. Let $M_{1}$ be a monotone branch of generation $k_{1}>k$. We say $M_{1}$ is subordinate to $M$ if (i) $M_{1} \subset M$, (ii) the ends of $M$ and $M_{1}$ are related as follows: Let $E$ and $E^{\prime}$ be the ends of $M$, and $E_{1}$ and $E_{1}^{\prime}$ the ends of $M_{1}$. Suppose $T^{-i} E=\hat{B}^{(k-i)}$; then $T^{-i} E_{1}=\hat{B}_{1}^{\left(k_{1}-i\right)}$ with $\hat{B}_{1}^{\left(k_{1}-i\right)} \subset \hat{B}^{(k-i)}$; and $E^{\prime}$ and $E_{1}^{\prime}$ are related the same way. A collection of monotone branches $\left\{M_{j}\right\}$ subordinate to $M$ is called a viable replacement for $M$ if $\left(M^{\circ} \cap \Omega\right) \subset \cup_{j} M_{j}$.

Lemma 8.2. There exists $K>0$ for which the following holds: Suppose $M \in \mathcal{T}_{k}, k \leq n$, has an end $E$ with the property that $T^{-i} E=\hat{B}^{(k-i)}, i \geq$ $K \alpha(k-i)$. Then $T^{-i} M$ is contained in a horizontal section $H$ of $R_{k-i}$ of length $<e^{-2 \alpha(k-i)}$ centered at $\hat{B}^{(k-i)}$.

Lemma 8.2 is proved in Appendix A.21. Let $M$ be as in Lemma 8.2 and assume, for definiteness, that $T^{-i} M$ lies in the right half of $H$ (it contains, needless to say, $\hat{B}^{(k-i)}$ ). To look for a viable replacement for $M$, we examine the structures inside $H$ more closely.

For $j=1,2, \cdots, i-1$, let $S_{j} \in \mathcal{T}_{k-i+j}$ be the ancestors of $M$, and let $\hat{S}_{0}=T^{-1} S_{1}$, so that $\hat{S}_{0}$ is a section of $R_{k-i}$ containing the right half of $H$. From Lemma 8.2, it follows that there exists $\ell$ with $\ell<K \alpha(k-i)$ such that $T^{-\ell} S_{\ell} \subset H$.

Consider now $P \in \mathcal{T}_{p},(k-i)<p<(1+2 \theta)(k-i+1)$, such that $P \cap H \neq \emptyset$. Then $P \cap H$ is the union of horizontal sections that run the entire length of $H$. (See Figure 3.) We fix one component of $H \cap P$, call it $\hat{H}$, and let $\hat{B}^{(p)}$ denote the $B^{(p)}$ in $\hat{H}$. Whenever possible, we define $P_{j} \in \mathcal{T}_{p+j}, j=1,2, \cdots, i$, as follows: $P_{1}$ is the child of $P$ such that $T^{-1} P_{1}$ contains the right half of $\hat{H}$; for $j>1, P_{j}$ is the child of $P_{j-1}$ one of whose ends is $T^{j} \hat{B}^{(p)}$. If well defined, $P_{j}$ depends on $M, P$ and $\hat{H}$; we write $P_{j}(M, P, \hat{H})$.


Figure 3. Replacement branches: $\cup P$ is used to replaced $M$
We observe that $P_{j}$ may not be defined: First, $P$ may be discontinued, in which case it has no children. If that is not the case, then $P_{1}$ is defined. Let $E_{1}^{\prime}$ denote the end of $P_{1}$ not equal to $T\left(\hat{B}^{(p)}\right)$. This is clearly the older of the two ends of $P_{1}$. It may cause $P_{1}$ to be discontinued. We may also have $P_{2}=T\left(P_{1}\right)$, with $T\left(E_{1}^{\prime}\right)$ causing $P_{2}$ to be discontinued, and so on.

Let $n$ be as in the statement of Proposition 8.2, and let $M, P$ and $H$ be as above.

Lemma 8.3. Assume $k \leq(1+2 \theta)^{-1} n$. Then the following hold for every component $\hat{H}$ of $P \cap H$ : If $P_{1}$ is well defined and $E_{1}^{\prime}$ remains active for $\ell$ generations for some $\ell$ with $T^{-\ell} S_{\ell} \subset H$, then
(i) $P_{j}$ is well defined for all $j \leq i$, and
(ii) $P_{i}$ is subordinate to $M$.

Lemma 8.3 is proved in Appendix A.21.
Proof of Proposition 8.2. Our strategy is to construct, for $m=1,2, \cdots$, $n(1+2 \theta)^{-1}$, a collection of monotone branches $\mathcal{S}_{m}$ with the following properties:
(i) For every $M \in \mathcal{S}_{m}$, if $M \in \mathcal{T}_{k}$ and $E$ is an end of $M$ with $T^{-i} E=B^{(k-i)}$, then $i \leq \frac{2}{3}(k-i) \theta^{-1}$;
(ii) $\mathcal{S}_{m} \subset \cup_{m \leq k<m(1+2 \theta)} \mathcal{T}_{k}$, and
(iii) $\cup\left\{M, M \in \mathcal{S}_{m}\right\} \supset R_{m(1+2 \theta)}$.

Proposition 8.2 follows immediately from (ii) and (iii). We say $M \in \mathcal{S}_{m}$ is at replacement time if equality holds in (i), i.e. $i=\frac{2}{3}(k-i) \theta^{-1}$, for one of its ends. For reasons to become clear momentarily, we have elected to define replacement time to occur somewhat before the branch is discontinued.

Let $\mathcal{S}_{1}=\left\{R_{1}\right\}$. We assume for all $k \leq m, \mathcal{S}_{k}$ has been constructed and has properties (i)-(iii).

Construction of $\mathcal{S}_{m+1}$ with property (i): We consider $M \in \mathcal{S}_{m}$ one at a time. If $M$ has not reached its replacement time, then we put all the children of $M$ into $\mathcal{S}_{m+1}$. If it has, then we put into $\mathcal{S}_{m+1}$ the children of $\left\{P^{\prime}\right\}$ where $\left\{P^{\prime}\right\}$ is defined as follows: Suppose $M \in \mathcal{T}_{k}$ and $E$ is an end with $T^{-i} E=\hat{B}^{(k-i)}$ and $i=\frac{2}{3}(k-i) \theta^{-1}$. Then $\left\{P^{\prime}\right\}=\left\{P_{i}(M, P, \hat{H}): P \in \mathcal{S}_{k-i+1}\right.$ and $\hat{H}$ is a component of $H \cap P\} .{ }^{9}$

First we show that $P^{\prime}$ is well defined as an element of $\mathcal{T}_{p+i}$ where $p$ is the generation of $P$. To do that, it suffices to verify the hypotheses of Lemma 8.3. To begin with, $P_{1}$ is well defined as an element of $\mathcal{T}_{p+1}$ because $P \in \mathcal{S}_{k-i+1}$ and by property (i) for $\mathcal{S}_{k-i+1}$, both ends of $P$ will remain active for some period of time. Let $E_{1}^{\prime}$ be as above, i.e. the "other" end of $P_{1}$. We claim that $E_{1}^{\prime}$ will last $\gg K \alpha(k-i)$ generations: Suppose it was created $\ell$ generations prior to $p+1$. Then $\ell \leq \frac{2}{3}(p+1-\ell) \theta^{-1}$, so that

$$
\begin{aligned}
K \alpha(k-i)+\ell & \leq K \alpha(k-i-\ell)+(K \alpha+1) \ell \\
& \leq K \alpha(k-i-\ell)+(K \alpha+1) \frac{2}{3}(p+1-\ell) \theta^{-1} \\
& <(p+1-\ell) \theta^{-1} .
\end{aligned}
$$

The hypothesis of Lemma 8.3 is verified and $P^{\prime}$ is defined.
To prove that the children of $P^{\prime}$ meet the condition in property (i) for $\mathcal{S}_{m+1}$, we let $\hat{E}$ and $\hat{E}^{\prime}$ be its two ends, $\hat{E}$ being the one contained in $E$. This end is created the same time $E$ is created. Clearly, $i<\frac{2}{3} p \theta^{-1}$ since $p>k-i$. As for $\hat{E}^{\prime}$, it follows from the analysis in Lemma 8.3 that this is the younger of the two ends. Thus it cannot have reached replacement time if $\hat{E}$ has not. This completes the construction of $\mathcal{S}_{m+1}$ with property (i).

Proof of property (ii) for $\mathcal{S}_{m+1}$ : Let $M \in \mathcal{S}_{m}$. If $M$ is not replaced, then the children of $M$ are obviously of acceptable generation. If replacement occurs, then the generation of $P^{\prime}$ is estimated as follows: Let all notation be as above. Since $P \in \mathcal{S}_{k-i+1}$, by inductive assumption, $p<(1+2 \theta)(k-i+1)$. Combining this with $\frac{3}{2} \theta i=k-i$, we have

$$
p+i \leq i+(1+2 \theta)(k-i+1)=i+(1+2 \theta)\left(\frac{3}{2} \theta i+1\right)<(i+1)(1+2 \theta) .
$$

To complete the proof, we show that $i \leq m$. First, it is true for $m=2$. In general, we claim that if $E$ is an end of $M$ and $T^{-i} E=B^{(k-i)}$, then $i \leq m$. This is obviously so if no replacement occurs. In a replacement procedure, observe that even though the generations of the new branches are higher, their ends are created exactly the same number of generations earlier as the branch replaced. (See the proof of Lemma 8.3.)

[^9]Proof of property (iii) for $\mathcal{S}_{m+1}$ : As before, it suffices to consider the case where $M \in \mathcal{S}_{m}$ is at replacement time. We claim that $\left\{P^{\prime}\right\}$ is a viable replacement for $M$. Let $H$ be as above. By the induction hypothesis, more specifically, by (iii) for $k-i+1$,

$$
H \cap R_{(k-i+1)(1+2 \theta)} \subset H \cap\left(\cup\left\{M^{\prime}, M^{\prime} \in \mathcal{S}_{k-i+1}\right\}\right) \subset H \cap(\cup P)
$$

that is to say, $H \backslash(\cup P)$ does not meet $R_{(k-i+1)(1+2 \theta)}$. Thus the part of phase space deleted as we replace $M$ by $\left\{P^{\prime}\right\}$ does not meet $R_{g}$ where

$$
g:=(k-i+1)(1+2 \theta)+i .
$$

The same computation as in the proof of property (ii) gives $g<(m+1)(1+2 \theta)$.
This completes the proof of Proposition 8.2.
Remarks. 1. As the proof shows, there is a natural time step for branch replacements. They occur very infrequently, roughly once every $\sim \theta^{-1}$ iterates. Thus when working with the replacement of a branch of generation $k$, it is mostly structures of generation up to $\sim \theta k$ - including assumption (A $6^{\prime}$ )(III) for these times - that count, although certain properties of critical orbits and the well-definedness of branches up to time $k-1$ are also needed.
2. In spite of the qualitative flavor of the statement, the existence of viable replacements in the setting above reflects the fact that monotone branches reproduce at rates much faster than the speed with which critical blobs are allowed to approach the critical set.

## 9. Construction of SRB Measures

The definition of SRB measures is given in Section 1. For more information on the subject, see [Y1] and [Y2]. The goal of this section is to prove

Proposition 9.1. Let $T \in \mathcal{G}$. Then $T$ has an SRB measure.

### 9.1. Generic part of construction

Our construction can be thought of as having a "generic" part, i.e. a part that can be used for many dynamical systems, and a "situation-dependent" part, i.e. a part that relies (seriously) on the properties of the map in question. The goal of this subsection is to treat the generic part. We pinpoint what specific information is needed and then assume it to complete the construction. For notational convenience, we give the proof in the setting of Part II of this paper, remarking that aside from $\operatorname{dim} W^{u}=1$, other properties of $T$ used below are inessential.

Step 1. Pushing forward Lebesgue measure on a $W^{u}$-leaf. Let $l_{0}$ be a piece of local unstable manifold through $\hat{z}$ where $\hat{z}$ is a hyperbolic periodic point or
belongs in a uniformly hyperbolic invariant set. We let $m_{0}$ be the Riemannian measure on $l_{0}$, and for $n=1,2, \cdots$, define

$$
\nu_{n}=\frac{1}{n} \sum_{i=0}^{n-1} T_{*}^{i}\left(m_{0}\right)
$$

where $T_{*}^{i}\left(m_{0}\right)$ is the measure with $T_{*}^{i}\left(m_{0}\right)(E)=m_{0}\left(T^{-i}(E)\right)$ for all Borel sets $E$. Let $\nu$ be a limit point of $\nu_{n}$ in the weak* topology. It is easy to see that $\nu$ is $T$-invariant.

To prove that $\nu$ is an SRB measure, it is necessary to show that it has absolutely continuous conditional measures ( accm ) on unstable manifolds. By design, this property is enjoyed by $\nu_{n}$ for all $n$. Whether it is passed from $\nu_{n}$ to $\nu$, however, depends on a number of factors that are situation-dependent. We describe in Step 2 a construction to facilitate this passage if the conditions are right. Our construction is based on the idea (also used in [BY]) that it suffices to control a small fraction of $\nu_{n}$.

Step 2. "Catching" a fraction of $\nu$ with accm on unstable curves. First we introduce some language convenient for our purposes. We call a curve $\gamma$ an unstable curve if there exist $\kappa<1$ and $K>1$ such that for all $z \in \gamma$ and $\tau \in X_{z}$ tangent to $\gamma,\left|D T_{z}^{-n} \tau\right| \leq K \kappa^{n}|\tau|$ for all $n \geq 0$. Next we introduce the objects used to "catch" a part of $\nu$. Let $L$ be an interval, and let $\Sigma$ be a compact set. We say $\Psi: L \times \Sigma \rightarrow R_{1}$ is a continuous family of unstable curves if
(a) $\Psi$ maps $L \times \Sigma$ homeomorphically onto its image;
(b) for each $\alpha \in \Sigma,\left.\Psi\right|_{L \times\{\alpha\}}$ is a $C^{1}$ embedding, and $D_{\alpha}:=\Psi(L \times\{\alpha\})$ is an unstable curve;
(c) $\left.\alpha \mapsto \Psi\right|_{L \times\{\alpha\}}$ is continuous as a map from $\Sigma$ to $C^{1}\left(L, R_{1}\right)$.

We will use the notation $\mathcal{N}=\Psi(L \times \Sigma)=\cup_{\alpha} D_{\alpha}$.
The following condition, the validity of which depends on the specifics of the map in question, is assumed for the rest of this subsection:
(S) There exist $c>0, K \geq 1$, a continuous family of unstable curves $\mathcal{N}=$ $\cup_{\alpha} \mathcal{D}_{\alpha}$, and a sequence of integers $n_{1}<n_{2}<\cdots$ for which the following hold. For each $i \geq 0$, there is a collection $\left\{\omega_{j}^{(i)}\right\}$ of subsegments of $l_{0}$ such that
(i) for each $i, j, T^{i}\left(\omega_{j}^{(i)}\right)=D_{\alpha}$ for some $\alpha$;
(ii) letting $\tau(z)$ denote a unit vector tangent to $l_{0}$ at $z \in l_{0}$, we have, for all $z, z^{\prime} \in \omega_{j}^{(i)}$,

$$
\frac{\left|D T_{z}^{i} \tau(z)\right|}{\left|D T_{z^{\prime}}^{i} \tau\left(z^{\prime}\right)\right|}<K
$$

(iii) $\frac{1}{n_{k}} \sum_{i=0}^{n_{k}-1} m_{0}\left(\cup_{j} \omega_{j}^{(i)}\right)>c m_{0}\left(l_{0}\right)$ for all $n_{k}$.

Let $\hat{\nu}_{n_{k}}=\frac{1}{n_{k}} \sum_{i=0}^{n_{k}-1} T_{*}^{i}\left(\left.m_{0}\right|_{\bigcup_{j} \omega_{j}^{(i)}}\right)$, and let $\hat{\nu}$ be a limit point of $\hat{\nu}_{n_{k}}$. It follows from (iii) in Condition (S) that $\hat{\nu}(\mathcal{N})>c m_{0}\left(l_{0}\right)>0$. From (i) and (ii), we see that for each $n_{k}, \hat{\nu}_{n_{k}}$ is supported on a finite number of $D_{\alpha}$, and its densities with respect to arclength measure on $D_{\alpha}$ are bounded between $K$ and $\frac{1}{K}$. The absolute continuity of conditional measures of $\hat{\nu}$ on $\left\{D_{\alpha}\right\}$ is now a simple exercise: Let $\eta_{1}<\eta_{2}<\cdots$ be any increasing sequence of finite partitions on $\Sigma$ such that $\bigvee_{i>0} \eta_{i}$ partitions $\Sigma$ into points, and let $\mathcal{E}_{i}$ be the partition on $\mathcal{N}$ given by $\left\{\Psi(L \times S): S \in \eta_{i}\right\}$. Then $\mathcal{E}_{\infty}:=\bigvee \mathcal{E}_{i}$ is the partition of $\mathcal{N}$ into $\left\{D_{\alpha}\right\}$. Let $\ell \subset L$ be an arbitrary interval, and let $A=\Psi(\ell \times \Sigma)$. Then there exists $K^{\prime}$ depending on the constant $K$ in (S)(ii) and on the norms of the embeddings $\left.\Psi\right|_{L \times\{\alpha\}}$ such that for all $n_{k}$ and $i$,

$$
\begin{equation*}
\frac{1}{K^{\prime}} \frac{|\ell|}{|L|} \leq\left(\hat{\nu}_{n_{k}} \mid \mathcal{E}_{i}\right)(A) \leq K^{\prime} \frac{|\ell|}{|L|} . \tag{10}
\end{equation*}
$$

Here $\hat{\nu}_{n_{k}} \mid \mathcal{E}_{i}$ denotes the conditional measure of $\hat{\nu}_{n_{k}}$ given $\mathcal{E}_{i}$. The relation in (10) is first passed to $\hat{\nu} \mid \mathcal{E}_{i}$ by letting $n_{k} \rightarrow \infty$. It is then passed, by the martingale convergence theorem, to $\hat{\nu} \mid \mathcal{E}_{\infty}$. Our assertion on the conditional measures of $\hat{\nu}$ follows as we let $\ell$ range over a countable basis of the Borel topology on $L$.

Step 3. Extracting an SRB measure from $\nu$. Let $\nu$ be as in Step 1. Then $\nu(\mathcal{N}) \geq \hat{\nu}(\mathcal{N})>0$, and $\nu_{\mathcal{N}}:=\left.\nu\right|_{\mathcal{N}}$ is $T_{\mathcal{N}}$-invariant where $T_{\mathcal{N}}: \mathcal{N} \rightarrow \mathcal{N}$ is the first return map of $T$ to $\mathcal{N}$. Let $R: \mathcal{N} \rightarrow \mathbb{Z}^{+}$be the first return time, and assume for the moment that $\nu_{\mathcal{N}}$ has accm on $\left\{D_{\alpha}\right\}$. We claim that

$$
\mu:=\sum_{n>0} \sum_{k=0}^{n-1} T_{*}^{k}\left(\left.\nu_{\mathcal{N}}\right|_{\{R=n\}}\right)
$$

normalized is an SRB measure. To prove this, it suffices to show (i) $T$ has a positive Lyapunov exponent $\mu$-a.e., and (ii) the $D_{\alpha}$-curve through a.e. $z \in \mathcal{N}$ is its local unstable manifold. Both are true by the construction in Step 2.

While $\hat{\nu}$ (which is not necessarily $T_{\mathcal{N}}$-invariant) has accm on $\left\{D_{\alpha}\right\}$, we do not know in general that $\nu_{\mathcal{N}}$ does as well. The following procedure is used to extract a part of $\nu_{\mathcal{N}}$ with the desired property: Let $\nu_{\Sigma}=\left(\pi_{\Sigma}\right)_{*}\left(\Psi_{*}^{-1}\left(\nu_{\mathcal{N}}\right)\right)$ where $\pi_{\Sigma}: L \times \Sigma \rightarrow \Sigma$ is projection. Let $\lambda$ denote Lebesgue measure on $L$, and let $\nu_{\text {Leb }}=\Psi_{*}\left(\lambda \times \nu_{\Sigma}\right)$. We then decompose $\nu_{\mathcal{N}}$ into $\nu_{\mathcal{N}}=\nu_{\text {ac }}+\nu_{\perp}$ where $\nu_{\text {ac }}$ is absolutely continuous with respect to $\nu_{\text {Leb }}$ and $\nu_{\perp}$ is singular with respect to it (written $\nu_{\text {ac }} \ll \nu_{\text {Leb }}$ and $\left.\mu_{\perp} \perp \nu_{\text {Leb }}\right)$. Since $\left(T_{\mathcal{N}}\right)_{*}\left(\nu_{\text {ac }}\right) \ll \nu_{\text {Leb }}$ and $\left(T_{\mathcal{N}}\right)_{*}\left(\nu_{\perp}\right) \perp \nu_{\text {Leb }}$, it follows that both $\nu_{\text {ac }}$ and $\nu_{\perp}$ are $T_{\mathcal{N}}$-invariant. By Condition $(\mathrm{S}), \nu_{\mathrm{ac}}(\mathcal{N}) \geq \hat{\nu}(\mathcal{N})>0$. The construction of $\mu$ above can now be carried out with $\nu_{\mathrm{ac}}$ in the place of $\nu_{\mathcal{N}}$.

Summary. We have shown that if for some $l_{0}=W_{\text {loc }}^{u}(\hat{z})$ there is a family of unstable curves $\mathcal{N}=\cup_{\alpha} \mathcal{D}_{\alpha}$ for which Condition (S) is satisfied, then $T$ admits an SRB measure.

For the rest of this section, we assume $T \in \mathcal{G}$.

### 9.2. Dynamics on unstable manifolds

Let

$$
\Omega_{\delta}:=\left\{z_{0} \in \Omega: d_{\mathcal{C}}\left(z_{i}\right) \geq \delta \forall i \in \mathbb{Z}\right\} .
$$

From Section 3.5 , we know that $\left.T\right|_{\Omega_{\delta}}$ is uniformly hyperbolic. Fix $\hat{z} \in \Omega_{\delta}$, and let $l_{0}=W_{r}^{u}(\hat{z})$ denote its local unstable curve of radius $r$. We assume $r$ is small enough that for all $\xi_{0} \in l_{0}, d_{\mathcal{C}}\left(\xi_{-i}\right)>\frac{1}{2} \delta$ for all $i \geq 0$. In the rest of this section, we let $\tau_{0}\left(\xi_{0}\right) \in X_{\xi_{0}}$ denote the positively oriented unit vector tangent to $l_{0}$, and use $\tau$ to denote generic unit vectors tangent to $l_{i}:=T^{i} l_{0}$.
A. Control of $\left(\xi_{0}, \tau_{0}\right)$ for $\xi_{0} \in l_{0}$. For $z \in \Omega \backslash \mathcal{C}$, a natural choice of $\phi(z)$ is $\phi(z)=z_{0}^{*}\left(Q^{(j)}\right)$ where $j$ is the largest $k$ such that $z \in Q^{(k)}$. Observe that $j=\infty$ corresponds exactly to $z \in \mathcal{C}$. Thus for all $\xi_{0} \in \Omega$ such that $\xi_{i} \notin \mathcal{C}$ for all $0 \leq i<k, \xi_{0}$ is controlled by $\cup_{j>0} \Gamma_{j}$ for $k$ iterates. Proof of control for $\tau_{0}\left(\xi_{0}\right)$ is obtained by leveraging ( $\mathrm{A}^{\prime}$ ); details are given in Appendix A.22.

Lemma 9.1. For all $\xi_{0} \in l_{0}$, the sequence $\left(\xi_{0}, \tau_{0}\right), \cdots,\left(\xi_{k}, \tau_{k}\right)$ is controlled by $\cup \Gamma_{j}$ provided $\xi_{i} \notin \mathcal{C}$ for all $i<k$.

Once control is established, the evolution of $l_{i}$ is very similar to that of $\mathcal{F}_{k}$-leaves. We record their geometric and dynamical properties in B and C below.
B. Geometry of $l_{i}$. The proof of the following is entirely parallel to that of Lemma 7.1: (i) For each $i>0, l_{i}$ is partitioned by $\left\{T^{i-k} z_{0}, k<i, z_{0} \in\left(l_{k} \cap \bar{\Gamma}\right)\right\}$ where $\bar{\Gamma}$ is the closure of $\cup_{j} \Gamma_{j}$ into a finite disjoint union of monotone segments $\{\sigma\}$.
(ii) The free part of $\sigma$, if nonempty, is connected, and the function on $\sigma$ giving the number of iterates before a point becomes free is $U$-shaped. (See the last paragraph of the proof of Lemma 7.1 in Appendix A.20.)
(iii) If $\sigma \cap Q^{(1)} \neq \emptyset$, then either (a) $\sigma \cap Q^{(1)}$ meets $\bar{\Gamma}$ in a single point $\hat{z}_{0}$, and $\sigma$ contains a $C^{2}(b)$ curve of length $e^{-\alpha i}$ centered at $\hat{z}_{0}$, or (b) $\sigma \cap Q^{(1)}$ lies strictly to one side of $\bar{\Gamma}$; in this case we let $\hat{z}_{0}=\phi(\xi)$ where $\xi$ is the point in $\sigma \cap Q^{(1)}$ closest to $\mathcal{C}$.
C. 1D behavior. Behavior near the "turns" excepted, the dynamics of $l_{0} \rightarrow l_{1} \rightarrow l_{2} \rightarrow \cdots$ bear a striking resemblance to those of iterated 1D maps. By this, we mean a qualitative resemblance rather than the existence of a specific map $f: I \rightarrow I$ with the property that $f^{k}\left(\pi_{x}\left(l_{0}\right)\right) \approx \pi_{x}\left(l_{k}\right)$ for all $k$. Here
$\pi_{x}$ denotes projection onto the $x$-axis. To make precise this qualitative resemblance, we formulate three properties in analogy with (P1)-(P3) in Section 2.2.

Proposition 9.2. $\left.T^{k}\right|_{l_{0}}, k=1,2, \cdots$, satisfy ( $\mathrm{P}^{\prime}$ )-( $\mathrm{P} 3^{\prime}$ ) below.
( $\mathrm{P}^{\prime}$ ) (Outside of $\mathcal{C}^{(1)}$ ). Let $f_{0} \in \mathcal{M}$ be as in Section 6.1, and let $\varepsilon=$ $\max (\mathcal{O}(a), \mathcal{O}(b))$. Then for all $z \in l_{i} \backslash \mathcal{C}^{(1)}$ in a free state, we have

$$
\left|\pi_{x}(T z)-f_{0}\left(\pi_{x}(z)\right)\right|, \quad| | D T_{z}(\tau)\left|-\left|f_{0}^{\prime}\left(\pi_{x}(z)\right)\right|\right|<\varepsilon .
$$

From this it follows that results analogous to (P1)(i),(ii) with $\left|\left(f^{n}\right)^{\prime}\right|$ replaced by $\left|D T^{n}(\tau)\right|$ hold with slightly weaker constants for segments of $l_{i}$ in free state.
( $\mathrm{P}^{\prime}$ ) (Bound periods and derivative recovery). Let $\omega$ be the maximal free segment in a component of $l_{i} \cap \mathcal{C}^{(1)}$, let $\hat{z}_{0}$ be as in $\mathrm{B}(\mathrm{iii})$ above, and let $\mathcal{P}^{z_{0}}$ be the partition in Section 2.2 centered at $\pi_{x}\left(\hat{z}_{0}\right)$. Then there exist $K_{0}$ and $K_{1}$ such that
(i) $K_{0}^{-1} \log \frac{1}{\left|z-\hat{z}_{0}\right|} \leq p(z) \leq K_{0} \log \frac{1}{\left|z-\hat{z}_{0}\right|}$ for all $z \in \omega$;
(ii) $\left|D T_{z}^{p}(\tau)\right|>e^{\frac{1}{4} \lambda p(z)}$ for all $z \in \omega$;
(iii) if $\pi_{x}(\omega) \approx I_{\mu j}$ for some $I_{\mu j} \in \mathcal{P}^{\hat{z}_{0}}$, then $\left|T^{p}(\omega)\right|>e^{-K_{1} \alpha|\mu|}$.

Let $\omega$ be a segment of $l_{i}$. We say all $z \in \omega$ have the same itinerary for $n-1$ iterates if there exist $t_{1}<t_{1}+p_{1} \leq t_{2}<t_{2}+p_{2} \leq \cdots \leq n$ such that for every $k, \pi_{x} \circ T^{t_{k}} \omega \subset P^{+}$for some $P \subset C_{\delta}, p_{k}=\min _{z \in \omega} p\left(T^{t_{k}} z\right)$, and for all $i \in[0, n) \backslash \cup_{k}\left[t_{k}, t_{k}+p_{k}\right), \pi_{x} \circ T^{i} \omega \subset P^{+}$for some $P \cap C_{\delta}=\emptyset$.
( $\mathrm{P}^{\prime}$ ) (Distortion estimate). There exists $K_{2}$ such that the following hold for all $i, n$ and $\omega \subset l_{i}$ satisfying (i) all $z \in \omega$ have the same itinerary for $n-1$ iterates and (ii) both $\omega$ and $T^{n}(\omega)$ are free. Then for all $z, z^{\prime} \in \omega$,

$$
\frac{\left|D T_{z}^{n}(\tau(z))\right|}{\left|D T_{z^{\prime}}^{n}\left(\tau\left(z^{\prime}\right)\right)\right|}<K_{2} .
$$

A proof of Proposition 9.2 is given in Appendix A.22.

### 9.3. Distribution of free segments of length $>\delta$

For definiteness, assume $l_{0}$ is such that either (i) $\pi_{x}\left(l_{0}\right)=I_{\mu_{0} j_{0}}$ where $I_{\mu_{0} j_{0}}$ is one of the outermost $I_{\mu j}$ or (ii) $l_{0} \cap \mathcal{C}^{(1)}=\emptyset$ and has length $>K^{-1} \delta$. Following the procedure in Section 2.3, we introduce on $l_{0}$ an increasing sequence of partitions $\mathcal{Q}_{0}<\mathcal{Q}_{1}<\mathcal{Q}_{2}<\cdots$ with $\mathcal{Q}_{i}$ representing a canonical subdivision by itinerary. This means in particular that $\mathcal{Q}_{0}=\left\{l_{0}\right\}$, and each $\omega \in \mathcal{Q}_{i-1}$ has the property that all $z \in \omega$ have the same itinerary through step $i-1$ in the sense of Section 9.2C. We are particularly interested in those $\omega \in \mathcal{Q}_{i-1}$ for which $T^{i} \omega$ is free and $\left|T^{i} \omega\right|>\delta$. These are the segments that will be used in our constructions in Section 9.1. Observe that ( $\mathrm{P} 3^{\prime}$ ) holds for $T^{i}{ }_{\omega}$.

As in Section 2.3, let $S$ be the stopping time on $l_{0}$ defined by $S(z)=i$ if and only if $T^{i}\left(\mathcal{Q}_{i-1}(z)\right)$ is free and has length $>\delta$. We introduce a sequence of stopping times $S_{0}<S_{1}<S_{2}<\cdots$ on $l_{0}$ as follows: Let $S_{0}=0$ and $S_{1}=S$. On $\omega \in \mathcal{Q}_{i-1}$ such that $\left.S_{k}\right|_{\omega}=i$, we define $S_{k+1}(z)$ to be the smallest $j>i$ such that $T^{j}\left(\mathcal{Q}_{j-1}(z)\right)$ is free and has length $>\delta$.

Lemma 9.2. There exists $K_{3}$ such that the following holds for all $k \geq 0$ and $n>K_{3} \log \frac{1}{\delta}$ : If $\omega \in \mathcal{Q}_{i-1}$ is such that $\left.S_{k}\right|_{\omega}=i$, then

$$
m_{0}\left(\omega \cap\left\{\left(S_{k+1}-S_{k}\right)>n\right\}\right)<e^{-K_{3}^{-1} n} m_{0}(\omega)
$$

Proof. This lemma corresponds to Lemma 2.5 and Corollary 2.1 in Section 2.3. The proofs in Appendix A. 2 for Lemma 2.5 and Corollary 2.1 depend only on (P1)-(P3), and they continue to hold as (P1)-(P3) are replaced by ( $\mathrm{P} 1^{\prime}$ )-( $\mathrm{P} 3^{\prime}$ ).

The next lemma locates sites suitable for the construction of $\mathcal{N}$.
Lemma 9.3. There exist an interval $\tilde{L} \subset I$, a number $\tilde{c}>0$, a sequence of integers $n_{1}<n_{2}<\cdots$, and a collection of segments $\left\{\tilde{\omega}_{j}^{(i)}\right\}$ of $l_{0}$ such that
(i) $\pi_{x}\left(T^{i}\left(\tilde{\omega}_{j}^{(i)}\right)\right)=\tilde{L}$;
(ii) $\frac{1}{n_{k}} \sum_{i=0}^{n_{k}-1} m_{0}\left(\cup_{j} \tilde{\omega}_{j}^{(i)}\right) \geq \tilde{c} m_{0}\left(l_{0}\right)$.

Proof. (1) Estimate from below of the total measure of $\omega \in \mathcal{Q}_{i-1}$ with $\left|T^{i}(\omega)\right|>\delta$. Let $\mathcal{R}_{i k}=\left\{\omega \in \mathcal{Q}_{i-1}:\left.S_{k}\right|_{\omega}=i\right\}$. By Lemma 9.2, there exists $K^{\prime \prime}$ such that

$$
\int_{\omega}\left(S_{k+1}-S_{k}\right) d m_{0} \leq K^{\prime \prime} m_{0}(\omega) \quad \text { for all } \omega \in \mathcal{R}_{i k}
$$

Writing $S_{n}=\sum_{k=0}^{n-1}\left(S_{k+1}-S_{k}\right)$ and summing over all $\omega \in \cup_{i} \mathcal{R}_{i k}$ for each $k$, we obtain

$$
\begin{equation*}
\int_{l_{0}} S_{n} d m_{0} \leq K^{\prime \prime} n \cdot m_{0}\left(l_{0}\right) \tag{11}
\end{equation*}
$$

Let $N$ be a large integer. Applying Chebychev's Inequality to (11), we obtain $m_{0}\left\{S_{\left[\frac{1}{2 K^{\prime \prime}} N\right]}>N\right\} \leq \frac{1}{N} \int S_{\left[\frac{1}{2 K^{\prime \prime}} N\right]} d m_{0} \leq \frac{1}{N}\left(K^{\prime \prime} \frac{1}{2 K^{\prime \prime}} N\right) m_{0}\left(l_{0}\right)=\frac{1}{2} m_{0}\left(l_{0}\right)$.
This implies

$$
\begin{equation*}
\sum_{i \leq N} m_{0}\left(\cup\left\{\omega \in \cup_{k} \mathcal{R}_{i k}\right\}\right) \geq \frac{1}{4 K^{\prime \prime}} N m_{0}\left(l_{0}\right) . \tag{12}
\end{equation*}
$$

(2) Selection of $\tilde{L}$. We partition $I$ into intervals $L_{1}, L_{2}, \cdots, L_{\frac{3}{\delta}}$ of length $\frac{1}{3} \delta$ each. For $\omega \in \mathcal{R}_{i k}$, since $\left|T^{i}(\omega)\right|>\delta$, there exists $q=\psi(\omega)$ such that
$\pi_{x}\left(T^{i}(\omega)\right) \supset L_{q}$. Let $\hat{\omega}=\omega \cap T^{-i}\left(\pi_{x}^{-1}\left(L_{q}\right)\right)$. By ( $\left.\mathrm{P} 3^{\prime}\right)$, there exists $K^{\prime \prime \prime}$ independent of $\omega$ such that $m_{0}(\hat{\omega})>K^{\prime \prime \prime}-1 m_{0}(\omega)$. Together with (12), this implies that for each $N$, there exists $q(N)$ such that

$$
\frac{1}{N} \sum_{i \leq N} m_{0}\left(\cup\left\{\hat{\omega}: \omega \in \cup_{k} \mathcal{R}_{i k}, \psi(\omega)=q\right\}\right) \geq \frac{\delta}{12 K^{\prime \prime} K^{\prime \prime \prime}} m_{0}\left(l_{0}\right) .
$$

Let $\tilde{L}=L_{q}$ where $q=q(N)$ for infinitely many $N$. For each $i$, then, the collection $\left\{\tilde{\omega}_{j}^{(i)}\right\}$ is $\left\{\hat{\omega}: \omega \in \cup_{k} \mathcal{R}_{i k}, \psi(\omega)=q\right\}$.

### 9.4. Completing the proof of Proposition 9.1

As explained in Section 9.1, it suffices to verify Condition (S). Let $l_{0}$ be as in Section 9.2. Using the notation in Lemma 9.3, we let $L$ be the middle $\frac{9}{10}$ of $\tilde{L}$, and let

$$
\mathcal{N}:=\operatorname{closure}\left\{\left(\cup_{i, j>0} T^{i}\left(\tilde{\omega}_{j}^{(i)}\right)\right) \cap\left(L \times D_{m-1}\right)\right\} .
$$

It remains to show that $\mathcal{N}$ is a continuous family of unstable curves.
Lemma 9.4. Let $M \in \mathcal{T}_{n}$ be such that $M \cap \mathcal{N} \neq \emptyset$. Then
(i) there exists $\tilde{\omega}_{j}^{(i)}$ such that $T^{i}\left(\tilde{\omega}_{j}^{(i)}\right) \subset M^{\circ}$;
(ii) $M \cap\left(L \times D_{m-1}\right)$ is a horizontal section.

Proof. (i) Since $\mathcal{N} \subset \Omega \subset \operatorname{int}\left(R_{n}\right), M \cap \mathcal{N} \neq \emptyset \Longrightarrow M \cap \gamma \neq \emptyset$ for some $\gamma=T^{i}\left(\tilde{\omega}_{j}^{(i)}\right)$. It remains to show that $\gamma$ cannot meet $E \cup E^{\prime}$, the ends of $M$. This is because all points in $E \cup E^{\prime}$ are in bound state, while $\gamma$ is free. (We need to know $\xi_{i} \in \gamma$ is free when we view the orbit as starting from $\xi_{i-n}$. For $n>i$, this is true because for $\xi_{0} \in l_{0}, d_{\mathcal{C}}\left(\xi_{j}\right)>\frac{1}{2} \delta$ for all $j<0$.) (ii) By Lemma 8.1, there is, through each $\xi \in \gamma$, a codimension-one stable manifold $V(\xi):=T^{n}\left(W_{n}^{s}\left(T^{-n} \xi\right)\right)$. Each $V(\xi)$ has diameter $<K b^{\frac{n}{5}}$ and spans the crosssection of $M$, i.e. $\cup_{\xi \in \gamma} V(\xi)$ is a section of $R_{n}$. All points in this section are free, so it is a horizontal section of length $\approx \frac{1}{3} \delta$, of which roughly the middle $\frac{9}{10}$ is occupied by $M \cap\left(L \times D_{m-1}\right)$.

Verification of Condition (S). Construction of $\left\{\mathcal{E}_{i}\right\}$. Let

$$
\mathcal{A}_{n}=\left\{M \cap\left(L \times D_{m-1}\right): M \in \mathcal{T}_{n}, M \cap \mathcal{N} \neq \emptyset\right\}
$$

By Lemma 9.4, elements of $\mathcal{A}_{n}$ are horizontal sections. We give an algorithm below that selects, for each $n$, a cover $\mathcal{E}_{n}^{\prime}$ of $\mathcal{N}$ by a finite number of pairwise disjoint elements of $\cup_{j \geq n} \mathcal{A}_{j}$. We then let $\mathcal{E}_{n}=\left\{E \cap \mathcal{N}: E \in \mathcal{E}_{n}^{\prime}\right\}$.

Let $\mathcal{E}_{1}^{\prime}=\mathcal{A}_{1}$, and assume $\mathcal{E}_{i-1}^{\prime}$ is constructed. To construct $\mathcal{E}_{i}^{\prime}$, we first put all $E \in \mathcal{E}_{i-1}^{\prime}$ of generation $\geq i$ in $\mathcal{E}_{i}^{\prime}$. Each $E \in \mathcal{E}_{i-1}^{\prime}$ of generation $i-1$ is then replaced systematically by elements of $\left\{A \in \mathcal{A}_{j}, j \geq i, A \subset E\right\}$ as follows: first pick all $F \in \mathcal{A}_{i}$, then pick all $F \in \mathcal{A}_{i+1}$ that cover some points in $\mathcal{N} \cap E$ not
yet covered, then pick all $F \in \mathcal{A}_{i+2}$ covering some points not covered before, and so on. Notice at each stage that the branches chosen are pairwise disjoint. Moreover, the process stops in finite time, for every $z \in \mathcal{N}$ lies in some $M \in \mathcal{T}$.

Properties of $\mathcal{E}_{\infty}=\bigvee_{n>0} \mathcal{E}_{n}$. First we show that the elements of $\mathcal{E}_{\infty}$ form a continuous family of $C^{1}$ curves. Since every $E \in \mathcal{E}_{\infty}$ is the nested intersection of a sequence of horizontal sections whose cross-sectional diameters tend to zero, it is the graph of a function $\varphi: L \rightarrow D_{m-1}$. By Lemma 9.4, $\varphi$ is the pointwise limit of a sequence of functions $\varphi_{k}$, the graph of each one of which is contained in $T^{i}\left(\tilde{\omega}_{j}^{(i)}\right)$ for some $\tilde{\omega}_{j}^{(i)}$. Since $T^{i}\left(\tilde{\omega}_{j}^{(i)}\right)$ is a $C^{2}(b)$ curve, $\left|\varphi_{k}^{\prime \prime}\right|$ is uniformly bounded for all $k$; therefore a subsequence $\varphi_{k_{i}}$ converges to $\varphi$ in the $C^{1}$ norm.

To see that the curves in $\mathcal{E}_{\infty}$ are unstable curves, we use the fact that $T^{i}\left(\tilde{\omega}_{j}^{(i)}\right)$ are unstable curves (Lemma 5.5). The uniform derivative estimates along these curves in backward time are passed to the graph of $\varphi$, and the distortion estimate in (S)(ii) is verified similarly.

Finally, (S)(iii) is given by Lemma 9.3. This completes the verification of Condition (S) and the proof of Proposition 9.1.

## PART III. PARAMETER ISSUES

Let $\mathcal{G}=\cap_{n \geq 0} \mathcal{G}_{n}$. The purpose of Part III is to prove the existence and abundance of maps in $\mathcal{G}$. More specifically, we will prove that for 1-parameter families $T_{a}: X \rightarrow X$ satisfying the Standing Hypotheses in Section 1, the set $\left\{a: T_{a} \in \mathcal{G}\right\}$ has positive Lebesgue measure. Our plan is to construct a set $\Delta \subset\left\{a: T_{a} \in \mathcal{G}\right\}$ with a generalized Cantor structure in which the gap ratios tend to zero exponentially fast.

We cannot overemphasize the dependence of Part III on earlier sections. Results from Part II on properties of individual $T \in \mathcal{G}$ are clearly relevant as we now seek to identify such maps from a given 1-parameter family. Since the criteria for belonging in $\mathcal{G}$ reside with the behavior of critical orbits, a major focus of the present study is on the evolution of critical curves, i.e. curves of the form $a \mapsto z_{i}(a), i=0,1,2, \cdots$, where $z_{0}$ is a critical point. We will show that $a \mapsto z_{i}(a)$ define processes that have a great deal in common with the 1D maps studied in Section 2. Part of the analysis involves adapting the results of Section 2 to the present context.

Each of the first three sections of Part III discusses one important aspect of the problem. These ideas culminate in Section 13, which contains the actual construction of $\Delta$.

Hypotheses for Part III. We assume
(1) $\left\{T_{a}, a \in\left[a_{0}, a_{1}\right]\right\}$ satisfies the Standing Hypotheses in Section 1;
(2) $\left[a_{0}, a_{1}\right]$ is in a sufficiently small neighborhood of $a^{*}$ that $T_{a}$ satisfies the hypotheses at the beginning of Part II for all $a \in\left[a_{0}, a_{1}\right]$.
The generic constant $K$ here depends on the family $T_{a}$ as well as our choice of $\lambda$.

## 10. Dependence of dynamical structures on parameter

Notation such as $\mathcal{T}_{k}(a)$ and $\mathcal{C}^{(k)}(a)$ are used to indicate dependence on the $\operatorname{map} T_{a}$.

### 10.1. Continuation of critical regions and critical points

Definition 10.1. Let $J \subset\left[a_{0}, a_{1}\right]$ be an interval, and assume that for some $\hat{a} \in J, T_{\hat{a}} \in \mathcal{G}_{n}$. We say $\left\{T_{a}, a \in J\right\}$ is a continuation of $T_{\hat{a}}$ in $\mathcal{G}_{n}$ if the following hold:
(1) For all $a \in J, T_{a} \in \mathcal{G}_{n}$, and there is a choice of $\Gamma_{\theta n}(a)$ with the following properties:
(2) The monotone branches of $T_{\hat{a}}$ of generation $\leq \theta n$ deform continuously on $J$; i.e., for each $k \leq \theta n$, there is a map $\Phi_{k}$ defined on $J \times \mathcal{T}_{k}(\hat{a})$ such that
(i) for each fixed $a, M \mapsto \Phi_{k}(a, M)$ is a bijection between $\mathcal{T}_{k}(\hat{a})$ and $\mathcal{T}_{k}(a)$;
(ii) for each fixed $M, a \mapsto \Phi_{k}(a, M)$ is continuous (in the Hausdorff metric).
(3) The critical regions of $T_{\hat{a}}$ of generation $\leq \theta n$ deform continuously on $J$; i.e., for each $k \leq \theta n$, there is a map $\Psi_{k}$ defined on $J \times\left\{Q^{(k)}(\hat{a})\right\}$ such that
(i) for each fixed $a, Q \mapsto \Psi_{k}(a, Q)$ is a bijection between $\left\{Q^{(k)}(\hat{a})\right\}$ and $\left\{Q^{(k)}(a)\right\}$;
(ii) for each fixed $Q, a \mapsto \Psi_{k}(a, Q)$ is continuous.
(4) The critical points of $T_{\hat{a}}$ of generation $\leq \theta n$ continue smoothly to all of $J$; i.e., for each $z_{0}(\hat{a})=z_{0}^{*}\left(Q^{(k)}(\hat{a})\right), k \leq \theta n, a \mapsto z_{0}(a)$ is a $C^{2}$ curve satisfying
(i) $z_{0}(a)=z_{0}^{*}\left(Q^{(k)}(a)\right)$ where $Q^{(k)}(a)=\Psi\left(a, Q^{(k)}(\hat{a})\right)$;
(ii) if $\xi_{1}(\hat{a})=T_{\hat{a}}^{-(k-1)} z_{0}(\hat{a})$ and $l$ is the $\mathcal{F}_{1}$-leaf containing $\xi_{1}(\hat{a}),{ }^{10}$ then there is a $C^{2}$-function $\xi_{1}: J \rightarrow l$ such that $z_{0}(a)=T_{a}^{k-1}\left(\xi_{1}(a)\right)$.

[^10]We refer to $a \mapsto \Gamma_{\theta n}(a)$ with property (4) as a coherent choice of $\Gamma_{\theta n}(a) .{ }^{11}$
Observe that if $\left\{T_{a}, a \in J\right\}$ is a continuation of $T_{\hat{a}}$, then (i) $\left\{\Phi_{k}\right\}$ "commutes" with the actions of $T_{a}$; i.e., if $M \in \mathcal{T}_{k}(\hat{a})$ is such that $T_{\hat{a}} M=\cup_{i=1}^{s} M_{i}$, $M_{i} \in \mathcal{T}_{k+1}(\hat{a})$, then $T_{a} \Phi_{k}(a, M)=\cup_{i=1}^{s} \Phi_{k}\left(a, M_{i}\right)$; and (ii) the partial order on $\left\{Q^{(k)}, k \leq \theta n\right\}$ is respected by $\left\{\Psi_{k}\right\}$. The validity of these statements is easily seen by comparing two nearby $a$.

Our first goal is to give sufficient conditions for the existence of continuations. To ensure that a nontrivial continuation exists, we choose $T_{\hat{a}}$ in the "interior" of $\mathcal{G}_{N}$. Let

$$
\mathcal{G}_{N}^{\#}=\left\{T \in \mathcal{G}_{N}: T \text { satisfies (A2) }{ }^{\#} \text { and (A4) }{ }^{\#}\right\}
$$

where (A2) ${ }^{\#}$ and (A4) $\#$ below require that for all $z_{0} \in \Gamma_{\theta N}$ of generation $k$, the following hold for all $i \leq k \theta^{-1}$ :

$$
\begin{aligned}
& \text { (A2) }{ }^{\#} \quad d_{\mathcal{C}}\left(z_{i}\right)>\min \left(\delta, 2 e^{-\alpha i}\right) ; \\
& \text { (A4) }{ }^{\#}\left|w_{i}^{*}\right|>c_{2} e^{\lambda^{*} i} \text { where } \lambda^{*}=\lambda+\frac{1}{100} \lambda_{0} .
\end{aligned}
$$

Clearly, $\mathcal{G}_{N}^{\#} \subset \mathcal{G}_{N}$.
Proposition 10.1 (Dynamical continuation). There exists $\rho>0$ depending only on $\left\|T_{a}\right\|_{C^{3}}$ and $c_{2}$ for which the following holds: Assume $T_{\hat{a}} \in \mathcal{G}_{N}^{\#}$. For $n \leq N$, let $J_{n}=\left[\hat{a}-\rho^{n}, \hat{a}+\rho^{n}\right]$. Then $\left\{T_{a}, a \in J_{n} \cap\left[a_{0}, a_{1}\right]\right\}$ is a continuation of $T_{\hat{a}}$ in $\mathcal{G}_{n}$.

Proof. We assume the following hold for $n<N$ and prove it for $n+1$ :
(i) $\left\{T_{a}, a \in J_{n}\right\}$ is a continuation of $T_{\hat{a}}$ in $\mathcal{G}_{n}$;
(ii) (a priori estimate on $\left.\left|\frac{d}{d a} z_{0}(a)\right|\right)$ there is a constant $K>0$ independent of $\rho$ such that for all $a \in J_{n}$, if $z_{0}(a) \in \Gamma_{[\theta n]}(a)$ is of generation $j$, then

$$
\left|\frac{d}{d a} z_{0}(a)\right|<K^{j} .
$$

There is nothing to do if $k \theta^{-1}<n+1<(k+1) \theta^{-1}$ : no new monotone branches or critical regions are constructed, and all critical points of generation $\leq k$ are treated in the previous step. We assume therefore that $n+1=(k+1) \theta^{-1}$ for some $k$.

1. Coherent choice of $\Gamma_{k+1}(a)$, construction of $\mathcal{C}^{(k+1)}(a)$ and $\mathcal{T}_{k+1}(a)$, and verification of (2) and (3) in Definition 10.1 for objects of generation $k+1$. For each individual $a \in J_{n}$, since $T=T_{a} \in \mathcal{G}_{n}$, we know by Proposition 8.1 that $\mathcal{C}^{(k+1)}$ and $\mathcal{T}_{k+1}$ can be constructed. Moreover, for each $M \in \mathcal{T}^{(k)}$ and

[^11]$Q=Q^{(j)},(k+1)(1+2 \theta)^{-1} \leq j \leq k$, there is the following dichotomy: either $T M \cap Q=\emptyset$, or a horizontal section of $T M$ pierces through the entire length of $Q$.

Now let $a$ vary over $J_{n}$. For $M \in \mathcal{T}_{k}(\hat{a})$, we know by inductive assumption that $M(a):=\Phi_{k}(a, M)$ varies continuously with $a$, as does $Q(a):=\Psi_{j}(Q)$. Since the dichotomy above holds for all $a$ - and there is no way to go from one scenario to the other in a continuous manner - it follows that exactly one of the two scenarios must prevail for all $a \in J_{n}$. Indeed, the number of times $T_{a} M(a)$ goes through $Q(a)$ is constant for all $a$. This proves properties (2)(i) and (3)(i) in Definition 10.1.

Suppose for $M$ and $Q$ as above, $T_{\hat{a}}$ has a critical point $\hat{z}_{0}(a) \in T_{\hat{a}} M(\hat{a}) \cap$ $Q(\hat{a})$. Let $l_{k}(\hat{a})$ be the connected component of $\mathcal{F}_{k}$-leave in $M(\hat{a})$ on which $T_{\hat{a}}^{-1} z_{0}(\hat{a})$ is located. Let $l=T_{\hat{a}}^{-k+1} l_{k}$. Then $l$ is a leaf of $\mathcal{F}_{1}$, and $\mathcal{F}_{1}$ does not depend on $a$. From the last paragraph, we know that $T_{a}^{k}(l)$ pierces through $Q(a)$ for all $a$. Let $z_{0}(a)$ be constructed on $T_{a}^{k}(l)$. This construction guarantees the continuity of $a \mapsto z_{0}(a)$ for all critical points of generation $k+1$ and consequently properties (2)(ii) and (3)(ii) in Definition 10.1.
2. Smoothness of $a \mapsto z_{0}(a)$ and estimate on $\left|\frac{d}{d a} z_{0}(a)\right|$ for $z_{0}$ of generation $k+1$. Continuing to use notation from the last paragraph, we let $x \mapsto$ $\gamma(x, a)=(x, \psi(x, a))$ be the curve $T_{a}^{k} l$ in $Q^{(k+1)}(a)$, and let $z_{0}(a)=(\bar{x}(a), \bar{y}(a))$. For each $(x, a)$, we consider in $X_{\gamma(x, a)}$ the 2D plane $S=S\left(\partial_{x} \gamma(x, a)\right.$, v) with orthonormal basis $\{u, v\}$ where $u=\partial_{x} \gamma /\left|\partial_{x} \gamma\right|$ and $v$ points in roughly the same direction as $\mathbf{v}$. Let $e_{k+1}$ be the most contracted direction of $D T_{a}^{k+1}$ in $S$. As in Section 3.6, we write $e_{k+1}$ as a linear combination of $u$ and $v$, and let $\eta_{k+1}$ denote its $v$-component. Then $\bar{x}(a)$ is defined implicitly by $\eta_{k+1}(\bar{x}(a), a)=0$, and therefore is $C^{2}$ as a function of $a$. Likewise, $\bar{y}(a)=\psi(\bar{x}(a), a)$ is a $C^{2}$ function of $a$.

The following lemma is proved in Appendix A.23.
Lemma 10.1. As functions of $x$ and $a$,
(a) $\|u\|_{C^{2}},\|v\|_{C^{2}}<K^{k+1}$;
(b) $\left\|\eta_{k+1}\right\|_{C^{2}}<K^{k+1}$.

Corollary 10.1.

$$
\begin{equation*}
\left|\frac{d z_{0}(a)}{d a}\right|,\left|\frac{d^{2} z_{0}(a)}{d a^{2}}\right| \leq K^{k+1} . \tag{13}
\end{equation*}
$$

Proof of Corollary 10.1. Differentiating $\eta_{k+1}(\bar{x}(a), a)=0$, we obtain

$$
\begin{equation*}
\frac{d \bar{x}(a)}{d a}=-\frac{\partial_{a} \eta_{k+1}}{\partial_{x} \eta_{k+1}}(\bar{x}(a), a) . \tag{14}
\end{equation*}
$$

Observe that $\left|\partial_{x} \eta_{k+1}\right|>K^{-1}$ : This follows from Lemma 3.7 and the fact that derivative growth along the orbit of $z_{0}^{*}\left(Q^{(k)}(a)\right)$ is passed on to that of $z_{0}(a)$
via Lemma 3.2. Our claim on the first derivative follows directly from Lemma 10.1 (b) and the fact that $\frac{d \bar{y}}{d a}=\partial_{x} \psi \frac{d \bar{x}}{d a}+\partial_{a} \psi$. To estimate the second derivative we differentiate (14) one more time with respect to $a$, and use again Lemma 10.1.
3. Proofs of $(\mathrm{A} 2)(n+1)$ and $(\mathrm{A} 4)(n+1)$. Let $z_{0}$ be a critical point of generation $k+1$. We give details only for step $n+1$ : By Corollary 10.1, $\left|z_{0}(a)-z_{0}(\hat{a})\right| \leq K^{k+1}\left(2 \rho^{n+1}\right)$ for all $a \in J_{n+1}$, so that

$$
\begin{equation*}
\left|z_{n+1}(a)-z_{n+1}(\hat{a})\right|<K\|D T\|^{n+1}\left|z_{0}(a)-z_{0}(\hat{a})\right| \ll e^{-\alpha(n+1)} \tag{15}
\end{equation*}
$$

provided $\rho$ is sufficiently small relative to $\|D T\|^{-1}$ and $K^{-1}$. To finish, we need to deal with the differences between $d_{\mathcal{C}(a)}(\cdot)$ and $d_{\mathcal{C}(\hat{a})}(\cdot)$. Suppose $z_{n+1}(\hat{a}) \in$ $\mathcal{C}^{(1)}$, and $\phi\left(z_{n+1}(\hat{a})\right)=z_{0}^{*}\left(Q^{(j)}(\hat{a})\right)$. Then $j \leq \alpha^{*}(n+1) \theta, z_{n+1}(\hat{a}) \in Q^{(j)}(\hat{a})$, and $T_{\hat{a}}$ has a horizontal section $H(\hat{a})$ extending considerably beyond $Q^{(j)}(\hat{a})$ on both sides. We conclude from the continuous deformation of structures of generation $j$, the estimate $\left|z_{0}^{*}\left(Q^{(j)}(\hat{a})\right)-z_{0}^{*}\left(Q^{(j)}(a)\right)\right| \leq K^{j}|\hat{a}-a|$ and (15) above that $z_{n+1}(a)$ is either in $Q^{(j)}(a)$ or it is in $H(a)$ and just outside of $Q^{(j)}(a)$, and $d_{\mathcal{C}_{a}}\left(z_{n+1}(a)\right)>\frac{1}{2} d_{\mathcal{C}_{\hat{a}}}\left(z_{n+1}(\hat{a})\right)>e^{-\alpha(n+1)}$.

To prove (A4)(n+1), we first convert the problem to one involving $\left|w_{i}\right|$, thereby picking up some factors of $e^{\alpha i}$. The comparability of $\left|w_{i}(a)\right|$ and $\left|w_{i}(\hat{a})\right|$ is given by the following lemma, the proof of which follows closely that of Lemma 3.2 and is omitted.

Lemma 10.2. Let $z_{0}(\hat{a})$ be of generation $j \leq \theta(n+1)$. For $i \leq n+1$, let $w_{i}(\hat{a})=\left(D T_{\hat{a}}^{i}\right)_{z_{0}(\hat{a})} \mathbf{v}$, and $w_{i}(a)=\left(D T_{a}^{i}\right)_{z_{0}(a)} \mathbf{v}, a \in J_{n+1}$. Then $\left|w_{i}(a)\right| \geq$ $\frac{1}{2}\left|w_{i}(\hat{a})\right|$.

This completes the proof of (i) and (ii) for step $n+1$.
10.2. Properties of $a \mapsto z_{0}(a), z_{0} \in \cup \Gamma_{j}$

Unlike the situation in 1D, $\cup_{j \geq 1} \Gamma_{j}$ is an infinite set, and the domains of definition of $a \mapsto z_{0}(a)$ decrease as the generation of $z_{0}$ increases. For $T_{\hat{a}} \in \mathcal{G}_{N}^{\#}$ and $z_{0} \in \Gamma_{\theta N}(\hat{a})$ of generation $\theta n, n \leq N$, we guarantee the continuation of $z_{0}(\hat{a})$ only to the interval $J_{n}=\left[\hat{a}-\rho^{n}, \hat{a}+\rho^{n}\right]$. We claim, however, that there is a uniform bound on $\frac{d}{d a} z_{0}(a)$ that is valid for all $z_{0}$ (independent of generation) on their respective intervals of continuation.

Lemma 10.3. Let $\hat{a}, J_{n}$ and $\Gamma_{\theta n}$ be as in Proposition 10.1, and let $a \mapsto$ $z_{0}^{(k)}(a), a \in J_{n}$, be a curve of critical points of generation $k \leq[\theta n]$. Then there is $z_{0}^{\left(k^{\prime}\right)}$ of generation $k^{\prime}, k^{\prime}<k \leq k^{\prime}(1+2 \theta)$, with $z_{0}^{(k)} \in Q^{\left(k^{\prime}\right)}\left(z_{0}^{\left(k^{\prime}\right)}\right)$ such that

$$
\left|\frac{d}{d a}\left(z_{0}^{(k)}(a)-z_{0}^{\left(k^{\prime}\right)}(a)\right)\right|<b^{\frac{k^{\prime}}{9}} .
$$

A proof of Lemma 10.3 is given in Appendix A.23.

Corollary 10.2. Under the hypotheses of Lemma 10.3, there exists $K_{1}$ such that for every curve of critical points $a \mapsto z_{0}^{(k)}(a), k \leq \theta n$, if $z_{0}^{(k)}=$ $\left(x_{0}, y_{0}\right)$, then

$$
\left|\frac{d}{d a} x_{0}(a)\right| \leq K_{1}, \quad\left|\frac{d}{d a} y_{0}(a)\right| \leq b^{\frac{1}{10}} .
$$

Proof. Let $z_{0}=z_{0}^{*}\left(Q^{(k)}\right)$, and suppose $Q^{(k)}=Q^{\left(k_{0}\right)} \subset Q^{\left(k_{1}\right)} \subset \cdots \subset$ $Q^{(1)}$ are consecutive critical regions. We obtain, by comparing $a \mapsto z_{0}^{\left(k_{i}\right)}$ and $a \mapsto z_{0}^{\left(k_{i-1}\right)}$, that $\left|\frac{d}{d a}\left(z_{0}^{(k)}(a)-z_{0}^{(1)}(a)\right)\right|<b^{\frac{1}{10}}$. The assertions in this corollary now follow immediately from properties of the finitely many critical points $z_{0}^{(1)}=\left(x_{0}^{(1)}, y_{0}^{(1)}\right)$ of generation one, namely $\frac{d}{d a} y_{0}^{(1)}=0$ and $\left|\frac{d}{d a} x_{0}^{(1)}\right|<\frac{1}{2} K_{1}$ for some $K_{1}$.

Remark. Lemma 10.3 and Corollary 10.2 together imply the following: (1) The speeds of movement of all critical points are uniformly bounded. (2) While $\Gamma_{n}$ as a whole moves with speed $\mathcal{O}(1)$, the relative speed of motion of $z_{0}^{\prime}$ and $z_{0}$ for $z_{0}^{\prime} \in Q^{(k)}\left(z_{0}\right)$ decreases exponentially fast with $k$.

### 10.3. Setting for the analysis to follow

Recall that for a single map $T \in \mathcal{G}_{N}$, whether or not $T$ is in $\mathcal{G}_{\frac{1}{\alpha^{*}} N}$ is determined by whether (A2) and (A4) are satisfied up to time $\frac{1}{\alpha^{*}} N$. We now consider a family $\left\{T_{a}, a \in J\right\}$ where $T_{a} \in \mathcal{G}_{N}$ for all $a$. In addition to asking whether $T_{a} \in \mathcal{G}_{\frac{1}{\sigma^{*}} N}$ for each individual $a$, we will also want to know for what fraction of $a \in \mathcal{J}^{*}$ is $T_{a} \in \mathcal{G}_{\frac{1}{\alpha^{*}} N}$.

This leads us to study the evolution of $\zeta_{i}: a \mapsto z_{i}(a)$ where $\zeta_{0}$ is a coherent choice of $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N}$. The analysis is highly inductive: For each $T_{a}$, critical points are defined inductively, and the presence of certain structures is needed to track their orbits. We wish now to track not single orbits but entire curves. Precise conditions under which this analysis will be carried out are as follows:

Assumptions in the inductive analysis of critical curves. Let $J \subset\left[a_{0}, a_{1}\right]$ be a parameter interval.
(C1) $\left\{T_{a}, a \in J\right\}$ is a continuation (of some $T_{\hat{a}}$ ) in $\mathcal{G}_{N}$.
(C2) A coherent choice of $\Gamma_{\frac{1}{\alpha^{*} \theta N}}(a), a \in J$, has been made.
Clarification 1. (C1) and (C2) are related as follows: If, in addition to (C1), we have $T_{\hat{a}} \in \mathcal{G}_{N}^{\#}$ for some $\hat{a} \in J$, then steps 1 and 2 in Proposition 10.1 can be carried out for critical points of generation $k$ for all $k \leq \frac{1}{\alpha^{*}} \theta N$. In Sections 11 and 12 we are not concerned with how (C2) comes about, but note that once a coherent choice of $\Gamma_{\frac{1}{\alpha^{*}} \theta N}(a)$ is made, the estimates in Lemma 10.3 and Corollary 10.2 are valid for the critical points in question. Justifications for this claim follow verbatim those in Sections 10.1 and 10.2.
2. Nothing is assumed or claimed at this point about the behavior of critical orbits beyond time $N$. That is the objective of the investigation in the pages to follow. More precisely, we will be concerned with the evolution of

$$
\zeta_{i}: a \mapsto z_{i}(a) \quad \text { for } \quad a \in J, \quad z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N} \backslash \Gamma_{\theta N}, \quad i \leq \frac{1}{\alpha^{*}} N
$$

with particular interest in the time range $N<i \leq \frac{1}{\alpha^{*}} N$.
3. As we will see, this analysis requires that all critical structures of generation $\leq \theta N$ exist and vary with parameter in a certain way. This is provided by (C1) and Corollary 10.2. Also, critical structures beyond generation $\theta N$ are not relevant for this analysis.

## 11. Dynamics of curves of critical points

The aim of this section is to bring to light a certain resemblance between the evolution of $\zeta_{i}$ and that of certain "horizontal" curves under the iteration of $T_{a}$ for a fixed $a$. The reason behind this resemblance is that $\left|\frac{d}{d a} \zeta_{i}\right|$ is comparable to $\left|D T^{i}(\mathbf{v})\right|$. We will show that for as long as $T_{a} \in \mathcal{G}_{N}$, this comparability selfperpetuates once we get it going, and the start-up mechanism is provided by the parameter transversality condition in the Standing Hypotheses in Section 1. This is discussed in Section 11.1. Basic properties in the evolution of $\zeta_{i}$, such as bound and free periods, are discussed in Section 11.2.

Conditions (C1) and (C2) in Section 10.3 are assumed throughout.

### 11.1. Equivalence of space- and $a$-derivatives

We use the notation $z_{i}(a)=T_{a}^{i}\left(z_{0}(a)\right), w_{i}(a)=\left(D T_{a}^{i}\right)_{z_{0}(a)}(\mathbf{v})$ and $\tau_{i}(a)=$ $\frac{d}{d a} z_{i}(a)$.

Proposition 11.1. There exist $\hat{K}>1$ and $i_{0} \in \mathbb{Z}^{+}$(both depending only on $\left.\left\{F_{a}\right\}\right)$, such that the following holds for all $(a, b)$ sufficiently close to $\left(a^{*}, 0\right)$ : Fix $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N} \backslash \Gamma_{\theta N}$. We assume that for some $n \leq \frac{1}{\alpha^{*}} N$, $d_{\mathcal{C}}\left(z_{i}(a)\right) \geq$ $\min \left(\delta, e^{-\alpha i}\right)$ for all $a \in J$ and $i<n$. Then

$$
\hat{K}^{-1} \leq \frac{\left|\tau_{i}\right|}{\left|w_{i}\right|} \leq \hat{K} \quad \text { for all } \quad i_{0}<i \leq n .
$$

The required proximity of $(a, b)$ to $\left(a^{*}, 0\right)$ depends also on $i_{0}$, and under the conditions above, the pair $\left(z_{0}, w_{0}\right)$ is controlled by $\Gamma_{\theta N}$ up to time $n$ (see Proposition 6.2).

Proposition 11.1 is a consequence of Standing Hypothesis (b) in Section 1 and can be viewed as the higher dimensional analog of Proposition 2.3 in

Section 2.4. Recall from Proposition 2.3 that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \sum_{s=1}^{k} \frac{\frac{d}{d a}\left(f_{a}\left(x_{s-1}\right)\right)\left(a^{*}\right)}{\left(f^{s-1}\right)^{\prime}\left(x_{1}\left(a^{*}\right)\right)}=\left[\frac{d}{d a} f_{a}(\hat{x}(a))-\frac{d}{d a} q(a)\right]_{a=a^{*}}:=\hat{c} \tag{16}
\end{equation*}
$$

The constant $\hat{K}$ in Proposition 11.1 is derived from $\hat{c}$ together with angle and other considerations.

Proof. Letting $\psi(z)=\frac{\partial}{\partial a}\left(T_{a} z\right)$, we write

$$
\tau_{i}=D T_{z_{i-1}} \tau_{i-1}+\psi\left(z_{i-1}\right)=D T_{z_{0}}^{i} \tau_{0}+\sum_{s=1}^{i} D T_{z_{s}}^{i-s} \psi\left(z_{s-1}\right):=I+I I
$$

where

$$
I=D T_{z_{0}}^{i} \tau_{0}+\sum_{s=1}^{i_{0}} D T_{z_{s}}^{i-s} \psi\left(z_{s-1}\right) \quad \text { and } \quad I I=\sum_{s=i_{0}+1}^{i} D T_{z_{s}}^{i-s} \psi\left(z_{s-1}\right)
$$

$i_{0}$ being a number to be determined. We will show there exist $K_{0}$ (depending only on $\left\{T_{a}\right\}$ ) and $i_{0}$ such that if $(a, b)$ is sufficiently near $\left(a^{*}, 0\right)$, then for $i_{0}<i \leq n$,
$-K_{0}^{-1}|\hat{c}|<\frac{|I|}{\left|w_{i}\right|}<K_{0}|\hat{c}| \quad$ and
$-\frac{|I I|}{\left|w_{i}\right|}$ is as small as we wish.
These estimates together give the desired result.
Estimate on $\frac{|\mathrm{III}|}{\left|w_{i}\right|}$. Since $|\psi(\cdot)|<K$, it follows from Lemma 5.6 that

$$
|\mathrm{II}| \leq K \sum_{s=i_{0}+1}^{i}\left\|D T_{z_{s}}^{i-s}\right\| \leq K \sum_{s=i_{0}+1}^{i} K e^{-\hat{\lambda} s}\left|w_{i}\right| .
$$

Choosing $i_{0}$ large enough, we can make $K^{2} \sum_{s=i_{0}+1}^{\infty} e^{-\hat{\lambda} s} \ll K_{0}^{-1}|\hat{c}|$.
Estimate on $\frac{|I|}{\left|w_{i}\right|}$. Increase $i_{0}$ if necessary so that with $k=i_{0}$, the sum on the left side of (16) is $<\frac{1}{2} \hat{c}$ from its limit; $i_{0}$ is fixed from here on. Let $V$ be such that $I=D T_{z_{i_{0}}}^{i-i_{0}} V$, i.e.

$$
V=D T_{z_{0}}^{i_{0}} \tau_{0}+\sum_{s=1}^{i_{0}} D T_{z_{s}}^{i_{0}-s} \psi\left(z_{s-1}\right)
$$

The verification of Lemma 11.1 is given in Appendix A. 24 .
Lemma 11.1. As $(a, b) \rightarrow\left(a^{*}, 0\right)$,

$$
\frac{\left|w_{1}\right|}{\left|w_{i_{0}}\right|} V \rightarrow\left( \pm \sum_{s=1}^{i_{0}} \frac{\frac{d}{d a}\left(f_{a}\left(x_{s-1}\right)\right)\left(a^{*}\right)}{\left(f^{s-1}\right)^{\prime}\left(x_{1}\left(a^{*}\right)\right)}, 0\right) .
$$

It is important to note that the convergence above is uniform among all critical curves. This is evident from the uniform bound on $\frac{d}{d a} z_{0}$ and the proof of Lemma 11.1.

To finish, we write

$$
\frac{|\mathrm{I}|}{\left|w_{i}\right|}=\frac{\left|D T_{z_{i_{0}}}^{i-i_{0}}(V)\right|}{\left|D T_{z_{i_{0}}}^{i-i_{0}}\left(w_{i_{0}}\right)\right|}=\frac{\left|D T_{z_{i_{0}}}^{i-i_{0}}\left(\frac{V}{|V|}\right)\right|}{\left|D T_{z_{i_{0}}}^{i-i_{0}}\left(\frac{w_{i_{0}}}{\mid w_{i_{0}}}\right)\right|} \cdot \frac{|V|}{\left|w_{i_{0}}\right|}
$$

Notice first that by Lemma 11.1, $\frac{1}{3} \frac{\hat{c}}{\left|w_{1}\right|}<\frac{|V|}{\left|w_{i_{0}}\right|}<2 \frac{\hat{c}}{\left|w_{1}\right|}$. We claim that

$$
\begin{equation*}
K^{-1}<\frac{\left|D T^{i-i_{0}}\left(\frac{V}{|V|}\right)\right|}{\left|D T^{i-i_{0}}\left(\frac{w_{i_{0}}}{\mid w_{i_{0}}}\right)\right|}<K \tag{17}
\end{equation*}
$$

provided $(a, b)$ is sufficiently near $\left(a^{*}, 0\right)$. Assume $\angle\left(V, w_{i_{0}}\right) \neq 0$, and let $e_{j-i_{0}}=$ $e_{j-i_{0}}(S)$ be the most contracted vector of order $j-i_{0}$ where $S=S\left(V, w_{i_{0}}\right)$. To prove (17), it suffices to show that $e_{j-i_{0}}$ is well-defined, $\angle\left(e_{j-i_{0}}, w_{i_{0}}\right)>K^{-1}$, and $\angle\left(V, w_{i_{0}}\right) \ll K^{-1}$. With $(a, b)$ sufficiently near $\left(a^{*}, 0\right)$, we may assume for some $s_{0} \gg i_{0}$ that $d_{\mathcal{C}}\left(z_{s}\right)>\frac{1}{2} \delta_{0}$ for all $s<s_{0}\left(\delta_{0}\right.$ is as in Definition 1.1). This together with our assumption that $d_{\mathcal{C}}\left(z_{j}\right) \geq \min \left(\delta, e^{-\alpha j}\right)$ implies that for all $j>i_{0},\left|w_{j}\right| /\left|w_{i_{0}}\right|>K^{-1} e^{\left(\lambda^{\prime}-2 \alpha\right)\left(j-i_{0}\right)}$, proving $e_{j-i_{0}}$ is well-defined. Since $w_{i_{0}}$ is $b$-horizontal, we have $\angle\left(e_{1}, w_{i_{0}}\right)>K^{-1}$ by Lemma 3.7. This together with $\angle\left(e_{1}, e_{j-i_{0}}\right)<(K b)^{j-i_{0}}$ (Lemma 3.1) gives $\angle\left(e_{j-i_{0}}, w_{i_{0}}\right)>K^{-1}$. As for $V$, Lemma 11.1 tells us its slope is as small as we wish. Hence $\angle\left(V, w_{i_{0}}\right) \ll K^{-1}$.

We remark that all the estimates in the proof above - and hence the constants in the statement of the proposition - are independent of $N$. For as long as both $\tau_{i}$ and $w_{i}$ grow in magnitude, the angles between them must shrink by rank-one arguments. The assumptions in the next lemma are as in Proposition 11.1. A detailed proof is given in Appendix A.24.

LEMMA 11.2. If $z_{i}$ is a free return, then $\angle\left(\tau_{i}, w_{i}\right)<\frac{K}{\left|\tau_{i}\right|}$.
The following are assumed for the rest of this paper: (i) $i_{0}$ is sufficiently large, (ii) $(a, b)$ is sufficiently close to $\left(a^{*}, 0\right)$, and (iii) all critical points stay at distances $>\frac{1}{2} \delta_{0}$ away from $\mathcal{C}$ for $\gg i_{0}$ iterates - where "sufficiently large", "sufficiently close" and " $\gg$ " are as required in the proof of Proposition 11.1.

### 11.2. Resemblance to phase-space dynamics

In addition to ( C 1 ) and $(\mathrm{C} 2)$, we now fix $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N} \backslash \Gamma_{\theta N}$ and impose on it
(C3) For some $i_{0}<n \leq \frac{1}{\alpha^{*}} N, d_{\mathcal{C}}\left(z_{i}\right)>\min \left(\delta, e^{-\alpha i}\right)$ for all $i \leq n$.

We discuss below four aspects of the dynamics of $\zeta_{i}: a \mapsto z_{i}(a), i=$ $0,1,2, \cdots$. The notation is as in Section 11.1; in particular, $w_{i}(a)=\left(D T_{a}^{i}\right)_{z_{0}(a)}(\mathbf{v})$ and $\tau_{i}(a)=\frac{d}{d a} z_{i}(a)$. Recall also that $s(u)=\frac{\left|u_{y}\right|}{\left|u_{x}\right|}$ where $u=\left(u_{x}, u_{y}\right)$ is a vector in $\mathbb{R}^{m}=\mathbb{R} \times \mathbb{R}^{m-1}$.
A. Outside of $\mathcal{C}^{(1)}$. For the first $i_{0}$ iterates, we do not have a great deal of information on $\tau_{i}$. Let $\hat{\varepsilon}:=\hat{K} c_{2} e^{-\frac{1}{4} \lambda_{0} i_{0}}+\mathcal{O}(b)$ where $\hat{K}$ is as in Proposition 11.1. We may assume $\hat{\varepsilon} \ll \delta$.

Lemma 11.3. The following hold for every $a$ :
(a) If $z_{n}$ is free, then $s\left(\tau_{n}\right)<\hat{\varepsilon}$.
(b) If $z_{n}$ is free, and $z_{n+j} \notin \mathcal{C}^{(1)} \forall 0 \leq j<j_{0}$, then
(i) $\left|\tau_{n+j}\right|>K^{-1} \delta e^{\frac{1}{4} \lambda_{0} j}\left|\tau_{n}\right|$ for $j \leq j_{0}$; and
(ii) if in addition $\gamma_{n+j_{0}} \in \mathcal{C}^{(1)}$, then $\left|\tau_{n+j_{0}}\right|>K^{-1} e^{\frac{1}{4} \lambda_{0} j_{0}}\left|\tau_{n}\right|$.

Proof. (a) follows from Lemma 11.2 and the $b$-horizontal property of $w_{n}$. As for (b), since $\left|\tau_{n+j}\right| \gg 1$, we have $\left|\tau_{n+j+1}\right| /\left|\tau_{n+j}\right| \approx\left|f^{\prime}\left(x_{n+j}\right)\right|, z_{i}=\left(x_{i}, y_{i}\right)$. The assertions follow by a proof similar to that of Lemma 3.5.

Remark. We do not claim that free segments of $\zeta_{i}$ are $C^{2}(b)$, only that they are roughly horizontal (because $\hat{\varepsilon} \ll 1$ ). No effort will be made to control $\left(z_{0}, \tau_{0}\right)$. Information on $\tau_{i}$ is obtained instead through comparisons with $w_{i}$ via Proposition 11.1 and Lemma 11.2.
B. Geometry of critical curves inside $Q^{(1)}$. Let $\omega$ be a subinterval of $J$. We assume $\zeta_{n}(\omega)$ is free (meaning $z_{n}(a)$ is free for each $a$ ), and $\zeta_{n}(\omega) \subset Q^{(1)}$. For each individual $a$, we have seen in Part II how $z_{n}(a)$ is related to the critical structure of $T_{a}$. We now describe the geometric relationship between the curve $\zeta_{n}$ and the 1-parameter family of critical structures.

By Lemma 11.3(a), $\zeta_{n}$ is roughly horizontal. Consider an arbitrary point $\hat{a} \in \omega$, and assume that $\zeta_{n}(\hat{a})$ lies in the interior of $Q^{(j)}(\hat{a})$ for some $j \leq n$. We also assume, for definiteness, that as $a$ increases, we move right along $\zeta_{n}$. By continuity, for all $a$ in a neighborhood of $\hat{a}, \zeta_{n}(a)$ also lies in the interior of $Q^{(j)}(a)$. Let $\bar{a}$ be the first $a$ for which the last statement is not valid. Then $\zeta_{n}(\bar{a}) \in \partial Q^{(j)}(\bar{a})$. Since it cannot be in $\partial R_{j}$ (because $n \geq j$ ), it has to lie in the right (vertical) boundary of $Q^{(j)}(\bar{a})$.

We claim that as $a$ increases, $\zeta_{n}$ crosses $Q^{(j)}(a)$ in exactly one point, i.e. for all $a>\bar{a}, \zeta_{n}(a) \notin Q^{(j)}(a)$. This is because for $a \in \omega$, $\left|\frac{d}{d a} \zeta_{n}\right|>$ $\hat{K}^{-1}\left|w_{n}\right|>K^{-1} e^{\lambda n}$ (Proposition 11.1 and Lemma 5.2), while $\left|\frac{d}{d a} z_{0}^{*}\left(Q^{(j)}\right)\right|<$ $K_{1}$ (Corollary 10.2). Since $z_{0}^{*}\left(Q^{(j)}\right)$ and the "vertical" boundaries of $Q^{(j)}$ move in the horizontal direction at the same speed, and we may assume $n$ is large
enough that $K^{-1} e^{\lambda n} \gg K_{1}, a \mapsto \zeta_{n}(a)$ crosses $\cup_{a} \partial Q^{(j)}(a)$ transversally in $R_{1} \times J$ in exactly one point.

The picture can therefore be summarized as follows. Let $j_{0}$ be the largest $j \leq \frac{1}{\alpha^{*}} \theta n$ such that $\zeta_{n}$ meets $Q^{(j)}$, and assume $\zeta_{n}(\hat{a})$ lies in the interior of $Q^{\left(j_{0}\right)}$. Let $Q^{\left(j_{0}\right)} \subset Q^{\left(j_{1}\right)} \subset Q^{\left(j_{2}\right)} \ldots$ be consecutive critical regions starting from $Q^{\left(j_{0}\right)}$ for $T_{\hat{a}}$. This structure, as we know, is identical for all the $T_{a}$. As $a$ increases (or decreases), this nested structure moves with speed $\mathcal{O}(1)$, which is very slow relative to the speed of $a \mapsto \zeta_{n}(a)$. Thus from the point of view of $\zeta_{n}$, the critical structure appears stationary, and the picture resembles that of a single map.
C. Bound period and recovery. The setting is as in B above. For each $a$, we have defined for the map $T_{a}$ the notion of $\phi(z)$ for $z=\zeta_{n}(a), d_{\mathcal{C}}(z):=$ $|\phi(z)-z|$, and $p(z)$. To emphasize their dependence on the map $T_{a}$, we now write $\phi_{a}(z), d_{\mathcal{C}(a)}(z)$ and $p_{a}(z)$. For purposes of studying the evolution of $\zeta_{n}$, we can, if we so choose, use the definitions associated with each individual $a$ for $\zeta_{n}(a)$. For aesthetic as well as practical reasons (to become clear in the next section), we prefer to have some coherence along $\zeta_{n}$, even if this involves some small modifications in the definitions above. We explain how this can be done:

Step 1. Choosing a common guiding critical point $\phi(\omega)$. The choice is quite arbitrary. Let $j_{0}$ be the largest $j \leq \frac{1}{\alpha^{*}} \theta n$ such that $\zeta_{n}$ meets $Q^{(j)}$, and pick $\hat{a}$ with $\zeta_{n}(\hat{a}) \in Q^{\left(j_{0}\right)}$. Let $\phi(\omega):=z_{0}^{*}\left(Q^{\left(j_{0}\right)}\right)(\hat{a})$, and define $d_{\mathcal{C}}(z)=|z-\phi(\omega)|$ for $z \in \zeta_{n}(a)$.

Lemma 11.4. For $z=\zeta_{n}(a)$, let $j$ be the generation of $\phi_{a}(z)$. Then

$$
\left|d_{\mathcal{C}(a)}(z)-d_{\mathcal{C}}(z)\right| \leq b^{\frac{j}{4}}+K_{1}|a-\hat{a}| .
$$

Proof. Let $j_{0}>j_{1}>\cdots$ be as in B. Then $z \in Q^{\left(j_{i}\right)}(a)$ for some $i \geq 0$. By definition, $j \leq j_{i}$, and so $Q^{\left(j_{0}\right)} \subset Q^{(j)}$. We then have

$$
\begin{aligned}
\left|d_{\mathcal{C}(a)}(z)-d_{\mathcal{C}}(z)\right| & \leq\left|z_{0}^{*}\left(Q^{(j)}\right)(a)-z_{0}^{*}\left(Q^{\left(j_{0}\right)}\right)(\hat{a})\right| \\
& \leq\left|z_{0}^{*}\left(Q^{(j)}\right)(a)-z_{0}^{*}\left(Q^{\left(j_{0}\right)}\right)(a)\right|+\left|z_{0}^{*}\left(Q^{\left(j_{0}\right)}\right)(a)-z_{0}^{*}\left(Q^{\left(j_{0}\right)}\right)(\hat{a})\right| \\
& \leq b^{\frac{j}{4}}+K_{1}|a-\hat{a}| ;
\end{aligned}
$$

the first term in the last inequality is by Lemma 4.1 and the second by Corollary 10.2.

The first term in the error above is innocuous. Since

$$
|a-\hat{a}|<\hat{K} e^{-\lambda n}\left|\zeta_{n}(a)-\zeta_{n}(\hat{a})\right|,
$$

the second term is negligible if $d_{\mathcal{C}(a)}(z)$ is > $K_{1} \hat{K} e^{-\lambda n} \delta$. In particular, for $z$ with $d_{\mathcal{C}(a)}(z)>e^{-\alpha n}$, we have $d_{\mathcal{C}(a)}(z) \approx d_{\mathcal{C}}(z)$.

Step 2. Definition of a new bound period $p(\cdot)$. Let $\mathcal{P}^{\omega}$ be the partition $\mathcal{P}$ as in Section 2.2, centered at $\pi_{x}(\phi(\omega))$, and let $\hat{\omega}$ be such that $\pi_{x}\left(\zeta_{n}(\hat{\omega})\right) \approx I_{\mu j}$. We define

$$
p(\hat{\omega})=\min _{a \in \hat{\omega}} p_{a}\left(\zeta_{n}(a)\right) .
$$

For this definition to be meaningful, we must verify that it has the properties of bound periods in the sense of $p_{a}(\cdot)$ for a single map. The next lemma, the proof of which is given in Appendix A.25, assures us this is the case.

Lemma 11.5. Let $\hat{\omega}$ be as above, and let $p=p(\hat{\omega})$. Then
(a) $K^{-1}|\mu|<p<K|\mu|$;
(b) for $a, a^{\prime} \in \hat{\omega}$ and $j<p,\left|\zeta_{n+j}(a)-\zeta_{n+j}\left(a^{\prime}\right)\right|<2 e^{-\beta j}$;
(c) $\left|\tau_{n+p}(a)\right|>K^{-1} e^{\frac{p}{4}}\left|\tau_{n}(a)\right|$ for all $a \in \hat{\omega}$;
(d) $\zeta_{n+p}(a)$ is out of all splitting periods, $s\left(\tau_{n+p}(a)\right)<\hat{\varepsilon}$;
(e) $\left|\zeta_{n+p}(\hat{\omega})\right| \geq \frac{1}{\mu^{2}} e^{-K \alpha|\mu|}$.
D. Decomposition into bound and free states. With bound periods defined, we may now assign a bound or free state to each $\zeta_{i}(a)$ in the evolution of $\zeta_{i}$, namely that $\zeta_{i}(a)$ is free if it is not in a bound period as defined in C. We remark that this notion is not necessarily consistent with the one for a single map $T_{a}$. Note that in Lemma 11.3, and in the setting of B and C, the word "free" as stated referred to "free" in the sense of individual maps. We leave it to the reader to check that these statements are, in fact, valid if "free" is given the meaning in this paragraph.

## 12. Derivative growth via statistics

This section is about how to deal with (A4) (see §4.1). We focus on one critical point at a time, and discuss (i) what it takes to maintain regular derivative growth along its orbit, and (ii) why one should expect the conditions guaranteeing this growth to be satisfied by a positive measure set of parameters.

### 12.1. Estimating $\left|w_{i}^{*}\right|$ in terms of itinerary

We return in this subsection to the dynamics of a single map to motivate a few ideas. As explained earlier, (A4) is not a self-perpetuating property. We now give a condition in terms of the itinerary of $z_{i}$ that guarantees sustained exponential growth of $\left|w_{i}^{*}\left(z_{0}\right)\right|$.

Consider for definiteness $T \in \mathcal{G}_{N}$, and assume that for some $z_{0} \in \Gamma_{\frac{1}{\alpha^{*} \theta} \theta N} \backslash$ $\Gamma_{\theta N}, d_{\mathcal{C}}\left(z_{i}\right)>e^{-\alpha i}$ for all $i \leq n, n \leq \frac{1}{\alpha^{*}} N$. We know from Corollary 6.1 that for $w_{0}=\mathbf{v} \in X_{z_{0}},\left|w_{i}^{*}\right|>K^{-1} e^{\left(\frac{1}{3} \lambda-2 \alpha\right) i}$ for all $i \leq n$. Let us examine more closely how this growth comes about. Let $t_{1}<t_{1}+p_{1} \leq t_{2}<t_{2}+p_{2} \leq t_{3}<$
$t_{3}+p_{3} \leq \cdots$ be such that $t_{i}$ are the consecutive free return times up to time $n$ and $p_{i}$ the lengths of the ensuing bound periods. Then we have

$$
\text { (i) } \frac{\left|w_{t_{i+1}}^{*}\right|}{\left|w_{t_{i}+p_{i}}^{*}\right|} \geq c_{2} e^{\frac{1}{4} \lambda_{0}\left(t_{i+1}-\left(t_{i}+p_{i}\right)\right)}, \quad \text { (ii) } \quad \frac{\left|w_{t_{i}+p_{i}}^{*}\right|}{\left|w_{t_{i}}^{*}\right|} \geq K^{-1} e^{\frac{1}{3} \lambda p_{i}} .
$$

Observe that in (i), the exponent is the "outside exponent" $\frac{1}{4} \lambda_{0}$, which is strictly $>\lambda$. As for (ii), the guaranteed growth rate of $\frac{1}{3} \lambda$ does not contribute much to maintaining a Lyapunov exponent of $\lambda$. It is no significant loss if we replace it by the weaker estimate $\frac{\left|w_{t_{i}+p_{i}}^{*}\right|}{\left|w_{t_{i}}^{*}\right|} \geq c_{2}^{-1}$, which is what we will do.

We continue to reason as follows: If the fraction of time $z_{i}$ spends in bound periods between time 0 and $n$ is $<\sigma$, and $z_{n}$ is not in a bound period, then $\left|w_{n}^{*}\right|>$ const $e^{(1-\sigma) \frac{1}{4} \lambda_{0} n}$. This number is $>e^{\lambda n}$ if $\sigma$ is sufficiently small. Since the "outside exponent" does not decrease as $\delta$ decreases (Lemma 2.1), it is logical to attempt to decrease $\sigma$ by decreasing $\delta$, the idea being that some of the time intervals that are bound periods for the original $\delta$ would no longer be counted as bound periods for a smaller $\delta$. We summarize the conclusion of this reasoning in the following lemma:

Let $B(\hat{\delta} ; 0, n)$ be the total time between 0 and $n$ during which $z_{i}$ spends in bound periods initiated from visits to the region $d_{\mathcal{C}}(\cdot)<\hat{\delta}$.

Lemma 12.1. Let $z_{0}$ be as above, and let $\sigma>0$ and $0<\hat{\delta} \leq \delta$. If $B(\hat{\delta} ; 0, n)<\sigma n$, then

$$
\left|w_{n}^{*}\right|>K^{-1} \hat{\delta} e^{\left[(1-\sigma) \frac{1}{4} \lambda_{0}-3 \alpha\right] n} .
$$

Proof. Consider first the case where $z_{n}$ is free. Let $\hat{t}_{1}<\hat{t}_{1}+\hat{p}_{1} \leq \hat{t}_{2}<$ $\hat{t}_{2}+\hat{p}_{2} \leq \cdots \leq \hat{t}_{k}+\hat{p}_{k} \leq n$ be such that $\hat{t}_{1}, \cdots, \hat{t}_{k}$ are the consecutive free return times to $\left\{d_{\mathcal{C}}(\cdot)<\overline{\hat{\delta}}\right\}$. Then

$$
\left|w_{n}\right|=\frac{\left|w_{n}\right|}{\left|w_{\hat{t}_{k}+\hat{p}_{k}}\right|} \cdots \frac{\left|w_{\hat{t}_{2}}\right|}{\left|w_{\hat{t}_{1}+\hat{p}_{1}}\right|} \frac{\left|w_{\hat{t}_{1}+\hat{p}_{1}}\right|}{\left|w_{\hat{t}_{1}}\right|}\left|w_{\hat{t}_{1}}\right| .
$$

We have $\frac{\left|w_{n}\right|}{\left|w_{\hat{t}_{k}+\hat{p}_{k}}\right|}>c_{2} \hat{\delta} e^{\frac{1}{4} \lambda_{0}\left(n-\hat{t}_{k}-\hat{p}_{k}\right)}$ by Lemma 3.5(i) with $\hat{\delta}$ in the place of $\delta$, and $\frac{\left|w_{\hat{t}_{i+1}}\right|}{\left|w_{t_{i}+\hat{p}_{i}}\right|}>c_{2} e^{\frac{1}{4} \lambda_{0}\left(\hat{t}_{i+1}-\hat{t}_{i}-\hat{p}_{i}\right)}$ by Lemma 3.5(ii). To cancel $c_{2}$ we use $\frac{\left|w_{\hat{e}_{i}+\hat{p}_{i}}\right|}{\left|w_{\hat{t}_{i}}\right|}>c_{2}^{-1}$, which is a trivial consequence of ( $\mathrm{P} 2^{\prime}$ )(ii). This gives $\left|w_{n}^{*}\right|>$ $K^{-1} \hat{\delta} e^{\left[(1-\sigma) \frac{1}{4} \lambda_{0}\right] n}$ since $\hat{p}_{1}+\cdots+\hat{p}_{k} \leq \sigma n$ by assumption. The factor $-3 \alpha n$ is needed if $n$ is not free; see Lemma 5.4.

In this lemma, we think of $\hat{\delta}$ as possibly $\ll \delta$, and the factor $\hat{\delta}$ as the price to pay due to the greater nonuniformity in growth properties "outside" - where "outside" now means $d_{\mathcal{C}}(\cdot)>\hat{\delta}$. As we will see, this factor will be absorbed into the initial growth if the critical orbit remains outside of $\mathcal{C}^{(1)}$ for a long time.

We point out that a quantity similar to $B_{n}(\hat{\delta} ; 0, n)$ has already appeared in the context of 1 D maps; see Proposition 2.2. Our next step is to make a connection to this proposition.

### 12.2. Processes defined by curves of critical points

In Section 9.3, we borrowed some results from Section 2.3 for the dynamics of $T \in \mathcal{G}$ on unstable manifolds - after establishing a strong resemblance between $\left.T^{k}\right|_{W^{u}}$ and iterated 1D maps. Section 11.2 suggests that this similarity can perhaps be extended to $\zeta_{i}: a \mapsto z_{i}(a)$. But the relation between the "dynamics" of critical curves and iterated 1D maps is a little more tenuous. For one thing, there is no reasonable description of global geometry for critical curves: even though $\zeta_{0}$ is defined on an interval, it is inevitable that one will lose control of $\zeta_{n}$ on parts of this interval as $n$ increases.

We claim, nevertheless, that the statistical results in Section 2.3 are valid. To see that, we return to Section 2 to examine the situation more closely:

Observations. 1. In order to apply the results in Section 2.3 to critical curves, we must verify for them estimates analogous to (P1)-(P3) in Section 2.2 and fix a definition of canonical subdivision by itinerary with properties identical to those in Section 2.3.
2. Once that is done, we may proceed, letting $\gamma_{i}=\zeta_{i}$ wherever it makes sense. It is not important in Section 2.3 for $\gamma_{i}$ to be $=f^{i}$ or the $i$ th iterate of any map; $\left\{\gamma_{i}\right\}$ could have been a process, meaning a sequence of maps from $J$ to $I$.
3. Finally, as noted at the end of Section 2.3 , the stated results are entirely unaffected if we choose to stop considering any element of $\mathcal{Q}_{i}$ at any time by simply setting $\zeta_{j}=*$ for all $j \geq i$. Here, the symbol $*$ will correspond to deleted parameters.

Definition of a process $\left\{\gamma_{i}\right\}$ associated with $\zeta_{i}: a \mapsto z_{i}(a)$. We assume (C1) and (C2) on an interval $J$, and fix $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N} \backslash \Gamma_{\theta N}$. Associated with $z_{0}$, we seek to define a sequence of maps

$$
\gamma_{i}: J \rightarrow R_{1} \cup\{*\} \quad \text { for } \quad 0 \leq i \leq \frac{1}{\alpha^{*}} N
$$

with the property that $\gamma_{i}(a)=\zeta_{i}(a)=z_{i}(a)$ whenever $\gamma_{i}(a) \neq *$. Here as in Section $2.3, *$ is the "garbage symbol": once $\gamma_{i}(a)=*$ for some $i, \gamma_{j}(a)=*$ for all $j>i$; that is to say, we stop considering $a \in J$ from that point on. To facilitate the description of $\gamma_{i}$, we first introduce the following language:

Let $\gamma: \omega \rightarrow R_{1}$ be such that $\frac{d}{d a} \gamma$ is nonzero and roughly horizontal. We introduce on $\omega$ a partition that will be referred to as the "canonical partition defined by $\gamma$ ". This partition is the pullback of the following partition on $\gamma(\omega)$ : First we divide $\gamma(\omega) \backslash \mathcal{C}^{(1)}$ into intervals of length $\approx \delta$ each (except possibly
for the end interval(s), which may be shorter), and partition each component of $\gamma(\omega) \cap \mathcal{C}^{(1)}$ into $\left\{\pi_{x}^{-1} I_{\mu j}\right\}$ using the guiding critical orbit $\phi(\omega)$ to center the partition $\mathcal{P}$ as is done in Section 11.2C. The final partition on $\gamma(\omega)$ is obtained by adjoining the end intervals in the partition above to their neighbors. We say a canonical partition is nontrivial if at least one partition point is introduced. Note that in a nontrivial partition, each element $\omega^{\prime}$ has the property that either $\gamma\left(\omega^{\prime}\right) \cap \mathcal{C}^{(1)}=\emptyset$ and $\delta \leq\left|\gamma\left(\omega^{\prime}\right)\right| \leq 3 \delta$, or $\gamma\left(\omega^{\prime}\right) \subset \mathcal{C}^{(1)}$ and $\pi_{x}\left(\gamma\left(\omega^{\prime}\right)\right) \approx I_{\mu j}$ for some $I_{\mu j} .{ }^{12}$

The canonical subdivision by itinerary for $\gamma_{i}$ proceeds as follows. We ignore the first $i_{0}$ iterates as they are not particularly meaningful. In general, for $\omega \in \mathcal{Q}_{i}$, we use the language "delete $\omega$ " and "set $\gamma_{i} \mid \omega=*$ " interchangeably. Let $\mathcal{Q}_{i_{0}}$ be the canonical partition on $J$ defined by $\gamma_{i_{0}}$. Here is how we go from one step to the next:

Case 1. Consider $\omega \in \mathcal{Q}_{i}$ where $\zeta_{i}(\omega)$ is free and outside of $\mathcal{C}^{(1)}$. We take the canonical partition defined by $\zeta_{i+1}$ on $\omega$, and set $\left.\zeta_{i+1}\right|_{\omega^{\prime}}=*$ for elements $\omega^{\prime}$ of this partition for which $d_{\mathcal{C}}\left(\zeta_{i+1}\left(\omega^{\prime}\right)\right)<e^{-\alpha(i+1)}$. On the part of $\omega$ not deleted, we set $\gamma_{i+1}=\zeta_{i+1}$, and call the restriction of the canonical partition on it $\mathcal{Q}_{i+1}$.

Case 2. Consider $\omega \in \mathcal{Q}_{i}$ for which $\gamma_{i}(\omega)$ is free and inside $\mathcal{C}^{(1)}$. It follows from the previous step that $\pi_{x}\left(\gamma_{i}(\omega)\right) \subset I_{\mu j}^{+}$. A bound period $p(\omega)$ is set as in Section 11.2C. We put $\omega \in \mathcal{Q}_{i+j}$ for all $j<p$, and at step $i+p$, we do as in Case 1, i.e. consider the canonical partition defined by $\zeta_{i+p}$, delete those elements with $d_{\mathcal{C}}(\cdot)<e^{-\alpha(i+p)}$, set $\gamma_{i+p}=\zeta_{i+p}$ on the rest and call the resulting partition $\mathcal{Q}_{i+p}$.

This completes the definition of the canonical subdivision by itinerary. We remark that by virtue of $(\mathrm{C} 1), d_{\mathcal{C}}\left(z_{i}(a)\right)>e^{-\alpha i}$ for all $i \leq N$ and $a \in J$, so that no deletions take place before time $N$. For $N<n \leq \frac{1}{\alpha^{*}} N$, the construction above is designed to guarantee that if $\gamma_{n}(a) \neq *$, then (C3) is satisfied for $z_{0}$ up to time $n$ on the parameter interval $\omega \in \mathcal{Q}_{n}$ containing $a$.

Verification of estimates analogous to (P1)-(P3) for $\gamma_{n}$. When $\tau_{n}=$ $\frac{d}{d a} \gamma_{n}$, Lemmas 11.3 and 11.5 play the role of (P1) and (P2). (P3) at $i_{0}$ follows from (a) the corresponding result for $w_{i_{0}},(\mathrm{~b})$ the fact that all $\zeta_{n}(a), a \in J$, remain very close to each other for all $n \leq i_{0}$, and (c) Proposition 11.1. The following distortion estimate for parameters is needed to take the place of (P3) for $n>i_{0}$. Its proof is given in Appendix A.26.

[^12]Lemma 12.2. There exists $K>0$ such that for any $\omega \in \mathcal{Q}_{n-1}$ such that $\gamma_{n}(\omega)$ is free,

$$
K^{-1}<\frac{\left|\tau_{n}(a)\right|}{\left|\tau_{n}\left(a^{\prime}\right)\right|}<K \quad \text { for all } a, a^{\prime} \in \omega
$$

### 12.3. Large deviation estimate

Having established the resemblance between $\gamma_{i}$ and iterated 1D maps, we now state the analog of Proposition 2.2 for parameters. For $i_{1}<i_{2}$, let $B\left(a, \hat{\delta} ; i_{1}, i_{2}\right)$ denote the number of $i \in\left(i_{1}, i_{2}\right]$ such that $\gamma_{i}(a)$ is in a bound period initiated at a previous time $j$ where $d_{\mathcal{C}}\left(z_{j}(a)\right)<\hat{\delta}$. Built into this definition is the implication that if $i$ is one of the times counted in $B\left(a, \hat{\delta} ; i_{1}, i_{2}\right)$, then $\gamma_{i}(a) \neq *$. The proof of Corollary 12.1 follows closely that of Proposition 2.2.

Corollary 12.1. Assume ( C 1 ) and ( C 2 ), and let $z_{0} \in \Gamma_{\frac{1}{\alpha^{*} \theta N}} \backslash \Gamma_{\theta N}$. Then Proposition 2.2 holds for $\left\{\gamma_{i}, i \leq \frac{1}{\alpha^{*}} N\right\}$, where $\left\{\gamma_{i}\right\}$ is the process associated with $a \mapsto z_{i}(a)$ defined above. More precisely, given any $\sigma>0$, there exists $\hat{\varepsilon}_{1}>0$ such that the following holds for all sufficiently small $\hat{\delta}>0$ : Let $t_{0}$ and $\omega \in \mathcal{Q}_{t_{0}}$ be such that $\gamma_{t_{0}}(\omega)$ is free and $\approx I_{\mu_{0} j_{0}}$ (in particular, $\left.\left.\gamma_{t_{0}}\right|_{\omega} \neq *\right)$, and let $n$ be such that (i) $t_{0}+n \leq \frac{1}{\alpha^{*}} \theta N$ and (ii) $n>K \sigma^{-1}\left|\mu_{0}\right|$. Then

$$
\left|\left\{a \in \omega: B\left(a, \hat{\delta} ; t_{0}, t_{0}+n\right)>\sigma n\right\}\right|<e^{-\hat{\varepsilon}_{1} n}|\omega| .
$$

This result holds also if $\gamma_{t_{0}}(\omega)$ is outside of $\mathcal{C}^{(1)}$, free and has length $\geq \delta$. In this case condition (ii) for $n$ is replaced by $n>K \sigma^{-1} \log \frac{1}{\delta}$.

### 12.4. In preparation for the selection of good parameters

The main ingredients for dealing with (A4) are treated in the last 3 subsections. The results as stated, however, are not quite in a form that can be applied directly. This subsection contains the adjustments needed to render Lemma 12.1 and Corollary 12.1 ready for use in the construction in Section 13. We also specify the desired relations between $\sigma, \hat{\delta}$ etc. and constants chosen earlier.

Setting. We assume (C1) and (C2), and continue to focus on a single critical point $z_{0} \in \Gamma_{\frac{1}{\alpha^{*}} \theta N} \backslash \Gamma_{\theta N}$. Let $\gamma_{i}, i \leq \frac{1}{\alpha^{*}} \theta N$, be the process associated with $z_{0}$. In what follows, we may assume also that $d_{\mathcal{C}}\left(z_{i}\right)>\delta_{0}$ for all $i \leq n_{0}$ where $n_{0}$ is as large as we need, and that at least one subdivision of the parameter interval takes place before $\gamma_{i}(J)$ meets $\mathcal{C}^{(1)}$.
(1) Measure of parameters deleted in connection with (A4). The procedure in Section 13 requires that we work with time intervals of the type [ $n, 2 n$ ].

Corollary 12.2. Assume $\sigma \gg \alpha$, and let $\hat{\varepsilon}_{1}=\hat{\varepsilon}_{1}\left(\frac{1}{2} \sigma\right)$ be given by Corollary 12.1 (with $\frac{1}{2} \sigma$ in the place of $\sigma$ ). Let $\hat{\delta}$ be small enough to satisfy the requirement in Corollary 12.1. Then for all $\omega \in \mathcal{Q}_{n}$ with $\left.\gamma_{n}\right|_{\omega} \neq *$,

$$
|\{a \in \omega: B(a, \hat{\delta} ; n, 2 n)>\sigma n\}|<e^{-\hat{\varepsilon}_{1} n}|\omega| .
$$

Proof. We explain the modifications necessary to apply Corollary 12.1. If $\gamma_{n}(\omega)$ is free and is either $\approx I_{\mu j}$ or is outside and has length $\geq \delta$, then we apply Corollary 12.1 directly with $t_{0}=n$. Note that the lower bound on $n$ in Corollary 12.1 is satisfied: if $\gamma_{n}(\omega) \approx I_{\mu j}$, then $\mu \leq \alpha n$, so that $n \geq \frac{1}{\alpha} \mu \gg \sigma^{-1} \mu$; we may assume $n>K \sigma^{-1} \log \frac{1}{\delta}$ since $n_{0}$ can be arbitrarily large. If $\gamma_{n}(\omega)$ is not free or is shorter than required, we back up to step $n_{1}$ when $\omega$ was first created as an element of some $\mathcal{Q}_{n_{1}}$. Notice first that there is such an $n_{1}$, for by assumption a subdivision occurred before time $n$. Moreover, $\gamma_{n_{1}}(\omega)$ is free, and it is either $\approx I_{\mu j}$ or is outside and has length $\geq \delta$. By the parameter version of Lemma $2.3, n \leq(1+K \alpha) n_{1}$. We may then apply Corollary 12.1 with $t_{0}=n_{1}$ and $\frac{1}{2} \sigma$ in the place of $\sigma$. Since $\sigma \gg \alpha$, a parameter $a$ with $B\left(a, \hat{\delta} ; n_{1}, 2 n\right)>\sigma\left(2 n-n_{1}\right)$ clearly satisfies $B(a, \hat{\delta} ; n, 2 n)>\frac{1}{2} \sigma n$.
(2) Growth of $\left|w_{i}^{*}\left(z_{0}\right)\right|$ for good parameters. Let $n=2^{j_{1}} n_{0}$, and consider the following procedure repeated on time intervals $\left[n_{0}, 2 n_{0}\right], \cdots,\left[2^{j} n_{0}, 2^{j+1} n_{0}\right]$, $\cdots,\left[2^{j_{1}-1} n_{0}, n\right]$ : On each time interval $\left[2^{j} n_{0}, 2^{j+1} n_{0}\right]$, in addition to the deletions corresponding to (C3) (see $\S 11.2$ ), we delete at time $2^{j+1} n_{0}$ all $\omega \in \mathcal{Q}_{2^{j+1} n_{0}}$ on which $B\left(a, \hat{\delta} ; 2^{j} n_{0}, 2^{j+1} n_{0}\right)>\sigma 2^{j} n_{0}$. The following corollary gives a lower bound on $\left|w_{i}^{*}\left(z_{0}\right)\right|$ for $T=T_{a}$ where $a$ survives these deletions up to time $n$.

Corollary 12.3. Assume
(a) $d_{\mathcal{C}}\left(z_{i}\right)>\min \left(\delta, e^{-\alpha i}\right)$ for all $i \leq n$;
(b) $B\left(\hat{\delta} ; 2^{j} n_{0}, 2^{j+1} n_{0}\right)<\sigma 2^{j} n_{0}$ for all $j<j_{1}$.

Then for all $i \leq n,\left|w_{i}^{*}\left(z_{0}\right)\right|>c_{2} e^{\left[\frac{1}{4}(1-2 \sigma) \lambda_{0}-3 \alpha\right] i}$.
Proof. Assumptions (a) and (b) together with Lemma 12.1 imply that at times $2^{j} n_{0},\left|w_{i}^{*}\right|>K^{-1} \hat{\delta} e^{\left[\frac{1}{4}(1-\sigma) \lambda_{0}-3 \alpha\right] i}$. Between times $i$ and $2 i$, the worst-case scenario is that all the close returns (i.e. returns to $\left\{d_{\mathcal{C}}(\cdot)<\hat{\delta}\right\}$ ) occur at the beginning of this time block. Even so, we guarantee easily that for all $k<i$, $\left|w_{i+k}^{*}\right|>K^{-1} \hat{\delta} e^{\left[\frac{1}{4}(1-2 \sigma) \lambda_{0}-3 \alpha\right](i+k)}$. Observe finally that the factor $K^{-1} \hat{\delta}$ can be replaced by $c_{2}$, i.e. it is absorbed into the initial stretch if $n_{0}$ is sufficiently large depending on $\hat{\delta}$.
(3) Choice of constants. The exponents directly related to derivative growth are $\lambda_{0}, \lambda, \lambda^{*}$ and $\alpha$. We review briefly what they represent. First, outside of $\mathcal{C}^{(1)}, b$-horizontal vectors grow at rate $\frac{1}{4} \lambda_{0}$; see Lemma 3.5. The first constant chosen in this paper, $\lambda<\frac{1}{5} \lambda_{0}$, is the minimum growth rate along critical orbits guaranteed by (A4). Because more stringent estimates are needed for reasons to be explained, we fix a slightly larger target Lyapunov exponent $\lambda^{*}=\lambda+\frac{1}{100} \lambda_{0}$; see Section 10.1. The constant $\alpha$ is then chosen to satisfy $\alpha \ll \min \{\lambda, 1\}$.

Next we come to $\sigma$, which is chosen so that $\frac{1}{4}(1-2 \sigma) \lambda_{0}-3 \alpha$, the exponent in Corollary 12.3, is $>\lambda^{*}$. For example, $\sigma=\frac{1}{100}$ will work. We may assume this is in agreement with the relation $\sigma \gg \alpha$ as required in Corollary 12.2. Once $\sigma$ is fixed, we choose $\hat{\delta}$ small enough to satisfy Corollary 12.1.

Summary. If $\sigma$ and $\hat{\delta}$ are as in the last paragraph, and the hypotheses of Corollary 12.3 are satisfied, then $\left|w_{i}^{*}\left(z_{0}\right)\right|>c_{2} e^{\lambda^{*} i}$ for all $i \leq n$. Moreover, between times $2^{j} n_{0}$ and $2^{j+1} n_{0}$, the measure of parameters in violation of Corollary $12.3(\mathrm{~b})$ is, by Corollary $12.2,<e^{-\hat{\varepsilon} 2^{j} n_{0}}|J|$.

## 13. Positive measure sets of good parameters

The purpose of this section is to construct, for a given family $\left\{T_{a}\right\}$ satisfying the Standing Hypotheses in Section 1 and with $b$ sufficiently small, a sequence of sets

$$
\Delta_{0} \supset \Delta_{n_{0}} \supset \Delta_{2 n_{0}} \supset \Delta_{2^{2} n_{0}} \supset \cdots
$$

in parameter space with the properties that
(i) $\left\{T_{a}, a \in \Delta_{2^{j} n_{0}}\right\} \subset \mathcal{G}_{2^{j} n_{0}}^{\#}$ (where $\mathcal{G}_{n}^{\#}$ is as defined in Section 10.1) and
(ii) $\Delta:=\cap_{j \geq 0} \Delta_{2^{j} n_{0}}$ has positive Lebesgue measure.

Together with the material in Section 9, this construction brings to completion the proof of our Main Theorem.

The construction in this section requires more stringent conditions on the global constants in Section 6.1 than are imposed in Part II. See, for instance, the end of Section 11.1 and Section 12.4(3).

### 13.1. Getting started

The two properties required of the start-up interval $\Delta_{0}$ are:
(1) For all $a \in \Delta_{0}$ and $z_{0} \in \Gamma_{1}, d_{\mathcal{C}}\left(z_{i}\right)>\delta_{0}$ for all $i \leq n_{0}$ where $n_{0}$ is a very large number to be prescribed.
(2) For each $z_{0} \in \Gamma_{1}$, a subdivision occurs in the process $a \mapsto z_{i}(a)$ before $\gamma_{i}\left(\Delta_{0}\right)$ meets $\mathcal{C}^{(1)}$.

Here $\delta_{0}$ is as in Definition 1.1; recall that $d\left(f_{a^{*}}^{i}(\hat{x}), C\right)>2 \delta_{0}$ for all $i>0$. Lower bounds have been placed on $n_{0}$ a finite number of times in previous sections: among the more important places where this condition appeared are (i) to provide time for hyperbolicity of $T_{a}$ to build up initially (in Part II);
(ii) to allow the comparability of space and $a$-derivatives to take hold (see Proposition 11.1); and (iii) to absorb the small constant $\hat{\delta}$ from Lemma 12.1 (see Corollary 12.3). A few more conditions on $n_{0}$ will be imposed in this section. The process referred to in (2) is the one in Section 12.2. The purpose of (2) is to ensure that the entire parameter interval is not lost in the first deletion: Let $n_{1}$ be the first time $\Delta_{0}$ is subdivided. Then $\left|\gamma_{n_{1}}(\omega)\right| \geq \delta$ for every $\omega \in \mathcal{Q}_{n_{1}}$, and if $\omega \in \mathcal{Q}_{j}$ is such that $\gamma_{i}(\omega) \cap \mathcal{C}^{(1)}=\emptyset$ for all $i<j$ and $\gamma_{j}(\omega) \cap \mathcal{C}^{(1)} \neq \emptyset$, then $\left|\gamma_{j}(\omega)\right|>\hat{K}^{-1} c_{2} \delta$. Thus assuming $n_{0}$ is large enough that $e^{-\alpha n_{0}} \ll \hat{K}^{-1} c_{2} \delta$, we are guaranteed that only a small fraction of the measure is deleted.

We claim that for any $\Delta_{0}$ containing a $a^{*}$ short enough for (1) to be satisfied, (2) is automatically satisfied if $b$ is sufficiently small. To see this, let $\left\{\hat{x}_{0}^{k}\right\}$ be the critical points of the 1D maps $f_{a}$, and let $a \mapsto \hat{x}_{i}^{k}(a), i=0,1, \cdots$, be the critical curves defined by the 1D maps. For each $k$, let $\hat{n}_{1}^{k} \geq n_{0}$ be the first time $\left|\hat{x}_{i}^{k}\left(\Delta_{0}\right)\right|>3 \delta$, and let $\hat{n}_{1}$ be the maximum of the $\hat{n}_{1}^{k}$. Now let $\zeta_{i}^{k}: a \mapsto z_{i}^{k}(a)$ where $z_{0}^{k}$ is the critical point of $T_{a}$ near ( $\left.\hat{x}_{0}^{k}, 0\right)$. We choose $b$ small enough that $\left|\hat{x}_{i}^{k}(a)-\pi_{x}\left(z_{i}^{k}(a)\right)\right| \ll \delta$ for $i \leq \hat{n}_{1}$ for all $a \in \Delta_{0}$. Then for $i<\hat{n}_{1}^{k}$, $\zeta_{i}^{k}\left(\Delta_{0}\right) \cap \mathcal{C}^{(1)}=\emptyset$, and $\left|\zeta_{\hat{n}_{1}^{k}}^{k}\left(\Delta_{0}\right)\right|>2 \delta$, so a subdivision occurs at or before time $\hat{n}_{1}^{k}$ in the process associated with $z_{0}^{k}$.

In the rest of this section, let $\lambda^{*}$ and $\mathcal{G}_{n}^{\#}$ be as defined in Section 10.1. For clarity of presentation, we first describe the construction up to time $\theta^{-1}$ (where the situation is simpler) before giving it in full generality.

### 13.2. Construction of $\Delta_{N}$ for $N \leq \theta^{-1}$

A. Outline of scheme. This time period is characterized by the fact that the only relevant critical points are those in $\Gamma_{1}:=\left\{z_{0}^{1}, \cdots, z_{0}^{q}\right\}$. Associated with each $z_{0}^{k}$, we construct a sequence of parameter sets $\Delta_{0}=\Delta_{0}^{k} \supset \Delta_{1}^{k} \supset$ $\Delta_{2}^{k} \supset \cdots \supset \Delta_{\theta^{-1}}^{k}$ with the property that for $a \in \Delta_{i}^{k}, z_{j}^{k}(a)$ has the desired properties for all $j \leq i$. The parameter sets $\Delta_{i}:=\cap_{1 \leq k \leq q} \Delta_{i}^{k}$ consist, therefore, of parameters for which all the critical orbits have the desired properties up to time $i$.

The sets $\left\{\Delta_{i}^{k}\right\}$ are constructed in the following order. First, we set $\Delta_{i}^{k}=$ $\Delta_{0}$ for all $i \leq n_{0}$ and all $k$. Then we proceed with an $N$-to- $2 N$ scheme, i.e. we go from step $n_{0}$ to step $2 n_{0}, 2 n_{0}$ to $4 n_{0}, 4 n_{0}$ to $8 n_{0}$, and so on, until $\theta^{-1}$, which we may assume is $=2^{\ell_{0}} n_{0}$, is reached. Within each stage, i.e. from $N=2^{\ell} n_{0}$ to $2 N$, we construct for each $k$ the parameter sets $\Delta_{i}^{k}, N<i \leq 2 N$. Which
$k$ goes first is immaterial, but it is important that all the critical orbits be treated up to time $2 N$ before we go to the next stage.

Remark. The number 2 in our $N$-to- $2 N$ scheme is somewhat arbitrary; the idea of updating all the critical orbits to order $\sim N$ simultaneously (as opposed to treating one to an arbitrarily large time before beginning on a second) is not. This is because the derivative estimate (A4) for the $q$ critical orbits cannot be developed independently of each other: when $z_{i}^{k}$ visits $Q^{(1)}\left(z_{0}^{k^{\prime}}\right)$, it relies on the orbit of $z_{0}^{k^{\prime}}$ to guide it through its derivative recovery, and parameters that are favorable for $z_{0}^{k}$ may have been deleted for $z_{0}^{k^{\prime}}$. As we will see, the times $2^{\ell} n_{0}, \ell=1,2, \cdots$, are designated times for different critical orbits to communicate to each other their selected parameter sets.
B. Processes $\left\{\gamma_{i}^{k}\right\}$ defined on $\Delta_{0}$. In Section 12.2 , we considered a parameter interval $J$ on which all $T_{a}$ are assumed to be in $\mathcal{G}_{N}$, and introduced for each critical point a process $\gamma_{i}$ defined up to time $\frac{1}{\alpha^{*}} N$. In a similar manner, we now wish to define for each $z_{0}^{k}$ a process

$$
\gamma_{i}^{k}: \Delta_{0} \rightarrow R_{1} \cup\{*\}, \quad i=0,1,2, \cdots, \theta^{-1}
$$

Section 12.2 does not guarantee that such a process is well defined, for it is not likely that $T_{a} \in \mathcal{G}_{\alpha^{*} \theta^{-1}}$ for all $a \in \Delta_{0}$. Here is how we circumvent the problem: we use the procedure in Section 12.2 to extend $\gamma_{i}^{k}$ from step $N=2^{\ell} n_{0}$ to step $2 N$ whenever it is feasible, and to set $\gamma_{i}=*$ whenever it is not. More precisely, for fixed $k$ and $N$, we assume $\gamma_{i}^{k}$ is defined on $\Delta_{0}$ for all $i \leq N$. Associated with $\gamma_{i}$ is its canonical subdivision by itinerary $\mathcal{Q}_{i}$. For each $\omega \in \mathcal{Q}_{N}$, we set $\left.\gamma_{N+1}\right|_{\omega}=*$ unless $T_{a} \in \mathcal{G}_{2 \alpha^{*} N}$ for all $a \in \omega$. Thus when $\left.\gamma_{N+1}\right|_{\omega} \neq *$, the construction in Section 12.2 can legitimately be carried out on $\omega$ up to time $2 N$.

There is one other difference between the construction here and that in Section 12.2 , where $\gamma_{i}=*$ is set only to achieve $d_{\mathcal{C}}\left(z_{i}\right)>e^{-\alpha i}$. Here we permit the setting of $\gamma_{i} \mid \omega=*, \omega \in \mathcal{Q}_{i}$, for a wider range of reasons as we will see in paragraph C.
C. Formal procedure from step $N=2^{\ell} n_{0}$ to step $2 N$. At time $N$, assume we are handed the following objects: For each $k=1,2, \cdots, q$, there is a process $\gamma_{i}^{k}: \Delta_{0} \rightarrow R_{1} \cup\{*\}$ well defined up to time $N$. The set $\Delta_{N}^{k}:=\left\{\gamma_{N}^{k} \neq *\right\}$ has the property that for all $a \in \Delta_{N}^{k}$,
(i) $d_{\mathcal{C}}\left(z_{i}^{k}(a)\right)>3 e^{-\alpha i}$ for all $i \leq N$;
(ii) $B\left(a, \hat{\delta} ; 2^{j} n_{0}, 2^{j+1} n_{0}\right)<\sigma 2^{j} n_{0}$ for all $j<\ell$.

Observe, by Corollary 12.3 , that $\Delta_{N}:=\cap_{k} \Delta_{N}^{k} \subset \mathcal{G}_{N}^{\#}$.
How to go from step $N$ to step $2 N$. The following steps are taken for each $k$.
(1) First we set $\left.\gamma_{N+1}^{k}\right|_{\hat{\omega}}=*$ on those $\hat{\omega} \in \mathcal{Q}_{N}^{k}$ with $\hat{\omega} \cap \Delta_{N}=\emptyset$.
(2) On the rest of the $\hat{\omega} \in \mathcal{Q}_{N}^{k}$, we extend the process $\gamma_{i}^{k}$ to $2 N$ (see justification below), deleting all $\left.\omega \in \mathcal{Q}_{i}\right|_{\omega}$ with $d_{\mathcal{C}}\left(\gamma_{i}(\omega)\right)<3 e^{-\alpha i}$.
(3) Set $\gamma_{2 N}^{k} \mid \omega=*$ on those $\omega \in \mathcal{Q}_{2 N}^{k}$ with the property that $B(a, \hat{\delta} ; N, 2 N)>$ $\sigma N$ for $a \in \omega$.

Step (1) stipulates that unless some $a \in \hat{\omega}$ is good for all $q$ critical points, the entire parameter interval will be abandoned.

Justification for step (2): We need to show that $T_{a} \in \mathcal{G}_{2 \alpha^{*} N}$ for all $a \in \hat{\omega}$. By assumption, there exists $\hat{a} \in \hat{\omega}$ such that $T_{\hat{a}} \in \mathcal{G}_{N}^{\#}$. It follows from Proposition 10.1 that $T_{\hat{a}}$ has a continuation in $\mathcal{G}_{2 \alpha^{*} N}$ on the interval $\left[\hat{a}-\rho^{-2 \alpha^{*} N}, \hat{a}+\rho^{-2 \alpha^{*} N}\right]$. On the other hand, Proposition 11.1 gives $|\hat{\omega}|<\hat{K} e^{-\hat{\lambda} N}$, which is $\ll \rho^{-2 \alpha^{*} N}$.

Note that steps (2) and (3) lead directly to (i) and (ii) above at time $2 N$.
D. Measure deleted from step $N$ to step $2 N$. Consider one $z_{0}^{k}$ at a time. We wish to estimate the contribution to $\Delta_{N} \backslash \Delta_{2 N}$ by the orbit of $z_{0}^{k}$ (this is not to be confused with $\left.\Delta_{N}^{k} \backslash \Delta_{2 N}^{k}\right)$.

Deletions in Step (1). We have no control on the total measure of all the $\hat{\omega} \in \mathcal{Q}_{N}^{k}$ removed in this step, but all the $\hat{\omega}$ removed have the property that $\hat{\omega} \cap \Delta_{N}=\emptyset$ : the very fact that $\hat{\omega} \cap \Delta_{N}=\emptyset$ means that all the parameters in $\hat{\omega}$ have been deleted earlier due to violations on the part of critical orbits other than that of $z_{0}^{k}$. Thus from the point of view of $\Delta_{N} \backslash \Delta_{2 N}$, no measure is deleted in this step.

Deletions in Step (2). For $i$ with $N<i \leq 2 N$, we consider $\omega \in \mathcal{Q}_{i-1}^{k}$, and give an upper bound on the fraction of $\omega$ that may be deleted at the $i$ th iterate. Let $i_{0}$ be the smallest $j<i$ such that $\omega \in \mathcal{Q}_{j}^{k}$, i.e. $i_{0}$ is the time when the partition interval $\omega$ is created. There are two possibilities:
(i) $\gamma_{i_{0}}^{k}(\omega)$ is outside of $\mathcal{C}^{(1)}$ and $\delta<\left|\gamma_{i_{0}}^{k}(\omega)\right|<3 \delta$. In this case, $\left|\gamma_{i}^{k}(\omega)\right|>K^{-1} \delta$, and not knowing the location of $\gamma_{i}^{k}(\omega)$, we assume the worstcase scenario, i.e. $\gamma_{i}^{k}(\omega)$ crosses entirely a forbidden region $d_{\mathcal{C}}(\cdot)<3 e^{-\alpha i}$. The fraction of $\omega$ deleted is then $<K \delta^{-1} 6 e^{-\alpha i}<K e^{-\frac{1}{2} \alpha i}$. Here $K$ is the distortion constant (Lemma 12.2) as one transfers the length ratio on $\gamma_{i}^{k}(\omega)$ back to $\omega$.
(ii) $\pi_{x}\left(\gamma_{i_{0}}^{k}(\omega)\right) \approx I_{\mu j}$. Let $p$ be the bound period initiated at time $i_{0}$. Then $p \leq K|\mu|$, so that $\left|\gamma_{i}^{k}(\omega)\right|>K^{-1}\left|\gamma_{i_{0}+p}^{k}(\omega)\right|>\frac{K^{-1}}{\mu^{2}} e^{-K \alpha|\mu|}>e^{-2 K \alpha|\mu|}$. For the first inequality above, we use first Proposition 11.1, then $\left|w_{i}\right| \geq c_{2}\left|w_{i_{0}+p}\right|$. For the second inequality, we use Lemma 11.5(e). Thus the fraction of $\omega$ deleted is $<K 6 e^{-\alpha i} e^{2 K \alpha^{2} n}<K e^{-\frac{1}{2} \alpha i}$. (We note here the significance of the rule that in canonical subdivisions no partition point is introduced that would result in
an element $\omega$ with $\pi_{x}\left(\gamma_{i_{0}}^{k}(\omega)\right) \subset I_{\mu j}$ and $\left|\gamma_{i_{0}}^{k}(\omega)\right| \ll\left|I_{\mu j}\right|$. For such an element, we would not be able to control the fraction of parameters deleted.)

We conclude that between times $N$ and $2 N$, the total measure deleted in the course of executing step (2) on the orbit of each $z_{0}^{k}$ is

$$
<\sum_{N<i \leq 2 N} K e^{-\frac{1}{2} \alpha i}\left|\Delta_{0}\right| .
$$

Deletions in Step (3). By Corollary 12.2, on each $\hat{\omega} \in \mathcal{Q}_{N}^{k}$, the total measure deleted is $<e^{-\hat{\varepsilon}_{1} N}|\hat{\omega}|$. Thus the total measure deleted at time $2 N$ on account of executing step (3) on the orbit of $z_{0}^{k}$ is $<e^{-\hat{\varepsilon}_{1} N}\left|\Delta_{0}\right|$.

Summary. Let $\mathcal{D}_{N, 2 N}^{k}$ denote the set of $a \in \Delta_{N}$ deleted on account of the orbit of $z_{0}^{k}$ as we carry out our procedure from time $N$ to time $2 N$. Then

$$
\left|\mathcal{D}_{N, 2 N}^{k}\right|<\left(K \sum_{N<i \leq 2 N} e^{-\frac{1}{2} \alpha i}+e^{-\hat{\varepsilon}_{1} N}\right)\left|\Delta_{0}\right| .
$$

Writing the quantity in parentheses as $K^{\prime} e^{-\varepsilon^{\prime} N}$, and letting $\mathcal{D}_{N, 2 N}$ denote the set of all parameters deleted from $\Delta_{N}$ between time $N$ and time $2 N$, we have the estimate

$$
\left|\mathcal{D}_{N, 2 N}\right| \leq q K e^{-\varepsilon^{\prime} N}\left|\Delta_{0}\right|
$$

### 13.3. Construction of $\Delta_{N}$ for $N>\theta^{-1}$

A. Outline of scheme. Our basic strategy is as before; i.e., we work with cycles that go from time $N$ to time $2 N$, treating all relevant critical orbits in each cycle before going to the next and making deletions with the aid of processes of the type in the last subsection. There are two new aspects in the situation: the number of distinguishable critical orbits grows with time, and the critical structures of the maps $T_{a}$ are not uniform for all $a \in \Delta_{0}$. The processes we consider must reflect this reality; they are discussed in part B below.

Parts C and D follow their counterparts in Section 13.2. Except for treating the new complexities brought to light in part B, they do not differ substantially from before.
B. Processes defined by critical orbits: two new aspects.
(1) Processes associated with critical blobs. Relabeling the processes $\left\{\gamma_{i}^{k}\right\}$ in Section 13.2B as $\left\{\gamma_{i}^{z_{0}}\right\}$, we seek to explain what is meant by a process $\left\{\gamma_{i}^{z_{0}}\right\}$ where $z_{0}$ is an arbitrary critical point. Questions surrounding the domains of definition of $\left\{\gamma_{i}^{z_{0}}\right\}$ are treated in item (2) below. We discuss here the more basic question of whether to set $\gamma_{i}^{z_{0}}=z_{i}$ when it is $\neq *$.

Setting $\gamma_{i}^{z_{0}}=z_{i}$ for all $i$ is only natural, but it has the following drawback: Let $z_{0}^{\prime} \in B^{(j)}\left(z_{0}\right)$. Even though the orbits of $z_{0}$ and $z_{0}^{\prime}$ stay together for a long time, if treated independently, the canonical subdivisions accompanying $a \mapsto z_{i}(a)$ and $a \mapsto z_{i}^{\prime}(a)$ are likely to produce slightly different partitions, and rules of deletion such as those in Section 13.2C may yield slightly different results. However inconsequential, these differences are a technical nuisance.

To avoid this technical nuisance, we have elected to view $\left\{\gamma_{i}^{z_{0}}\right\}$ as associated with critical blobs. More precisely, let $z_{0}$ be a critical point of generation $j$. We let $B^{\left(j_{1}\right)} \supset B^{\left(j_{2}\right)} \supset \cdots \supset B^{\left(j_{n}\right)}$ be the complete chain of critical blobs containing $B^{(j)}\left(z_{0}\right)$; i.e., $j_{1}=1, j_{n}=j$, and for each $i$, there is no $B^{(k)}$, $j_{i}<k<j_{i+1}$, such that $B^{\left(j_{i}\right)} \supset B^{(k)} \supset B^{\left(j_{i+1}\right)}$. (See $\S 7.3$ for the geometry of critical regions.) We say $B^{\left(j_{i}\right)}$ is visible on the time interval $\left(j_{i-1} \theta^{-1}, j_{i} \theta^{-1}\right]$, thinking of it as "hidden" inside lower-generation critical blobs before time $j_{i-1} \theta^{-1}$ and no longer active after time $j_{i} \theta^{-1}$. During the time period when $B^{\left(j_{i}\right)}$ is visible, we set $\gamma_{\ell}^{z_{0}}=z_{\ell}^{\left(j_{i}\right)}$ or $*$ where $z_{0}^{\left(j_{i}\right)}=z_{0}^{*}\left(B^{\left(j_{i}\right)}\right)$. That is to say, $B^{(1)}\left(z_{0}\right)$ is visible in the first $\theta^{-1}$ iterates, and for $\ell \leq \theta^{-1}, \gamma_{\ell}^{z_{0}}=z_{\ell}^{(1)}$ or $*$. Since $z_{0}^{(1)}=z_{0}^{k}$ for some $k,\left\{\gamma_{i}^{z_{0}}\right\}$ is identical to one of the processes defined in the last subsection for $i \leq \theta^{-1}$. At time $\theta^{-1}$, the critical blob $B^{(1)}$ retires, and $B^{(2)}$ becomes visible. We have $\gamma_{\ell}^{z_{0}}=z_{\ell}^{(2)}$ or $*$ for $\theta^{-1}<\ell \leq 2 \theta^{-1}$, and so on. (Note that some of the orbit segments are visible for more than $\theta^{-1}$ steps due to the "skipping of generations"; see Section 7.3).

Assuming for the moment that all definitions are legitimate and all rules for deletion are as in Section 13.2C, we make the following observation: During the period when $\gamma_{\ell}^{z_{0}}=z_{\ell}^{\left(j_{i}\right)}$ or $*, d_{\mathcal{C}}\left(z_{\ell}^{\left(j_{i}\right)}\right)>3 e^{-\alpha \ell}$ implies that $d_{\mathcal{C}}\left(\xi_{\ell}\right)>2 e^{-\alpha \ell}$ for all $\xi_{0} \in B^{\left(j_{i}\right)}$. This follows from earlier estimates on the sizes of critical blobs; see Section 4.2. Moreover, if $d_{\mathcal{C}}\left(\gamma_{\ell}^{z_{0}}\right)>3 e^{-\alpha \ell}$ for all $\ell \leq j \theta^{-1}$, then all $\xi_{0} \in B^{(j)}\left(z_{0}\right)$ satisfy $d_{\mathcal{C}}\left(\xi_{\ell}\right)>2 e^{-\alpha \ell}$. In particular, all critical points inside $B^{(j)}\left(z_{0}\right)$ obey (A2) ${ }^{\#}$ up to time $j \theta^{-1}$. The same conclusion is valid for (A4) ${ }^{\#}$ since up to time $j \theta^{-1}$, all $\xi_{0} \in B^{(j)}\left(z_{0}\right)$ can be regarded as having the same itinerary. Hence they have the same fraction of "bad iterates" in the sense of $B(\cdot, \hat{\delta} ; n, 2 n)$.
(2) Stabilization of critical structures and extending the processes $\left\{\gamma_{i}^{z_{0}}\right\}$. In Section 13.2, we considered processes defined by $z_{0} \in \Gamma_{1}$, which has a continuation on all of $\Delta_{0}$. Critical structures of higher generations do not have such continuations. To stabilize these structures, we introduce an increasing sequence of partitions $\mathcal{J}_{\theta^{-1}}<\mathcal{J}_{2 \theta^{-1}}<\mathcal{J}_{4 \theta^{-1}}<\cdots$ on $\Delta_{0}$ with the following properties: $\mathcal{J}_{\theta^{-1}}=\left\{\Delta_{0}\right\}$; for each $N=2^{\ell} \theta^{-1}, \ell \geq 1, \mathcal{J}_{N}$ is a refinement of $\mathcal{J}_{\frac{1}{2} N}$ and partitions $\Delta_{0}$ into intervals of length $\approx \rho^{\alpha^{*} N}$. Leaving precise rules of deletion to part C, we explain here the relation between these partitions and the processes defined in (1).

For $N=2^{\ell} \theta^{-1}, \ell=1,2, \cdots$, the picture is as follows:
(i) There is a decreasing sequence of "good sets" $\Delta_{N}$ with the property that $T_{a} \in \mathcal{G}_{N}^{\#}$ for all $a \in \Delta_{N}$. (This is not the definition of $\Delta_{N}$, however.)
(ii) There are subcollections of "good" intervals $\mathcal{J}_{N}^{*} \subset \mathcal{J}_{N}$. For each $J_{N} \in$ $\mathcal{J}_{N}^{*}$,

- $J_{N} \subset J_{\frac{1}{2} N}$ for some $J_{\frac{1}{2} N} \in \mathcal{J}_{\frac{1}{2} N}^{*}$,
- $J_{N} \cap \Delta_{\frac{1}{2} N} \neq \emptyset$.
(iii) For each $J_{N} \in \mathcal{J}_{N}^{*}$, let $\Gamma^{N}\left(J_{N}\right)$ be the set of critical points of $T_{a}, a \in J_{N}$, of generations between $\frac{1}{2} \theta N$ and $\theta N$. Then for each $z_{0} \in \Gamma^{N}\left(J_{N}\right)$ of generation $k$, there is a well defined process $\left\{\gamma_{i}^{z_{0}}, i \leq \min \left(k \theta^{-1}, N\right)\right\}$, the domains of definition of which are as follows: Let $\Delta_{0}=J_{\theta^{-1}} \supset J_{2 \theta^{-1}} \supset$ $\cdots \supset J_{N}$ be the elements of $\mathcal{J}_{2^{\ell} \theta^{-1}}$ containing $J_{N}$. Then
- $\gamma_{i}^{z_{0}}, i \leq \theta^{-1}$, is defined on $\Delta_{0}=J_{\theta^{-1}}$ (this is what is constructed in Section 13.2);
- the process above is extended from $i=\theta^{-1}$ to $i=2 \theta^{-1}$ on $J_{2 \theta^{-1}}$, then from $i=2 \theta^{-1}$ to $i=4 \theta^{-1}$ on $J_{4 \theta^{-1}}$, and so on, up to $i=\frac{1}{2} N$;
- the product of this last extension is extended from $i=\frac{1}{2} N$ to $i=$ $\min \left(k \theta^{-1}, N\right)$ on $J_{N}$ where $k$, as we recall, is the generation of $z_{0}$.

We explain how to go from step $N$ to step $2 N$, clarifying along the way what exactly is meant by some of the statements in (iii) and how they can be achieved:

Elements of $\mathcal{J}_{N}$ not in $\mathcal{J}_{N}^{*}$ are discarded since all parameters in them have been deleted in a previous step (second property of $J_{N}$ in assumption (ii) above). Let $J_{N} \in \mathcal{J}_{N}^{*}$ be fixed. We consider $\left.\mathcal{J}_{2 N}\right|_{J_{N}}$, and put into $\mathcal{J}_{2 N}^{*}$ those elements of $\mathcal{J}_{2 N}$ that meet $\Delta_{N}$ (as required by (ii)). Consider a (fixed) $J \subset J_{N}$ such that $J \in \mathcal{J}_{2 N}^{*}$. Since there exists $\hat{a} \in J$ such that $T_{\hat{a}} \in \mathcal{G}_{N}^{\#}$ (assumption (ii)), $T_{a} \in \mathcal{G}_{2 \alpha^{*} N}$ for all $a \in J$ (Proposition 10.1). Thus on $J$ there is a coherent choice of $\Gamma_{2 \theta N}$ whose orbits can be treated up to time $\min \left(k \theta^{-1}, 2 N\right)$ where $k$ is the generation of the critical point.

Fix $z_{0} \in \Gamma^{2 N}(J)$. Since $k$, the generation of $z_{0}$, is $\geq \theta N, \gamma_{i}^{z_{0}}$ is defined for all $i \leq N$ in the sense of (iii). As mandated by (iii), we now seek to extend this process to all $i \leq \min \left(k \theta^{-1}, 2 N\right)$ on the interval $J$. Such an extension is carried out on one $\omega \in \mathcal{Q}_{N}^{z_{0}}$ at a time. Fix $\omega$ such that $\gamma_{N}^{z_{0}} \mid \omega \neq *$. If $\omega \subset J$, then we consider $\gamma_{i}^{z_{0}}$ for $i=N+1, N+2, \cdots$ starting from $\omega$ as explained in Section 13.2B. If $\omega \cap J=\emptyset$, then $\omega$ is not our concern. It remains to consider the case $\omega \cap \partial J \neq \emptyset$.

Observation. If for all $a \in \omega, z_{0}(a)$ satisfies the hypotheses of Proposition 11.1 up to time $N$, then $|\omega| \ll|J|$.

Indeed, $\omega \in \mathcal{Q}_{N}^{z_{0}}$ and $J \in \mathcal{J}_{N}$ have exponentially different length scales. This is because by Proposition 11.1, $|\omega|<\hat{K} e^{-\hat{\lambda} N}$, which is $\ll \rho^{\alpha^{*} N}=|J|$, and our rules of deletion, which are stated precisely in part C below, are built to ensure that the hypotheses of Proposition 11.1 are satisfied for $z_{0}(a)$ for every $a \in\left\{\gamma_{N}^{z_{0}} \neq *\right\}$. To deal with those $\omega$ that intersect some $J \in \mathcal{J}_{N}^{*}$ but are not completely contained in it, we let $J^{+}$be, say, $10 \%$ longer than $J$, and treat all $\omega \in \mathcal{Q}_{N}^{z_{0}}$ that are completely contained in $J^{+}$. The properties of $J$ continue to be valid in $J^{+}$, and these overlapping intervals lead to an overcount by a factor of at most 2 . This completes the qualitative description of the extension of (iii).

We finish with the inductive definition of $\Delta_{N}$, even though the following acquires meaning only after the deletion rules are specified. We let

$$
\Delta_{2 N}^{z_{0}}:=\left\{\gamma_{\min \left(k \theta^{-1}, 2 N\right)}^{z_{0}} \neq *\right\}, \quad \Delta_{2 N}(J):=\cap_{z_{0} \in \Gamma^{2 N}(J)} \Delta_{2 N}^{z_{0}}
$$

and

$$
\Delta_{2 N}:=\Delta_{N} \cap\left(\cup_{J \in \mathcal{J}_{2 N}^{*}} \Delta_{2 N}(J)\right)
$$

Our deletion rules are designed to ensure that $\Delta_{N}$ as defined above has the property in assumption (i), and that $\cap_{N} \Delta_{N}$ has positive measure.

Remarks. (1) The intersection with $\Delta_{N}$ in the definition of $\Delta_{2 N}$ may seem redundant, for on $J \in \mathcal{J}_{2 N}^{*}$, all critical blobs corresponding to $z_{0} \in \Gamma^{2 N}(J)$ have already been treated up to time $\min \left\{k^{-1} \theta, 2 N\right\}$. Consequently, it is tempting to claim that $T_{a} \in \mathcal{G}_{2 N}^{\#}$ for all $a \in \cup_{J \in \mathcal{J}_{2 N}^{*}} \Delta_{2 N}(J)$. This is not true in general, for not every critical blob has offspring (meaning smaller critical blobs inside), and the definition of $\cup_{J \in \mathcal{J}_{2 N}^{*}} \Delta_{2 N}(J)$ does not take into consideration the behavior of critical blobs that expired without reproducing before time $N$. To ensure that $T_{a} \in \mathcal{G}_{2 N}^{\#}$ for all $a \in \Delta_{2 N}$, we require that $\Delta_{2 N} \subset \Delta_{N}$, and use the inductively obtained fact that for all $a \in \Delta_{N}$, all $z_{0} \in \Gamma^{N}(J)$ are well behaved.
(2) To deal with the phenomenon called "skipping of generations", we need to work with slightly overlapping intervals of generations to ensure that all critical behaviors are represented. For example, we should have included in the definition of $\Gamma^{N}(J)$ all critical points from generation $\theta N$ to generation $2 \theta N(1+2 \theta)$, and the elements of $\mathcal{J}_{N}$ should have been taken to be of length $\approx \frac{1}{2} \rho^{\alpha^{*} N(1+2 \theta)}$, and so on. We have omitted - and will continue to omit - all of these factors of $(1+2 \theta)$ to slightly simplify the discussion. The problem is easy to rectify (and should probably be ignored on first pass).
C. Formal procedure from step $N=2^{\ell} \theta^{-1}$ to step $2 N$. We now give the formal procedure at a generic step $N$. The following should not be thought of as induction hypotheses, but rather as a summary of the situation as we arrive at step $N$ following the procedure described in part B. At time $N$ we assume we have the following:
(a) A subcollection $\mathcal{J}_{N}^{*}$ of $\mathcal{J}_{N}$ with the property that on each $J \in \mathcal{J}_{N}^{*}$, there is a coherent choice of $\Gamma_{\theta N} ; \Gamma^{N}(J)$ and $\Delta_{N}(J)$ are as defined above.
(b) On each $J \in \mathcal{J}_{N}^{*}$, associated with each $z_{0} \in \Gamma^{N}(J)$ of generation $k$, is a process $\gamma_{i}^{z_{0}}: J \rightarrow R_{1} \cup\{*\}$ for which the following hold: For all $a \in \Delta_{N}^{z_{0}}$ and $i \leq \min \left(k \theta^{-1}, N\right)$,
(i) $d_{\mathcal{C}}\left(\gamma_{i}^{z_{0}}(a)\right)>3 e^{-\alpha i}$;
(ii) $B\left(a, \hat{\delta} ; 2^{j} \theta^{-1}, 2^{j+1} \theta^{-1}\right)<\sigma 2^{j} \theta^{-1}$.
(c) A subset of $\cup_{J \in \mathcal{J}_{N}^{*}} \Delta_{N}(J)$ in $\left\{a: T_{a} \in \mathcal{G}_{N}^{\#}\right\}$ is called $\Delta_{N}$.

As noted in $\mathrm{B}(1)$ above, all $z_{0} \in \Gamma^{N}(J)$ of generation $k$ obey (A2) \# and (A4) ${ }^{\#}$ up to time $\min \left(k \theta^{-1}, N\right)$. Step (c) is needed because $\cup_{J \in \mathcal{J}_{N}^{*}} \Delta_{N}(J)$ is not necessarily in $\left\{a: T_{a} \in \mathcal{G}_{N}^{\#}\right\}$ (see Remark (1) in $\mathrm{B}(2)$ above).

What is done from time $N$ to time $2 N$.
(0) First we introduce the partition $\mathcal{J}_{2 N}$, and let $\mathcal{J}_{2 N}^{*} \subset \mathcal{J}_{2 N}$ be the collection of $J$ with $J \cap \Delta_{N} \neq \emptyset$. Elements of $\mathcal{J}_{2 N} \backslash \mathcal{J}_{2 N}^{*}$ are excluded from further consideration.

We then treat one $J \in \mathcal{J}_{2 N}^{*}$ at a time, carrying out for it steps (1)-(4) below. Steps (1)-(3) are carried out for each $z_{0} \in \Gamma^{2 N}(J)$, beginning with the $z_{0}$ of the lowest generations.
(1) Set $\left.\gamma_{N+1}^{z_{0}}\right|_{\omega}=*$ on those $\omega \in \mathcal{Q}_{N}^{z_{0}}$ with $\omega \cap \Delta_{N}=\emptyset$.
(2) On the rest of the $\omega \in \mathcal{Q}_{N}^{z_{0}}$, we continue the process to time $2 N$ in the manner described above, deleting along the way all $\omega^{\prime} \in \mathcal{Q}_{i}^{z_{0}}$ with $d_{\mathcal{C}}\left(\gamma_{i}^{z_{0}}\left(\omega^{\prime}\right)\right)<3 e^{-\alpha i}$.
(3) Set $\left.\gamma_{2 N}^{z_{0}}\right|_{\omega}=*$ on those $\omega \in \mathcal{Q}_{2 N}^{z_{0}}$ with the property that $B(a, \hat{\delta}, N, 2 N)>$ $\sigma N$ for $a \in \omega$.
(4) Define $\Delta_{2 N}^{z_{0}}=\left\{\gamma_{\min \left(k \theta^{-1}, 2 N\right)}^{z_{0}} \neq *\right\}$ and $\Delta_{2 N}(J)=\cap_{z_{0} \in \Gamma^{2 N}(J)} \Delta_{2 N}^{z_{0}}$ as in Part B.

Finally, after all the $J \in \mathcal{J}_{2 N}^{*}$ are treated, we set

$$
\Delta_{2 N}=\Delta_{N} \cap\left(\cup_{J \in \mathcal{J}_{2 N}^{*}} \Delta_{2 N}(J)\right) .
$$

Step (0) is to ensure the existence of a coherent choice of $\Gamma_{2 \theta N}$ on each selected $J$. Step (1) is to ensure that the process can legitimately be extended on those $\omega$ on which $\gamma_{N+1}^{z_{0}} \neq *$. Note also that every $z_{0}$ has an ancestor, so all $\gamma_{i}^{z_{0}}$ are extensions of previously constructed processes. Since many of the $z_{0} \in \Gamma^{2 N}(J)$ are related to each other via ancestry, the steps above in fact contain many duplications.

It is evident that the steps above lead to (a)-(c) at the beginning of Part B for time $2 N$.
D. Measure deleted from time $N$ to time $2 N$. First we estimate the measure deleted on account of a fixed $J \in \mathcal{J}_{2 N}^{*}$ and a fixed $z_{0} \in \Gamma^{2 N}(J)$ : Step (0) does not contribute to $\Delta_{N} \backslash \Delta_{2 N}$ since no $a \in \cup \mathcal{J}_{N}^{*} \backslash \cup \mathcal{J}_{2 N}^{*}$ belongs in $\Delta_{N}$. The same remark holds for step (1). Explanations and estimates for steps (1)-(3) are exactly as before, except that $\left|\Delta_{0}\right|$ should be replaced by $\left|J^{+}\right|$. Thus we have

$$
\left|\mathcal{D}_{N, 2 N}^{z_{0}}\right| \leq K^{\prime} e^{-\varepsilon^{\prime} N} \cdot 2|J|
$$

Since the cardinality of $\Gamma^{2 N}$ is $\leq 2 N \theta K_{2}^{2 \theta N(1+2 \theta)}$ (Corollary 7.1), we have

$$
\left|\mathcal{D}_{N, 2 N}\right| \leq \sum_{J \in \mathcal{J}_{2 N}^{*}} \sum_{z_{0} \in \Gamma^{2 N}}\left|\mathcal{D}_{N, 2 N}^{z_{0}}\right| \leq 2 N \theta K_{2}^{2 \theta N(1+2 \theta)} \cdot K^{\prime} e^{-\varepsilon^{\prime} N} \cdot 2\left|\Delta_{0}\right|
$$

### 13.4. The final count

From Sections $13.1,13.2 \mathrm{C}$ and 13.3 D , we see that the total measure deleted at the end of the procedure is

$$
\leq\left(K \delta_{0}^{-1} e^{-\alpha n_{0}}+q K^{\prime} \sum_{N=2^{\ell} n_{0}, \ell \in \mathbb{Z}^{+}} N \theta K_{2}^{2 \theta N(1+2 \theta)} \cdot e^{-\varepsilon^{\prime} N}\right)\left|\Delta_{0}\right|
$$

As $(a, b) \rightarrow\left(a^{*}, 0\right), n_{0} \rightarrow \infty$ and $\theta \rightarrow 0$, but none of the other constants is affected. Thus with $(a, b)$ sufficiently near $\left(a^{*}, 0\right)$, the quantity in parenthesis can be made arbitrarily small. In other words, $|\Delta|$ can be made as large a fraction of $\left|\Delta_{0}\right|$ as we wish. This completes the proof of (ii) in the first paragraph of this section.

## Appendices

## A.1. Properties of "good" 1 D maps ( $\S \S 2.1$ and 2.2)

Proof of Lemma 2.1. Let $x$ be such that $f^{i}(x) \notin C_{\delta}$ for $i \in[0, n)$. We divide $[0, n]$ into maximal time intervals $[i, i+k]$ such that $f^{i+j}(x) \notin C_{\delta_{0}}$ for $0<j<k$, and estimate $\left|\left(f^{k}\right)^{\prime}\left(f^{i}(x)\right)\right|$ as follows:

Case 1. $f^{i}(x), f^{i+k}(x) \in C_{\delta_{0}}$. Definitions 1.1(b)(ii) and (c)(ii) together guarantee that $\left|\left(f^{k}\right)^{\prime}\left(f^{i}(x)\right)\right| \geq e^{\frac{1}{3} \lambda_{0} k}$.

Case 2. $f^{i}(x) \notin C_{\delta_{0}}, f^{i+k}(x) \in C_{\delta_{0}}$, same as Definition 1.1(b)(ii).
Case 3. $f^{i}(x), f^{i+k}(x) \notin C_{\delta_{0}}$. If $k \geq M_{0}$, then $\left|\left(f^{k}\right)^{\prime}\left(f^{i}(x)\right)\right|>e^{\lambda_{0} k}$ from Definition 1.1(b)(i). If $k<M_{0}$, we let $\hat{k}$ be the smallest integer $>k$ such that $f^{i+\hat{k}}(x) \in C_{\delta_{0}}$. Using Definition 1.1(b)(i) for $\hat{k} \geq M_{0}$ and Definition 1.1(b)(ii) for $\hat{k}<M_{0}$, we conclude that $\left|\left(f^{k}\right)^{\prime}\left(f^{i}(x)\right)\right|>c_{0} K_{0}^{-M_{0}} e^{\lambda_{0} k}$ where $K_{0}=\max \left|f^{\prime}(x)\right|$.

Case 4. $f^{i}(x) \in C_{\delta_{0}}, f^{i+k}(x) \notin C_{\delta_{0}}$, same as Case 3, with an extra factor $\geq\left(\min _{y \in C_{\delta_{0}}}\left|f^{\prime \prime}(y)\right|\right) \delta$.

Cases 3 and 4 are relevant only for part (a).
Proof of Lemma 2.2. Proceed as in the proof of Lemma 2.1. From Definitions 1.1(b)(i) and (c)(ii), we see that for $f$, the estimates in all four cases are determined by $\left|\left(f^{j}\right)^{\prime}(y)\right|$ for $y \notin C_{\delta}$ and $j \leq N:=\max \left(M_{0}, K \log \frac{1}{\delta}\right)$. Choose $g$ sufficiently near $f$ that $\left|g^{j}(y)-f^{j}(y)\right|$ is sufficiently small for all $y \notin C_{\delta}$ and $j \leq N$.

We will use the notation $x_{i}=f^{i}(x)$.
Proof of Proposition 2.1. (P1) is Lemma 2.2. Let $x \in C_{\delta}(\hat{x})$.
Sublemma A.1.1. For all $y \in[\hat{x}, x]$ and $k<p$,

$$
\frac{1}{2} \leq \frac{\left(f^{k}\right)^{\prime}\left(y_{1}\right)}{\left(f^{k}\right)^{\prime}\left(\hat{x}_{1}\right)} \leq 2
$$

provided that $\delta$ and $\varepsilon$ are sufficiently small.
Proof. We write

$$
\log \frac{\left(f^{k}\right)^{\prime}\left(y_{1}\right)}{\left(f^{k}\right)^{\prime}\left(\hat{x}_{1}\right)} \leq \sum_{j=1}^{k} \frac{\left|f^{\prime}\left(y_{j}\right)-f^{\prime}\left(\hat{x}_{j}\right)\right|}{\left|f^{\prime}\left(\hat{x}_{j}\right)\right|} \leq K \sum_{j=1}^{k} \frac{\left|y_{j}-\hat{x}_{j}\right|}{d\left(\hat{x}_{j}, C\right)}
$$

We first choose $h_{0}$ large enough that $\frac{1}{\delta_{0}} \sum_{i=h_{0}+1}^{\infty} e^{-2 \alpha j} \ll 1$, followed by $\delta$ small enough that $\delta \sum_{j=1}^{h_{0}} \frac{1}{\delta_{0}} K^{j} \ll 1$. We then require $\varepsilon$ to be sufficiently small so that $d\left(\hat{x}_{j}, C\right)>\delta_{0} \forall j<n_{0}$ for some $n_{0}$ satisfying $e^{-\alpha n_{0}}<\delta$. These choices ensure that

$$
\sum_{j=1}^{k} \frac{\left|y_{j}-\hat{x}_{j}\right|}{d\left(\hat{x}_{j}, C\right)}<\sum_{j=1}^{h_{0}} \frac{1}{\delta_{0}} K^{j} \delta+\sum_{j=h_{0}+1}^{n_{0}} \frac{1}{\delta_{0}} e^{-2 \alpha j}+\sum_{j=n_{0}+1}^{k} e^{-(2 \alpha-\alpha) j} \ll 1
$$

Proof of (P2). Suppose $|x-\hat{x}|=e^{-h}$. Then (G2) together with the sublemma above imply that

$$
\left|x_{p}-\hat{x}_{p}\right|=\left|\left(f^{p-1}\right)^{\prime}\left(y_{1}\right)\right|\left|x_{1}-\hat{x}_{1}\right| \geq K^{-1} e^{\lambda(p-1)}(x-\hat{x})^{2}
$$

From $\left|x_{p}-\hat{x}_{p}\right|<1$, we read off the upper bound $p<\frac{3}{\lambda} h$ for $h$ sufficiently large. For the lower bound, we write $\left|x_{p}-\hat{x}_{p}\right|<K^{p-1} e^{-2 h}$ and recall that $p$ is defined such that $\left|x_{p}-\hat{x}_{p}\right| \geq e^{-2 \alpha p}$. That $p>$ constant $\cdot h$ follows directly from $K^{p-1} e^{-2 h}>e^{-2 \alpha p}$.

To prove (P2)(ii) we again write $\left|x_{p}-\hat{x}_{p}\right|<K\left|\left(f^{p-1}\right)^{\prime}\left(y_{1}\right)\right|(x-\hat{x})^{2}$, so that

$$
\begin{equation*}
K\left|\left(f^{p-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|^{\frac{1}{2}}|x-\hat{x}|>e^{-\alpha p} \tag{18}
\end{equation*}
$$

We also have

$$
\left|\left(f^{p}\right)^{\prime}(x)\right|=\left|\left(f^{p-1}\right)^{\prime}\left(x_{1}\right)\right|\left|f^{\prime}(x)\right|>\left(K^{-1}\left|\left(f^{p-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|\right) \cdot\left(K^{-1}|x-\hat{x}|\right) .
$$

Combined with (18) this gives

$$
\left|\left(f^{p}\right)^{\prime}(x)\right|>K^{-3}\left|\left(f^{p-1}\right)^{\prime}\left(\hat{x}_{1}\right)\right|^{\frac{1}{2}} e^{-\alpha p}>\hat{c}_{1}^{\frac{1}{2}} K^{-3} e^{\frac{1}{2} \lambda(p-1)} e^{-\alpha p}
$$

which we may assume is $>e^{-\frac{1}{3} \lambda p}$ if $p$ is sufficiently large, or equivalently, $\delta$ is sufficiently small.

It remains to prove (P2)(iii). From (P2)(i), (ii) and Sublemma A.1.1, for $\left.I_{\mu j} \in \mathcal{P}\right|_{C_{\delta}(\hat{x})}$,
$\left|f^{p}\left(I_{\mu j}\right)\right| \geq K^{-1} \frac{\left|f\left(I_{\mu j}\right)\right|}{\left|f\left(\left[\hat{x}, \hat{x}+e^{-\mu}\right]\right)\right|}\left|f^{p}\left(\left[\hat{x}, \hat{x}+e^{-\mu}\right]\right)\right| \geq K^{-1} \frac{1}{\mu^{2}} e^{-2 \alpha p}>e^{-K \alpha|\mu|}$.

Proof of (P3). We write $\sigma_{0}=[x, y], \sigma_{k}=f^{t_{k}} \sigma_{0}$, and assume for definiteness that $\sigma_{0} \subset C_{\delta}$ and $n \geq t_{q}+p_{q}$. Then

$$
\log \frac{\left(f^{n}\right)^{\prime}(x)}{\left(f^{n}\right)^{\prime}(y)} \leq \sum_{j=0}^{n-1} \frac{\left|f^{\prime}\left(y_{j}\right)-f^{\prime}\left(x_{j}\right)\right|}{\left|f^{\prime}\left(y_{j}\right)\right|} \leq K \sum_{k=1}^{q}\left(S_{k}^{\prime}+S_{k}^{\prime \prime}\right)
$$

where

$$
S_{k}^{\prime}=\sum_{j=t_{k}}^{t_{k}+p_{k}-1} \frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)} \quad \text { and } \quad S_{k}^{\prime \prime}=\sum_{t_{k}+p_{k}}^{t_{k+1}-1} \frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)}
$$

except for $S_{q}^{\prime \prime}$ which ends at index $n-1$.
I. Bound on $\sum_{k=1}^{q} S_{k}^{\prime \prime}$. For $k<q$ and $t_{k}+p_{k} \leq j<t_{k+1}-1$, we have, by (P1)(ii), $\left|\sigma_{k+1}\right| \geq c_{1} e^{\lambda\left(t_{k+1}-j\right)}\left|x_{j}-y_{j}\right|$, so that $S_{k}^{\prime \prime} \leq K \frac{\left|\sigma_{k+1}\right|}{\delta}$. Also, by combining (P2)(ii) and (P1)(ii), we have $\left|\sigma_{k+1}\right| \geq e^{\frac{1}{3} \lambda\left(t_{k+1}-t_{k}\right)}\left|\sigma_{k}\right| \geq \tau\left|\sigma_{k}\right|$ for some $\tau>1$, so that $\sum_{t=0}^{q-1} S_{k}^{\prime \prime} \leq K \frac{\left|\sigma_{q}\right|}{\delta}$.

The term $S_{q}^{\prime \prime}$ is treated differently because $x_{n}$ may not be a return. Observe the following: (i) If $\left[x_{n}, y_{n}\right] \subset C_{\delta_{0}}$, then (P1)(ii) gives, as before, $S_{q}^{\prime \prime} \leq$ $\frac{1}{\delta} K\left|y_{n}-x_{n}\right|$ which is $\leq K$ since $\left|y_{n-1}-x_{n-1}\right| \lesssim \delta$ by definition. (ii) If for $t_{q}+p_{q} \leq j<n,\left[x_{j}, y_{j}\right] \cap C_{\delta_{0}}=\emptyset$, then (P1)(i) with $\delta_{0}$ in the place of $\delta$ gives $S_{q}^{\prime \prime} \leq \frac{1}{\delta_{0}} \frac{K}{\delta_{0}}\left|y_{n}-x_{n}\right| \leq K \delta \delta_{0}^{-2} \leq K$. In general, if there is $\hat{n} \geq t_{q}+p_{q}$ such that $\hat{n}$ is the last return to $C_{\delta_{0}}$ before time $n$, then we apply (i) to $\sum_{t_{q}+p_{q}}^{\hat{n}}$ and (ii) to $\sum_{\hat{n}+1}^{n-1}$.
II. Bound on $\sum_{k=1}^{q} S_{k}^{\prime}$. First we estimate $S_{k}^{\prime}$. Suppose $y_{t_{k}} \in C_{\delta}(\hat{x})$. For $t_{k}<j<t_{k}+p_{k}$ we write

$$
\frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)}=\frac{\left|y_{j}-x_{j}\right|}{\left|y_{j}-\hat{x}_{j-t_{k}}\right|} \cdot \frac{\left|y_{j}-\hat{x}_{j-t_{k}}\right|}{d\left(y_{j}, C\right)} .
$$

By Sublemma A.1.1 and the usual estimates near $\hat{x}$, the first factor on the right is

$$
<K \frac{\left|y_{t_{k}+1}-x_{t_{k}+1}\right|}{\left|y_{t_{k}+1}-\hat{x}_{1}\right|}<K \frac{\left|f^{\prime}\left(x_{t_{k}}\right)\right|\left|y_{t_{k}}-x_{t_{k}}\right|}{\left|y_{t_{k}}-\hat{x}\right|^{2}}<K \frac{\left|\sigma_{k}\right|}{d\left(y_{t_{k}}, C\right)} .
$$

Thus

$$
\begin{aligned}
S_{k}^{\prime} & =\frac{\left|y_{t_{k}}-x_{t_{k}}\right|}{d\left(y_{t_{k}}, C\right)}+\sum_{j=t_{k}+1}^{t_{k}+p_{k}-1} \frac{\left|y_{j}-x_{j}\right|}{d\left(y_{j}, C\right)} \\
& \leq K \frac{\left|\sigma_{k}\right|}{d\left(y_{t_{k}}, C\right)}\left(1+\sum_{j=t_{k}+1}^{t_{k}+p_{k}-1} \frac{\left|y_{j}-\hat{x}_{j-t_{k}}\right|}{d\left(y_{j}, C\right)}\right)<K \frac{\left|\sigma_{k}\right|}{d\left(y_{t_{k}}, C\right)} .
\end{aligned}
$$

Now let $\mathcal{K}_{\mu}=\left\{k \leq q: \sigma_{k} \subset I_{\mu, j}\right.$ for some $\left.j\right\}$. Then

$$
\sum_{k \in \mathcal{K}_{\mu}} S_{k}^{\prime}<\sum_{k \in \mathcal{K}_{\mu}} K \frac{\left|\sigma_{k}\right|}{e^{-|\mu|}}<K \frac{1}{\mu^{2}}
$$

The first inequality is from above. The second follows from the next two facts: (i) $\left|\sigma_{k+1}\right| \geq \tau\left|\sigma_{k}\right|$ for $k<q$, and (ii) the term with the largest index is bounded above by $\left|I_{\mu, j}^{+}\right|$, which is $<K \frac{1}{\mu^{2}} e^{-\mu}$. To finish, we sum over all $\mu$ to obtain $\sum S_{k}^{\prime}<K$.

## A.2. Growth estimates and large deviations (§2.3)

To avoid cumbersome notation, we write $\mu$ instead of $|\mu|$ in all estimates.
Proof of Lemma 2.3. Since points in $\omega$ are assumed to have the same itinerary up to time $n,[0, n]$ is divided into bound intervals $\left(t_{k}, t_{k}+p_{k}\right)$ and free intervals $\left[t_{k}+p_{k}, t_{k+1}\right]$. From (P2)(ii), we have $\left|\gamma_{t_{k}+p_{k}}(\omega)\right|>e^{\frac{1}{4} \lambda p_{k}}\left|\gamma_{t_{k}}(\omega)\right|$ and from (P1)(ii), we have $\left|\gamma_{t_{k+1}}(\omega)\right|>c_{1} e^{\frac{1}{4} \lambda_{0}\left(t_{k+1}-t_{k}-p_{k}\right)}\left|\gamma_{t_{k}+p_{k}}(\omega)\right|$. Thus for any time $j$ such that $\gamma_{j}(\omega)$ is free, $\left|\gamma_{j}(\omega)\right|>e^{\frac{1}{5} \lambda j}|\omega|$. Now $\left|\gamma_{n}(\omega)\right|<1$ forces $n$ to be $<K \mu_{0}$.

Proof of Lemma 2.5. Let $s \in \omega$ be such that $S(s)>n$. We define the essential return times $t_{1}<t_{2}<\cdots$ and corresponding return addresses $I_{\mu_{1} j_{1}}^{i_{1}}, I_{\mu_{2} j_{2}}^{i_{2}}, \cdots$ for $s$ as follows: Let $t_{1}$ be the smallest $i>0$ when either (a) $\gamma_{i}(\omega)$ is out of bound period and $\left|\gamma_{i}(\omega)\right|>\delta$ or (b) $i$ is the extended bound period of $\gamma_{0}(\omega)$, whichever happens first. If (a) happens first, then $\left.S\right|_{\omega}=t_{1}$, and we stop iterating. If not, then we may assume $\gamma_{t_{1}}(\omega) \subset C_{\delta}$, and the return address of $s$ at time $t_{1}$ is $I_{\mu_{1} j_{1}}^{i_{1}}$ if $\gamma_{t_{1}}\left(\mathcal{Q}_{t_{1}}(s)\right) \approx I_{\mu_{1} j_{1}} \subset C_{\delta}\left(\hat{x}_{i_{1}}\right)$. Similarly, $t_{2}(s)$ is the first $i>t_{1}(s)$ when either (a) $\gamma_{i}\left(\mathcal{Q}_{t_{1}}(s)\right)$ is out of bound period and $\left|\gamma_{i}\left(\mathcal{Q}_{t_{1}}(s)\right)\right|>\delta$ or (b) $i$ is the extended bound period of $\mathcal{Q}_{t_{1}}(s)$, whichever happens first. Again if (a) happens first, then $\left.S\right|_{\mathcal{Q}_{t_{1}}(s)}=t_{2}$ and we stop considering $\mathcal{Q}_{t_{1}}(s)$; otherwise $\gamma_{t_{2}}\left(\mathcal{Q}_{t_{2}}(s)\right) \approx I_{\mu_{2} j_{2}} \subset C_{\delta}\left(\hat{x}_{i_{2}}\right)$, and so on.

Let $A_{q}=\left\{s \in \omega: S(s)>n\right.$, and $\gamma_{i}(s)$ makes a total of exactly $q$ essential returns before time $n\}$. Then $|\{S>n\}|=\sum_{q}\left|A_{q}\right|$. We write $A_{q}=\cup_{R} A_{q, R}$ where $A_{q, R}=\left\{s \in A_{q}\right.$ : if $\left(\mu_{1}, \cdots, \mu_{q}\right)$ are the $\mu$-coordinates of its first $q$ return addresses, then $\left.\left|\mu_{1}\right|+\left|\mu_{2}\right|+\cdots+\left|\mu_{q}\right|=R\right\}$. We further decompose $A_{q, R}$ into intervals $\sigma$ consisting of points whose first $q$ return addresses are identical. For $\sigma$ with return addresses $\left(I_{\mu_{1} j_{1}}, \cdots, I_{\mu_{q} j_{q}}\right)$, we let $Q_{t_{i}}=\mathcal{Q}_{t_{i}}(s)$ for $s \in \sigma$. Then

$$
|\sigma|=\frac{\left|Q_{t_{q}}\right|}{\left|Q_{t_{q-1}}\right|} \frac{\left|Q_{t_{q-1}}\right|}{\left|Q_{t_{q-2}}\right|} \cdots \frac{\left|Q_{t_{1}}\right|}{|\omega|}|\omega| \leq K^{q} \frac{\left|\gamma_{t_{q}}\left(Q_{t_{q}}\right)\right|}{\left|\gamma_{t_{q}}\left(Q_{t_{q-1}}\right)\right|} \cdots \frac{\left|\gamma_{t_{1}}\left(Q_{t_{1}}\right)\right|}{\left|\gamma_{t_{1}}(\omega)\right|}|\omega|
$$

where $K$ is the distortion constant in (P3). Now $\left|\gamma_{t_{q+1}}\left(Q_{t_{q}}\right)\right|<1$, and by (P2)(iii) and (P1)(ii), we have

$$
\frac{\left|\gamma_{t_{k}}\left(Q_{t_{k}}\right)\right|}{\mid \gamma_{t_{k+1}}\left(Q_{\left.t_{k}\right)} \mid\right.} \leq K \frac{\left|I_{\mu_{k} j_{k}}\right|}{\left|\gamma_{t_{k}+p_{k}}\left(Q_{t_{k}}\right)\right|} \leq K e^{-(1-K \alpha) \mu_{k}}
$$

Thus

$$
|\sigma|<K^{q} e^{-\sum_{k=1}^{q} \frac{9}{10} \mu_{k}+K \alpha \mu_{0}}|\omega|=K^{q} e^{-\frac{9}{10} R+K \alpha \mu_{0}}|\omega|:=|\sigma|_{R}
$$

(For $q=0$, this estimate presumes that $\gamma_{i}(\omega)$ has completed its initial bound period, i.e. $n>K \mu_{0}$.) We estimate $|\{S>n\}|$ by

$$
|\{S>n\}|=\sum_{q, R}\left|A_{q, R}\right| \leq \sum_{R}\left(\text { number of } \sigma \text { in } \cup_{q} A_{q, R}\right) \cdot|\sigma|_{R}
$$

There are $\binom{R-1}{q}$ ways of decomposing $R$ into a sum of $q+1$ integers. For a fixed $q$-tuple $\left(\mu_{1}, \cdots, \mu_{q}\right)$, we claim there are $\leq 2^{q} \mu_{1}^{2} \mu_{2}^{2} \cdots \mu_{q}^{2}$ possibilities for $\sigma$ with these data. This is because $\gamma_{t_{k}}\left(\mathcal{Q}_{t_{k}}(\sigma)\right)$ is short enough that it can meet at most one $C_{\delta}(\hat{x})$, which contains $\leq 2 \mu_{k}^{2}$ intervals of the form $I_{\mu_{k} j}$. Furthermore, for $\left(\mu_{1}, \cdots, \mu_{q}\right)$ with $\mu_{1}+\mu_{2}+\cdots+\mu_{q}=R$, we have $\mu_{1}^{2} \mu_{2}^{2} \cdots \mu_{q}^{2} \leq\left(\frac{R}{q}\right)^{2 q}$.

There is one other piece of information that is crucial to us, namely that all bound periods are $\geq \Delta:=K^{-1} \log \frac{1}{\delta}$. This means that for a given $R$, the only feasible $q$ are $\leq \frac{R}{\Delta}$. For a fixed $R$, then, the number of $\sigma$ in $\cup_{q} A_{q, R}$ is

$$
\leq \sum_{q}\binom{R-1}{q} \cdot 2^{q}\left(\frac{R}{q}\right)^{2 q} \leq \frac{R}{\Delta} \cdot\binom{R}{\frac{R}{\Delta}} \cdot 2^{\frac{R}{\Delta}} \Delta^{2 \frac{R}{\Delta}}
$$

which, by Sterling's formula, is $\sim \frac{R}{\Delta}\left(e^{\varepsilon\left(\frac{1}{\Delta}\right)} 2^{\frac{1}{\Delta}} \Delta^{\frac{2}{\Delta}}\right)^{R}$ where $\varepsilon\left(\frac{1}{\Delta}\right) \rightarrow 0$ as $\delta \rightarrow 0$. Calling the expression above $(1+\eta(\delta))^{R}$, we have $\eta(\delta) \rightarrow 0$ as $\delta \rightarrow 0$. Observe also that $n \leq K R+K \mu_{0}$ by Lemma 2.3 and so $R \geq K^{-1} n-\mu_{0}$. Thus

$$
\begin{aligned}
|\{S>n\}|< & \sum_{R \geq K^{-1} n-\mu_{0}} K^{q}(1+\eta(\delta))^{R} e^{-\frac{9}{10} R+K \alpha \mu_{0}}|\omega| \\
& <e^{-\frac{4}{5} K^{-1} n+\mu_{0}}|\omega|<e^{-\frac{1}{2} K^{-1} n}|\omega|
\end{aligned}
$$

provided that $n>3 K \mu_{0}$.

Proof of Corollary 2.1. Let $t_{1} \geq 0$ be the smallest $i$ such that there are points $s, s^{\prime} \in \omega$ with different itineraries in their first $i$ iterates. Then either $t_{1}=0$, or $t_{1}<K \log \delta^{-1}$ and $\left|\gamma_{t_{1}}(\omega)\right|>K^{-1} \delta$ by (P1). Let $n$ be an arbitrary integer $>t_{1}$. We partition $\omega$ into $\{\hat{\omega}\} \cup\left\{\omega_{\mu j}\right\}$ where $s \in \hat{\omega}$ if $\gamma_{t_{1}}\left(\mathcal{Q}_{t_{1}}(s)\right)$ is outside of $C_{\delta}$ and $\gamma_{t_{1}}\left(\omega_{\mu j}\right) \approx I_{\mu j}$. Then $\left.S\right|_{\hat{\omega}}=t_{1}$. For $\mu$ with $n>3 K \mu$, $\left|\omega_{\mu j} \cap\left\{S>n+t_{1}\right\}\right|<K e^{-\frac{1}{2} K^{-1} n}\left|\omega_{\mu j}\right|$ by Lemma 2.4; $K$ here is the distortion constant in (P3). Note also that the total length of $I_{\mu j}$ with $n \leq 3 K \mu$ is $\leq 2 e^{-\frac{1}{3} K^{-1} n}$. It follows therefore that
$\left|\left\{s \in \omega, S(s)>n+t_{1}\right\}\right|<K e^{-\frac{1}{2} K^{-1} n}|\omega|+K\left(\frac{2 e^{-\frac{1}{3} K^{-1} n}}{K^{-1} \delta}\right)|\omega|<e^{-\hat{K}^{-1}\left(n+t_{1}\right)}|\omega|$
provided $\hat{K}$ is sufficiently large and $n+t_{1}>\hat{K} \log \delta^{-1}$.
In the next proof, it is advantageous to take a probabilistic viewpoint, with $(\omega, P), P$ being normalized Lebesgue measure, as the underlying probability space.

Proof of Proposition 2.2. Let $\hat{\delta}>0$ be a small number to be determined, and let $B_{n}$ be as in the statement of the proposition. The idea of this proof is to introduce random variables $\hat{X}_{i}, i=0,1, \cdots$, with the property that (i) $B_{n} \leq \sum_{i \leq n} \hat{X}_{i}$ and (ii) the conditional expectations of $\hat{X}_{i}$ are dominated by certain exponential random variables.

Step I. Reformulation of problem as one involving $\sum_{i \leq n} \hat{X}_{i}$. Define a sequence of random variables $t_{1}<t_{2}<\cdots$ marking certain intersection times with $C_{\hat{\delta}}$ as follows: If $I_{\mu_{0} j_{0}} \subset C_{\hat{\delta}}$, let $t_{1}=0$, and let $S_{1}$ be the stopping time $S$ defined in Section 2.3. If $I_{\mu_{0} j_{0}} \cap C_{\hat{\delta}}=\emptyset$, let $t_{1}$ be the smallest $i$ for which $\gamma_{i}\left(\mathcal{Q}_{i-1}(s)\right) \cap C_{\hat{\delta}} \neq \emptyset$, and define $S_{1}$ on each element of $\mathcal{Q}_{t_{1}}$ as follows: If $\gamma_{t_{1}}\left(\mathcal{Q}_{t_{1}}(s)\right) \cap C_{\hat{\delta}}=\emptyset$, set $S_{1}(s)=0$. If $\gamma_{t_{1}}\left(\mathcal{Q}_{t_{1}}(s)\right) \approx I_{\mu j} \subset C_{\hat{\delta}}$, let $S_{1}$ be the stopping time $S$ on $\mathcal{Q}_{t_{1}}(s)$ for the sequence $\gamma_{t_{1}}, \gamma_{t_{1}+1}, \cdots\left(\right.$ instead of $\left.\gamma_{0}, \gamma_{1}, \cdots\right)$; that is to say, $S_{1}(s)$ is the smallest $i$ such that $\gamma_{t_{1}+i}\left(\mathcal{Q}_{t_{1}+i-1}(s)\right)$ is not in a bound period and $\left|\gamma_{t_{1}+i}\left(\mathcal{Q}_{t_{1}+i-1}(s)\right)\right|>\delta$. Then on each element of $\mathcal{Q}_{t_{1}}$, we define $t_{2}$ to be the smallest $i \geq t_{1}+S_{1}$ such that $\gamma_{i}\left(\mathcal{Q}_{i-1}(s)\right) \cap C_{\hat{\delta}} \neq \emptyset$, and on each $\mathcal{Q}_{t_{2}}(s)$, define $S_{2}$ to be either $S$ or 0 as before depending on whether $\gamma_{t_{2}}\left(\mathcal{Q}_{t_{2}}(s)\right) \subset C_{\hat{\delta}}$ or not, and so on.

Before proceeding further, we record the following lower estimate on $\left|\gamma_{t_{i}}\left(\mathcal{Q}_{t_{i}-1}(s)\right)\right|$. Let $t$ be the time $\mathcal{Q}_{t_{i}-1}(s)$ is created. By definition, $t_{i-1}+$ $S_{i-1} \leq t<t_{i}$, and $\gamma_{t}\left(\mathcal{Q}_{t_{i}-1}(s)\right) \approx P$ for some $P \in \mathcal{P}$. Moreover, there are only two possibilities: either $P$ is outside and $|P| \geq \delta$, or it is $\approx I_{\mu j}$ for some $I_{\mu j} \subset C_{\delta} \backslash C_{\hat{\delta}}$. By (P1) and (P2)(iii), $\left|\gamma_{t_{i}}\left(\mathcal{Q}_{t_{i}-1}(s)\right)\right|>\hat{\delta}^{\prime}:=\min \left(c_{2} \delta, \hat{\delta}^{K_{1} \alpha}\right)$. Note that if $\hat{\delta} \ll \delta$, then $\hat{\delta}^{\prime} \gg \hat{\delta}$.

We now head toward the promised random variables. For $i=0,1,2, \cdots$, let $X_{i}(s)=1$ for $i \in\left[t_{k}, t_{k}+S_{k}\right)$, any $k$, and $=0$ otherwise. Then $B_{n} \leq$
$\sum_{i \leq n} X_{i}$; in fact, this is likely to be an overcount, for $S_{k}$ goes beyond bound periods. It is thus sufficient to show that $P\left\{\sum_{i \leq n} X_{i}>\sigma n\right\}$ decreases exponentially with $n$.

As we will see, it simplifies the discussion slightly if we "speed up time" to skip over the intervals $\left[t_{k}, t_{k}+S_{k}\right)$. Let $T_{-1}=-1$. With $T_{i}$ defined, we let $T_{i+1}=T_{i}+1$ except when $T_{i}+1=t_{k}$, in which case we let $T_{i+1}=T_{i}+1+S_{k}$. We let $\hat{X}_{0}=S$ if $\omega \subset C_{\hat{\delta}}, 0$ otherwise, and let $\hat{X}_{i+1}=S_{k}$ if $T_{i}+1=t_{k}, 0$ otherwise. Let $\mathcal{Q}_{T_{i}}$ be the partition defined by $\mathcal{Q}_{T_{i}}(s)=\mathcal{Q}_{T_{i}(s)}(s)$, and note that $\hat{X}_{i}$ is measurable with respect to $\mathcal{Q}_{T_{i}}$. Since $X_{i} \leq \hat{X}_{i}$, it is all the more true that $B_{n} \leq \sum_{i \leq n} \hat{X}_{i}$.

Step II. Conditional distribution of $\hat{X}_{i+1}$ given $\mathcal{Q}_{T_{i}}$. Let $i \geq 0$, and consider $Q \in \mathcal{Q}_{T_{i}}$. On most $Q, \hat{X}_{i+1}$ is identically equal to 0 . The only time when this is not the case is when $\gamma_{T_{i}(Q)+1}(Q)$ meets $C_{\hat{\delta}}$. We note that
(1) for all $s, s^{\prime} \in Q, \gamma_{T_{i}(Q)+1}^{\prime}(s) / \gamma_{T_{i}(Q)+1}^{\prime}\left(s^{\prime}\right)<K$;
(2) $\left|\gamma_{T_{i}(Q)+1}(Q)\right|>\hat{\delta}^{\prime}$.
(1) follows from (P3); (2) is from Step I. From (1) and (2), we deduce that (i) $P\left(\hat{X}_{i+1}=0 \mid Q\right) \geq 1-K \hat{\delta} \hat{\delta}^{\prime}-1$ and (ii) $P\left(\hat{X}_{i+1}>n \mid Q \cap\left\{\gamma_{T_{i}+1} \in I_{\mu j}\right\}\right)<$ $K e^{-\frac{1}{2} K^{-1} n}$ if $n \geq 3 K \mu$ (Lemma 2.5); for $n<3 K \mu$, there is no information. It follows that for all $n \geq 0$,

$$
\begin{equation*}
P\left(\hat{X}_{i+1}>n \mid Q\right)<K \hat{\delta}^{\prime-1} \min \left(\hat{\delta}, e^{-(3 K)^{-1} n}\right)+K \hat{\delta} \hat{\delta}^{\prime-1} e^{-\frac{1}{2} K^{-1} n} . \tag{19}
\end{equation*}
$$

A simple computation shows that if $\varepsilon<\frac{1}{6} K^{-1}$ (where $K$ is as in the exponents above), then $E\left[e^{\varepsilon \hat{X}_{i+1}} \mid Q\right]<\infty$. We note further that by decreasing $\hat{\delta}$ (keeping $\varepsilon$ fixed), $E\left[e^{\varepsilon X_{i+1}} \mid Q\right]$ can be made arbitrarily close to 1 . Let $\eta>0$ be a number to be determined shortly, and choose $\hat{\delta}=\hat{\delta}(\eta)$ sufficiently small that $E\left[e^{\varepsilon \hat{X}_{i+1}} \mid Q\right]<e^{\eta}$. Observing that the upper bound in (19) and hence that for $E\left[e^{\varepsilon \hat{X}_{i+1}} \mid Q\right]$ do not depend on $i$ or on $Q$, we conclude that with $\hat{\delta}=\hat{\delta}(\eta)$ as above, $E\left[e^{\varepsilon \hat{X}_{i+1}} \mid \mathcal{Q}_{T_{i}}\right]<e^{\eta}$ for every $i \geq 0$.

Step III. Large deviation estimate for $\sum_{i \leq n} \hat{X}_{i}$. To finish, we write

$$
\begin{aligned}
E\left[e^{\varepsilon \sum_{i \leq n} \hat{X}_{i}}\right] & =E\left[E\left[e^{\varepsilon \sum_{i \leq n} \hat{X}_{i}} \mid \mathcal{Q}_{T_{n-1}}\right]\right] \\
& =E\left[e^{\varepsilon \sum_{i<n} \hat{X}_{i}} E\left[e^{\varepsilon \hat{X}_{n}} \mid \mathcal{Q}_{T_{n-1}}\right]\right]<e^{\eta} E\left[e^{\varepsilon \sum_{i<n} \hat{X}_{i}}\right],
\end{aligned}
$$

giving inductively $E\left[e^{\varepsilon \sum_{i \leq n} \hat{X}_{i}}\right]<e^{n \eta} E\left[e^{\varepsilon \hat{X}_{0}}\right]$. Since $E\left[e^{\varepsilon \hat{X}_{0}}\right]<e^{K \varepsilon \mu_{0}}$, we arrive at

$$
P\left\{B_{n}>\sigma n\right\}<P\left\{\sum_{i \leq n} \hat{X}_{i}>\sigma n\right\}<e^{\eta n-\varepsilon \sigma n+K \varepsilon \mu_{0}} .
$$

This is $<e^{-\frac{1}{2} \varepsilon \sigma n}$ if $\eta$ is chosen $<\frac{1}{4} \varepsilon \sigma$ and $n$ is $>K \mu_{0} \sigma^{-1}$.

## A.3. Parameter transversality (§2.4)

Proof of Lemma 2.6. First we show that $\cup_{i \geq 0} f^{-i} C$ is dense in $I$. If not, there would be an interval $\omega$ with the property that $\phi(x)$ is identical for all $x \in \omega$. Let $\omega$ be a maximal interval of this type. Then either (i) $f^{n+k}(\omega) \subset f^{n}(\omega)$ for some $n, k$, or (ii) $f^{k}(\omega), k=0,1, \cdots$, are pairwise disjoint. Case (i) cannot happen since it implies the presence of a periodic point $x$ with $\left|\left(f^{k}\right)^{\prime} x\right| \leq 1$. Case (ii) is equally absurd, for it implies the existence of $\left\{k_{i}\right\}$ where $f^{k_{i}}(\omega)$ are arbitrarily short and arbitrarily close to $C$, a scenario not permitted by Definition 1.1 (c)(ii) and Lemma 2.1.

For each $n$, let $l_{n}(\hat{x})$ and $r_{n}(\hat{x})$ be the two points in $\cup_{0<i \leq n} f^{-i} C$ closest to $\hat{x} \in C$. In the case $I=S^{1}$, let

$$
\Lambda^{(n)}=\left\{x \in I: f^{i} x \notin \cup_{\hat{x} \in C}\left(l_{n}(\hat{x}), r_{n}(\hat{x})\right) \forall i \geq 0\right\}
$$

If $I$ is an interval, we may assume $n$ is large enough that $f(I) \subset\left(z_{n}^{1}, z_{n}^{2}\right)$ where $z_{n}^{1}$ and $z_{n}^{2}$ are the two points in $\cup_{0 \leq i \leq n} f^{-i} C$ closest to the ends of $I$. We then define $\Lambda^{(n)}$ as in the circle case but with $I$ replaced by $\left[z_{n}^{1}, z_{n}^{2}\right]$. In both cases, $\Lambda^{(n)}$ is compact and $f\left(\Lambda^{(n)}\right) \subset \Lambda^{(n)}$. Clearly, $\cup \Lambda^{(n)}$ is dense in $I$ since $\cup_{i \geq 0} f^{-i} C$ is dense in $I$ and the gaps in $\Lambda^{(n)}$ decrease in size as $n$ increases.

For part (a), it remains to show that $\left.f\right|_{\Lambda^{(n)}}$ is conjugate to a shift of finite type. Let $\mathcal{J}^{(n)}=\left\{J_{i}^{(n)}\right\}$ be the partition of $I$ by $\cup_{0 \leq i \leq n} f^{-i} C$. Observe that for $J_{i}(n) \neq\left(l_{n}(\hat{x}), \hat{x}\right)$ or $\left(\hat{x}, r_{n}(\hat{x})\right), f\left(J_{i}^{(n)}\right)$ is equal to the union of a finite number of elements of $\mathcal{J}^{(n)}$. Let $\Lambda_{i}^{(n)}=\Lambda^{(n)} \cap J_{i}^{(n)}$. Then the alphabet of the shift in question is $\left\{i: \Lambda_{i}^{(n)} \neq \emptyset\right\}$, and the transition $i \rightarrow j$ is admissible if $f\left(\Lambda_{i}^{(n)}\right) \supset \Lambda_{j}^{(n)}$.

Assertion (b) follows from our construction.
Proof of Corollary 2.2. Fix $n$ large enough that for all $i \geq 0, f^{i}(q) \notin$ $\left(l_{n}(\hat{x}), r_{n}(\hat{x})\right)$ for all $\hat{x} \in C$, and let $\Lambda=\Lambda^{(n)}$. Let $B=\cup_{i} \partial \bar{\Lambda}_{i}$ where $\bar{\Lambda}_{i}$ is the shortest interval containing $\Lambda_{i}$. Since $B$ is a finite set with $f(B) \subset B$, it consists of pre-periodic points. From Lemma 2.1, these periodic points are repelling. Thus if $g$ is sufficiently near $f$, there is a unique set $B_{g}$ with $g\left(B_{g}\right) \subset B_{g}$ such that $\left.g\right|_{B_{g}}$ is conjugate to $\left.f\right|_{B}$. Using $B_{g}$, we recover a set $\Lambda_{g}$ on which $g$ is conjugate to $\left.f\right|_{\Lambda}$. The uniqueness of $q_{g}$ follows from the expanding property of $g$ away from $C$ (Lemma 2.2).

Proof of Proposition 2.3. (i) We prove $a \mapsto q(a)$ is differentiable with

$$
\begin{equation*}
\frac{d}{d a} q(a)=-\sum_{i=1}^{\infty} \frac{\partial_{a} f_{a}\left(f_{a}^{i-1}(q)\right)}{\left(f_{a}^{i}\right)^{\prime}(q)} \tag{20}
\end{equation*}
$$

Here all objects depend on $a$, mention of which is often suppressed (e.g. $f=f_{a}$, $q=q(a))$. Continuing to use the notation in Corollary 2.2, we let $\Lambda_{i_{0}, i_{1}, \cdots, i_{n}}=$ $\left\{x \in I: f^{j}(x) \in \Lambda_{i_{j}}, 0 \leq j \leq n\right\}$, and let $\Lambda_{i_{0}, i_{1}, \cdots, i_{n}}(q)$ be the cylinder set
containing $q$. For each $n$, choose $q_{n} \in \partial \bar{\Lambda}_{i_{0}, i_{1}, \cdots, i_{n}}(q)$. We will show that as functions of $a, \frac{d}{d a} q_{n}$ converge uniformly to the right side of (20). This requires in particular the uniform bound $\max _{i_{0}, i_{1}, \cdots, i_{n}}\left|\bar{\Lambda}_{i_{0}, i_{1}, \cdots, i_{n}}\right|<K e^{-\frac{1}{4} \lambda_{0} n}$ for all $n>0$ (see Lemma 2.1).

Introduce $G(x, a)=\left(f_{a}(x), a\right)$, and let $G^{n}\left(q_{n}, a\right)=\left(p_{n}, a\right)$. Then $p_{n} \in B$. Differentiating, we obtain

$$
\frac{d}{d a} p_{n}=\sum_{i=1}^{n} \partial_{x} f^{n-i}\left(f^{i} q_{n}\right) \partial_{a} f\left(f^{i-1} q_{n}\right)+\partial_{x} f^{n}\left(q_{n}\right) \frac{d}{d a} q_{n}
$$

Hence we have

$$
\begin{equation*}
\frac{d}{d a} q_{n}=\frac{\frac{d}{d a} p_{n}}{\partial_{x} f^{n}\left(q_{n}\right)}-\sum_{i=1}^{n} \frac{\partial_{a} f\left(f^{i-1} q_{n}\right)}{\partial_{x} f^{i}\left(q_{n}\right)} \tag{21}
\end{equation*}
$$

Since $B$ is a finite set, $\frac{d}{d a} p_{n}$ is uniformly bounded for all $n$. With $\left|\left(f^{n}\right)^{\prime}\left(q_{n}\right)\right|$ growing exponentially, the first term on the right is exponentially small. It remains to check that the second term converges uniformly to the right side of (20). In addition to the growth of $f^{n}$, this uses our estimates on max $\left|\bar{\Lambda}_{i_{0}, i_{1}, \cdots, i_{n}}\right|$ above and a distortion estimate for $f^{i}$. We leave it as an exercise.
(ii) This is a direct application of (20) to $q$ defined by $q\left(a^{*}\right)=f_{a^{*}}(\hat{x})$ :

$$
\begin{aligned}
\frac{\frac{d}{d a} f^{n}(\hat{x})}{\partial_{x} f^{n-1}\left(\hat{x}_{1}\right)} & =\frac{1}{\partial_{x} f^{n-1}\left(\hat{x}_{1}\right)}\left(\partial_{x} f^{n-1}\left(\hat{x}_{1}\right) \frac{d \hat{x}_{1}}{d a}+\partial_{a} f^{n-1}\left(\hat{x}_{1}\right)\right) \\
& =\frac{d \hat{x}_{1}}{d a}+\frac{\partial_{a} f^{n-1}\left(\hat{x}_{1}\right)}{\partial_{x} f^{n-1}\left(\hat{x}_{1}\right)} .
\end{aligned}
$$

Thus the limit as $n \rightarrow \infty$ at $a=a^{*}$ differs from $\frac{d q}{d a}\left(a^{*}\right)$ by $\frac{d \hat{x}_{1}}{d a}\left(a^{*}\right)$, which is also easily seen to be the term corresponding to $i=0$ in the sum in Proposition 2.3(ii).

## A.4. Most contracted directions: Preliminaries (§3.1)

We record in this appendix some elementary estimates in preparation for the proof of Lemma 3.1.
I. Area growth. Let $\left\{E_{1}, E_{2}, \cdots, E_{m}\right\}$ denote the usual basis of $\mathbb{R}^{m}$. Recall that if $u=\sum u_{i} E_{i}$ and $v=\sum v_{i} E_{i}$, then $u \wedge v=\sum_{i<j} I_{i j} E_{i} \wedge E_{j}$ where $I_{i j}=u_{i} v_{j}-v_{i} u_{j}$, and the area of the parallelogram spanned by $u$ and $v$ is equal to

$$
|u \wedge v|=\sqrt{|u|^{2}|v|^{2}-\langle u, v\rangle^{2}}=\left(\sum_{i<j} I_{i j}^{2}\right)^{\frac{1}{2}}
$$

Sublemma A.4.1. Let $M$ and $\tilde{M}$ have the properties in (H1) in Section 3.1. Then
(a) $|M u \wedge M v|<K b|u \wedge v|$;
(b) $|M u \wedge \tilde{M} v| \leq K b|u||v|$.

Consider next linearly independent unit vectors $u^{(0)}$ and $v^{(0)}$ parametrized by $s=\left(s_{1}, s_{2}\right)$. For $n=1,2, \cdots$, let $u^{(n)}=M_{n} u^{(n-1)}$ and $v^{(n)}=M_{n} v^{(n-1)}$, and let $u^{(n)} \wedge v^{(n)}=\sum_{i<j} I_{i j}^{(n)} E_{i} \wedge E_{j}$.

Sublemma A.4.2. For $k=0,1,2$, assume the respective versions of (H2), in which the $C^{2}$-norms are replaced by the corresponding $C^{k}$-norms. Then

$$
\left|\partial^{k}\left(u^{(n)} \wedge v^{(n)}\right)\right|<(K b)^{n} .
$$

It follows that if $\Delta_{n}=\left|u^{(n)} \wedge v^{(n)}\right|$, then $\left|\partial^{k} \Delta_{n}^{2}\right|<(K b)^{2 n}$.
II. Formulas for $e$ and $f$. We fix $M \in \mathrm{~L}(m, \mathbb{R})$ and $S=S(u, v)$ where for simplicity we assume $u$ and $v$ are unit vectors with $u \perp v$. The formulas below all pertain to $\left.M\right|_{S}$; mention of $S$ is suppressed (e.g. we write $e=e(S)$ ) except where ambiguity arises. The following formulas are results of elementary computations:

First, we write down the squares of the singular values of $\left.M\right|_{S}$ :

$$
|M e|^{2}=\frac{1}{2}\left(B-\sqrt{B^{2}-4 C}\right):=\lambda, \quad|M f|^{2}=\frac{1}{2}\left(B+\sqrt{B^{2}-4 C}\right)
$$

where $B=|M u|^{2}+|M v|^{2}, C=|M u \wedge M v|^{2}$. (Note that the formulas above are in agreement with $|M u \wedge M v|=|M e||M f|$.) We then write $e=\alpha_{0} u+\beta_{0} v$, and solve for $|M e|=\sqrt{\lambda}$ subject to $\alpha_{0}^{2}+\beta_{0}^{2}=1$. There are two solutions (a vector and its negative): either $e= \pm v$, or the solution with a positive $u$-component is given by

$$
\begin{equation*}
e=\frac{1}{Z}(\alpha u+\beta v) \tag{22}
\end{equation*}
$$

with $\alpha=|M v|^{2}-\lambda, \beta=-\langle M v, M u\rangle$ and $Z=\sqrt{\alpha^{2}+\beta^{2}}$. From this we deduce that a solution for $f$ is $f=\frac{1}{Z}(-\beta u+\alpha v)$.

## A.5. Most contracted directions: Proof of Lemma 3.1 (§3.1)

I. Proof of Lemma 3.1(a). We assume $M_{i}$ satisfies (H1) and let $S=$ $S(u, v)$. As before, mention of $S$ is suppressed. Recall that $\Delta_{i}:=\mid M^{(i)} u \wedge$ $M^{(i)} v \mid$.

Sublemma A.5.1. (i) $\Delta_{i}<(K b)^{i}$;
(ii) $\left|M^{(i)} e_{i}\right|<\left(\frac{K b}{\kappa}\right)^{i}$;
(iii) $\left|M^{(i+1)} f_{i}\right|=\left|M^{(i+1)} f_{i+1}\right| \pm \mathcal{O}\left(\left(\frac{K b}{\kappa}\right)^{i}\right) \gtrsim K_{0}^{-1} \kappa^{i}$;
(iv) If we substitute $u=e_{i}, v=f_{i}$ and $M=M^{(i+1)}$ into the formulas in Appendix A.4, part II, and let $\alpha_{i+1}, \beta_{i+1}$ and $Z_{i+1}$ be the resulting quantities, then $Z_{i+1} \approx\left|\alpha_{i+1}\right| \approx\left|M^{(i+1)} f_{i}\right|^{2}$.

Proof. (i) is the $k=0$ case of Sublemma A.4.2. For (ii), write $\left|M^{(i)} e_{i}\right|=$ $\frac{\Delta_{i}}{\mid M^{(i)} f_{i}}$; the assertion follows from (i) and our assumption on $\left|M^{(i)} f_{i}\right|$. Now make the substitution in (iv). From the formula for $|M f|$, we see that

$$
\left|M^{(i+1)} f_{i+1}\right|^{2}=B_{i+1}+\mathcal{O}\left(C_{i+1}\right)=\left|M^{(i+1)} f_{i}\right|^{2}+\left|M_{i+1} M^{(i)} e_{i}\right|^{2}+\mathcal{O}\left(C_{i+1}\right)
$$

estimates for the last two terms are given by (ii) and (i). This proves (iii). (iv) is now obvious.

We now prove Lemma 3.1(a). Continuing to substitute $u=e_{i}$ and $v=f_{i}$ in the formulas in Appendix A.4, we have, from (22),

$$
\begin{equation*}
e_{i+1}-e_{i}=\frac{1}{Z_{i+1}}\left(\frac{-\beta_{i+1}^{2}}{\alpha_{i+1}+Z_{i+1}} e_{i}+\beta_{i+1} f_{i}\right) . \tag{23}
\end{equation*}
$$

To estimate $\left|e_{i+1}-e_{i}\right|$, then, we need to obtain a suitable upper bound for $\left|\beta_{i+1}\right|$ and lower bounds for $\left|\alpha_{i+1}\right|$ and $Z_{i+1}$. Sublemma A.5.1 gives

$$
\begin{equation*}
\left|\beta_{i+1}\right| \leq\left|M^{(i+1)} e_{i}\right|\left|M^{(i+1)} f_{i}\right|<\left(\frac{K b}{\kappa}\right)^{i} \sqrt{Z_{i+1}} \tag{24}
\end{equation*}
$$

and $\left|\alpha_{i+1}\right| \approx Z_{i+1}$. These estimates together with $Z_{i+1}>K_{0}^{-2} \kappa^{2 i}$ tell us $\left|e_{i+1}-e_{i}\right| \approx \frac{\left|\beta_{i+1}\right|}{Z_{i+1}}<\left(\frac{K b}{\kappa^{2}}\right)^{i}$. The second assertion follows easily from $\left|M^{(i)} e_{n}\right| \leq$ $\left|M^{(i)}\left(e_{n}-e_{n-1}\right)\right|+\cdots+\left|M^{(i)}\left(e_{i+1}-e_{i}\right)\right|+\left|M^{(i)} e_{i}\right|<\left(\frac{K b}{\kappa^{2}}\right)^{i}$.
II. Proof of Lemma 3.1(b): First derivative estimates. For this part we assume $M_{i}$ and $S$ satisfy (H2) with $C^{2}$ norms replaced by $C^{1}$ norms. Let $\partial$ denote a fixed partial derivative.

Sublemma A.5.2. $\left|\partial e_{1}\right|,\left|\partial f_{1}\right|<K_{1}$ for some $K_{1}$.
Proof. Switching $u$ and $v$ in (H2) if necessary, we may assume $\left|M_{1} v\right| \geq$ $\left|M_{1} u\right|$. Then from Appendix A.4.II we have $Z_{1}>\alpha \geq\left|M_{1} v\right|^{2}-K b>\frac{1}{2} B-$ $K b>\frac{1}{4} K_{0}^{-2}$. Differentiating (22) gives the desired result.

Our plan of proof is as follows: For $k=1,2, \cdots$, we assume for all $i \leq k$

$$
\begin{equation*}
\left|\partial e_{i}\right|,\left|\partial f_{i}\right|<2 K_{1} \text { where } K_{1} \text { is as in Sublemma A.5.2, } \tag{*}
\end{equation*}
$$

and prove for all $i \leq k$ :
(A) $\left|\partial\left(M^{(i)} f_{i}\right)\right|<K^{i},\left|\partial\left(M^{(i)} e_{i}\right)\right|<\left(\frac{K b}{\kappa^{2}}\right)^{i}$;
(B) $\left|\partial\left(e_{i+1}-e_{i}\right)\right|, \quad\left|\partial\left(f_{i+1}-f_{i}\right)\right|<\left(\frac{K b}{\kappa^{3}}\right)^{i}$.

Observe that for $i=1,\left({ }^{*}\right)$ is given by Sublemma A.5.2. It is easy to see that (B) above implies (*) with $i=k+1$, namely $\left|\partial f_{k+1}\right| \leq\left|\partial\left(f_{k+1}-f_{k}\right)\right|+$ $\cdots+\left|\partial\left(f_{2}-f_{1}\right)\right|+\left|\partial f_{1}\right|$. From (B), we have $\left|\partial\left(f_{i+1}-f_{i}\right)\right|<\left(\frac{K b}{\kappa^{3}}\right)^{i}$, and from Sublemma A.5.2, we have $\left|\partial f_{1}\right|<K_{1}$. Hence $\left|\partial f_{k+1}\right|<\frac{K b}{\kappa^{3}}+K_{1}$, which, for $b$ sufficiently small, is $<2 K_{1}$. The computation for $e_{k+1}$ is identical.

Proof that $(*) \Longrightarrow$ (A). First we prove the estimate for $\partial\left(M^{(i)} f_{i}\right)$. Writing $\partial\left(M^{(i)} f_{i}\right)=\sum_{j=1}^{i} M_{i} \cdots\left(\partial M_{j}\right) \cdots M_{1} f_{i}+M^{(i)} \partial f_{i}$, we obtain easily $\left|\partial\left(M^{(i)} f_{i}\right)\right| \leq \sum_{j=1}^{i}\left|M_{i} \cdots\left(\partial M_{j}\right) \cdots M_{1} f_{i}\right|+\left\|M^{(i)}\right\|\left|\partial f_{i}\right| \leq i K^{i}+K^{i}\left(2 K_{1}\right)$.

This inequality is used to estimate $\partial\left(M^{(i)} e_{i}\right) \cdot{ }^{13}$ Write $\partial\left(M^{(i)} e_{i}\right)=(\mathrm{I})+$ (II) where (I) is its component in the direction of $M^{(i)} f_{i}$ and (II) is its component orthogonal to $M^{(i)} f_{i}$. Recall that $\partial\left\langle M^{(i)} e_{i}, M^{(i)} f_{i}\right\rangle=0$. We have
$|(\mathrm{I})|=\left|\left\langle\partial\left(M^{(i)} e_{i}\right), \frac{M^{(i)} f_{i}}{\left|M^{(i)} f_{i}\right|}\right\rangle\right|=\frac{1}{\left|M^{(i)} f_{i}\right|}\left|\left\langle M^{(i)} e_{i}, \partial\left(M^{(i)} f_{i}\right)\right\rangle\right|<\frac{1}{\kappa^{i}}\left(\frac{K b}{\kappa}\right)^{i} K^{i} ;$ $|(\mathrm{II})|\left|M^{(i)} f_{i}\right|=\left|\partial\left(M^{(i)} e_{i}\right) \wedge M^{(i)} f_{i}\right| \leq\left|\partial\left(M^{(i)} e_{i} \wedge M^{(i)} f_{i}\right)\right|+\left|M^{(i)} e_{i} \wedge \partial\left(M^{(i)} f_{i}\right)\right|$.
The first term in the last line is $<(K b)^{i}$ by Sublemma A.4.2, since we have established $\left|\partial e_{i}\right|,\left|\partial f_{i}\right|<2 K_{1}$; the second term is $<\left(\frac{K b}{\kappa}\right)^{i} \cdot K^{i}$. This completes the proof of (A).

To prove (B), we first compute some quantities associated with the next iterate. Substitute $u=e_{i}, v=f_{i}, M=M^{(i+1)}$ into the formulas in Appendix A.4, and let $B_{i+1}, C_{i+1}, \lambda_{i+1}$ etc. be the resulting quantities. The following is a straightforward computation.

Sublemma A.5.3. Assume (*) and (A). Then for all $i \leq k$ :
(a) $\left|\partial \lambda_{i+1}\right|<\left(\frac{K b}{\kappa^{2}}\right)^{2(i+1)}$;
(b) $\left|\partial \beta_{i+1}\right|<\left(\frac{K b}{\kappa^{2}}\right)^{i} \sqrt{Z_{i+1}}$;
(c) $\left|\partial \alpha_{i+1}\right|,\left|\partial Z_{i+1}\right|<K^{i} \sqrt{Z_{i+1}}$.

Proof that $(*),(\mathrm{A}) \Longrightarrow(\mathrm{B})$. We work with $e_{i}$; the computation for $f_{i}$ is similar. From (23) we have $\partial\left(e_{i+1}-e_{i}\right)=(\mathrm{III})+(\mathrm{IV})+(\mathrm{V})$ where

$$
\begin{aligned}
& |(\mathrm{III})|=\left|\frac{1}{Z_{i+1}}\left(e_{i+1}-e_{i}\right) \partial_{1} Z_{i+1}\right|<\frac{K^{i} \sqrt{Z_{i+1}}}{Z_{i+1}} \cdot\left(\frac{K b}{\kappa^{2}}\right)^{i}<\left(\frac{K b}{\kappa^{3}}\right)^{i} \\
& |(\mathrm{IV})|=\left|\frac{1}{Z_{i+1}} \partial\left(\beta_{i+1} f_{i}\right)\right|<\frac{1}{Z_{i+1}}\left(\left|\partial \beta_{i+1}\right|+\left|\beta_{i+1}\right|\left|\partial f_{i}\right|\right)<\left(\frac{K b}{\kappa^{3}}\right)^{i} ; \\
& |(\mathrm{V})|=\left|\frac{1}{Z_{i+1}} \partial\left(\frac{\beta_{i+1}^{2}}{\alpha_{i+1}+Z_{i+1}} e_{i}\right)\right| \ll\left(\frac{K b}{\kappa^{3}}\right)^{i} .
\end{aligned}
$$

[^13]To estimate (III), we have used Sublemmas A.5.1, A.5.3(c) and part (a) of Lemma 3.1. To estimate (IV), we have used Sublemma A.5.3(b), (*) and $\left|\beta_{i+1}\right|<\left(\frac{K b}{\kappa}\right)^{i}$. The estimate for (V) is easy.

This completes the proofs of the first derivative estimates in Lemma 3.1(b).
III. Proof of Lemma 3.1(b): Second derivative estimates. We now assume the full force of (H2). The proof proceeds in a manner entirely analogous to that for first derivatives: We first prove $\left|\partial^{2} e_{1}\right|,\left|\partial^{2} f_{1}\right|<K_{1}^{\prime}$ for some $K_{1}^{\prime}$. Then for $k=1,2, \cdots$, we assume for all $i \leq k$

$$
\left(*^{\prime}\right)\left|\partial^{2} e_{i}\right|,\left|\partial^{2} f_{i}\right|<2 K_{1}^{\prime},
$$

and prove for all $i \leq k$ :
(A') $\left|\partial^{2}\left(M^{(i)} f_{i}\right)\right|<K^{i},\left|\partial^{2}\left(M^{(i)} e_{i}\right)\right|<\left(\frac{K b}{\kappa^{3}}\right)^{i} ;$
$\left(\mathrm{B}^{\prime}\right)\left|\partial^{2}\left(e_{i+1}-e_{i}\right)\right|, \quad\left|\partial^{2}\left(f_{i+1}-f_{i}\right)\right|<\left(\frac{K b}{\kappa^{4}}\right)^{i}$.
Details are left to the reader.

## A.6. A perturbation lemma (§3.2)

Proof of Lemma 3.2. Assume inductively that $\angle\left(w_{i}, w_{i}^{\prime}\right)<\eta^{\frac{i+1}{4}}$ for all $i<n$. Let $n=2 j$ (or $2 j \pm 1$ ). Let $u_{j}=\frac{w_{j}}{\mid w_{j}}, A=D T_{z_{j}}^{j}$, and let $u_{j}^{\prime}$ and $A^{\prime}$ be the corresponding quantities for $\left(z_{0}^{\prime}, w_{0}^{\prime}\right)$. Since $\left|w_{j}\right|<K^{j}$ and $\left|w_{2 j}\right|>K_{0}^{-1} \kappa^{2 j-1}$ by hypothesis, we have

$$
\begin{equation*}
\left|A u_{j}\right|=\frac{\left|w_{2 j}\right|}{\left|w_{j}\right|}>\left(\frac{\kappa^{2}}{K}\right)^{j} \tag{25}
\end{equation*}
$$

We observe first that $\left|A^{\prime} u_{j}^{\prime}\right| \gtrsim\left(\frac{\kappa^{2}}{K}\right)^{j}$ : Clearly, $\left|A^{\prime} u_{j}^{\prime}\right| \geq\left|A u_{j}\right|-\|A\|\left|u_{j}-u_{j}^{\prime}\right|-$ $\left\|A-A^{\prime}\right\|\left|u_{j}^{\prime}\right|$. The desired estimate follows from the fact that $\|A\|\left|u_{j}-u_{j}^{\prime}\right| \approx$ $\|A\| \angle\left(u_{j}, u_{j}^{\prime}\right) \leq K^{j} \eta^{\frac{j+1}{4}},\left\|A-A^{\prime}\right\|=\left|D T_{z_{j}}^{j}-D T_{z_{j}^{\prime}}^{j}\right| \leq j K^{j} \eta^{j+1}$, and both of these quantities are $\ll\left(\frac{\kappa^{2}}{K}\right)^{j}$ by the relation imposed on $\eta$ and $\kappa$.

We estimate $\angle\left(u_{2 j}, u_{2 j}^{\prime}\right) \approx\left|u_{2 j} \wedge u_{2 j}^{\prime}\right|$ by

$$
\left|u_{2 j} \wedge u_{2 j}^{\prime}\right| \leq \frac{\left|A u_{j} \wedge A u_{j}^{\prime}\right|}{\left|A u_{j}\right| \cdot\left|A_{j}^{\prime} u_{j}^{\prime}\right|}+\frac{\left|A u_{j} \wedge\left(A-A^{\prime}\right) u_{j}^{\prime}\right|}{\left|A u_{j}\right| \cdot\left|A_{j}^{\prime} u_{j}^{\prime}\right|} .
$$

The first term is $<(K b)^{j} \eta^{\frac{j+1}{4}}\left(\frac{K}{\kappa^{2}}\right)^{2 j}$. The second term is $<K^{j}(K \eta)^{j+1}\left(\frac{K}{\kappa^{2}}\right)^{2 j}$. Both are $<\frac{1}{2} \eta^{\frac{n+1}{4}}$ by the relations we imposed on $b, \eta$ and $\kappa$ with $K_{1}$ appropriately chosen. This completes the proof of (b) for $n=2 j$.

To prove (a), we write

$$
\frac{\left|w_{2 j}^{\prime}\right|}{\left|w_{j}^{\prime}\right|} \geq \frac{\left|w_{2 j}\right|}{\left|w_{j}\right|}\left(1-\frac{\left|w_{j}\right|}{\left|w_{2 j}\right|}\left(\left\|A^{\prime}-A\right\|+\|A\|| | u_{j}^{\prime}-u_{j} \mid\right)\right) .
$$

Using the same bounds as before, we see that the factor inside parentheses is $>1-\left(K \kappa^{-2}\right)^{j}\left[(K \eta)^{j+1}+K^{j} \eta^{\frac{j+1}{4}}\right]>1-\left(\frac{1}{4}\right)^{j}$. This proves $\left|w_{2 j}^{\prime}\right| \geq$ $\left|w_{2 j}\right|\left(\sum_{1 \leq i \leq j} \frac{1}{4^{2}}\right)$.

## A.7. Temporary stable curves (§3.3)

Proof of Proposition 3.1. Let $B_{0}$ be the ball of radius $\eta$ in $S$ centered at $z_{0}$. Then on $B_{0}$ we have, by Lemma 3.2, $\left\|\left.D T\right|_{S}\right\| \geq \frac{1}{2} K_{0}^{-1}$, so that $e_{1}(S)$ is well defined. Let $\gamma_{1}$ be the integral curve to $e_{1}(S)$ defined for $s \in(-\eta, \eta)$ with $\gamma_{1}(0)=z_{0}$. Note that $\left|D T\left(e_{1}\right)\right|<K b$.

To construct $\gamma_{2}$, let $B_{1}$ be the $\frac{\eta^{2}}{2 K_{0}}$-neighborhood of $\gamma_{1}$ in $S$ where $K_{0}$ is a constant related to $\|T\|_{C^{2}}$. For $\xi \in B_{1}$, let $\xi^{\prime}$ be a point in $\gamma_{1}$ with $\left|\xi-\xi^{\prime}\right|<\frac{\eta^{2}}{2 K_{0}}$. Then $\left|T \xi-T z_{0}\right| \leq\left|T \xi-T \xi^{\prime}\right|+\left|T \xi^{\prime}-T z_{0}\right| \leq \frac{\eta^{2}}{2}+\frac{K b}{\kappa^{2}} \eta<\eta^{2}$. Thus by Lemma 3.2, $\left|\left|D T_{\xi}^{2}\right|_{S} \| \geq \frac{1}{2} K_{0}^{-1} \kappa\right.$. This ensures that $e_{2}(S)$ is defined on all of $B_{1}$. Let $\gamma_{2}$ be the integral curve to $e_{2}(S)$ with $\gamma_{2}(0)=z_{0}$. We verify that $\gamma_{2}$ is defined on $(-\eta, \eta)$ and runs alongside $\gamma_{1}$. More precisely,

$$
\begin{aligned}
\left|\frac{d}{d s}\left(\gamma_{2}(s)-\gamma_{1}(s)\right)\right| & \leq\left|e_{2}\left(\gamma_{2}(s)\right)-e_{1}\left(\gamma_{2}(s)\right)\right|+\left|e_{1}\left(\gamma_{2}(s)\right)-e_{1}\left(\gamma_{1}(s)\right)\right| \\
& \leq\left|e_{2}-e_{1}\right|+\left|\partial e_{1}\right|\left|\gamma_{2}(s)-\gamma_{1}(s)\right| \leq \frac{K b}{\kappa^{2}}+\frac{K}{\kappa^{3}}\left|\gamma_{2}(s)-\gamma_{1}(s)\right|
\end{aligned}
$$

by Lemma 3.1. By Gronwall's inequality, $\left|\gamma_{2}(s)-\gamma_{1}(s)\right| \leq \frac{K b}{\kappa^{2}}|s| e^{\frac{K}{\kappa^{3}}|s|}$, which is $\ll \frac{\eta^{2}}{2 K_{0}}$ for $|s|<\eta$. This ensures that $\gamma_{2}$ remains in $B_{1}$ and hence is well defined for all $s<\eta$.

In general, we inductively construct $\gamma_{i}$ by letting $B_{i-1}$ be the $\frac{\eta^{i}}{2 K_{0}^{i-1}}$ neighborhood of $\gamma_{i-1}$ in $S$. Then for all $\xi \in B_{i-1},\left|T^{j} \xi-T^{j} z_{0}\right|<\eta^{j+1}$ for $k<i$. Thus by Lemma 3.2, $\left\|\left.D T_{\xi}^{i}\right|_{S}\right\| \geq \frac{1}{2} K_{0}^{-1} \kappa^{i-1}$, and so $e_{i}$ is well defined. Integrating and arguing as above, we obtain $\gamma_{i}$ with $\left|\gamma_{i}(s)-\gamma_{i-1}(s)\right|<$ $K\left(\frac{K b}{\kappa^{2}}\right)^{i-1}|s| \ll \frac{\eta^{i}}{2 K_{0}^{i-1}}$ for all $s$ with $|s|<\eta$.

## A.8. A curvature estimate ( $\S 3.4$ )

Proof of Lemma 3.3. Recall that

$$
k_{i}(s)=\frac{\left|\gamma_{i}^{\prime}(s) \wedge \gamma_{i}^{\prime \prime}(s)\right|}{\left|\gamma^{\prime}(s)\right|^{3}}
$$

Since $\gamma_{i}^{\prime}=D T_{\gamma_{i-1}}\left(\gamma_{i-1}^{\prime}\right)$, we have $\gamma_{i}^{\prime \prime}=\left(\frac{d}{d s} D T_{\gamma_{i-1}}\right)\left(\gamma_{i-1}^{\prime}\right)+D T_{\gamma_{i-1}}\left(\gamma_{i-1}^{\prime \prime}\right)$. Thus $k_{i} \leq \frac{1}{\left|\gamma_{i}^{\prime}\right|^{3}}(\mathrm{I}+\mathrm{II})$ where

$$
\mathrm{I}=\left|D T\left(\gamma_{i-1}^{\prime}\right) \wedge D T\left(\gamma_{i-1}^{\prime \prime}\right)\right|, \quad \mathrm{II}=\left|D T\left(\gamma_{i-1}^{\prime}\right) \wedge\left(\frac{d}{d s} D T_{\gamma_{i-1}}\right)\left(\gamma_{i-1}^{\prime}\right)\right| .
$$

Since $D T=D T_{\gamma_{i-1}}$ has the form in (H1) in Section 3.1, we have $\mathrm{I}<K b \mid \gamma_{i-1}^{\prime} \wedge$ $\gamma_{i-1}^{\prime \prime} \mid$ (see Sublemma A.4.1). Observe that $\frac{d}{d s} D T_{\gamma_{i-1}}$ has the same form with $K_{0}$ replaced by $K_{0}\left|\gamma_{i-1}^{\prime}\right|$, i.e. if $\frac{d}{d s} D T_{\gamma_{i-1}}=M=\left(\hat{M}^{1}, \cdots, \hat{M}^{m}\right)$, then $\left\|\hat{M}^{1}\right\|<$ $K_{0}\left|\gamma_{i-1}^{\prime}\right|$ and $\left\|\hat{M}^{j}\right\|<K_{0}\left|\gamma_{i-1}^{\prime}\right| b$ for $j \geq 2$. Thus II $<K b\left|\gamma_{i-1}^{\prime}\right|^{3}$, and so

$$
k_{i} \leq\left(K b \cdot k_{i-1}+K b\right) \frac{\left|\gamma_{i-1}^{\prime}\right|^{3}}{\left|\gamma_{i}^{\prime}\right|^{3}}=\sum_{j=1}^{i-1}(K b)^{j} \frac{\left|\gamma_{i-j}^{\prime}\right|^{3}}{\left|\gamma_{i}^{\prime}\right|^{3}}+(K b)^{i} \frac{\left|\gamma_{0}^{\prime}\right|^{3}}{\left|\gamma_{i}^{\prime}\right|^{3}} \kappa_{0} \leq \frac{K b}{\kappa^{3}}
$$

## A.9. Properties of $e_{1}$ in $\mathcal{C}^{(1)}$ (§3.6)

Proof of Lemma 3.7. Consider $T_{0}=\left(\hat{T}^{1}, 0, \cdots, 0\right)$ acting at $(x, 0), x \in C_{\delta}$, and let $e_{1}=e_{1}\left(T_{0},(x, 0) ; S\left(\frac{\partial}{\partial x}, \mathbf{v}\right)\right)$ be the most contracted direction of $D T_{0}$ at $(x, 0)$ on the plane indicated. Since $\operatorname{det}\left(D T_{0}\left(e_{1}\right)\right)=0$, it is easy to see that

$$
\left\langle e_{1}, \mathbf{v}\right\rangle= \pm \frac{f^{\prime}(x)}{\sqrt{\left|D T_{0}(\mathbf{v})\right|^{2}+\left(f^{\prime}(x)\right)^{2}}}
$$

the sign depending on the orientation of $D T_{0}(\mathbf{v})$. Using the facts that $\left|f^{\prime \prime}\right|>$ $K^{-1}$ and $\left|D T_{0}(\mathbf{v})\right|>K^{-1}$, one verifies readily that $\left|\frac{d}{d x} e_{1}\left(T_{0},(x, 0) ; S\left(\frac{\partial}{\partial x}, \mathbf{v}\right)\right)\right|>$ $K^{-1}$.

Observe that if $b$ is sufficiently small, then by continuity,

$$
e_{1}\left(T, \gamma(x) ; S\left(\gamma^{\prime}(x), \mathbf{v}\right)\right)
$$

is defined everywhere on $\gamma$. We compare it to $e_{1}\left(T_{0},(x, 0) ; S\left(\frac{\partial}{\partial x}, \mathbf{v}\right)\right)$ : First, we continue to focus on $(x, 0)$ and $S=S\left(\frac{\partial}{\partial x}, \mathbf{v}\right)$, and interpolate between $T_{0}$ and $T$ by introducing $T_{s}:=\left(\hat{T}^{1}, \frac{s}{\sqrt{b}} \hat{T}^{2}, \cdots, \frac{s}{\sqrt{b}} \hat{T}^{m}\right), s \in[0, \sqrt{b}]$. More precisely, we consider the 2-parameter family $M(s, x):=\left(D T_{s}\right)_{(x, 0)}$. Observing that $M$ satisfies (H2) in Section 3.1 with $\sqrt{b}$ in the place of $b$, we obtain, by Lemma 3.1, $\left|\frac{\partial^{2}}{\partial x \partial s} e_{1}\right|<K$. From this we conclude

$$
\left|\frac{d}{d x} e_{1}\left(T,(x, 0) ; S\left(\frac{\partial}{\partial x}, \mathbf{v}\right)\right)-\frac{d}{d x} e_{1}\left(T_{0},(x, 0) ; S\left(\frac{\partial}{\partial x}, \mathbf{v}\right)\right)\right|=\mathcal{O}(\sqrt{b})
$$

Next we consider $T$ and interpolate between the $x$-axis and $\gamma$. Write $\gamma(x)=\left(x, \gamma_{y}(x)\right)$. For $s \in[0, b]$, let $z(s, x)=\left(x, \frac{s}{b} \gamma_{y}(x)\right)$, and let $M(s, x)=$ $D T_{z(s, x)}, S=S(u(s, x), \mathbf{v})$ where $u(s, x)=\left(1, \frac{s}{b} \gamma_{y}^{\prime}(x)\right)$. Another application of Lemma 3.1 gives

$$
\left|\frac{d}{d x} e_{1}\left(T, \gamma(x) ; S\left(\gamma^{\prime}(x), \mathbf{v}\right)\right)-\frac{d}{d x} e_{1}\left(T,(x, 0) ; S\left(\frac{\partial}{\partial x}, \mathbf{v}\right)\right)\right|=\mathcal{O}(b)
$$

The inequality in (3) now follows from $\left|\frac{d}{d s} e_{1}\right|>K^{-1}$ and the fact that both $\left|\frac{d}{d s} \gamma^{\prime}\right|$, which is equal to the curvature of $\gamma$, and $\left|\frac{d}{d s} S\left(\gamma^{\prime}, \mathbf{v}\right)\right|$ are $\ll 1$.

## A.10. Critical points on $C^{2}(b)$-curves (§3.7)

Proof of Corollary 3.1. Let $\gamma:[\hat{x}-\delta, \hat{x}+\delta] \rightarrow R_{1}$ be the $C^{2}(b)$-curve in question, with $\hat{x} \in C$ and $\gamma(x)=\left(x, \gamma_{y}(x)\right)$. Let $\eta=\left\langle e_{1}(S), v\right\rangle$ be as defined in Section 3.6. Since $\left|\frac{d \eta}{d x}\right|>K_{1}^{-1}$ (Lemma 3.7), there can be at most one $x \in[\hat{x}-\delta, \hat{x}+\delta]$ with $\eta(x)=0$. Observe that if we show $\eta(\hat{x})=\mathcal{O}(b)$, that will force $\eta(x)=0$ for some $x$ with $|x-\hat{x}|<K_{1}|\eta(\hat{x})|$. The claim on $\eta(\hat{x})$ follows by interpolating between $\left(T_{0},(x, 0), S\left(\frac{\partial}{\partial x}, \mathbf{v}\right)\right)$ and $\left(T, \gamma(x), S\left(\gamma^{\prime}(x), \mathbf{v}\right)\right)$ as detailed in Appendix A.9.

Proof of Lemma 3.8. We obtain by using Lemma 3.2 that for all $i<n$, $D T_{z}^{i}(\mathbf{v})>\hat{K}_{0}^{-1}$ for all $z$ with $|z-\gamma(0)|<2 b^{\frac{n}{5}}$. This guarantees that $e_{n}(S)$ with $S=S\left(\hat{\gamma}^{\prime}, \mathbf{v}\right)$ is defined at $\hat{\gamma}(s)$ for all $s \in\left[-b^{\frac{n}{5}}, b^{\frac{n}{5}}\right]$.

Let $\eta_{n}$ be defined by using $e_{n}$ instead of $e_{1}$ in the definition of $\eta$ in Section 3.6. We have $\left|\frac{d}{d s} \eta_{n}\right|=\left|\frac{d}{d s} \eta_{1}\right|+\mathcal{O}(b)>\frac{1}{2} K_{1}^{-1}$ from Lemmas 3.1 and 3.7. This shows that there is at most one point at which $\eta_{n}=0$, i.e. a critical point of order $n$. To see there exists one such point, we first interpolate between $\left(\gamma(0), S\left(\gamma^{\prime}(0), \mathbf{v}\right)\right)$ and ( $\left.\hat{\gamma}(0), S\left(\hat{\gamma}^{\prime}(0), \mathbf{v}\right)\right)$. By Lemma 3.1 and assumption (b) in this lemma,

$$
\begin{equation*}
\left|e_{n}(\hat{\gamma}(0))-e_{n}(\gamma(0))\right|<K b^{\frac{n}{4}} . \tag{26}
\end{equation*}
$$

We have

$$
\left|\eta_{n}(0)\right| \leq\left|e_{n}(\hat{\gamma}(0))-e_{n}(\gamma(0))\right|+\left|e_{n}(\gamma(0))-\gamma^{\prime}(0)\right|+\left|\gamma^{\prime}(0)-\hat{\gamma}^{\prime}(0)\right|<K b^{\frac{n}{4}}
$$

because $\left|e_{n}(\gamma(0))-\gamma^{\prime}(0)\right|=0$, and $\left|\gamma^{\prime}(0)-\hat{\gamma}^{\prime}(0)\right|<b^{\frac{n}{4}}$ from assumption (b) of this lemma. This estimate on $\eta_{n}(0)$ forces $\eta_{n}(s)=0$ for some $s$ with $|s|<K b^{\frac{n}{4}}$.

Proof of Lemma 3.9. Since $e_{n+1}$ is defined on a neighborhood of $\gamma(0)$ of radius $\gg(K b)^{n}$, and $\left|\eta_{n+1}(0)\right|<(K b)^{n}$ by Lemma 3.1(a), we proceed as in the proof of Lemma 3.8 to obtain a critical point of order $n+1$. This argument is then repeated to obtain successively critical points of order $n+2, n+3$, and so on. The distances between critical points of consecutive orders decrease geometrically. (It is not necessary to increase the order by 1 each time, but we may not be able to construct a critical point of order $n+m$ in a single step: for $m$ large, $e_{n+m}$ may not be defined in a neighborhood of order $(K b)^{n}$.)

## A.11. Splitting algorithm (§3.8)

Proof of Lemma 3.10. Consider first $I_{j}$ with the property that $I_{j} \not \supset I_{j^{\prime}}$ for any $j^{\prime}$. We observe that (i) for $i=t_{j}+1, \cdots, t_{j}+\ell_{t_{j}}-1, w_{i}^{*}$ is $b$-horizontal by Lemma 3.4, and (ii) $w_{t_{j}+\ell_{t_{j}}}^{*}$ is $b$-horizontal by assumption (a) in the lemma and the single-return argument in Section 3.8. We emphasize that the preceding discussion is entirely independent of what happens before time $t_{j}$, for assump-
tion (a) guarantees that whatever happens before, $w_{t_{j}}^{*}$ splits in a desirable manner.

Consider next $I_{j}$ with the property that all $I_{j^{\prime}} \subset I_{j}$ are of the type in the last paragraph. For definiteness, we label these inner intervals as $I_{j_{1}}, \cdots, I_{j_{k}}$ with $j_{1}<\cdots<j_{k}$. Then applying the observations in the last paragraph to each of the inner intervals and Lemma 3.4 to the times in between, we see that the only time $i$ we need to be concerned with is $i=t_{j}+\ell_{t_{j}}$. There are two cases: $t_{j_{k}}+\ell_{t_{k}}<i$, and $=i$.

If $t_{j_{k}}+\ell_{t_{k}}<i$, the $b$-horizontal property of $w_{i}^{*}$ follows from an argument identical to that of the single-return case applied to the time interval $I_{j}$; note that when making this argument, one is entirely oblivious to whether or not $\hat{w}_{t_{j}}$ is split and recombined between times $t_{j}$ and $i$.

If $t_{j_{k}}+\ell_{t_{k}}=i$, we argue first that the rejoining of $D T^{t_{j_{k}}}\left(\hat{E}_{j_{t_{k}}}\right)$ increases the slope of $D T_{z_{i-1}}\left(w_{i-1}^{*}\right)$ by at most $(K b)^{\frac{1}{2} \ell_{j_{k}}}$. Then we apply the single-return argument to $I_{j}$ (ignoring the splitting and re-combinations that occurred in between), and note that with the rejoining of $D T^{\ell_{t_{j}}} \hat{E}_{t_{j}}$, the slope deteriorates by an additional $(K b)^{\frac{1}{2} \ell_{t_{j}}}$. Since $s\left(D T_{z_{i-1}}\left(w_{i-1}^{*}\right)\right)<\frac{3 K_{0}}{2 \delta} b$, the resulting vector $w_{i}^{*}$ is still $b$-horizontal.

Inducting on the number of layers inside an $I_{j}$, we see that the only question that remains to be treated is the following: Suppose there exist $j_{1}<$ $\cdots<j_{k}$ such that $j_{1}+\ell_{j_{1}}=\cdots=j_{k}+\ell_{j_{k}}=i$. Can we be assured of the $b$-horizontal property of $w_{i}^{*}$ for arbitrary $k$ ? We answer in the affirmative, on the grounds that the deterioration in slope caused by recombining $D T^{\ell_{j_{t}}}\left(\hat{E}_{j_{t}}\right)$ is a geometric series of the form $\sum(K b)^{q}$. To see this, one must start from the rejoining of the vector that is split off last, and work backwards one step at a time in the estimation of additional deterioration in slope.

## A.12. Estimates on $B^{(k)}$ and $\mathcal{F}_{k}(\S 4.2)$

Sublemma A.12.1. For $\varepsilon, a>0$, let $J$ be an interval containing $\left[0, \frac{\varepsilon}{a}\right]$, and let $\psi: J \rightarrow \mathbb{R}$ be a $C^{2}$ function with $\left|\psi^{\prime \prime}\right| \leq a$ and $\left|\psi\left(\frac{\varepsilon}{a}\right)-\psi(0)\right| \leq \frac{1}{2} \frac{\varepsilon^{2}}{a}$. Then $\left|\psi^{\prime}(0)\right| \leq \varepsilon$.

Proof. Suppose $\left|\psi^{\prime}(0)\right|=\varepsilon^{\prime}>\varepsilon$. Then $\left|\psi\left(\frac{\varepsilon}{a}\right)-\psi(0)\right| \geq \varepsilon^{\prime} \frac{\varepsilon}{a}-\frac{1}{2} a\left(\frac{\varepsilon}{a}\right)^{2}>\frac{1}{2} \frac{\varepsilon^{2}}{a}$.
Proof of Lemma 4.2. Between $Q^{(k)}$ and $Q^{(\hat{k})}$ we have $Q^{(k)} \supset Q^{(k+1)} \supset$ $\cdots \supset Q^{(k+n)}=Q^{(\hat{k})}$. Denote $\hat{z}$ as $\hat{z}_{(k+n)}$, and for $0 \leq i<n$, choose $z_{(k+i)}$ so that
$-z_{(k+i)} \in Q^{(k+i)}$,

- $z_{(k+i)}$ has the same $x$-coordinates as $\hat{z}$, and
- $z_{(k)}$ lies on the $\mathcal{F}_{k}$-leaf containing $z$.

Let $\gamma_{k+i}$ be the $\mathcal{F}_{k+i}$-leaf containing $z_{(k+i)}$ and let $\tau_{k+i}$ be the tangent to $\gamma_{k+i}$ at $z_{(k+i)}$. We claim that $\angle\left(\tau_{k+i}, \tau_{k+i+1}\right) \leq K \delta^{-\frac{3}{2}} b^{\frac{k+i}{4}+\frac{1}{2}}$. To see this, regard $\gamma_{k+i}$ and $\gamma_{k+i+1}$ as graphs of functions defined on the $x$-axis, fix $l$ with $1 \leq l<m$ (where $m=\operatorname{dim}(X))$, and let $\psi(x)=y^{l}$ coordinate of $\gamma_{k+i+1}(x)-$ $y^{l}$ coordinate of $\gamma_{k+i}(x)$. Since the diameter of $Q^{(k+i)}$ is $<b^{\frac{k+i}{2}}$ by (A1)(ii), and the $\gamma_{j}$ are $C^{2}(b)$-curves, we wish to use Sublemma A.12.1 with $a=\frac{K b}{\delta^{3}}$ and $\frac{1}{2} \frac{\varepsilon^{2}}{a}=b^{\frac{k+i}{2}}$ to conclude that $\left|\psi^{\prime}\right| \leq \varepsilon=K \delta^{-\frac{3}{2}} b^{\frac{k+i}{4}+\frac{1}{2}}$. To do this, we need to first verify that $\frac{\varepsilon}{a} \ll$ the length of $Q^{(k+i+1)}$, i.e. $K \delta^{\frac{3}{2}} b^{\frac{k+i}{4}-\frac{1}{2}} \ll$ $\min \left\{\delta, e^{-\lambda(k+i+1)}\right\}$. This is true for $k+i>1$. The claim is also valid when $k+i=1$, for $\left|\psi^{\prime}\right|<\frac{K b}{\delta}$ by the $C^{2}(b)$-property of the curves in question.

Thus we have

$$
\angle(\tau, \hat{\tau}) \leq \angle\left(\tau, \tau_{k}\right)+\sum_{i=0}^{n-1} \angle\left(\tau_{i+k}, \tau_{i+k+1}\right)<K \delta^{-3} b \cdot|z-\hat{z}|_{h}+b^{\frac{k}{4}}
$$

the first term in the last inequality following again from the fact that $\gamma_{k}$ is $C^{2}(b)$.

Proof of Lemma 4.1. It suffices for us to prove $\left|z_{0}^{*}\left(Q^{(k)}\right)-z_{0}^{*}\left(Q^{(k+1)}\right)\right|<$ $K b^{\frac{k}{4}}$. The rest follows immediately. To prove the estimate on $z_{0}^{*}$, let $\gamma$ and $\hat{\gamma}$ be the leaves of $\mathcal{F}_{k}$ and $\mathcal{F}_{k+1}$ containing $z_{0}^{*}\left(Q^{(k)}\right)$ and $z_{0}^{*}\left(Q^{(k+1)}\right)$ respectively, parametrized so that $\gamma(0)=z_{0}^{*}\left(Q^{(k)}\right)$ and $\hat{\gamma}(0)$ has the same $x$-coordinate as $\gamma(0)$. We apply Lemma 3.8 to obtain a critical point $\hat{z}$ of order $k$ on $\hat{\gamma}$ : the bound for $|\gamma(0)-\hat{\gamma}(0)|$ comes from the diameter bound for $Q^{(k)}$ given by (A1); the one for $\left|\gamma^{\prime}(0)-\hat{\gamma}^{\prime}(0)\right|$ comes from Lemma 4.2. Lemma 3.8 tells us also that $|\hat{z}-\hat{\gamma}(0)|<K b^{\frac{k}{4}}$. Lemma 3.9 says that $\left|z_{0}^{*}\left(Q^{(k+1)}\right)-\hat{z}\right|<(K b)^{k}$.

## A.13. Correct alignment implies correct splitting (§4.4)

Proof of Lemma 4.7. Assume that $\phi\left(z_{i}\right)=z_{0}^{*}\left(Q^{(j)}\right)$, and let $\gamma$ and $\hat{\gamma}$ be the $\mathcal{F}_{j}$-leaves parametrized by arc-length through $\phi\left(z_{i}\right)$ and $z_{i}$ respectively. Both $\gamma$ and $\hat{\gamma}$ are $C^{2}(b)$-curves by (A1)(ii). Let $\ell_{i}$ be the splitting period of $z_{i}$.

From Lemma 4.3 and Lemma 3.8, there exists a critical point $\hat{z}$ of order $j$ on $\hat{\gamma}$ with $\left|\hat{z}-\phi\left(z_{i}\right)\right|<K b^{\frac{j}{4}} \ll d_{\mathcal{C}}\left(z_{i}\right)$. From Lemma 3.2, the most contracted direction of order $\ell_{i}$ in $S$ is well defined on $\hat{\gamma}$ from $\hat{z}$ to $z_{i}$ where $S=S\left(\hat{\gamma}^{\prime}, \mathbf{v}\right)$. By Lemma 3.7, $\left|\frac{d}{d s} e_{1}\right|>K_{1}^{-1}$. By Lemma 3.1, $\left|\frac{d}{d s}\left(e_{\ell_{i}}-e_{1}\right)\right|<K b$. Together we have

$$
\begin{equation*}
\angle\left(e_{\ell_{i}}(S), \hat{\gamma}^{\prime}\right)\left(z_{i}\right)>K_{1}^{-1}\left|\hat{z}-z_{i}\right|-\left|e_{j}(\hat{z})-e_{\ell_{i}}(\hat{z})\right|>\frac{3}{4} K_{1}^{-1} d_{\mathcal{C}}\left(z_{i}\right) \tag{27}
\end{equation*}
$$

For the last inequality we used $\left|e_{j}(\hat{z})-e_{\ell_{i}}(\hat{z})\right|<\min \left((K b)^{\frac{1}{4} j},(K b)^{\ell_{i}}\right) \ll d_{\mathcal{C}}\left(z_{i}\right)$.
Next we pass from $e_{\ell_{i}}(S)$ to $e_{\ell_{i}}\left(S^{*}\right)$ at $z_{i}$ where $S^{*}=S\left(w_{i}^{*}, \mathbf{v}\right)$. This is a straightforward interpolation between $u=\hat{\gamma}^{\prime}$ and $u=w_{i}^{*}$ by Lemma 3.1.

Since the difference in $u$ is $\leq \varepsilon d_{\mathcal{C}}\left(z_{i}\right)$ by the assumption of correct alignment, we obtain

$$
\begin{equation*}
\left|\angle\left(e_{\ell_{i}}(S), e_{\ell_{i}}\left(S^{*}\right)\right)\right|<K \varepsilon d_{\mathcal{C}}\left(z_{i}\right) \tag{28}
\end{equation*}
$$

where $K$ is the constant in Lemma 3.1. Finally, $\left|B_{i}\right| /\left|A_{i}\right|$ in Lemma 4.7 is $\approx\left|\angle\left(e_{\ell_{i}}\left(S^{*}\right), w_{i}^{*}\right)\right|$, and from (27) and (28), we have the following estimates at $z_{i}$ :

$$
\begin{aligned}
\left|\angle\left(e_{\ell_{i}}\left(S^{*}\right), w_{i}^{*}\right)\right| & >\left|\angle\left(e_{\ell_{i}}(S), \hat{\gamma}^{\prime}\right)\right|-\left|\angle\left(e_{\ell_{i}}(S), e_{\ell_{i}}\left(S^{*}\right)\right)\right|-\left|\angle\left(w_{i}^{*}, \hat{\gamma}^{\prime}\right)\right| \\
& >\frac{3}{4} K_{1}^{-1} d_{\mathcal{C}}\left(z_{i}\right)-K \varepsilon d_{\mathcal{C}}\left(z_{i}\right)-\varepsilon d_{\mathcal{C}}\left(z_{i}\right)>\frac{1}{2} K_{1}^{-1} d_{\mathcal{C}}\left(z_{i}\right) .
\end{aligned}
$$

## A.14. Comparison of derivatives during bound periods (§5.1)

The following sublemma is used in a number of places. Its proof is easy and left to the reader:

Sublemma A.14.1. Let $z_{0} \in \Gamma_{\theta N}$ be of generation $k$. Then for all $i \leq$ $\theta^{-1} k$, the size of the longest splitting period $z_{i}$ is in is $<K \theta \alpha i$.

Proof of Proposition 5.1. The proof proceeds by induction. Let $i<N$ be the inductive index. We assume that (6) and (7) hold for all triples $\left(z_{0}, \xi_{0}, \xi_{0}^{\prime}\right)$ in the same component of $\mathcal{C}^{(1)}$ and all $j \leq \min \left(p\left(z_{0} ; \xi_{0}, \xi_{0}^{\prime}\right), i-1\right)$. We then fix a specific triple $\left(z_{0}, \xi_{0}, \xi_{0}^{\prime}\right)$ and prove for it step $i$ of these two assertions assuming $i \leq p\left(z_{0} ; \xi_{0}, \xi_{0}^{\prime}\right)$. Note that $K_{1}$, the constant in the statement of the proposition, must not be allowed to increase from step to step. It is larger than any other generic constant $K$ that appears in the proof. In particular, $K$ does not depend on $K_{1}$.

$$
\text { Let } M_{i}=\left|w_{i}^{*}\left(\xi_{0}\right)\right|, M_{i}^{\prime}=\left|w_{i}^{*}\left(\xi_{0}^{\prime}\right)\right|, \text { and } \theta_{i}\left(\xi_{0}, \xi_{0}^{\prime}\right)=\angle\left(w_{i}^{*}\left(\xi_{0}\right), w_{i}^{*}\left(\xi_{0}^{\prime}\right)\right)
$$

Case 1. No splitting period expires at $z_{i}$ and $i-1$ is not a return time. In this case $w_{i}^{*}=D T w_{i-1}^{*}$. Writing $C=D T_{\xi_{i-1}}, C^{\prime}=D T_{\xi_{i-1}^{\prime}}, u=\frac{w_{i-1}^{*}\left(\xi_{0}\right)}{\left|w_{i-1}^{*}\left(\xi_{0}\right)\right|}$ and $u^{\prime}=\frac{\hat{w}_{i-1}^{*}\left(\xi_{0}^{\prime}\right)}{\left|\hat{w}_{i-1}^{*}\left(\xi_{0}^{\prime}\right)\right|}$, we have

$$
\theta_{i}\left(\xi_{0}, \xi_{0}^{\prime}\right) \approx \frac{\left|C u \wedge C^{\prime} u^{\prime}\right|}{|C u|\left|C^{\prime} u^{\prime}\right|} \leq K^{-1} \delta^{-2}\left(\left|C u \wedge C u^{\prime}\right|+\left|C u \wedge\left(C^{\prime}-C\right) u^{\prime}\right|\right)
$$

By Sublemma A.4.1, $\left|C u \wedge C u^{\prime}\right|<K b \theta_{i-1}$. This together with

$$
\left|C u \wedge\left(C^{\prime}-C\right) u^{\prime}\right|<K b\left|\xi_{i-1}-\xi_{i-1}^{\prime}\right|
$$

gives $\theta_{i} \leq \frac{K b}{\delta^{2}}\left(\left|\xi_{i-1}-\xi_{i-1}^{\prime}\right|+\theta_{i-1}\right)<b^{\frac{1}{2}} \Delta_{i-1}\left(\xi_{0}, \xi_{0}^{\prime}\right)$, proving (7) for step $i$.
To compare magnitudes, we have

$$
\begin{aligned}
\frac{M_{i}^{\prime}}{M_{i}}=\frac{M_{i-1}^{\prime}}{M_{i-1}} \cdot \frac{\left|C^{\prime} u^{\prime}\right|}{|C u|} & \leq \frac{M_{i-1}^{\prime}}{M_{i-1}}\left(1+\frac{\left|C^{\prime} u^{\prime}-C u\right|}{|C u|}\right) \\
& \leq \frac{M_{i-1}^{\prime}}{M_{i-1}}\left(1+\frac{\left\|C^{\prime}-C\right\|}{|C u|}+\frac{\left|C\left(u-u^{\prime}\right)\right|}{|C u|}\right) .
\end{aligned}
$$

Since $|C u|>K^{-1} d_{\mathcal{C}}\left(z_{i-1}\right),\left\|C-C^{\prime}\right\|<K\left|\xi_{i-1}-\xi_{i-1}^{\prime}\right|$ and $\left|u-u^{\prime}\right| \approx \theta_{i-1}\left(\xi_{0}, \xi_{0}^{\prime}\right)<$ $b^{\frac{1}{2}} \Delta_{i-2}\left(\xi_{0}, \xi_{0}^{\prime}\right)$,

$$
\begin{equation*}
\frac{M_{i}^{\prime}}{M_{i}} \leq \frac{M_{i-1}^{\prime}}{M_{i-1}}\left(1+K \frac{\Delta_{i-1}\left(\xi_{0}, \xi_{0}^{\prime}\right)}{d_{\mathcal{C}}\left(z_{i-1}\right)}\right) \leq \exp \left\{K_{1} \sum_{n=1}^{i-1} \frac{\Delta_{n}\left(\xi_{0}, \xi_{0}^{\prime}\right)}{d_{\mathcal{C}}\left(z_{n}\right)}\right\} \tag{29}
\end{equation*}
$$

the last inequality following from (6) for step $i-1$ and the fact that $K<K_{1}$.
Case 2. $\quad i-1$ is a return time. In this case the angle estimate is trivial since $\theta_{i}\left(\xi_{0}, \xi_{0}^{\prime}\right)=\angle\left(C \mathbf{v}, C^{\prime} \mathbf{v}\right)$. To compare magnitudes, we first recall that

$$
w_{i-1}^{*}\left(\xi_{0}\right)=A\left(\xi_{i-1}\right) \cdot e\left(\xi_{i-1}\right)+B\left(\xi_{i-1}\right) \cdot \mathbf{v}
$$

where $e=e\left(\xi_{i-1}\right)=e_{\ell\left(z_{i-1}\right)}\left(\xi_{i-1}, S\left(w_{i-1}^{*}\left(\xi_{0}\right), \mathbf{v}\right)\right) ; w_{i-1}^{*}\left(\xi_{0}^{\prime}\right)$ and $e^{\prime}=e\left(\xi_{i-1}^{\prime}\right)$ are defined similarly. From Lemma 3.1,

$$
\begin{equation*}
\left|e-e^{\prime}\right| \leq K\left(\left|\xi_{i-1}-\xi_{i-1}^{\prime}\right|+\theta_{i-1}\left(\xi_{0}, \xi_{0}^{\prime}\right)\right) \tag{30}
\end{equation*}
$$

Let $B_{0}=\frac{B\left(\xi_{i-1}\right)}{\left|w_{i-1}^{*}\left(\xi_{0}\right)\right|}$. Since $w_{i}^{*}\left(\xi_{0}\right)=B\left(\xi_{i-1}\right) \cdot C \mathbf{v}$,

$$
\begin{equation*}
\frac{M_{i}^{\prime}}{M_{i}}=\frac{M_{i-1}^{\prime}}{M_{i-1}} \cdot \frac{\left|B_{0}^{\prime}\right|}{\left|B_{0}\right|} \cdot \frac{\left|C^{\prime} \mathbf{v}\right|}{|C \mathbf{v}|} . \tag{31}
\end{equation*}
$$

To estimate $\frac{B_{0}^{\prime}}{B_{0}}$, we let $u=\frac{w_{i-1}\left(\xi_{0}\right)}{\left|w_{i-1}\left(\xi_{0}\right)\right|}$, and let $e^{\perp}$ denote the unit vector orthogonal to $e$ in $S(u, \mathbf{v})$. Then a straightforward computation (using the fact that $\left\langle\mathbf{v}, e^{\perp}\right\rangle \approx 1$ ) gives

$$
\begin{align*}
\left|B_{0}-B_{0}^{\prime}\right| & =\left|\frac{\left\langle u, e^{\perp}\right\rangle}{\left\langle\mathbf{v}, e^{\perp}\right\rangle}-\frac{\left\langle u^{\prime}, e^{\prime} \perp\right.}{\left\langle\mathbf{v}, e^{\prime} \perp\right\rangle}\right|  \tag{32}\\
& \leq 2\left(\left|u-u^{\prime}\right|+\left|e^{\perp}-e^{\prime \perp}\right|\right) \leq 2\left(\left|u-u^{\prime}\right|+\left|e-e^{\prime}\right|\right)
\end{align*}
$$

This together with $\left|B_{0}\right| \sim d_{\mathcal{C}}\left(z_{i-1}\right)$ gives

$$
\begin{equation*}
\left|\frac{B_{0}^{\prime}}{B_{0}}-1\right|<\frac{K}{\left|B_{0}\right|}\left(b^{\frac{1}{2}} \Delta_{i-2}+\left|\xi_{i-1}-\xi_{i-1}^{\prime}\right|\right)<\frac{K \Delta_{i-1}\left(\xi_{0}, \xi_{0}^{\prime}\right)}{d_{\mathcal{C}}\left(z_{i-1}\right)} \tag{33}
\end{equation*}
$$

For the last ratio,

$$
\begin{equation*}
\left|\frac{\left|C^{\prime} \mathbf{v}\right|}{|C \mathbf{v}|}-1\right| \leq K\left|\xi_{i-1}-\xi_{i-1}^{\prime}\right|<\frac{K \Delta_{i-1}\left(\xi_{0}, \xi_{0}^{\prime}\right)}{d_{\mathcal{C}}\left(z_{i-1}\right)} \tag{34}
\end{equation*}
$$

Thus (6) is proved for step $i$ by substituting (33) and (34) into (31) and taking $K_{1}>2 K$.

Case 3. At least one splitting period initiated previously expires at time $i$. Among the splitting periods expiring at this time, let $j$ be the time when the first one is initiated. Then

$$
\begin{equation*}
w_{i}^{*}\left(\xi_{0}\right)=B\left(\xi_{j}\right) \cdot D T_{\xi_{j}}^{i-j} \mathbf{v}+A\left(\xi_{j}\right) \cdot D T_{\xi_{j}}^{i-j} e\left(\xi_{j}\right) \tag{35}
\end{equation*}
$$

Let

$$
B_{0}=\frac{B\left(\xi_{j}\right)}{\left|w_{j}^{*}\left(\xi_{0}\right)\right|}, \quad A_{0}=\frac{A\left(\xi_{j}\right)}{\left|w_{j}^{*}\left(\xi_{0}\right)\right|}, \quad V=D T_{\xi_{j}}^{i-j} \mathbf{v}, \quad E=D T_{\xi_{j}}^{i-j} e\left(\xi_{j}\right)
$$

As before, all corresponding quantities for $\xi_{0}^{\prime}$ carry a prime. We shall use $\theta_{i}\left(\xi_{0}, \xi_{0}^{\prime}\right) \leq(\mathrm{I})+(\mathrm{II})$ where

$$
\text { (I) }:=\left|\frac{V^{\prime}}{\left|V^{\prime}\right|}-\frac{V}{|V|}\right| \quad \text { and } \quad \text { (II) }:=\left|\frac{A_{0}^{\prime} E^{\prime}}{B_{0}^{\prime}\left|V^{\prime}\right|}-\frac{A_{0} E}{B_{0}|V|}\right| .
$$

Assume that $z_{j}$ is bound to $\eta_{0} \in \Gamma_{\theta N}$. We apply our inductive hypotheses to the triple $\left(\eta_{0}, \xi_{j}, \xi_{j}^{\prime}\right)$ for time $i-j$. From (7), we get (I) $<b^{\frac{1}{2}} \hat{\Delta}_{i-j-1}$ where

$$
\begin{equation*}
\hat{\Delta}_{n}=\sum_{s=1}^{n} b^{\frac{s}{4}} 2^{\hat{l}_{n-s}}\left|\xi_{j+n-s}-\xi_{j+n-s}^{\prime}\right| \tag{36}
\end{equation*}
$$

and $\hat{\ell}_{n-s}$ is the longest splitting period $\eta_{n-s}$ finds itself in. Clearly we have $\ell_{j+n-s} \geq \hat{\ell}_{n-s}+n_{0}$ where $n_{0}$ is the minimum number of iterations between returns to $\mathcal{C}^{(1)}$. Set $n_{0}=2$ if there is no return to $\mathcal{C}^{(1)}$ between time $j+1$ to $i$. Then $\hat{\Delta}_{n}<2^{-n_{0}} \Delta_{j+n}$ and (I) $<\frac{1}{2} b^{\frac{1}{2}} \Delta_{i-1}$.

For (II), we first write

$$
\begin{equation*}
(\mathrm{II}) \leq \frac{\left|A_{0}\right|}{\left|B_{0}\right|} \cdot \frac{\left|E^{\prime}-E\right|}{|V|}+\left|\frac{A_{0}}{B_{0}|V|}-\frac{A_{0}^{\prime}}{B_{0}^{\prime}\left|V^{\prime}\right|}\right|\left|E^{\prime}\right| . \tag{37}
\end{equation*}
$$

From $\frac{\left|A_{0}\right|}{\left|B_{0}\right|} \sim \frac{1}{d_{C}\left(z_{j}\right)},\left|E^{\prime}-E\right| \leq(K b)^{i-j}\left(\left|\xi_{j}-\xi_{j}^{\prime}\right|+\theta_{j}\left(\xi_{0}, \xi_{0}^{\prime}\right)\right)$ (Lemma 3.1), and $|V|>1$, we obtain

For the second term on the right side of (37), we write

$$
\begin{aligned}
\left|\frac{A_{0}}{B_{0}|V|}-\frac{A_{0}^{\prime}}{B_{0}^{\prime}\left|V^{\prime}\right|}\right|\left|E^{\prime}\right| & \leq(K b)^{i-j} \frac{\left|A_{0}^{\prime}\right|}{\left|B_{0}^{\prime}\right|} \cdot \frac{1}{|V|} \cdot\left(\left|\frac{A_{0}}{A_{0}^{\prime}} \cdot \frac{B_{0}^{\prime}}{B_{0}}-1\right|+\left|1-\frac{|V|}{\left|V^{\prime}\right|}\right|\right) \\
& \leq \frac{(K b)^{i-j}}{d_{\mathcal{C}}\left(z_{j}\right)}\left(\frac{\left|A_{0}\right|}{\left|A_{0}^{\prime}\right|}\left|\frac{B_{0}^{\prime}}{B_{0}}-1\right|+\left|\frac{A_{0}}{A_{0}^{\prime}}-1\right|+\left|1-\frac{|V|}{\left|V^{\prime}\right|}\right|\right)
\end{aligned}
$$

This is estimated term by term: For the first term,

$$
\begin{equation*}
\frac{(K b)^{i-j}}{d_{\mathcal{C}}\left(z_{j}\right)} \cdot \frac{\left|A_{0}\right|}{\left|A_{0}^{\prime}\right|} \cdot\left|\frac{B_{0}^{\prime}}{B_{0}}-1\right| \leq \frac{(K b)^{i-j}}{d_{\mathcal{C}}\left(z_{j}\right)} \cdot \frac{\Delta_{j}}{d_{\mathcal{C}}\left(z_{j}\right)}<(K b)^{\frac{i-j}{3}} \Delta_{j} \ll b^{\frac{1}{2}} \Delta_{i-1} \tag{39}
\end{equation*}
$$

because $b^{\frac{i-j}{3}} \approx d_{\mathcal{C}}\left(z_{j}\right)$ by the definition of splitting period. For the second term,

$$
\begin{equation*}
\frac{(K b)^{i-j}}{d_{\mathcal{C}}\left(z_{j}\right)} \cdot\left|\frac{A_{0}-A_{0}^{\prime}}{A_{0}^{\prime}}\right| \leq \frac{(K b)^{i-j}}{d_{\mathcal{C}}\left(z_{j}\right)} K \Delta_{j} \ll b^{\frac{1}{2}} \Delta_{i-1} \tag{40}
\end{equation*}
$$

because $A_{0}^{\prime} \approx 1$ and $\left|A_{0}-A_{0}^{\prime}\right| \leq K\left(\left|\xi_{j}-\xi_{j}^{\prime}\right|+\theta_{j}\left(\xi_{0}, \xi_{0}^{\prime}\right)\right)$. Finally, for the third term, we have

$$
\begin{align*}
\frac{(K b)^{i-j}}{d_{\mathcal{C}}\left(z_{j}\right)}\left|1-\frac{\left|V^{\prime}\right|}{|V|}\right| & \leq \frac{(K b)^{i-j}}{d_{\mathcal{C}}\left(z_{j}\right)} \sum_{k=1}^{i-j-1} \frac{K_{1} \hat{\Delta}_{k}}{d_{\mathcal{C}}\left(z_{k+j}\right)}  \tag{41}\\
& <K_{1}(K b)^{\frac{2}{3}(i-j)} \sum_{k=j}^{i-1} e^{\alpha(k-j)} \Delta_{k} \ll b^{\frac{1}{2}} \Delta_{i-1} .
\end{align*}
$$

Here we have used our inductive assumption (6) for $\frac{\left|V^{\prime}\right|}{|V|}$. Putting (38)-(41) together, we conclude (II) $\ll b^{\frac{1}{2}} \Delta_{i-1}$. Hence $\theta_{i}\left(\xi_{0}, \xi_{0}^{\prime}\right)<(\mathrm{I})+(\mathrm{II})<b^{\frac{1}{2}} \Delta_{i-1}$.

To compare magnitudes, we write

$$
\frac{M_{i}^{\prime}}{M_{i}}=\frac{M_{j}^{\prime}}{M_{j}} \cdot \frac{\left|B_{0}^{\prime} V^{\prime}+A_{0}^{\prime} E^{\prime}\right|}{\left|B_{0} V+A_{0} E\right|} \leq \frac{M_{j}^{\prime}}{M_{j}} \cdot \frac{\left|V^{\prime}\right|}{|V|} \cdot \frac{\left|B_{0}^{\prime}\right|}{\left|B_{0}\right|} \cdot\left(1+\frac{\left|\frac{V^{\prime}}{\left|V^{\prime}\right|}-\frac{V}{|V|}+\frac{A_{0}^{\prime} E^{\prime}}{B_{0}^{\prime}\left|V^{\prime}\right|}-\frac{A_{0} E}{B_{0}|V|}\right|}{\left|\frac{V}{|V|}+\frac{A_{0} E}{B_{0}|V|}\right|}\right) .
$$

Since $\frac{\left|A_{0}\right|}{\left|B_{0}\right|} \sim \frac{1}{d_{\mathcal{C}}\left(z_{j}\right)}$ and $\frac{|E|}{|V|} \leq d_{\mathcal{C}}^{3}\left(z_{j}\right)$ (by the definition of the splitting period at $z_{j}$ ), it follows that $\frac{\left|A_{0} E\right|}{\left|B_{0} V\right|} \ll 1$, giving

$$
\begin{equation*}
\frac{M_{i}^{\prime}}{M_{i}} \leq \frac{M_{j}^{\prime}}{M_{j}} \cdot \frac{\left|V^{\prime}\right|}{|V|} \cdot \frac{\left|B_{0}^{\prime}\right|}{\left|B_{0}\right|} \cdot\left(1+2\left|\frac{V^{\prime}}{\left|V^{\prime}\right|}-\frac{V}{|V|}\right|+2\left|\frac{A_{0}^{\prime} E^{\prime}}{B_{0}^{\prime}\left|V^{\prime}\right|}-\frac{A_{0} E}{B_{0}|V|}\right|\right) \tag{42}
\end{equation*}
$$

We estimate the contributions from the first three ratios on the right side. Applying our inductive assumption (6) to the first ratio, we obtain an upper bound of $\exp \left\{K_{1} \sum_{n=1}^{j-1} \frac{\Delta_{n}\left(\xi_{0}, \xi^{\prime}\right)}{d_{\mathcal{C}}\left(z_{n}\right)}\right\}$. Let $z_{j}$ be bound to $\eta_{0} \in \Gamma_{\theta N}$. Applying inductive assumption (6) to ( $\eta_{0}, \xi_{j}, \xi_{j}^{\prime}$ ), we obtain

$$
\begin{equation*}
\frac{\left|V^{\prime}\right|}{|V|} \leq 1+2^{-n_{0}+1} K_{1} \sum_{k=1}^{i-j-1} \frac{\Delta_{j+k}}{d_{\mathcal{C}}\left(z_{k+j}\right)}=1+\frac{K_{1}}{2^{n_{0}-1}} \sum_{k=j+1}^{i-1} \frac{\Delta_{k}}{d_{\mathcal{C}}\left(z_{k}\right)} . \tag{43}
\end{equation*}
$$

Observe that without the factor $2^{-n_{0}+1}$ in front of $K_{1}$, there would be no room for contributions from the remaining terms (unless we allow $K_{1}$ to increase). $\frac{\left|B_{0}^{\prime}\right|}{\left|B_{0}\right|}$ is estimated as in (33), giving

$$
\begin{equation*}
\left|\frac{B_{0}^{\prime}}{B_{0}}-1\right|<\frac{K \Delta_{j}}{d_{\mathcal{C}}\left(z_{j}\right)} . \tag{44}
\end{equation*}
$$

Since $K$ is independent of $K_{1}$, this term is easily absorbed. Finally the terms inside parentheses in (42) sum up to

$$
\begin{equation*}
<1+2(\mathrm{I})+2(\mathrm{II})<1+2 b^{\frac{1}{2}} \Delta_{i-1}\left(\xi_{0}, \xi_{0}^{\prime}\right) \tag{45}
\end{equation*}
$$

Substituting (43)-(45) into (42), we complete the proof of (6) for the triple $\left(z_{0}, \xi_{0}, \xi_{0}^{\prime}\right)$ at step $i$.

Proof of Lemma 5.1. From Sublemma A.14.1, we have $\ell_{n-s}<K \alpha \theta(n-s)$, from which it follows that $\Delta_{n}<\sum_{s=0}^{n} b^{\frac{s}{4}} e^{-\frac{1}{2} \beta(n-s)}<2 e^{-\frac{1}{2} \beta n}$. We now choose $b$ small enough, and follow the last part of the proof of Sublemma A.1.1 in Appendix A. 1 to finish.

## A.15. Properties of $w_{i}^{*}$ along controlled orbits ( $\S 5.3 \mathrm{~A}, \mathrm{~B}$ )

Proof of Lemma 5.2. We may assume that $\xi_{i}$ is in a splitting period; otherwise there is nothing to prove. Let $i_{1}<i \leq i_{2}$ be the longest splitting period containing $i$. By Sublemma A. 14.1 we have $i_{2}-i_{1} \leq K \alpha \theta i$. Let $w_{i_{1}}=A e+B \mathbf{v}$ be the usual splitting. An upper bound for $\left|w_{i}^{*}\right|$ in terms of $\left|w_{i}\right|$ is then given by

$$
\left|w_{i}^{*}\right| \leq K^{i-i_{1}}|B| \leq K^{i-i_{1}}\left|w_{i_{2}}^{*}\right|=K^{i-i_{1}}\left|w_{i_{2}}\right| \leq K^{i-i_{1}}\left(K^{i_{2}-i}\left|w_{i}\right|\right) \leq K^{\varepsilon i}\left|w_{i}\right| .
$$

The first " $\leq$ " uses the fact that $\frac{\left|w_{j+1}^{*}\right|}{\left|w_{j}^{*}\right|} \leq$ some $K$, the second uses $\left|w_{i_{2}}^{*}\right|>$ $K^{-1}|B|$, and the third $\|D T\| \leq K$.

To obtain an upper bound for $\left|w_{i}\right|$ in terms of $\left|w_{i}^{*}\right|$, we let $j_{1}<\cdots<j_{n}$ be the return times between $i_{1}$ and $i$ with the property that the splitting period initiated at each $j_{k}$ extends beyond $i$. Using the nested structure of splitting periods, and from the way $w_{i}^{*}$ is defined, we have
$\left|w_{i}\right| \leq K^{i-i_{1}}\left|w_{i_{1}}\right| \leq K^{i-i_{1}}\left|\xi_{i_{1}}-\phi\left(\xi_{i_{1}}\right)\right|^{-1}\left|\xi_{j_{1}}-\phi\left(\xi_{j_{1}}\right)\right|^{-1} \cdots\left|\xi_{j_{n}}-\phi\left(\xi_{j_{n}}\right)\right|^{-1}\left|w_{i}^{*}\right|$.
From Sublemma A.14.1 and (A2), we have

$$
i-j_{k}<K \alpha \theta\left(i-j_{k-1}\right)<\cdots<(K \alpha \theta)^{k+1} i_{1}
$$

and

$$
\left|\xi_{j_{k}}-\phi\left(\xi_{j_{k}}\right)\right|>e^{-\alpha\left(i-j_{k-1}\right)}
$$

Hence $\left|w_{i}\right|<K^{i-i_{1}} e^{\alpha(1+2 K \alpha \theta) i_{1}}\left|w_{i}^{*}\right|<K^{\varepsilon i} e^{2 \alpha i}\left|w_{i}^{*}\right|$.
Proof of Lemma 5.3. Observe first that if $t$ is any return, and $\ell_{t}$ is its splitting period, then by Corollary 5.1, $w_{i}^{*}$ aligns correctly at all returns in the time interval $\left(t, t+\ell_{t}\right)$ with $2 \varepsilon_{0}$-error. This is because before the rejoining of the vector split off at time $t$, the situation is identical to that in Proposition 5.1.

To prove the lemma, we consider, in the notation of Proposition 5.2, one bound interval $\left[n_{i}, n_{i}+p_{i}\right)$ at a time. At time $n_{i}$, we have correct alignment by assumption. From the observation in the first paragraph, it suffices to consider returns at time $t \in\left(n_{i}, n_{i}+p_{i}\right)$ where $\xi_{t}$ is not in any splitting period. Write $w_{n_{i}}^{*}=B \mathbf{v}+A e_{p_{i}}$. Then $w_{t}=B \cdot D T_{\xi_{n_{i}}}^{t-n_{i}} \mathbf{v}+A \cdot D T_{\xi_{n_{i}}}^{t-n_{i}} e_{p_{i}}$. The first of these two vectors has length $>K^{-1} d_{\mathcal{C}}\left(\xi_{n_{i}}\right) e^{\lambda\left(t-n_{i}\right)}$ and aligns correctly with $<2 \varepsilon_{0}$-error at time $t$ by Corollary 5.1. Addition of the second, which has length $<(K b)^{t-n_{i}}$, changes the angle of alignment by an insignificant amount relative to $d_{\mathcal{C}}\left(\xi_{t}\right)>e^{-\alpha\left(t-n_{i}\right)}$ because $t-n_{i}$ is larger than the splitting period initiated at $\xi_{n_{i}}$ by assumption. Thus $w_{t}$ aligns correctly with $<3 \varepsilon_{0}$-error.

## A.16. Derivative growth along controlled orbits (§5.3C)

Proof of Lemma 5.4. We give a proof in the case where $j$ exists; the other case is simpler. Let $k \leq i_{1}<i_{1}+p_{1} \leq i_{2}<i_{2}+p_{2} \leq \cdots \leq i_{r}=j<n$ be defined as follows: we let $i_{1}$ be the first return to $\mathcal{C}^{(1)}$ at or after time $k, p_{1}$ the bound period of $z_{i_{1}}, i_{2}$ the first return after $i_{1}+p_{1}$, and so on until $i_{r}=j$. Writing $k=i_{0}+p_{0}$, we have that $\frac{\left|w_{n}^{*}\right|}{\left|w_{k}^{*}\right|}$ is a product of factors of the following three types:

$$
\text { I }:=\frac{\left|w_{i_{s+1}}^{*}\right|}{\left|w_{i_{s}+p_{s}}^{*}\right|}, \quad \text { II }:=\frac{\left|w_{i_{s+1}+p_{s}}^{*}\right|}{\left|w_{i_{s+1}}^{*}\right|} \quad \text { and } \quad \text { III }:=\frac{\left|w_{n}^{*}\right|}{\left|w_{j}^{*}\right|} .
$$

First we prove the lemma assuming that no splitting period initiated before time $k$ expires between times $k$ and $n$. By Lemma 3.5,

$$
\mathrm{I} \geq \frac{1}{2} c_{2} e^{\frac{1}{4} \lambda_{0}\left(i_{s+1}-\left(i_{s}+p_{s}\right)\right)}
$$

By Proposition 5.2(ii), II $\geq K^{-1} e^{\frac{\lambda}{3} p_{s+1}}$. Moreover, $K^{-1}$ can be easily absorbed into the exponential estimate for the bound period $\left[i_{s}, i_{s}+p_{s}\right]$. For III, let $\ell$ be the splitting period initiated at time $j$. If $\ell>n-j$, then III $\geq K^{-1} d_{\mathcal{C}}\left(\xi_{j}\right) e^{\lambda(n-j)}$. If not, we split $w_{j}^{*}$ into $w_{j}^{*}=A e_{n-j}+B \mathbf{v}$ where $e_{n-j}$ is the most contracted direction of order $n-j$ at $\xi_{j}$ in $S=S\left(\mathbf{v}, w_{j}^{*}\right)$. (Note that $e_{n-j}$ is well-defined.) Then III $\geq K^{-1} d_{\mathcal{C}}\left(\xi_{j}\right) e^{\lambda(n-j)}-(K b)^{n-j}$. The last term is negligible because $d_{\mathcal{C}}\left(\xi_{j}\right) \sim b^{\frac{\ell}{3}} \gg(K b)^{n-j}$. Altogether, this gives $\frac{\left|w_{n}^{*}\right|}{\left|w_{k}^{*}\right|} \geq K^{-1} d_{\mathcal{C}}\left(\xi_{j}\right) e^{\lambda^{\prime}(n-k)}$ for some $\lambda^{\prime}>0$ as claimed.

To finish this proof, we view contributions from splitting periods initiated before time $k$ as perturbations of the estimates above, and verify that they are in fact inconsequential.

Proof of Lemma 5.5. The case where $\xi_{k}$ is not in a splitting period is contained in Lemma 5.4. Let $j$ be the starting point of the largest splitting period covering $\xi_{k}$. We claim that its length $\ell$ is $<K \theta(n-j)$. If not, then we would have $\left|\xi_{j}-z_{0}\right|<b^{K \theta(n-j)}$ where $z_{0}=\phi\left(\xi_{j}\right)$, so that for all $\hat{m}$ with $j<\hat{m} \leq n$,

$$
\left|\xi_{\hat{m}}-z_{\hat{m}-j}\right|<\|D T\|^{\hat{m}-j} b^{K \theta(n-j)}<e^{-\beta(\hat{m}-j)},
$$

contradicting our assumption that $\xi_{n}$ is free. Since $n-j>p\left(\xi_{j}\right) \gg \ell\left(\xi_{j}\right)>$ $k-j$, it follows that $\ell<2 K \theta(n-k)$. By Lemma 5.4, $\left|w_{n}\right|>K^{-1} \delta e^{\lambda^{\prime}(n-j)}\left|w_{j}\right| \geq$ $K^{-1} \delta e^{\lambda^{\prime}(n-k)} K^{-\ell}\left|w_{k}\right|$.

$$
\text { A.17. }\left\|D T_{\xi_{s}}^{i-s}\right\| \text { and } w_{i}\left(\xi_{0}\right)(\S 5.3 \mathrm{C})
$$

It suffices to show for any (fixed) unit vector $u \in \mathbb{R}^{m}$ that $\left|D T_{\xi_{s}}^{i-s} u\right| \leq$ $K e^{-\hat{\lambda} s}\left|w_{i}\right|$. Since this involves only two vectors, the problem is a 2 -dimensional one.

To understand the result, recall that in 2D, we have, by simple linear algebra,

$$
\begin{equation*}
\left\|D T_{\xi_{0}}^{i}\right\|=\left\|D T_{\xi_{0}}^{s}\right\|\left\|D T_{\xi_{s}}^{i-s}\right\| \cdot \angle\left(e_{s}\left(D T_{\xi_{s}}^{-s}\right), e_{i-s}\left(D T_{\xi_{s}}^{i-s}\right)\right) \tag{46}
\end{equation*}
$$

Note that $\left\|D T_{\xi_{0}}^{i}\right\| \sim\left|w_{i}\right|$ and $\left\|D T_{\xi_{0}}^{s}\right\| \sim\left|w_{s}\right|$. This is because $\left|w_{j}\right|>e^{\lambda^{\prime \prime} j}$ for $j=1,2, \cdots$ (Lemmas 5.2 and 5.4), so that Lemma 3.1 applies, and since $w_{0}$ makes a definite angle with $e_{1}=e_{1}(D T)$, it makes a definite angle with $e_{j}=e_{j}\left(D T^{j}\right)$ for all $j$. Plugging these estimates into (46), we obtain

$$
\frac{\left|w_{i}\right|}{\left|w_{s}\right|} \sim\left\|D T_{\xi_{s}}^{i-s}\right\| \cdot \angle\left(e_{s}\left(D T_{\xi_{s}}^{-s}\right), e_{i-s}\left(D T_{\xi_{s}}^{i-s}\right)\right)
$$

The key, therefore, is to understand the angle in the displayed formula above, and to compare it to $\left|w_{s}\right|$, which is $>e^{\lambda^{\prime \prime} s}$. This angle is clearly more delicate during or around splitting periods.

Sublemma A.17.1. Let $t$ be a return time to $\mathcal{C}^{(1)}$ for $\xi_{0}$. We denote its splitting period by $\ell_{t}$, and let $I_{t}:=\left(t-5 \ell_{t}, t+\ell_{t}\right)$. Then modifying $I_{t}$ slightly to $\tilde{I}_{t}=\left(t-(5 \pm \varepsilon) \ell_{t}, t+(1 \pm \varepsilon) \ell_{t}\right)$, we may assume $\left\{\tilde{I}_{t}\right\}$ has a nested structure.

Proof. We consider $t=0,1,2, \cdots$ in this order, and determine, if $t$ is a return time, what $\tilde{I}_{t}$ will be. The right end point of $\tilde{I}_{t}$ is determined by the following algorithm: Go to $t+\ell_{t}$, and look for the largest $t^{\prime}$ inside the bound period initiated at time $t$ with the property that $t^{\prime}-5 \ell_{t^{\prime}}<t+\ell_{t}$. If no such $t^{\prime}$ exists, then $t+\ell_{t}$ is the right end point of $\tilde{I}_{t}$. If $t^{\prime}$ exists, then the new candidate end point is $t^{\prime}+\ell_{t^{\prime}}$, and the search continues. For the same reasons as in the proof of Lemma 4.6, the increments in length are exponentially small and the process terminates.

As for the left end point of $\tilde{I}_{t}$, it is possible that $t-5 \ell_{t} \in \tilde{I}_{t^{\prime}}$ for some $t^{\prime}$, the bound period initiated at which time does not extend to time $t$. This means that $\ell_{t^{\prime}} \ll \ell_{t}$, and since we assume a nested structure has been arranged for $\tilde{I}_{t^{\prime}}$ for all $t^{\prime}<t$, we simply extend the left end of $\tilde{I}_{t}$ to include the largest $\tilde{I}_{t^{\prime}}$ that it meets.

Let us assume this nested structure and write $I_{t}$ instead of $\tilde{I}_{t}$ from here on.
Sublemma A.17.2. For $s \notin \cup I_{t}$, for all $j$ with $1 \leq j<i-s,\left|w_{s+j}\right| \geq$ $b^{\frac{j}{9}}\left|w_{s}\right|$.

Proof. We fix $j$ and let $r$ be such that $\xi_{r}$ makes the deepest return between times $s$ and $s+j$. Let $j^{\prime}$ be the smallest integer $\geq j$ such that $\xi_{s+j^{\prime}}$ is outside of all splitting periods. Then, from Lemma 5.4, it follows that

$$
\begin{equation*}
\left|w_{s+j}\right| \geq K^{-K \theta\left(j^{\prime}-j\right)}\left|w_{s+j^{\prime}}\right| \geq K^{-K \theta\left(j^{\prime}-j\right)} d_{\mathcal{C}}\left(z_{r}\right)\left|w_{s}\right| \approx K^{-K \theta\left(j^{\prime}-j\right)} b^{\frac{\ell_{r}}{3}}\left|w_{s}\right| \tag{47}
\end{equation*}
$$

Case 1. $s+j \notin I_{r}$. In this case, $6 \ell_{r}<j$ since $I_{r}$ is sandwiched between $s$ and $s+j$, and $j^{\prime}-j \leq \ell_{r}$ because $r$ is the deepest return. The rightmost quantity in (47) is therefore $>K^{-\ell_{r}} b^{\frac{\ell_{r}}{3}}\left|w_{s}\right|>b^{\frac{j}{9}}\left|w_{s}\right|$.

Case 2. $s+j \in I_{r}$. The argument is as above, except we only have $5 \ell_{r}<j$.

This completes the proof of the sublemma.
Proof of Lemma 5.6. Consider the case $s \notin \cup I_{t}$, and assume for the moment that $d_{\mathcal{C}}\left(\xi_{s}\right) \geq \delta_{0}$. Then by Sublemma A.17.2 $e_{i-s}\left(\xi_{s}\right)$ is well defined, and since $w_{s}$ is $b$-horizontal, we have $\angle\left(w_{s}, e_{i-s}\left(\xi_{s}\right)\right)>K^{-1}$ by Lemma 3.6. Thus $\left\|D T_{\xi_{s}}^{i-s}\right\|\left|w_{s}\right| \leq K\left|D T^{i-s}\left(\xi_{s}\right) w_{s}\right|=K\left|w_{i}\right|$, which together with $\left|w_{s}\right|>$ $e^{\lambda^{\prime \prime} s}$ gives the desired estimate. Here $\lambda^{\prime \prime}$ is slightly less than $\lambda^{\prime}$. For $s$ with $s \notin \cup I_{t}$ and $d_{\mathcal{C}}\left(\xi_{s}\right)<\delta_{0}$, consider $\xi_{s+1}$. It remains to prove the lemma for $s \in \cup I_{t}$. Let $I_{r}$ be the maximal $I_{t}$-interval containing $s$. Observe that $6 \ell_{r}<$ $K \alpha \theta s$ (recall that $\xi_{0}$ obeys (A2)). If $i \in I_{r}$, then $\left\|D T_{\xi_{s}}^{i-s}\right\|<K^{6 \ell_{r}} \ll e^{\frac{1}{2} \lambda^{\prime \prime} i}<$ $e^{-\frac{1}{2} \lambda^{\prime \prime} s} e^{\lambda^{\prime \prime} i}<e^{-\frac{1}{2} \lambda^{\prime \prime} s}\left|w_{i}\right|$. If $i \notin I_{r}$, let $s^{\prime}=r+\ell_{r}$. Then $s^{\prime} \notin \cup I_{t}$. This case having been dealt with, we have

$$
\left\|D T_{\xi_{s}}^{i-s}\right\| \leq\left\|D T_{\xi_{s}}^{s^{\prime}-s}\right\| \cdot\left\|D T_{\xi_{s^{\prime}}}^{i-s^{\prime}}\right\| \leq K^{6 \ell_{r}} \cdot K e^{-\lambda^{\prime \prime} s^{\prime}}\left|w_{i}\right|
$$

In summary, Lemma 5.6 holds with $\hat{\lambda}=\frac{1}{2} \lambda^{\prime \prime}$.

## A.18. Quadratic turn estimates ( $£ 5.3 \mathrm{D})$

Proof of Proposition 5.3. We fix $s_{1}>0$, and let

$$
p^{*}=\min _{0<s \leq s_{1}}\left\{p\left(\xi_{0}(s), z_{0}\right), M\right\}
$$

All time indices $i$ considered are $\leq p^{*}$, and all $s$ considered are in $\left(0, s_{1}\right)$, with further restrictions indicated where necessary. For $s \in\left(0, s_{1}\right)$ and $S=S\left(\gamma^{\prime}, \mathbf{v}\right)$, $e_{p^{*}}=e_{p^{*}}(S)$ is well defined by Proposition 5.1. Let

$$
\begin{equation*}
\gamma^{\prime}(s)=A_{0}(s) e_{p^{*}}(s)+B_{0}(s) \mathbf{v} \tag{48}
\end{equation*}
$$

Then

$$
\gamma_{i}^{\prime}(s)=A_{0}(s) D T^{i} e_{p^{*}}(s)+B_{0}(s) w_{i}(s)
$$

All splitting periods below are determined by the orbit of $z_{0}$; we use them for all the $\xi_{0}(s)$ in question. Writing

$$
w_{i}(s)=w_{i}(0)+\left(w_{i}^{*}(s)-w_{i}^{*}(0)\right)+\left(E_{i}(s)-E_{i}(0)\right)
$$

where

$$
E_{i}(s)=\sum_{k \in \Lambda_{i}} A_{k}(s) D T^{i-k} e_{\ell_{k}}
$$

and $\Lambda_{i}$ is the collection of $k>0$ such that the splitting period begun at time $k$ extends beyond time $i$, we arrive at the formula

$$
\begin{equation*}
\xi_{i}(s)-z_{i}=\int_{0}^{s} \gamma^{\prime}(u) d u=w_{i}(0) \int_{0}^{s} B_{0}(u) d u+\mathrm{I}+\mathrm{II}+\mathrm{III} \tag{49}
\end{equation*}
$$

where

$$
\begin{gathered}
\mathrm{I}=\int_{0}^{s} A_{0}(u) D T^{i}(u) e_{p^{*}}(u) d u, \quad \mathrm{II}=\int_{0}^{s} B_{0}(u)\left(w_{i}^{*}(u)-w_{i}^{*}(0)\right) d u \\
\mathrm{III}=\int_{0}^{s} B_{0}(u)\left(E_{i}(u)-E_{i}(0)\right) d u
\end{gathered}
$$

Plan of proof. We will prove that for $i$ and $s$ satisfying $i \in\left[\ell(s), p^{*}\right]$, the first term on the right side of (49) dominates, so that assuming $s_{1}$ is sufficiently small,

$$
\xi_{i}(s)-z_{i} \approx w_{i}(0) \int_{0}^{s} B_{0}(u) d u \approx \frac{1}{2} B_{0}^{\prime}(0) s^{2} w_{i}(0)
$$

The following estimate then completes the proof: Differentiating (48), we obtain $\gamma^{\prime \prime}=A_{0}^{\prime} e_{p^{*}}+A_{0} \frac{d}{d s} e_{p^{*}}+B_{0}^{\prime} \mathbf{v}$. On the left side, $\left|\gamma^{\prime \prime}\right|=\mathcal{O}(b)$ since $\gamma$ is $C^{2}(b)$. On the right side, $\left|A_{0} \frac{d}{d s} e_{p^{*}}\right| \approx\left|\frac{d}{d s} e_{1}\right|>K^{-1}$ (Lemma 3.7) and $\left\langle e_{p^{*}}, \frac{d}{d s} e_{p^{*}}\right\rangle=0$. It follows therefore that $B_{0}^{\prime} \mathbf{v} \approx \frac{d}{d s} e_{1}$.

We divide the main argument of the proof into the following two steps:
Step I. Estimates on $|\mathrm{I}|,|\mathrm{II}|$ and $|\mathrm{III}|$. Beginning with $|\mathrm{I}|$, we have $A_{0}(s) \approx 1$, so that $|\mathrm{I}| \leq(K b)^{i} s \ll s^{2}$ provided $b^{i}<b^{\ell(s)}:=s^{2}$. This is where the lower bound on $i$ is used for each $s$.

By assumption, $z_{0}$ is a critical point of order $M$. If $p^{*}=M$, then $B(0)=0$. For $p^{*}<M, B(0)$ may not be zero but we have $|B(0)|<(K b)^{p^{*}}$. Since this error is negligible, we will write $B(0)=0$ in the computation that follows. For |II|, then,

$$
\begin{aligned}
|\mathrm{II}| & \leq K\left|w_{i}^{*}(0)\right| \int_{0}^{s} u\left|\frac{\left|w_{i}^{*}(u)\right|}{\left|w_{i}^{*}(0)\right|}-1\right| d u \leq K\left|w_{i}^{*}(0)\right| \int_{0}^{s} u\left(\sum_{j<i} \frac{\Delta_{j}}{d_{\mathcal{C}}\left(z_{j}\right)}\right) d u \\
& \leq K\left|w_{i}^{*}(0)\right| \int_{0}^{s} u\left(\sum_{j<i}\left(2^{K \alpha \theta j} d_{\mathcal{C}}\left(z_{j}\right)\right)^{-1}\right) \sup _{j<i}\left|z_{j}-\xi_{j}(u)\right| d u \\
& \leq K e^{2 \alpha i}\left|w_{i}^{*}(0)\right| \int_{0}^{s} u \sup _{j<i}\left|z_{j}-\xi_{j}(u)\right| d u .
\end{aligned}
$$

Here we have used Proposition 5.1 for the second inequality, and assumption (2) in Section 5.3D and Sublemma A.14.1 for the next two.

To estimate |III|, we have, for each $k \in \Lambda_{i}$,

$$
\begin{aligned}
& \left|A_{k} D T_{\xi_{k}}^{i-k} e\left(\xi_{k}\right)-A_{k}(0) D T_{z_{k}}^{i-k} e\left(z_{k}\right)\right| \\
\leq & (K b)^{i-k}\left|A_{k}-A_{k}(0)\right|+\left|A_{k}(0)\right|\left|D T_{\xi_{k}}^{i-k} e\left(\xi_{k}\right)-D T_{z_{k}}^{i-k} e\left(z_{k}\right)\right|
\end{aligned}
$$

We claim that the first term can be estimated by

$$
\left|A_{k}(u)-A_{k}(0)\right|<K\left|w_{k}^{*}(0)\right| e^{4 \alpha k} \sup _{j<k}\left|z_{j}-\xi_{j}\right|<K\left|w_{i}^{*}(0)\right| e^{5 \alpha i} \sup _{j<i}\left|z_{j}-\xi_{j}\right| .
$$

For the first inequality, we use

$$
\left|\frac{A_{k}(u)}{A_{k}(0)}-1\right| \approx\left|\frac{\left|w_{k}^{*}(s)\right|}{\left|w_{k}^{*}(0)\right|}\left(1+\mathcal{O}\left(\left|z_{k}-\xi_{k}\right|+\theta_{k}\left(\xi_{0}(u), z_{0}\right)\right)\right)-1\right|,
$$

and $\left|A_{k}(0)\right| \leq K\left|w_{k}^{*}(0)\right| e^{\alpha k}$ because $w_{k}^{*}(0)$ aligns correctly at time $k$. For the second inequality, we use $\left|w_{k}^{*}(0)\right| \leq K e^{\alpha k}\left|w_{i}^{*}(0)\right|$ by virtue of Lemma 5.4 and assumption (2) in Section 5.3D. Summing over all $k \in \Lambda_{i}$ is not problematic because of the factor $(K b)^{i-k}$ in front. For the second term we use

$$
\left|D T_{\xi_{k}}^{i-k} e\left(\xi_{k}\right)-D T_{z_{k}}^{i-k} e\left(z_{k}\right)\right| \leq(K b)^{i-k}\left(\left|\xi_{k}-z_{k}\right|+\theta_{k}\left(\xi_{0}(u), z_{0}\right)\right)
$$

This inequality is derived from Lemma 3.1. Altogether, we have proved

$$
\begin{equation*}
|\mathrm{II}|,|\mathrm{III}|<K\left|w_{i}^{*}(0)\right| e^{5 \alpha i} \int_{0}^{s} u \sup _{j<i}\left|z_{j}-\xi_{j}(u)\right| d u . \tag{50}
\end{equation*}
$$

Step II. Proof of formula for $\left|\xi_{i}(s)-z_{i}\right|$. Fix $i_{0}$ so that $e^{9 \alpha i_{0}} e^{-\beta i_{0}} \ll 1$. We define

$$
U_{i}:=K e^{5 \alpha i} \sup _{j \leq i}\left|w_{j}^{*}(0)\right|
$$

where $K$ is the constant in the bound for $|\mathrm{II}|$ and $|\mathrm{III}|$ above. By choosing $\delta$ sufficiently small, we may assume $U_{i_{0}} s^{2} \leq U_{i_{0}} \delta^{2} \ll 1$. We now prove inductively (and in tandem) the following two statements:
(i) $K e^{2 \alpha i} U_{i} s^{2} \ll 1$,
(ii) $\left|\xi_{i}(s)-z_{i}\right| \approx \frac{1}{2}\left|\frac{d}{d s} e_{1}(0)\right|\left|w_{i}(0)\right| s^{2}$, or, equivalently, $|\mathrm{II}|,|\mathrm{III}| \ll\left|w_{i}(0)\right| s^{2}$. The first $n_{0}$ steps, where the entire action takes place away from $\mathcal{C}^{(1)}$, are trivial. We assume now that (i) and (ii) have been proved for all $j<i$, and prove (i) for step $i$. Using Lemmas 5.2 and 5.4, one has

$$
\sup _{j \leq i}\left|w_{j}^{*}(0)\right|<e^{\alpha i}\left|w_{i}^{*}(0)\right| \leq e^{2 \alpha i}\left|w_{i}(0)\right| \leq K e^{2 \alpha i}\left|w_{i-1}(0)\right| .
$$

This combined with (ii) for step $i-1$ gives

$$
U_{i} s^{2} \leq K e^{5 \alpha i}\left(K e^{2 \alpha i}\left|w_{i-1}(0)\right|\right) s^{2} \approx K^{2} e^{7 \alpha i}\left|\frac{d}{d s} e_{1}(0)\right|^{-1}\left|\xi_{i-1}(s)-z_{i-1}\right| .
$$

Thus

$$
\begin{equation*}
K e^{2 \alpha i} U_{i} s^{2} \leq K^{2} e^{9 \alpha i}\left|\frac{d}{d s} e_{1}(0)\right|^{-1} e^{-\beta(i-1)} \ll 1 \tag{51}
\end{equation*}
$$

proving (i). To prove (ii), first observe that from Lemmas 5.2 and 5.4 we have

$$
\begin{equation*}
\left|w_{j}(0)\right|<e^{2 \alpha i}\left|w_{i}(0)\right| \tag{52}
\end{equation*}
$$

for all $j<i$. We claim that

$$
\begin{aligned}
|\mathrm{II}|+|\mathrm{III}| & \leq 2 U_{i} \int_{0}^{s} u \sup _{j<i}\left|z_{j}-\xi_{j}(u)\right| d u \leq K U_{i} \sup _{j<i}\left|w_{j}(0)\right| s^{4} \\
& \leq\left(K e^{2 \alpha i} U_{i} s^{2}\right)\left|w_{i}(0)\right| s^{2} \ll\left|w_{i}(0)\right| s^{2} .
\end{aligned}
$$

The first inequality above is the conclusion of Step I, the second is obtained by using (ii) for $j<i$, the third is by (52), and the last is step $i$ of (i). This completes the proof of Step II.

Step III. Monotonicity of $s \mapsto p(s)$. To prove that the distance formula (Step II(ii)) holds for all $i \in\left[\ell(s), p\left(s_{1}\right)\right]$ and $0<s<s_{1}$, it remains to show that $p(s)$ is monotone in $s$ so that $p^{*}$ introduced at the beginning of this proof is equal to $\min \left\{p\left(\xi_{0}\left(s_{1}\right), z_{0}\right), M\right\}$. To do this, we check that for $s$ and $i$ with $i>\ell(s)$,

$$
\left|\frac{d}{d s} I\right|,\left|\frac{d}{d s} I I\right|,\left|\frac{d}{d s} I I I\right| \ll\left|B_{0}(s)\right|\left|w_{i}(0)\right| .
$$

Thus $\frac{d}{d s}\left(\xi_{i}(s)-z_{i}\right) \approx B_{0}(s) w_{i}(0) \approx B_{0}^{\prime}(0) s w_{i}(0)$; i.e., $\left|\xi_{i}(s)-z_{i}\right|$ increases monotonically with $s$. It follows by definition that $p(s)$ is monotone in $s$.

## A.19. Sectional diameter of $Q^{(k)}(\S 6.3)$

Proof of Lemma 6.2. Let $\xi_{k} \in Q^{(k)}$ be fixed. We argue as before that $d_{\mathcal{C}}\left(\xi_{i}\right)>2 b^{\frac{i}{5}}$ for $i=1, \cdots, k-1$. Let $S$ be a 2D subspace through $\xi_{1}$ containing $\xi_{1}$ and $\xi_{1}+\tau_{1}$. All constructions are in $S$ until the very end of the proof. There is clearly a stable curve $\gamma_{1}$ of order one in $S$ passing through $\xi_{1}$. Since $d_{\mathcal{C}}\left(\xi_{1}\right)>$ $2 b^{\frac{1}{5}}, \gamma_{1}$ makes an angle $\gtrsim b^{\frac{1}{5}}$ with the $x$-axis by Lemma 3.7; thus it connects the two components of $\partial\left(R_{1} \cap S\right)$. We wish to borrow the argument in Appendix A. 7 to construct inductively stable curves $\gamma_{i}$ of order $i, i=2,3, \cdots, k$, through $\xi_{1}$, but are prevented from doing so due to the following technical problem: with $\kappa=b^{\frac{1}{5}}$, Lemma 3.2 (a general perturbative result) does not apply. We seek instead to use Lemma 5.4, which relies on the control of $\left(\xi_{1}, \tau_{1}\right)$ for $k$ iterates, to estimate the growth of $\tau_{i}$. Details of the argument are as follows:

Assume that $\gamma_{i}=\gamma_{i}(S)$ with the following properties has been constructed: (i) $\gamma_{i}$ is tangent to $e_{i}$, passes through $\xi_{1}$ and connects the two components of $\partial R_{1} \cap S$; and (ii) for $\xi \in \gamma_{i}, d_{\mathcal{C}}\left(T^{j} \xi\right)>\frac{3}{2} b^{\frac{j}{5}}$ for $1 \leq j \leq i$.

To construct $\gamma_{i+1}$ we let $U_{i}$ be the $b^{\frac{i+1}{4}}$-neighborhood of $\gamma_{i}$. Then the following hold for all $\xi \in U_{i}$ : First, $d_{\mathcal{C}}\left(T^{j} \xi\right)>b^{\frac{j}{5}}$, so if $\tau_{1}$ is tangent to $\mathcal{F}_{1}$ at $\xi$, then $\left(\xi, \tau_{1}\right)$ is provisionally controlled by $\Gamma_{k}$ for $i$ iterates.

Claim. $\left|\tau_{j}\right|>(K b)^{\frac{j}{5}}$ for $j \leq i$.
Proof. If $T^{j} \xi$ is out of all splitting periods, then $\left|\tau_{j}\right|>(K b)^{\frac{j}{5}}$ by Lemma 5.4. If not, let $j_{1}<j$ be the time at which the longest splitting period extending
beyond $j$ is initiated. Since $d_{\mathcal{C}}\left(T^{j_{1}} \xi\right)>b^{\frac{j_{1}}{5}}$, it follows that $l_{j_{1}}$, the splitting period initiated at $j_{1}$, is $<\frac{3 j_{1}}{5}$. Thus $\left|\tau_{j}\right|>(\|D T\|)^{-\frac{3}{5} j_{1}}\left|\tau_{j_{1}+\ell_{j_{1}}}\right|>(\|D T\|)^{-\frac{3}{5} j_{1}} b^{\frac{j}{5}}$; the second inequality is obtained by applying Lemma 5.4 to $\tau_{j_{1}+l_{j_{1}}}$.

By Lemma 3.1, $e_{i+1}(\xi)$ is well-defined, $\left|D T_{\xi}^{j}\left(e_{i+1}\right)\right|<\left(K b^{\frac{3}{5}}\right)^{j}$ for all $j \leq$ $i+1,\left|e_{i+1}-e_{i}\right|<\left(K b^{\frac{3}{5}}\right)^{i}$, and $\left|\frac{d}{d s}\left(e_{i+1}-e_{i}\right)\right|<\left(K b^{\frac{2}{5}}\right)^{i}$. Let $\gamma_{i+1}$ be the integral curve of $e_{i+1}$ through $\xi_{1}$. We verify following the computation in Appendix A. 7 that $\left|\gamma_{i+1}-\gamma_{i}\right|<K b^{\frac{3}{5} i}$, so $\gamma_{i+1}$ stays in $U_{i}$ until it meets $\partial R_{1} \cap S$. Properties (i) and (ii) are again valid for $\gamma_{i+1}$.

To finish, we let $W=T^{k-1} W_{1}$ where $W_{1}=\cup_{S} \gamma_{k}(S)$, the union being taken over all 2D planes $S$ containing $\xi_{1}$ and $\xi_{1}+\tau_{1}$.

## A.20. Geometry of monotone branches (§7.3)

The proof of Lemma 7.1 uses material in Sections 7.3, 8.1 and 8.2.
Proof of Lemma 7.1. Let $T \in \mathcal{G}_{N}$. For $k \leq \theta N$, let $\hat{R}_{1, k}=\left\{\xi_{1} \in R_{1}\right.$ : $\left.\xi_{k} \in \cup_{M \in \mathcal{T}_{k}} M^{\circ}\right\}$. Then for $\xi_{1} \in \hat{R}_{1, k}$ and time indices $\leq k$, bound periods $p\left(\xi_{i}\right)$ for $\xi_{i} \in \mathcal{C}^{(1)}$ are well defined and $\left\{p\left(\xi_{i}\right)\right\}$ has a nested structure; i.e., $i+p\left(\xi_{i}\right) \geq j+p\left(\xi_{j}\right)$ for $\xi_{i}, \xi_{j} \in \mathcal{C}^{(1)}$ satisfying $i<j<i+p\left(\xi_{i}\right)$. We introduce a function $\mathfrak{b}_{k}\left(\xi_{1}\right)$ on $\hat{R}_{1, k}$ as follows:

- if $\xi_{k}$ is free, then $\mathfrak{b}_{k}\left(\xi_{1}\right)=0 ;$
- if $\xi_{k}$ is bound to some point and the bound period lasts beyond time $\theta N$, then having no knowledge of events beyond time $\theta N$, we set $\mathfrak{b}_{k}\left(\xi_{1}\right)=\infty$;
- if $\left(j, j+p\left(\xi_{j}\right)\right)$ is the longest bound period $\xi_{k}$ finds itself in, and $j+p\left(\xi_{j}\right) \leq$ $\theta N$, then we set $\mathfrak{b}_{k}\left(\xi_{1}\right)=j+p\left(\xi_{j}\right)-k$.

That is to say, $\mathfrak{b}_{k}\left(\xi_{1}\right)$ gives the number of iterates it takes for $\xi_{k}$ to become free - without knowledge of events after time $\theta N$. We observe immediately that due to the nested structure of $\left\{p\left(\xi_{i}\right)\right\}$, if $\mathfrak{b}_{k-1}\left(\xi_{1}\right)=i, 0<i \leq \infty$, then $\mathfrak{b}_{k}\left(\xi_{1}\right)=i-1$.

Let $l$ be an arbitrary $\mathcal{F}_{1}$-leaf parametrized by $s$. Then $\mathfrak{b}_{k}$ is defined on $l_{k}:=l \cap \hat{R}_{1, k}$, and the $T^{k-1}$-images of the connected components of $l_{k}$ are exactly the maximal $\mathcal{F}_{k}$-segments in $M^{\circ}$ for $M \in \mathcal{T}_{k}$. We say $\mathfrak{b}_{k}$ restricted to $\omega=l\left(s_{1}, s_{2}\right) \subset l_{k}$ is a $U$-shaped function if there exists $s^{*} \in\left(s_{1}, s_{2}\right)$ such that $\mathfrak{b}_{k}$ is non-increasing on $l\left(s_{1}, s^{*}\right]$ and nondecreasing on $l\left[s^{*}, s_{2}\right)$. Lemma 7.1 is reduced to the following. We claim that on all connected components of $l_{k}, \mathfrak{b}_{k}$ is a $U$-shaped function, and leave the proof to the reader as an exercise.

Proof of Corollary 7.1. Corollary 7.1 follows immediately from the arguments above. The numbers $K_{1}$ and $K_{2}$ are determined from $f_{0}$ as follows: Let $\hat{x}_{1}<\hat{x}_{2}<\cdots<\hat{x}_{r}=\hat{x}_{0}$ be the critical points of $f_{0}$, and let $I_{i}=\left(\hat{x}_{i-1}, \hat{x}_{i}\right)$.

Then

$$
K_{1}=\max _{1 \leq i \leq r} N_{i} \quad \text { and } \quad K_{2}=\max _{1 \leq i \leq r} \sum_{1 \leq j \leq r} L_{i j}
$$

where $N_{i}$ is the number of $I_{j}$-intervals counted with multiplicity $f_{0}\left(I_{i}\right)$ intersects at least partially, and $L_{i j}$ is the cardinality of $f_{0}\left(I_{j}\right) \cap\left\{\hat{x}_{i}\right\}$.

## A.21. Branch replacement ( $\S 8.3$ )

Proof of Lemma 8.2. Let $\frac{1}{3} H$ be the middle third of $H$. We will show $T^{-i} M^{\circ}$ is inside $\frac{1}{3} H$ so that $T^{-i} M \subset H$. To prove $T^{-i} M^{\circ} \subset \frac{1}{3} H$, it suffices to show that if we start from $B^{(k-i)} \subset H$ and move right along any $\mathcal{F}_{k-i}$-segment $\gamma$, we will get out of $T^{-i} M^{\circ}$ before we reach the end of $\frac{1}{3} H$. Suppose, to derive a contradiction, that this is not true for some $\gamma$. Then every point in $\gamma$, which we may assume runs from $B^{(k-i)}$ to the right end of $\frac{1}{3} H$, is controlled for the next $i$ iterates. We will show if this is the case then there exists $j<K \alpha(k-i)$ such that $T^{j} \gamma$ crosses some $Q^{(1)}$. It follows that $T^{j} \gamma$ crosses some $B^{(k-i+j)}$.

Let $\gamma_{0}$ be a segment of $\gamma$ such that $\pi_{x}\left(\gamma_{0}\right)=I_{\mu j}(\hat{x})$ for some $I_{\mu j}(\hat{x})$ with $\mu \sim 2 \alpha(k-i)$ (see $\S 2.2$ for a formal definition of $I_{\mu j}$ ), and let $\gamma_{i}=T^{i} \gamma_{0}$. We follow the argument in Section 9.2 to conclude that $\gamma_{i}$ obeys the rules ( $\mathrm{P} 1^{\prime}$ ) and ( $\mathrm{P} 2^{\prime}$ ) in Section 9.2C (leaving details as an exercise to the reader). That $T^{j} \gamma$ crosses $Q^{(1)}$ for some $j<K \alpha(k-i)$ is then proved by repeating the proof of Lemma 2.4 using ( $\mathrm{P} 1^{\prime}$ ) and ( $\mathrm{P} 2^{\prime}$ ).

Proof of Lemma 8.3. Let $n_{1}$ be the smallest $\ell$ for which $T^{-\ell} S_{\ell} \subset H$. First, we observe that if $P_{1}$ is well defined and $E_{1}^{\prime}$ remains active for at least $n_{1}$ generations, then all the offspring of $P_{1}$ survive; i.e., they are not discontinued, for at least $n_{1}$ generations. This is because the new ends created as the offspring of $P_{1}$ reproduce are younger than the end originating from $\hat{B}^{(k-i)}$, and so will last longer than it. It follows that $P_{j}, j \leq n_{1}$, are well defined, and $T^{n_{1}-1} P_{1}$ is a union of branches in $\mathcal{T}_{k-i+n_{1}}$ with adjacent ones overlapping in critical blobs.

We prove next that $P_{n_{1}}$ is subordinate to $S_{n_{1}}$. Observe that since $n_{1}$ is the smallest $\ell$ with the property that $T^{-\ell} S_{\ell} \subset H$, it follows that $S_{n_{1}-1}$ contains a $B^{\left(k-i+n_{1}-1\right)}$ and $S_{n_{1}}$ is the image of the subset of $S_{n_{1}-1}$ between $T^{n_{1}-1} \hat{B}^{(k-i)}$ and this $B^{\left(k-i+n_{1}-1\right)}$. We claim that $T^{n_{1}-2} P_{1}$ meets the $Q^{\left(k-i+n_{1}-1\right)}$ contain$\operatorname{ing} B^{\left(k-i+n_{1}-1\right)}$, so that $T^{n_{1}-2} P_{1} \cap B^{\left(k-i+n_{1}-1\right)}$ contains a $B^{\left(p+n_{1}-1\right)}$. To prove this, let $z \in B^{\left(k-i+n_{1}-1\right)}$. By Lemma 6.2, we know there exists a stable manifold $W=W_{k-i+n_{1}-1}^{s}$ whose $T^{k-i+n_{1}-1}$-image contains $z$ and along which $T$ contracts at a rate $\sim b^{\frac{1}{2}}$. Since $T^{k-i-1} W$ meets every fiber in $H$ and has diameter $<b^{\frac{k-i-1}{2}}$, the desired result follows from the relation between $T^{-n_{1}} S_{n_{1}}$, $H$ and $P$.

We explain why $T^{n_{1}-1} P_{1}$ is a monotone branch: From the observation in the first paragraph of this proof, we see that it suffices to show there are no critical blobs between $T^{n_{1}} \hat{B}^{(p)}$ and $T B^{\left(p+n_{1}-1\right)}$ : If there was one, look at when
and where it was created, and argue that the corresponding iterate of $T^{-n_{1}} S_{n_{1}}$ also crosses some $Q^{(j)}$ containing it, leading to the existence of a monotone branch of generation $k-i+n_{1}$ contained properly in $S_{n_{1}}$, which is absurd.

For $n_{1}<j \leq i$, the arguments are as above, namely that if $S_{j}$ subdivides, then so does $P_{j}$ in corresponding locations; and that no other subdivisions of $P_{j}$ are possible.

## A.22. Dynamics on unstable manifolds (§9.2)

Proof of Lemma 9.1. As before, it suffices to prove correct alignment at free returns. Inducting on $k$, we let $\xi_{k}$ be a free return, and let $\tau^{*}$ denote the tangent to $\mathcal{F}_{j}$ at $\xi_{k}$ where $\phi\left(\xi_{k}\right)$ is of generation $j$. We need to show $\angle\left(\tau_{k}, \tau^{*}\right)<$ $\varepsilon_{1} d_{\mathcal{C}}\left(\xi_{k}\right)$, and our plan is to deduce that from the control of foliations proved earlier.

Let $n \geq k$ be a sufficiently large number to be determined. We let $\xi_{k}=\zeta_{n}$, so that $\tau_{k}$ is a multiple of $D T_{\zeta_{0}}^{n} \tau, \tau \in X_{\zeta_{0}}$ being a unit vector tangent to $T^{k-n} l_{0}$. Let $\hat{\tau}_{1} \in X_{\zeta_{0}}$ be a unit vector tangent to $\mathcal{F}_{1}$. By the bound on $|\operatorname{det}(D T)|$, we have

$$
\angle\left(D T_{\zeta_{0}}^{n} \tau, D T_{\zeta_{0}}^{n} \hat{\tau}_{1}\right) \leq(K b)^{n} \frac{1}{\left|D T_{\zeta_{0}}^{n} \tau\right|} \frac{1}{\left|D T_{\zeta_{0}}^{n} \hat{\tau}_{1}\right|}
$$

Observe that $\left|D T_{\zeta_{0}}^{n} \tau\right|>K^{-1} e^{\lambda^{\prime \prime} n}$ : for the first $n-k$ iterates, Lemma 3.5 applies since $\zeta_{i}$ is essentially outside of $\mathcal{C}^{(1)}$; for the next $k$ iterates, use Lemma 5.4 and the fact that $\xi_{k}$ is a free return. The constraints on $n$ are as follows: First, $\zeta_{n}$ must be in a monotone branch in $\mathcal{T}_{n+1}$, so that $\left(\zeta_{0}, \hat{\tau}_{1}\right)$ is controlled through this time, giving $\left|D T_{\zeta_{0}}^{n} \hat{\tau}_{1}\right|>K^{-1} e^{\lambda^{\prime \prime} n}$. This is not a problem since $\xi_{0} \in \Omega$. Second, we assume $n \geq j$, so that our choice of $\phi\left(\zeta_{n}\right)$ in the control of foliations is compatible with the definition of $\phi\left(\xi_{k}\right)$. ( $\left.\mathrm{A} 6^{\prime}\right)$ then guarantees $\angle\left(D T_{\zeta_{0}}^{n} \hat{\tau}_{1}, \tau^{*}\right)<\varepsilon_{1} d_{C}\left(\zeta_{n}\right)$, and the desired conclusion follows if $n$ is large enough that $\angle\left(D T_{\zeta_{0}}^{n} \tau, D T_{\zeta_{0}}^{n} \hat{\tau}_{1}\right)$ is negligible.

Proof of Proposition 9.2. ( $\mathrm{P} 1^{\prime}$ ) is an easy exercise. ( $\mathrm{P} 2^{\prime}$ ) (iii) and ( $\mathrm{P} 3^{\prime}$ ) require the following extensions of Proposition 5.1.

Sublemma A.22.1. The setting is as in Proposition 5.3. Let $\xi_{0}, \xi_{0}^{\prime} \in \gamma$ be such that $\left|\xi_{0}-\xi_{0}^{\prime}\right|<\frac{1}{10} d_{\mathcal{C}}\left(\xi_{0}\right)$, and let $\ell=\ell\left(\xi_{0}\right), p=p\left(\xi_{0}\right)$. Then
(a) for $\ell<i \leq p$,

$$
\frac{\left|\xi_{i}-\xi_{i}^{\prime}\right|}{\left|\xi_{i}-z_{i}\right|}<K^{6 \alpha i+1} \frac{\left|\xi_{0}-\xi_{0}^{\prime}\right|}{d_{\mathcal{C}}\left(\xi_{0}\right)}
$$

(b) with $w_{i}=D T^{i} \mathbf{v}$, we have

$$
\frac{\left|w_{p}\left(\xi_{0}\right)\right|}{\left|w_{p}\left(\xi_{0}^{\prime}\right)\right|}<\exp \left\{K \frac{\left|\xi_{0}-\xi_{0}^{\prime}\right|}{d_{\mathcal{C}}\left(\xi_{0}\right)}\right\}, \quad \angle\left(w_{p}\left(\xi_{0}\right), w_{p}\left(\xi_{0}^{\prime}\right)\right) \leq K b^{\frac{1}{2}} \frac{\left|\xi_{0}-\xi_{0}^{\prime}\right|}{d_{\mathcal{C}}\left(\xi_{0}\right)}
$$

Proof. (a) We remark that this is a rough a priori bound in which factors of $K^{\alpha i}$ are allowed to accumulate. Let $s \mapsto \xi_{0}(s)$ be the parametrization of the segment from $\xi_{0}$ to $\xi_{0}^{\prime}$. We write $S=S(\tau, \mathbf{v}), e_{p}=e_{p}(S)$, and decompose $\tau$ into $\tau=A e_{p}+B \mathbf{v}$. For $\ell<i \leq p$, since $\left|D T^{i}\left(e_{p}\right)\right|$ is negligible, we have

$$
\begin{equation*}
\left|D T_{\xi_{0}}^{i} \tau\right| \approx\left|B \| w_{i}\left(\xi_{0}\right)\right| \quad \text { where } \quad|B| \approx\left|\frac{d}{d s} e_{1}\right| d_{\mathcal{C}}\left(\xi_{0}\right) \tag{53}
\end{equation*}
$$

The combined use of Proposition 5.1, Lemma 5.1 and Lemma 5.2 gives, on the other hand,

$$
\begin{equation*}
\frac{\left|w_{i}\left(\xi_{0}\left(s_{1}\right)\right)\right|}{\left|w_{i}\left(\xi_{0}\left(s_{2}\right)\right)\right|}, \frac{\left|w_{i}\left(\xi_{0}\right)\right|}{\left|w_{i}\left(z_{0}\right)\right|} \leq K^{3 \alpha i} \tag{54}
\end{equation*}
$$

where $s_{1}, s_{2}$ are any parameters and $z_{0}$ is the guiding critical point. Clearly, $\left|B\left(s_{1}\right)\right| /\left|B\left(s_{2}\right)\right|<K$. We have thus shown that

$$
\begin{equation*}
\frac{\left|D T_{\xi_{0}\left(s_{1}\right)}^{i} \tau\right|}{\left|D T_{\xi_{0}\left(s_{2}\right)}^{i} \tau\right|} \approx \frac{\left|B\left(s_{1}\right)\right|\left|\frac{\left|w_{i}\left(\xi_{0}\left(s_{1}\right)\right)\right|}{\left|B\left(s_{2}\right)\right|}\right| \frac{w_{i}\left(\xi_{0}\left(s_{2}\right)\right) \mid}{}<K^{1+3 \alpha i} . . . .}{} \tag{55}
\end{equation*}
$$

Using " $\sim$ " to denote omitted factors of $K^{3 \alpha i}$ so the main terms show up more clearly, we then have for $\ell<i \leq p$ :
(i) $\left|\xi_{i}-\xi_{i}^{\prime}\right| \lesssim\left|D T_{\xi_{0}}^{i} \tau\right| \cdot\left|\xi_{0}-\xi_{0}^{\prime}\right|$;
(ii) $\left|\xi_{i}-z_{i}\right| \sim\left(\left|w_{i}\left(z_{0}\right)\right| d_{\mathcal{C}}\left(\xi_{0}\right)\right) \cdot d_{\mathcal{C}}\left(\xi_{0}\right)$.
(i) comes from $\left|\xi_{i}-\xi_{i}^{\prime}\right| \leq \int\left|D T_{\xi_{0}(s)}^{i} \tau\left(\xi_{0}(s)\right)\right| d s$ together with (55); (ii) is (A5)(iii). The assertions in this sublemma are immediate upon comparing (i) and (ii), substituting in (53), and using the comparison of $\left|w_{i}\left(\xi_{0}\right)\right|$ and $\left|w_{i}\left(z_{0}\right)\right|$ in (54).
(b) Proposition 5.1 can be written as $\left|w_{p}\left(\xi_{0}\right)\right| /\left|w_{p}\left(\xi_{0}^{\prime}\right)\right|<\exp \left\{\sum_{i=1}^{p-1} K \mathcal{D}_{i}\right\}$ where

$$
\mathcal{D}_{i}=2^{K \alpha \theta i}\left|\xi_{i}-\xi_{i}^{\prime}\right|\left(\frac{1}{d_{\mathcal{C}}\left(z_{i}\right)}+\frac{b^{\frac{1}{4}}}{d_{\mathcal{C}}\left(z_{i+1}\right)}+\cdots+\frac{b^{\frac{p-i}{4}}}{d_{\mathcal{C}}\left(z_{p}\right)}\right)<2 e^{2 \alpha i}\left|\xi_{i}-\xi_{i}^{\prime}\right|
$$

The upper bound for $\left|\xi_{i}-\xi_{i}^{\prime}\right|$ in (a) is used in the estimates below.
Case 1. $i \geq \ell$. Using $\left|\xi_{i}-z_{i}\right|<e^{-\beta i}$, we obtain $\mathcal{D}_{i}<K e^{-(\beta-K \alpha) i} \cdot \frac{\left|\xi_{0}-\xi_{0}\right|}{d_{c}\left(\xi_{0}\right)}$.
Case 2. $i<\ell$. Using $\left|\xi_{i}-z_{i}\right| \leq\|D T\|^{i} d_{\mathcal{C}}\left(\xi_{0}\right)$ and $b^{3 \ell} \approx d_{\mathcal{C}}\left(\xi_{0}\right)$, we obtain

$$
\frac{d_{\mathcal{C}}\left(\xi_{0}\right)}{\left|\xi_{0}-\xi_{0}^{\prime}\right|} \sum_{i<\ell} \mathcal{D}_{i}<\sum_{i<\ell} K^{i} d_{\mathcal{C}}\left(\xi_{0}\right)<K^{\ell} d_{\mathcal{C}}\left(\xi_{0}\right)<\left(d_{\mathcal{C}}\left(\xi_{0}\right)\right)^{1-K \theta} .
$$

The angle estimate is similar.
Remark. For $i=p$, the inequality sign in (i) in the proof of (a) becomes $\sim$ because $T^{p} \gamma$ is roughly horizontal. The same argument then gives

$$
\begin{equation*}
\frac{\left|\xi_{p}-\xi_{p}^{\prime}\right|}{\left|\xi_{p}-z_{p}\right|}>K^{-6 \alpha p-1} \frac{\left|\xi_{0}-\xi_{0}^{\prime}\right|}{d_{\mathcal{C}}\left(\xi_{0}\right)} . \tag{56}
\end{equation*}
$$

Sublemma A.22.2. Letting $\tau$ and $\tau^{\prime}$ be unit tangent vectors to $\gamma$ at $\xi_{0}$ and $\xi_{0}^{\prime}$ respectively, we have

$$
\frac{\left|D T_{\xi_{0}}^{p} \tau\right|}{\left|D T_{\xi_{0}^{\prime}}^{p} \tau^{\prime}\right|}<\exp \left\{K \frac{\left|\xi_{0}-\xi_{0}^{\prime}\right|}{d_{\mathcal{C}}\left(\xi_{0}\right)}\right\}
$$

Proof. Splitting $\tau=A e_{p}+B \mathbf{v}$ where $e_{p}$ is the most contracted direction of $D T^{p}$ in $S(\tau, \mathbf{v})$, and letting $V=D T_{\xi_{0}}^{p} \mathbf{v}, V^{\prime}=D T_{\xi_{0}}^{p} \mathbf{v}, E=D T_{\xi_{0}}^{p} e_{p}$, and $E^{\prime}=D T_{\xi_{0}^{\prime}}^{p} e_{p}^{\prime}$, we obtain

$$
\begin{equation*}
\frac{\left|D T_{\xi_{0}}^{p} \tau^{\prime}\right|}{\left|D T_{\xi_{0}}^{p} \tau\right|}<\frac{\left|V^{\prime}\right|}{|V|} \cdot \frac{\left|B^{\prime}\right|}{|B|} \cdot(1+2(I)+2(I I)) \tag{57}
\end{equation*}
$$

where

$$
(I)=\left|\frac{V}{|V|}-\frac{V^{\prime}}{\left|V^{\prime}\right|}\right|, \quad(I I)=\left|\frac{A^{\prime} E^{\prime}}{B^{\prime}\left|V^{\prime}\right|}-\frac{A E}{B|V|}\right| .
$$

To obtain (57), we have used $|A E| \ll|B V|$ and $\left|A^{\prime} E^{\prime}\right| \ll\left|B^{\prime} V^{\prime}\right|$. For $\frac{\left|V^{\prime}\right|}{|V|}$, see Sublemma A.22.1(b). Since $\gamma$ is $C^{2}(b)$, we have $\left|\tau-\tau^{\prime}\right| \leq \frac{K b}{\delta^{3}}\left|\xi_{0}-\xi_{0}^{\prime}\right|$. Lemma 3.1 then gives $\left|e_{p}-e_{p}^{\prime}\right| \leq K\left|\xi_{0}-\xi_{0}^{\prime}\right|$. The remaining estimates resemble those in the proof of Proposition 5.1 in Appendix A.14. As in (32), we have

$$
\begin{equation*}
\left|B-B^{\prime}\right|, \quad\left|A-A^{\prime}\right|<K\left(\left|\xi_{0}-\xi_{0}^{\prime}\right|+\left|e_{p}-e_{p}^{\prime}\right|+\left|\tau^{\prime}-\tau\right|\right)<K\left|\xi_{0}-\xi_{0}^{\prime}\right| . \tag{58}
\end{equation*}
$$

This gives $\frac{\left|B^{\prime}\right|}{|B|} \leq 1+\frac{\left|B^{\prime}-B\right|}{|B|}<1+\frac{K\left|\xi_{0}-\xi_{0}^{\prime}\right|}{d_{\mathcal{C}}\left(\xi_{0}\right)}$. (I) is the angle part of Sublemma A.22.1(b). As in Case 3 in the proof of Proposition 5.1, (II) is bounded by the sum of a collection of terms of the form

$$
\begin{array}{rlrl}
(\text { i }) & =K \frac{\left|E-E^{\prime}\right|}{d_{\mathcal{C}}\left(\xi_{0}\right)}, & (\text { ii }) & =\frac{|E|}{d_{\mathcal{C}}\left(\xi_{0}\right)}\left|\frac{B^{\prime}}{B}-1\right| \\
\text { (iii) } & =\frac{|E|}{d_{\mathcal{C}}\left(\xi_{0}\right)}\left|\frac{A^{\prime}}{A}-1\right|, & \text { (iv) }=\frac{|E|}{d_{\mathcal{C}}\left(\xi_{0}\right)}\left|\frac{\left|V^{\prime}\right|}{|V|}-1\right| .
\end{array}
$$

For (i), Lemma 3.1 gives $\left|E-E^{\prime}\right|<(K b)^{p}\left|\xi_{0}-\xi_{0}^{\prime}\right|$. Observing that $|E| \ll d_{\mathcal{C}}\left(\xi_{0}\right)$, we estimate (ii) using the bound on $\frac{\left|B^{\prime}\right|}{|B|}$ above, (iii) is similar, and (iv) is given by Sublemma A.22.1(b).

Proof of ( $\mathrm{P} 2^{\prime}$ ). Extending $\omega$ as a $C^{2}(b)$ curve to $B^{(j)}\left(\hat{z}_{0}\right)$ if necessary, we obtain (P2')(i) from (A5)(i). For (ii), the desired bound follows from (A5)(ii) and (54) above. For (iii), (56) gives

$$
\left|T^{p}(\omega)\right| \geq K^{-1} e^{-6 \alpha p}\left|\xi_{p}-\hat{z}_{p}\right| \cdot \frac{\left|I_{\mu j}\right|}{e^{-|\mu|}}>K^{-1} \frac{1}{\mu^{2}} e^{-(\beta+6 \alpha) p}>\frac{1}{\mu^{2}} e^{-K_{1} \alpha|\mu|} .
$$

Proof of ( $\mathrm{P} 3^{\prime}$ ). Follow the proof of (P3) in Appendix A. 1 and use Sublemma A.22.2.
A.23. Bounds for $\frac{d}{d a} z_{0}(\S \S 10.1$ and 10.2)

Proof of Lemma 10.1. (a) To bound the first derivatives of $u$, it suffices to estimate $\partial_{a} \psi$ and $\partial_{a x} \psi$; bounds for $\partial_{x} \psi$ and $\partial_{x x} \psi$ are known since $x \mapsto \gamma(x, a)$ is $C^{2}(b)$. To pass between $\gamma$ and $l=T_{a}^{-k} \gamma$, we use the notation

$$
t(x, a):=\pi_{x}\left(T_{a}^{-k} \gamma(x, a)\right) \quad \text { and } \quad(X(t, a), Y(t, a)):=T_{a}^{k}\left(t, y_{0}\right),
$$

assuming $l \subset\left\{y=y_{0}\right\}$. Differentiating $\psi(x, a)=Y(t(x, a), a)$, we obtain

$$
\begin{equation*}
\partial_{a} \psi=\partial_{t} Y(t, a) \partial_{a} t(x, a)+\partial_{a} Y(t, a) \tag{59}
\end{equation*}
$$

Here and in the rest of the proof, we use the fact that all first and second partial derivatives of $Y$ are bounded above by $K^{k} b$, and corresponding partials of $X$ are bounded by $K^{k}$. Partials in $t$, however, are potentially problematic and must be treated with care. To bound $\partial_{a} t(x, a)$, we write it as

$$
\begin{equation*}
\partial_{a} t(x, a)=-\frac{\partial_{a} X(t, a)}{\partial_{t} X(t, a)} . \tag{60}
\end{equation*}
$$

Since $\left.T_{a}^{k}\right|_{l}$ is controlled, $\left|\partial_{t} X(t, a)\right|>1$, and so this term is $<K^{k}$. Thus $\left|\partial_{a} \psi\right|<K^{k}$.

To estimate $\partial_{a x} \psi$, we take one more derivative to obtain

$$
\partial_{a x} \psi=\partial_{t t} Y \partial_{x} t \partial_{a} t+\partial_{t} Y \partial_{a x} t+\partial_{a t} Y \partial_{x} t
$$

Since $t=t(x, a)$ is implicitly defined by $x=X(t, a)$, we have $\left|\partial_{x} t(x, a)\right|=$ $\left|\partial_{t} X(t, a)\right|^{-1}<1$, and finally

$$
\left|\partial_{x a} t(x, a)\right|=\frac{1}{\left|\partial_{t} X(t, a)\right|^{2}}\left|\partial_{a t} X \partial_{x} t \partial_{t} X-\partial_{t t} X \partial_{x} t \partial_{a} X\right|<K^{k+1}
$$

This completes the proof of $\left|\partial_{a x} \psi\right|<K^{k+1}$.
To bound the second derivatives of $u$, we need to bound the third derivatives of $\psi$. These are estimated similarly and are left to the reader. Since $v=(\mathbf{v}-\langle u, \mathbf{v}\rangle u) /|\mathbf{v}-\langle u, \mathbf{v}\rangle u|$, bounds for $\|v\|_{C^{2}}$ follow from those of $u$. This completes the proof of (a).

For (b) we cannot appeal simply to Lemma 3.1 because the bound on the $C^{2}$-norms of $u$ and $v$ in Lemma 10.1 is not a single number depending on the family $T_{a}$; it increases with the generation of the critical point. We go directly instead to the formulas for the most contracted directions in Appendix A.4II. Since $\eta_{k+1}$ is the quantity $\beta$ in Appendix A.4II with $S=S(u, v)$ and $M=D T_{a}^{k+1}(\gamma(x, a))$, we have $\eta_{k+1}=\langle M u, M v\rangle$. (b) follows now from (a) and the $C^{2}$-norm of $M$.

Proof of Lemma 10.3. From Proposition 10.1, $z_{0}^{(k)}(a)$ is well-defined on $J_{n}$ with $n=k \theta^{-1}$. Let $k^{\prime}<k$ be the largest integer such that $Q^{\left(k^{\prime}\right)}(a) \supset$ $Q^{(k)}(a)$, and let $z_{0}^{\left(k^{\prime}\right)}(a)=z_{0}^{*}\left(Q^{\left(k^{\prime}\right)}(a)\right)$. Then $(1+2 \theta)^{-1} k \leq k^{\prime}\left(\right.$ see $\left.\left(\mathrm{Al}^{\prime}\right)\right)$ and $\left|z_{0}^{\left(k^{\prime}\right)}(a)-z_{0}^{(k)}(a)\right|<K b^{\frac{k^{\prime}}{4}}$ (Lemma 4.1).

The following calculus estimate will be used: Let $g$ be a real valued $C^{2}$ function defined on an interval of length $L$, and assume that $|g| \leq M_{0}$ and $\left|g^{\prime \prime}\right|<M_{2}$. If $4 M_{0}<L^{2}$, then $\left|g^{\prime}\right| \leq \sqrt{M_{0}}\left(1+M_{2}\right)$.

To apply this estimate, we write $z_{0}^{(k)}(a)=\left(x_{0}^{(k)}(a), y_{0}^{(k)}(a)\right)$, and let $g(a)=$ $x_{0}^{(k)}(a)-x_{0}^{\left(k^{\prime}\right)}(a)$. Then $g$ is defined on $J_{n}$ and so $L=2 \rho^{n}=2 \rho^{k \theta^{-1}}$. Here $M_{0}=K b^{\frac{k^{\prime}}{4}}$, and $M_{2}=K^{k}$ from Corollary 10.1. When $b^{\theta}<\rho^{5}, 4 M_{0}<L^{2}$ holds. Therefore

$$
\left|\frac{d}{d a}\left(x_{0}^{(k)}(a)-x_{0}^{\left(k^{\prime}\right)}(a)\right)\right|<b^{\frac{k^{\prime}}{8}} K^{k}<b^{\frac{k^{\prime}}{9}} .
$$

A similar estimate holds for $\frac{d}{d a} y_{0}^{(k)}$.

## A.24. Equivalence of $\tau$ - and $a$-derivatives (§11.1)

Proof of Lemma 11.1. In this proof we fix $i_{0}$ and let $(a, b) \rightarrow\left(a^{*}, 0\right)$. Recall that if $\tau_{0}=\left(\tau_{0, x}, \tau_{0, y}\right)$, then by Corollary 10.2, $\tau_{0, y} \rightarrow 0$ as $b \rightarrow 0$. The two terms of $V$ are estimated as follows:
(i) Writing $T_{a^{*}, 0}^{i_{0}}=\left(T^{1}, 0\right)$, we have, as $b \rightarrow 0$,

$$
\left(D T_{a}^{i_{0}}\right)_{z_{0}} \tau_{0} \rightarrow\left(\frac{\partial T^{1}}{\partial x}\left(x_{0}, 0\right) \tau_{0, x}+\frac{\partial T^{1}}{\partial y}\left(x_{0}, 0\right) \tau_{0, y}, 0\right)=(0,0) .
$$

(ii) Assume $z_{s}$ stays out of $\mathcal{C}^{(1)}$ for $>i_{0}$ iterates. Then as $(a, b) \rightarrow\left(a^{*}, 0\right)$,

$$
\begin{aligned}
\frac{\sum_{s=1}^{i_{0}} D T_{z_{s}}^{i_{0}-s} \psi\left(z_{s-1}\right)}{\left|w_{i_{0}}\right| /\left|w_{1}\right|} & \rightarrow\left(\frac{\sum_{s=1}^{i_{0}}\left(f^{i_{0}-s}\right)^{\prime}\left(x_{s}\left(a^{*}\right)\right) \frac{d}{d a}\left(f_{a}\left(x_{s-1}\right)\right)\left(a^{*}\right)}{ \pm\left(f^{i_{0}-1}\right)^{\prime}\left(x_{1}\left(a^{*}\right)\right)}, 0\right) \\
& =\left( \pm \sum_{s=1}^{i_{0}} \frac{\frac{d}{d a}\left(f_{a}\left(x_{s-1}\right)\right)\left(a^{*}\right)}{\left(f^{s-1}\right)^{\prime}\left(x_{1}\left(a^{*}\right)\right)}, 0\right) .
\end{aligned}
$$

Proof of Lemma 11.2.

$$
\begin{aligned}
\left|\angle\left(w_{i}, \tau_{i}\right)\right| & \approx \frac{\left|w_{i} \wedge \tau_{i}\right|}{\left|w_{i}\right|\left|\tau_{i}\right|} \leq \frac{1}{\left|\tau_{i}\right|}\left(\sum_{s=1}^{i} \frac{1}{\left|w_{i}\right|}\left|w_{i} \wedge D T_{z_{s}}^{i-s} \psi\left(z_{s-1}\right)\right|+\frac{\left|w_{i} \wedge D T_{z_{0}}^{i} \tau_{0}\right|}{\left|w_{i}\right|}\right) \\
& \leq \frac{1}{\left|\tau_{i}\right|}\left(\sum_{s=1}^{i} \frac{\left|w_{s}\right|}{\left|w_{i}\right|}\left|\frac{w_{s}}{\left|w_{s}\right|} \wedge \psi\left(z_{s-1}\right)\right| b^{i-s}+\frac{\left|\tau_{0}\right|}{\left|w_{i}\right|} b^{i}\right) \leq \frac{K}{\left|\tau_{i}\right|} \sum_{s=0}^{\infty} b^{s} .
\end{aligned}
$$

The last inequality is valid if $\left|w_{s}\right| \leq\left|w_{i}\right|$ for all $s \leq i$, which is the case at free returns.

## A.25. Bound period estimates for parameters (§11.2)

We begin with some estimates on derivative comparisons. Let $a, a^{\prime} \in \hat{\omega}$ be as in Lemma 11.5. We let $\xi_{0}=\zeta_{n}(a), \xi_{0}^{\prime}=\zeta_{n}\left(a^{\prime}\right), w_{i}\left(\xi_{0}\right)=\left(D T_{a}^{i}\right) \xi_{0} \mathbf{v}, w_{i}\left(\xi_{0}^{\prime}\right)=$ $\left(D T_{a^{\prime}}^{i}\right)_{\xi_{0}^{\prime} \mathbf{v}}$ and $p=p(\hat{\omega})$. We wish to compare $w_{i}\left(\xi_{0}\right)$ and $w_{i}\left(\xi_{0}^{\prime}\right)$ for $i \leq p$.

Sublemma A.25.1 (Parameter version of Proposition 5.1). There exists $K>0$ such that

$$
\begin{align*}
& \frac{\left|w_{i}\left(\xi_{0}\right)\right|}{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|}, \frac{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|}{\left|w_{i}\left(\xi_{0}\right)\right|} \leq K^{3 \alpha i} \cdot \exp \left\{\sum_{j=0}^{i-1} K e^{2 \alpha j}\left|\xi_{j}-\xi_{j}^{\prime}\right|+K^{i}\left|a-a^{\prime}\right|\right\} \text { for } i \leq p  \tag{61}\\
& \text { (62) } \quad \angle\left(w_{p}\left(\xi_{0}\right), w_{p}\left(\xi_{0}^{\prime}\right)\right)<b^{\frac{1}{2}} \sum_{j=0}^{p-1} b^{\frac{j}{4}} 2^{\ell_{p-1-j}}\left|\xi_{p-j-1}-\xi_{p-j-1}^{\prime}\right|+K^{p}\left|a-a^{\prime}\right| . \tag{62}
\end{align*}
$$

Remarks. (i) The factor $K^{3 \alpha i}$ in (61) can be dropped if $\xi_{i}$ is out of all splitting periods (see the proof below). (ii) We may assume the quantity inside brackets in (61) is $\ll 1$ (cf. Lemma 5.1). This is because $p<K \alpha n$ and $\left|a-a^{\prime}\right|<\hat{K} e^{-\lambda^{\prime} n}$ (Proposition 11.1).

Proof. Let $\eta_{j}=T_{a}^{j} \xi_{0}^{\prime}, v_{j}=\left(D T_{a}^{j}\right)_{\xi_{0}^{\prime}} \mathbf{v}$. Then $\frac{\left|w_{i}\left(\xi_{0}\right)\right|}{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|}=\frac{\left|w_{i}\left(\xi_{0}\right)\right|}{\left|v_{i}\right|} \frac{\left|v_{i}\right|}{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|}$. Since $\left|\xi_{j}-\eta_{j}\right| \leq\left|\xi_{j}-\xi_{j}^{\prime}\right|+\left|\xi_{j}^{\prime}-\eta_{j}\right|<\left|\xi_{j}-\xi_{j}^{\prime}\right|+K^{j}\left|a-a^{\prime}\right|$, we have, by Proposition 5.1 (see the proof of Sublemma A.22.1 in Appendix A.22),

$$
\begin{align*}
\frac{\left|w_{i}\left(\xi_{0}\right)\right|}{\left|v_{i}\right|} & \leq K^{3 \alpha i} \exp \left\{\sum_{j=0}^{i-1} K e^{2 \alpha j}\left|\xi_{j}-\eta_{j}\right|\right\}  \tag{63}\\
& \leq K^{3 \alpha i} \exp \left\{\sum_{j=0}^{i-1} K e^{2 \alpha j}\left|\xi_{j}-\xi_{j}^{\prime}\right|+K^{i}\left|a-a^{\prime}\right|\right\},
\end{align*}
$$

the $K^{3 \alpha i}$ factor being there to account for the discrepancy between $w_{i}\left(\xi_{0}\right)$ and $w_{i}^{*}\left(\xi_{0}\right)$. Since $\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|>K^{-1}$ and $\left|v_{i}-w_{i}\left(\xi_{0}^{\prime}\right)\right|<K^{i}\left|a-a^{\prime}\right|$, we have

$$
\begin{equation*}
\frac{\left|v_{i}\right|}{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|}<1+\frac{\left|v_{i}-w_{i}\left(\xi_{0}^{\prime}\right)\right|}{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|}<1+K^{i}\left|a-a^{\prime}\right| \tag{64}
\end{equation*}
$$

completing the proof of (61).
For (62), we write $\angle\left(w_{p}\left(\xi_{0}\right), w_{p}\left(\xi_{0}^{\prime}\right)\right) \leq \angle\left(w_{p}\left(\xi_{0}\right), v_{p}\right)+\angle\left(v_{p}, w_{p}\left(\xi_{0}^{\prime}\right)\right)$. By Proposition 5.1,

$$
\begin{align*}
\angle\left(w_{p}\left(\xi_{0}\right), v_{p}\right) & <b^{\frac{1}{2}} \sum_{j=0}^{p-1} b^{\frac{j}{4}} 2^{\ell_{p-j-1}}\left|\xi_{p-j-1}-\eta_{p-j-1}\right|  \tag{65}\\
& <b^{\frac{1}{2}}\left(\sum_{j=0}^{p-1} b^{\frac{j}{4}} 2^{\ell_{p-j-1}}\left|\xi_{p-j-1}-\xi_{p-j-1}^{\prime}\right|+K^{p}\left|a-a^{\prime}\right|\right) .
\end{align*}
$$

To estimate $\angle\left(v_{p}, w_{p}\left(\xi_{0}^{\prime}\right)\right)$, note that by (64) and Remark (ii) above, $\frac{\left|v_{i}\right|}{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|} \approx 1$ for all $i \leq p$. Proceeding inductively, we let $i \leq \frac{p}{2}, u=\frac{v_{i}}{\left|v_{i}\right|}$, and $\hat{u}=\frac{w_{i}\left(\xi_{0}^{\prime}\right)}{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|}$.

Since $\left|w_{2 i}\left(\xi_{0}^{\prime}\right)\right|>K^{-1}$, we have

$$
\begin{aligned}
& \angle\left(v_{2 i}, w_{2 i}\left(\xi_{0}^{\prime}\right)\right)<\left|\left(D T_{a}^{i}\right)_{\eta_{i}} u \wedge\left(D T_{a^{\prime}}^{i}\right)_{\xi_{i}^{\prime}} \hat{u}\right| \frac{\left|v_{i}\right|}{\left|v_{2 i}\right|} \frac{\left|w_{i}\left(\xi_{0}^{\prime}\right)\right|}{\left|w_{2 i}\left(\xi_{0}^{\prime}\right)\right|} \\
& \quad \leq\left(\left|\left(D T_{a}^{i}\right)_{\eta_{i}} u \wedge\left(D T_{a}^{i}\right)_{\eta_{i}} \hat{u}\right|+\left|\left(D T_{a}^{i}\right)_{\eta_{i}} u \wedge\left(\left(D T_{a}^{i}\right)_{\eta_{i}}-\left(D T_{a^{\prime}}^{i}\right)_{\xi_{i}^{\prime}}\right) \hat{u}\right|\right) K^{2 i} \\
& \quad \leq(K b)^{2} \angle\left(v_{i}, w_{i}\left(\xi_{0}^{\prime}\right)\right)+K^{4 i}\left|a-a^{\prime}\right| .
\end{aligned}
$$

We conclude inductively that $\angle\left(v_{p}, w_{p}\left(\xi_{0}^{\prime}\right)\right)<K^{2 p}\left|a-a^{\prime}\right|$, completing the proof of (62).

Sublemma A. 25.2 (Parameter version of Sublemma A.22.1 in Appendix A.22). Let $\hat{z}_{0}=\phi\left(\xi_{0}(a)\right)$. Then
(a) for $\ell<i<p$ where $\ell$ is the splitting period of $\xi_{0}$, we have

$$
\frac{\left|\xi_{i}-\xi_{i}^{\prime}\right|}{\left|\xi_{i}-\hat{z}_{i}\right|}<K^{6 \alpha i+1} \frac{\left|\xi_{0}-\xi_{0}^{\prime}\right|}{\left|\xi_{0}-\hat{z}_{0}\right|}, \quad \frac{\left|\xi_{p}-\xi_{p}^{\prime}\right|}{\left|\xi_{p}-\hat{z}_{p}\right|}>K^{-6 \alpha p-1} \frac{\left|\xi_{0}-\xi_{0}^{\prime}\right|}{\left|\xi_{0}-\hat{z}_{0}\right|} ;
$$

(b) $\frac{\left|w_{p}\left(\xi_{0}\right)\right|}{\left|w_{p}\left(\xi_{0}^{\prime}\right)\right|}<\exp \left\{K \frac{\left|\xi_{0}-\xi_{0}^{\prime}\right|}{d_{\mathcal{C}(a)}\left(\xi_{0}\right)}+K^{p}\left|a-a^{\prime}\right|\right\}$.

Proof. The proof follows closely that of Sublemma A. 22.1 with the following modifications: In part (a), we consider the parametrization of the critical curve $\zeta_{n_{k}}$ from $\xi_{0}$ to $\xi_{0}^{\prime}$ by arclength, and split its tangent vectors $\tau$. The correctness of this splitting is a consequence of Lemma 11.2 and the fact that $w_{n_{k}}$ splits correctly. (53) is a statement about individual parameters. To prove (54), we use Sublemma A.25.1 instead of Proposition 5.1. The rest of the proof then proceeds as before. The term $K^{p}\left|a-a^{\prime}\right|$ in (b) is from the corresponding term in (61).

Proof of Lemma 11.5. Let $\tilde{a} \in \hat{\omega}$ be the parameter at which the minimum in the definition of $p(\hat{\omega})$ is attained. Then (a) is an immediate consequence of (A5)(i) for $T_{\tilde{a}}$.

Let $\hat{z}_{0}(a)=\phi_{a}\left(\zeta_{n}(a)\right)$. (b) follows from the fact that for all $a \in \hat{\omega}$ and $j<p(\hat{\omega}),\left|\hat{z}_{j}(a)-\hat{z}_{j}(\tilde{a})\right| \leq K^{j}|\hat{\omega}| \leq K^{\alpha n} \hat{K} e^{-\lambda^{\prime} n} \ll e^{-\beta j}$. In the second inequality we have used $p(\hat{\omega}) \leq \alpha n$ and $|\hat{\omega}| \leq \hat{K} e^{-\lambda^{\prime} n}$, which follows from $\left|\zeta_{n}(\hat{\omega})\right| \leq 1$ and Proposition 11.1.
(c) is proved via the following string of inequalities:

$$
\frac{\left|\tau_{n+p}(a)\right|}{\left|\tau_{n}(a)\right|}>\hat{K}^{-2} \frac{\left|w_{n+p}(a)\right|}{\left|w_{n}(a)\right|}>K^{-1} \hat{K}^{-2} \frac{\left|w_{n+p}(\tilde{a})\right|}{\left|w_{n}(\tilde{a})\right|}>K^{-1} \hat{K}^{-2} e^{\frac{p}{3} \lambda} .
$$

The first inequality above is based on Proposition 11.1. For the second inequality, first recall that for both of the maps $T_{a}$ and $T_{\tilde{a}}$, since $w_{n}$ splits correctly, we have $\left|w_{n+p}\right| \approx \frac{1}{2} \frac{d e_{1}}{d s} d_{\mathcal{C}}\left(z_{n}\right)\left|D T_{z_{n}}^{p}(\mathbf{v})\right| \cdot\left|w_{n}\right|$. We then use Sublemma A.25.1 and the remarks following it to compare $\left|\left(D T_{a}^{p}\right)_{z_{n}(a)}(\mathbf{v})\right|$ and $\left|\left(D T_{\tilde{a}}^{p}\right)_{z_{n}(\tilde{a})}(\mathbf{v})\right|$,
and note that the other factors are comparable up to a fixed constant. The last inequality follows from Proposition $5.2(2)$ for $T_{\tilde{a}}$.
(d) is a simple consequence of the bound on $\angle\left(\tau_{n+i}, w_{n+i}\right)$ (Lemma 11.2) and the fact that $\angle\left(w_{n+i}, D T_{z_{n}}^{i}(\mathbf{v})\right) \ll 1$ outside of splitting periods.
(e) is an application of the second inequality in Sublemma A.25.2(a).

## A.26. Distortion estimates for parameters (§12.2)

We prove Lemma 12.2 in this appendix. Let $\omega \in \mathcal{Q}_{n-1}$ be as in Lemma 12.2 , and let $a, a^{\prime} \in \omega$. Where no ambiguity arises, we will omit mention of the parameters and write $z_{i}=z_{i}(a), z_{i}^{\prime}=z_{i}\left(a^{\prime}\right), w_{i}=w_{i}\left(z_{0}\right)=\left(D T_{a}^{i}\right)_{z_{0}}(\mathbf{v})$, $w_{i}^{\prime}=w_{i}\left(z_{0}^{\prime}\right)=\left(D T_{a^{\prime}}^{i}\right)_{z_{0}^{\prime}}(\mathbf{v})$, and similarly for $\tau_{i}$ and $\tau_{i}^{\prime}$.

Plan of proof. Since the formula for the evolution of $\tau_{i}$ is more involved, we again invoke Proposition 11.1 and prove $\frac{\left|w_{n}\right|}{\left|w_{n}^{\prime}\right|}<K$. Let $0<n_{1}<n_{1}+p_{1} \leq$ $n_{2}<n_{2}+p_{2} \leq n_{3}<\cdots<n_{q}+p_{q} \leq n$ be such that $n_{k}$ is a free return, $p_{k}$ is the ensuing bound period, and $n_{k+1}$ is the first return following $n_{k}+p_{k}$. We write

$$
\begin{equation*}
\left|w_{n}\right|=\left|w_{n_{1}}\right| \cdot \frac{\left|w_{n_{1}+p_{1}}\right|}{\left|w_{n_{1}}\right|} \cdot \frac{\left|w_{n_{2}}\right|}{\left|w_{n_{1}+p_{1}}\right|} \cdot \frac{\left|w_{n_{2}+p_{2}}\right|}{\left|w_{n_{2}}\right|} \cdot \frac{\left|w_{n_{3}}\right|}{\left|w_{n_{2}+p_{2}}\right|} \cdots \tag{66}
\end{equation*}
$$

and estimate the factors in (66) separately. These factors are of two types, the more complicated of which being $\frac{\left|w_{n_{k}+p_{k}}\right|}{\left|w_{n_{k}}\right|}$. In the proof of Lemma $11.5(\mathrm{c})$ in Appendix A.25, we reduced the comparison of $\frac{\left|w_{n_{k}+p_{k}}\right|}{\left|w_{n_{k}}\right|}$ to that of $\left|D T_{z_{n_{k}}}^{p_{k}}(\mathbf{v})\right|$ with bounded error. A more refined estimate is needed here to control the cumulative effect of these errors over time intervals that may contain arbitrarily large numbers of bound periods. Such a comparison involves the difference in slopes between $w_{n_{k}}$ and $w_{n_{k}}^{\prime}$. Let $\theta_{i}:=\angle\left(w_{i}, w_{i}^{\prime}\right)$.

Sublemma A.26.1. (i) Let $i_{0}$ be as in Proposition 11.1. Then $\theta_{i_{0}}<$ $K^{i_{0}}\left|a-a^{\prime}\right|$.
(ii) For all $k \geq 1$,
(a) $\theta_{n_{k}}<K b^{\frac{1}{2}}\left|z_{n_{k}-1}-z_{n_{k}-1}^{\prime}\right|+2 b^{\frac{1}{2}}\left|a-a^{\prime}\right|+b^{\frac{1}{2}\left(n_{k}-\left(n_{k-1}+p_{k-1}\right)\right)} \theta_{n_{k-1}+p_{k-1}}$, with " $n_{0}+p_{0}$ " in the inequality above replaced by " $i_{0}$ " in the case $k=1$;
(b) $\theta_{n_{k}+p_{k}} \leq 2 b^{\frac{1}{2}} \sum_{j=0}^{p_{k}-1} b^{\frac{j}{4}} 2^{\ell_{n_{k}+p_{k}-j-1}}\left|z_{n_{k}+p_{k}-j-1}-z_{n_{k}+p_{k}-j-1}^{\prime}\right|$

$$
+K^{p_{k}}\left|a-a^{\prime}\right|+b^{\frac{p_{k}}{4}} \theta_{n_{k}}
$$

Proof. (i) is straightforward since $\left|z_{0}-z_{0}^{\prime}\right|<K\left|a-a^{\prime}\right|$ from Corollary 10.2. (ii) follows from estimates very similar to those in the proof of Proposition 5.1 (Appendix A.14). More precisely:
(a) Let $\Theta_{j}=\angle\left(\left(D T_{a}^{j}\right)_{z_{n_{k}+p_{k}}} w_{n_{k}+p_{k}},\left(D T_{a^{\prime}}^{j}\right)_{z_{n_{k}+p_{k}}^{\prime}} w_{n_{k}+p_{k}}^{\prime}\right)$. The assertion is proved inductively by showing, as in Case 1 of Proposition 5.1,

$$
\begin{equation*}
\Theta_{j} \leq b^{\frac{1}{2}}\left(\Theta_{j-1}+\left|a-a^{\prime}\right|+\left|z_{n_{k}+p_{k}+j-1}-z_{n_{k}+p_{k}+j-1}^{\prime}\right|\right) \tag{67}
\end{equation*}
$$

(b) Let $p=p_{k}$, and write $u=\frac{w_{n_{k}}\left(\xi_{0}\right)}{\left|w_{n_{k}}\left(\xi_{0}\right)\right|}$. Let $e=e_{p}(S)$ where $S=S(u, \mathbf{v})$. As usual, we split $u$ into $u=B \mathbf{v}+A e$. Following the computation in the proof of Proposition 5.1, Case 2, we obtain

$$
\begin{equation*}
\left|B^{\prime}-B\right|,\left|A^{\prime}-A\right|<K\left(\theta_{n_{k}}+\left|z_{n_{k}}-z_{n_{k}}^{\prime}\right|+\left|a-a^{\prime}\right|\right) \tag{68}
\end{equation*}
$$

Here $\left|u-u^{\prime}\right|=\theta_{n_{k}}$, and by Lemma 3.1, $\left|e-e^{\prime}\right| \leq K\left(\theta_{n_{k}}+\left|a-a^{\prime}\right|+\left|z_{n_{k}}-z_{n_{k}}^{\prime}\right|\right)$. The rest of the proof follows Case 3 in the same proof. As usual, we write $V=\left(D T_{a}^{p}\right)_{z_{n}} \mathbf{v}, V^{\prime}=\left(D T_{a^{\prime}}^{p}\right)_{z_{n}^{\prime}} \mathbf{v}, E=\left(D T_{a}^{p}\right)_{z_{n}} e, E^{\prime}=\left(D T_{a^{\prime}}^{p}\right)_{z_{n}^{\prime}} e^{\prime}$. Then we have $\theta_{n_{k}+p_{k}} \leq(I)+(I I)$ where $(I)=\angle\left(V, V^{\prime}\right),(I I)=\left|\frac{A^{\prime} E^{\prime}}{B^{\prime}\left|V^{\prime}\right|}-\frac{A E}{B|V|}\right|$. For $(I)$ we use the angle part of Sublemma A.25.1. The other estimates involve the same terms as in the proof of Proposition 5.1, Case 3, and are carried out similarly.

Corollary A.1. (i) $\theta_{n_{k}} \leq b^{\frac{1}{2}}\left|z_{n_{k}-1}-z_{n_{k}-1}^{\prime}\right|+K^{K \alpha n}\left|a-a^{\prime}\right|$.
(ii) Letting $u=\frac{w_{n_{k}}}{\left|w_{n_{k}}\right|}$ and $p=p_{k}$, we have

$$
\frac{\left|\left(D T_{a}^{p}\right)_{z_{n_{k}}} u\right|}{\mid\left(D T_{a^{\prime}}^{p} z_{z_{n_{k}}} u^{\prime} \mid\right.}<\exp \left\{K \frac{\left|z_{n_{k}}-z_{n_{k}}^{\prime}\right|}{d_{\mathcal{C}(a)}\left(z_{n_{k}}\right)}+K^{K \alpha n}\left|a-a^{\prime}\right|\right\} .
$$

Proof. (i) follows inductively from Sublemma A.26.1(ii). We use $K \alpha n$ to dominate $p_{k}$, and assume $n$ is sufficiently large that $b^{\frac{1}{4}\left(n-i_{0}\right)} K^{i_{0}}<K^{K \alpha n}$. For (ii), we split $u$ as in part (ii) of Sublemma A.26.1, obtaining $\frac{\left|\left(D T^{p}\right)_{n_{k}} u\right|}{\mid\left(D T_{a^{\prime}}^{o} z_{z_{k}}^{\prime} u^{\prime} \mid\right.}<$ $\frac{|B|}{\left|B^{\prime}\right|} \left\lvert\, \frac{V}{\left|V^{\prime}\right|}(1+2(I)+2(I I))\right.$. From the estimates in Sublemma A.26.1(ii) and the bound on $\theta_{n_{k}}$ in (i) above, we see that the right side of this inequality is bounded by terms of the form as claimed.

Proof of Lemma 12.2. This proof follows that of (P3) in Appendix A.1. Letting $u_{i}=\frac{w_{i}}{\left|w_{i}\right|}$, we write $\log \frac{\left|w_{n}\right|}{\left|w_{n}^{\prime}\right|} \leq K \sum_{k=1}^{q}\left(S_{k}^{\prime}+S_{k}^{\prime \prime}\right)$ where

$$
S_{k}^{\prime}=\log \frac{\left|\left(D T_{a}^{p_{k}}\right)_{z_{n_{k}}} u_{n_{k}}\right|}{\mid\left(D T_{a^{\prime}}^{p_{k}} z_{n_{k}^{\prime}}^{\prime} u_{n_{k}}^{\prime} \mid\right.} \quad \text { and } \quad S_{k}^{\prime \prime}=\log \frac{\mid\left(D T_{a}^{n_{k+1}-\left(n_{k}+p_{k}\right)}\right)_{z_{n_{k}+p_{k}}} u_{n_{k}+p_{k} \mid}}{\left|\left(D T_{a^{\prime}}^{n_{k+1}-\left(n_{k}+p_{k}\right)}\right)_{z_{n_{k}+p_{k}}^{\prime}} u_{n_{k}+p_{k}}^{\prime}\right|}
$$

except for $S_{q}^{\prime \prime}$ which ends at index $n-1$.
To estimate $S_{k}^{\prime}$, we let $\sigma_{k}=\left|z_{n_{k}}-z_{n_{k}}^{\prime}\right|$. Then Corollary A.1(ii) gives $S_{k}^{\prime}<K \frac{\left|\sigma_{k}\right|}{d_{\mathcal{C}}\left(z_{n_{k}}\right)}+K^{K \alpha n}\left|a-a^{\prime}\right|$. The sum $\sum_{k} K \frac{\left|\sigma_{k}\right|}{d_{\mathcal{C}}\left(z_{n_{k}}\right)}$ is estimated as in Appendix A.1. The additional term representing parameter contributions sums to $<n K^{K \alpha n}\left|a-a^{\prime}\right|<n K^{K \alpha n} e^{-\lambda^{\prime} n}$, which is uniformly bounded in $n$.

Now, $\sum S_{k}^{\prime \prime}$, which treats iterates outside of $\mathcal{C}^{(1)}$, is easily estimated to be $<K \frac{\left|\sigma_{q}\right|}{\delta}+n K^{K \alpha n} e^{-\lambda^{\prime} n}$.

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[^1]:    ${ }^{1}$ Here $q(a)$ is the continuation of $q\left(a^{*}\right)$ viewed as a point whose orbit is bounded away from $C$; it is not to be confused with $f_{a}(\hat{x}(a))$.

[^2]:    ${ }^{2}$ We will, in fact, assume $f$ is sufficiently close to $f_{0}$ that $f^{n}(\hat{x}) \notin C_{\delta_{0}}$ for all $n$ with $e^{-\alpha n}>\delta$.

[^3]:    ${ }^{3}$ In particular, if $I_{\mu_{0} j_{0}}$ is one of the outermost $I_{\mu j}$ in $C_{\delta}$, then $I_{\mu_{0} j_{0}}^{+}$contains an interval of length $\delta$ just outside of $C_{\delta}$.

[^4]:    ${ }^{4}$ Quantities such as $\frac{K_{1}}{\delta^{3}} b, \frac{3 K_{0}}{\delta} b$ appearing in this definition will be denoted as $\mathcal{O}(b)$.

[^5]:    ${ }^{5}$ By this, we mean the condition $d_{\mathcal{C}}\left(z_{i}\right)>\min \left(\delta, e^{-\alpha i}\right)$ is assumed but only for the orbit of $z_{0}$; i.e., no assumptions are made on the behavior of other critical orbits beyond time $N$.

[^6]:    ${ }^{6}$ This is the reason why we require $j \leq \alpha^{*} \theta i$ in our definition of $\phi\left(z_{i}\right)$ in Section 4.1.

[^7]:    ${ }^{7}$ By "maximal" we mean no other $Q^{(k)}$ can be squeezed between between $Q^{\left(k_{i}\right)}$ and $Q^{\left(k_{i+1}\right)}$.

[^8]:    ${ }^{8}$ If the degree of $T$ is zero, then it can happen that $Q \cap\left(\cup_{j>k} \mathcal{C}^{(j)}\right)=\emptyset$.

[^9]:    ${ }^{9}$ If $H \cap P=\emptyset$ for all $P \in \mathcal{S}_{k-i+1}$, then there is no need for replacement. Also, where $I$ is an interval (see $\S 3.9$ ), the two special branches in $\mathcal{T}_{m}$ having $T^{m-1} V_{i}$ as one of their ends are always in $\mathcal{S}_{m}$.

[^10]:    ${ }^{10}$ We assume $\mathcal{F}_{1}$ is independent of $a$.

[^11]:    ${ }^{11}$ For $T_{a} \in \mathcal{G}_{n}, \Gamma_{\theta n}(a)$ is determined only up to a finite precision; the exact location of $\Gamma_{\theta n}(a)$ depends on choices of $\mathcal{F}_{k}$-leaves on which critical points are constructed.

[^12]:    ${ }^{12}$ Some fuzziness is allowed in boundary situations due to the adjoining of end intervals.

[^13]:    ${ }^{13}$ We thank O. Lanford for showing us this argument.

