# Hyperbolic components of the complex exponential family 

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#### Abstract

We describe the structure of the hyperbolic components of the parameter plane of the complex exponential family, as started in [1]. More precisely, we label each component with a parameter plane kneading sequence, and we prove the existence of a hyperbolic component for any given such sequence. We also compare these sequences with the more commonly used dynamical kneading sequences.


1. Introduction. Our goal in this paper is to describe the structure of the hyperbolic components in the parameter plane for the complex exponential family.

Let $E_{\lambda}(z)=\lambda e^{z}$ with $\lambda \in \mathbb{C}$. The map $E_{\lambda}$ has a unique singular value at 0 (the omitted value). As is well known, the fate of the orbit of 0 determines much of the dynamical behavior of $E_{\lambda}$. For example, if $E_{\lambda}$ admits an attracting cycle, then the orbit of 0 must tend to this cycle. As a consequence, $E_{\lambda}$ has at most one attracting cycle.

The parameter space of the exponential family was first studied in [1] and $[6,7]$ and later on in [3-5] and [10].

Let $\mathcal{H}_{n}$ denote the set of $\lambda$-values for which $E_{\lambda}$ admits an attracting cycle of period $n$. The connected components of $\mathcal{H}_{n}$ are called hyperbolic components and it is conjectured that they are dense in the parameter plane. As shown in [1] and [6], any hyperbolic component is simply connected and unbounded, with the exception of $\mathcal{H}_{1}$ which is a cardioid-shaped region containing 0 . The region $\mathcal{H}_{2}$ consists of a single component which occupies

[^0]

Fig. 1. The parameter plane of $E_{\lambda}$. White regions correspond to hyperbolic components. Black smooth regions are due to numerics. Dotted lines have been drawn on the imaginary axis and on the horizontal lines with imaginary parts $-3 \pi,-\pi, \pi$ and $3 \pi$.
a large portion of the left half plane. Each $\mathcal{H}_{n}$ for $n>2$ consists of infinitely many distinct components, each of which extends to $\infty$ in the right half plane.

The arrangement of these hyperbolic components in the $\lambda$-plane is quite complicated. A partial description can be found in [1] where the authors show the existence of infinitely many hyperbolic components of period $n$ in between two hyperbolic components of period $n-1$. Our goal in this paper is to give a more precise description by using the dynamics of the corresponding maps. In particular we shall give a label to each of the components which will describe the dynamical behaviour of the critical orbit for those parameters in the given component. We shall see that this label also determines the position of the component in the right half plane. See Figure 1.

With this goal in mind, there is a choice to be made, for there are two ways to identify the various hyperbolic components in the $\lambda$-plane. Each of these involves the association of a kneading sequence to the component. This sequence is a string of $n-2$ integers. For technical reasons we precede the string with a 0 and end the string with a $*$. That is, a kneading sequence
assumes the form $0 s_{1} \ldots s_{n-2} *$ with $s_{j} \in \mathbb{Z}$. The $*$ denotes a "wild card" that will be described below.

One of the two kneading sequences is a dynamical kneading sequence ( $K$ kneading sequence) which is useful mainly in the dynamical plane (see [2]), since it determines the topological structure of the Julia set of $E_{\lambda}$ for any $\lambda$ in the hyperbolic component. The other kneading sequence is a parameter plane kneading sequence ( $S$-kneading sequence) and, as we shall see, is more useful for describing the structure of the $\lambda$-plane. The main result in this paper is as follows.

Theorem A. Fix $n \geq 3$ and let $s_{1}, \ldots, s_{n-2} \in \mathbb{Z}$. There exists a hyperbolic component $\Omega_{0 s_{1} \ldots s_{n-2} *}$ that extends to $\infty$ in the right half plane and such that if $\lambda \in \Omega_{0 s_{1} \ldots s_{n-2}}$, the map $E_{\lambda}$ has an attracting cycle of period $n$ with parameter plane kneading sequence $s=0 s_{1} \ldots s_{n-2} *$. Moreover, the components $\Omega_{0 s_{1} \ldots s_{n-2}}$ are ordered lexicographically.

From the proof of this theorem one obtains the following corollary (see Figure 2).


Fig. 2. Magnification of Figure 1 showing infinitely many period 4 components in between two period 3 components

Corollary B. Let $\Omega_{0 s_{1} \ldots s_{n-2} *}$ be as in Theorem A. Then between this hyperbolic component and the hyperbolic component $\Omega_{0 s_{1} \ldots\left(s_{n-2}+1\right) *}$ there exist hyperbolic components $\Omega_{0 s_{1} \ldots\left(s_{n-2}+1\right) k *}$ for each $k \in \mathbb{Z}$.

In this statement the word "between" refers to the ordering given by the imaginary part, since all hyperbolic components extend to infinity in the right half plane.

These results give a description of the ordering of the hyperbolic components in the far right half plane as a function of their kneading sequence. Note that these are existence type results. Although uniqueness is most likely true, this fact does not follow directly from our work in this paper. D. Schleicher [10] has announced some results in this direction using the coding of hairs in parameter space.

In Section 2 below we define each of these kneading sequences and discuss several of their properties. We also derive an algorithm for obtaining one sequence given the other. In Section 3 we prove Theorem A; that is, we show the existence of hyperbolic components corresponding to any $S$-kneading sequence.
2. Kneading sequences. Let us consider a hyperbolic component $\Omega$ of period $n>2$. The main goal of this section is to define two different kneading sequences associated to the parameter value $\lambda \in \Omega$. We shall also study the relation between the two sequences and give an algorithm that transforms one into the other.

We start by giving a topological description of the dynamical plane of $E_{\lambda}(z)=\lambda e^{z}$ that holds for any parameter $\lambda$ in the hyperbolic component $\Omega$.
2.1. The fingers and the glove.. If $\lambda \in \Omega$, the map $E_{\lambda}(z)=\lambda \exp (z)$ has an attracting periodic orbit of period $n>2$. This orbit varies analytically with $\lambda$ as long as $\lambda$ lies in the hyperbolic component. Let $z_{0}(\lambda), z_{1}(\lambda)=$ $E_{\lambda}\left(z_{0}\right), \ldots, z_{n-1}(\lambda)=E_{\lambda}\left(z_{n-2}\right)$ be the points of the periodic orbit. To simplify notation we will omit the dependence on $\lambda$ if it does not lead to confusion.

Let $A^{*}$ denote the immediate basin of attraction of the periodic orbit and, for $0 \leq i \leq n-1$, define $A^{*}\left(z_{i}\right)$ to be the connected component of $A^{*}$ which contains $z_{i}$. We name the points in the orbit so that the asymptotic value 0 belongs to $A^{*}\left(z_{0}\right)$.

We now construct geometrically and define what we call fingers. More details can be found in [2]. For $\nu \in \mathbb{R}$, let $H_{\nu}=\{z \mid \operatorname{Re} z>\nu\}$.

Definition. An unbounded simply connected $F \in \mathbb{C}$ is called a finger of width $c$ if
(a) $F$ is bounded by a single simple curve $\gamma \subset \mathbb{C}$.
(b) There exists $\nu$ such that $F \cap H_{\nu}$ is simply connected, extends to infinity, and satisfies

$$
F \cap H_{\nu} \subset\{z \mid \operatorname{Im} z \in[a-d / 2, a+d / 2]\} \quad \text { for some } a \in \mathbb{R}
$$

and $c$ is the infimum value for $d$.
Observe that the preimage of any finger which does not contain 0 consists of infinitely many fingers of width smaller than $2 \pi$ which are $2 \pi i$-translates of each other.

We begin the construction by choosing $B=B(\lambda)$ to be a topological disc in $A^{*}\left(z_{0}\right)$ that contains both 0 and $z_{0}$, and having the property that $B$ is mapped strictly inside itself under $E_{\lambda}^{n}$. This set can be defined precisely using linearizing coordinates, and one can show that it moves holomorphically with $\lambda$. Although this is not crucial for this work, we have included the details in the Appendix.

We now take successive preimages of the disc $B$. More precisely, let $B_{n-1}$ be the open set in $\mathbb{C}$ which is mapped to $B$. Note that, since $0 \in B$, it follows that $B_{n-1}$ has a single connected component which contains a left half plane, and whose image under $E_{\lambda}$ wraps infinitely many times over $B \backslash\{0\}$. Note that the point $z_{n-1}$ belongs to the set $B_{n-1}$, which lies inside $A^{*}\left(z_{n-1}\right)$.

We now consider the preimage of $B_{n-1}$. It is easy to check (by looking at the image of vertical lines with increasing real part) that this preimage consists of infinitely many disjoint fingers of width less than $2 \pi$ which are $2 \pi i$-translates of each other. We define $B_{n-2} \subset A^{*}\left(z_{n-2}\right)$ to be the connected component for which $z_{n-2} \in B_{n-2}$. The map $E_{\lambda}$ takes $B_{n-2}$ conformally onto $B_{n-1}$.

Similarly, we define the sets $B_{n-3}, \ldots, B_{0}$, by setting $B_{i}$ to be the connected component of $E_{\lambda}^{-1}\left(B_{i+1}\right)$ that contains the point $z_{i}$. These inverses are all well defined and the map $E_{\lambda}$ sends $B_{i}$ conformally onto $B_{i+1}$. Each $B_{i}$ belongs to the immediate basin $A^{*}\left(z_{i}\right)$. The following characterization of the sets $B_{i}, i=0, \ldots, n-2$, is proved in [2].

Proposition 2.1. Let $n>2$. For $i=0, \ldots, n-2, B_{i}$ is a finger of width $c_{i}<2 \pi$.

It follows immediately from the above construction that the width of the finger $B_{n-2}$ that is mapped by $E_{\lambda}$ conformally onto $B_{n-1}$ is $\pi$, while the widths of the other fingers are 0 . So we will refer to $B_{n-2}$ as the big finger.

We proceed to the final step, by defining the set

$$
G=\left\{z \in \mathbb{C} \mid E_{\lambda}(z) \in B_{0}\right\},
$$

which we call the glove. We observe from the above construction that $G$ is a connected set and $B_{n-1} \subset G \subset A^{*}\left(z_{n-1}\right)$. See Figure 3. Moreover, the complement of $G$ consists of infinitely many fingers, all of which are


Fig. 3. Sketch of the sets $B_{0}$ to $B_{n-1}, G$ and $V_{j}$ for $j \in \mathbb{Z}$. Points in grey belong to the basin of attraction of the periodic orbit.
$2 \pi i$-translates of each other. We index these infinitely many connected components by $V_{j}, j \in \mathbb{Z}$, so that $2 \pi i j \in V_{j}$.

In fact, these $V_{j}$ form a set of fundamental domains for the Julia set of $E_{\lambda}$ in the following sense:

- $J\left(E_{\lambda}\right) \subset \bigcup_{j \in \mathbb{Z}} V_{j}$.
- $E_{\lambda}$ maps each $V_{j}$ conformally onto $\mathbb{C} \backslash B_{0}$, and so $E_{\lambda}\left(V_{j}\right) \supset J\left(E_{\lambda}\right)$.

Hence, for each $j \in \mathbb{Z}$ we have a well defined inverse branch of $E_{\lambda}$ :

$$
L_{j}=L_{\lambda, j}: \mathbb{C} \backslash B_{0} \rightarrow V_{j} .
$$

Note that $B_{0}$ lies inside $V_{0}$ since $0 \in B_{0}$. The other fingers $B_{1}, \ldots, B_{n-2}$ may lie inside any of the fundamental domains $V_{j}$, depending on the value of $\lambda$. In particular, several $B_{i}$ may lie in the same $V_{j}$.
2.2. $K$-kneading sequences and $S$-kneading sequences.. We first introduce the kneading sequence given by the fundamental domains $V_{j}$. We define the $K$-kneading sequence of $\lambda \in \Omega$ to be

$$
K(\lambda)=0 k_{1} k_{2} k_{3} \ldots k_{n-2} *
$$

where $B_{j} \subset V_{k_{j}}$ for all $1 \leq j \leq n-2$. We use $*$ for the position of the point $z_{n-1}$, since this point does not belong to any of the $V_{j}$. We claim that this kneading sequence is constant throughout the entire hyperbolic component $\Omega$. To see this we first notice that the function $K: \Omega \rightarrow \Sigma$, where $\Sigma$ denotes the set of all sequences with integer terms, is locally constant. Hence, since $\Omega$ is connected, $K$ must be constant through the entire hyperbolic component. An alternative proof of this fact can be deduced from the Appendix.

We define the $K$-itinerary of any point $z \in J\left(E_{\lambda}\right)$ to be

$$
K(z)=k_{0} k_{1} k_{2} k_{3} \ldots
$$

where $E_{\lambda}^{j}(z) \in V_{k_{j}}$ for any $j \geq 0$.
One can then use these itineraries together with the kneading sequence to give a complete description of the structure of the Julia set for $E_{\lambda}$ in terms of symbolic dynamics. See [2].

We now define the $S$-kneading sequence of a value $\lambda \in \Omega$. This sequence has been independently introduced in [10], in the same context as ours. If we look at the dynamical plane very far to the right, we see that any finger is basically a straight horizontal band; therefore it makes sense to define the order of fingers in terms of their imaginary part. In this fashion, we can speak about fingers sitting above or below each other. Likewise, we can talk about the upper boundary and the lower boundary of a finger, as long as we look in the far right half plane.

Consider the half plane $H_{\mu}=\{z \in \mathbb{C} \mid \operatorname{Re} z>\mu\}$ for a fixed $\mu$ large enough. Define the family of fingers $F_{j}, j \in \mathbb{Z}$, to be the infinitely many connected components of the preimage of $B_{n-1}$. We observe that the fingers $F_{j}$ are the $2 k \pi i$-translates of the big finger for any $k \in \mathbb{Z}$. We index these sets consecutively so that $F_{0}$ is the one immediately above $B_{0}$. For any $j \in \mathbb{Z}$, let $T_{j}$ be the region in $H_{\mu}$ that lies between the upper boundaries of $F_{j-1}$ and $F_{j}$ (so, we have $F_{j} \cap H_{\mu} \subset T_{j}$ ). See Figure 4.

Finally, we define the $S$-kneading sequence of $\lambda \in \Omega$ to be

$$
S(\lambda)=0 s_{1} s_{2} s_{3} \ldots s_{n-2} *
$$

where $B_{j} \cap H_{\mu} \subset T_{s_{j}}$ for all $1 \leq j \leq n-2$. Equivalently, $B_{j}$ tends to infinity between $F_{s_{j}}$ and $F_{s_{j}-1}$. It is clear that this definition does not depend on the choice of $\mu$. See Figure 4. Moreover, since $S(\lambda)$ is again the same for all $\lambda \in \Omega$, we will use the notation $\Omega_{0 s_{1} \ldots s_{n-2} *}$ to label the hyperbolic component with such dynamical behavior with respect to the $T_{j}$ 's.

We observe that the regions $T_{i}$ do not define a family of fundamental domains in the sense described above. Consequently, the $S$-itinerary (defined in the obvious way) is not well defined for all points in the Julia set, but only for those whose orbits have sufficiently large real part. However, $E_{\lambda}$ preserves the orientation in each $T_{i}$, a feature that will prove to be useful


Fig. 4. The families $F_{i}$ and $T_{i}$
later on. The $S$-kneading sequences and itineraries are not suitable for use in the dynamical plane, but we shall see that they are very convenient when the parameter plane is considered. Therefore, it is of interest to be able to use both of these kneading sequences.
2.3. Translation algorithm.. In this section we describe an algorithm that relates the $K$ - and $S$-kneading sequences. Let us denote the $S$-kneading sequence of $E_{\lambda}$ by

$$
S=0 s_{1} \ldots s_{n-2} *
$$

We will show how to compute the $K$-kneading sequence

$$
K=0 k_{1} \ldots k_{n-2^{*}}
$$

associated to $\lambda$.
The algorithm consists of two steps. The first step is to attach a sign $(+$ or -$)$ to each of the zero entries of $S$ (with the exception of the first entry of the sequence that will remain as 0 ). This sign indicates that the corresponding $B_{i}$ is above $\left(0^{+}\right)$or below $\left(0^{-}\right) B_{0}$, at least far to the right.

The second step will determine each of the $k_{i}$ based on $s_{i}$ and $s_{i+1}$, except for the last entry $k_{n-2}$ which will be determined by $s_{n-2}$ and $s_{1}$.

Step 1: Deciding on $0^{+}$or $0^{-}$. Let $s_{i}=0$. Then $B_{i} \subset T_{0}$ and $B_{i}$ lies either above or below $B_{0}$ in the far right half plane. We will attach the superscript + or - to 0 depending on whether $B_{i}$ is above $\left(0^{+}\right)$or below $\left(0^{-}\right) B_{0}$.

To determine the sign, consider the words $s_{1} s_{2} \ldots$ and $s_{i+1} s_{i+2} \ldots$ Compare these two words until the minimal $j \geq 1$ is found such that $s_{j} \neq s_{i+j}$.

Then set

$$
s_{i}= \begin{cases}0^{+} & \text {if } s_{j}<s_{i+j} \\ 0^{-} & \text {if } s_{j}>s_{i+j}\end{cases}
$$

We write $*=\infty$ for ordering purposes in the above criterion.
We now show that this rule gives the correct superscript. Since $s_{i}=0$, $B_{i}$ meets $T_{0}$ as does $B_{0}$. If $s_{1}>s_{i+1}$ (resp. $s_{1}<s_{i+1}$ ) then $B_{i+1}$ is below (resp. above) $B_{1}$. Since the order is preserved inside each of the $T_{k}$ 's we deduce that $B_{i}$ is below (resp. above) $B_{0}$. Hence $s_{i}=0^{-}$(resp. $0^{+}$). Observe that having defined $*=\infty$ takes care of the case $s_{i+1}=*$, i.e., the case of the big finger.

Now we use induction. Let us assume $s_{j}=s_{i+j}$ for $j=1, \ldots, k$ but $s_{k+1} \neq s_{i+k+1}$. Then $B_{j}$ and $B_{i+j}$ are contained in $T_{s_{j}}, j=1, \ldots, k$, and hence, their relative order can be decided by looking at their respective images $B_{k+1}$ and $B_{i+k+1}$. There are two cases.

If $s_{k+1}>s_{i+k+1}$ then $B_{i+k+1}$ is below $B_{k+1}$, and consequently, $B_{i+j}$ is below $B_{i}$ for all $j=1, \ldots, k$. Therefore $B_{i}$ is below $B_{0}$ and so $s_{i}=0^{-}$. If $s_{k+1}<s_{i+k+1}$ we substitute "above" for "below" in the previous paragraph and conclude that $s_{i}=0^{+}$.

In particular, we remark that there are two cases that never occur: (a) $s_{i}=0^{+}$and $s_{i+1} \leq 0^{-}$in the case $s_{1} \geq 0^{+}$, and (b) $s_{i}=0^{-}$and $s_{i+1} \geq 0^{+}$in the case $s_{1} \leq 0^{-}$. More generally, by arguments similar to the above, any $0^{+}$(respectively, $0^{-}$) must be followed by entries larger than or equal to (respectively, less than or equal to) $s_{1}$.

Step 2: Obtaining $k_{i}$. Let $S$ be a signed $S$-kneading sequence obtained by replacing each 0 with the corresponding $0^{+}$and $0^{-}$symbols. There are two completely symmetric cases: $s_{1} \geq 0^{+}$and $s_{1} \leq 0^{-}$. We adopt the conventions that $1-1=0^{+}$and $-1+1=0^{-}$. Now, for any $i$ with $1 \leq i \leq n-2$,
(a) If $s_{1} \geq 0^{+}$then $\quad k_{i}= \begin{cases}s_{i} & \text { if } i=n-2 \text { or } s_{i+1} \geq 0^{+}, \\ s_{i}-1 & \text { if } s_{i+1} \leq 0^{-} .\end{cases}$
(b) If $s_{1} \leq 0^{-}$then $\quad k_{i}= \begin{cases}s_{i}+1 & \text { if } i=n-2 \text { or } s_{i+1} \geq 0^{+}, \\ s_{i} & \text { if } s_{i+1} \leq 0^{-} .\end{cases}$

We now prove that for a given $\lambda \in \Omega$ the above rule translates any signed $S$ to a unique $K$. We consider the case $s_{1} \geq 0^{+}$, the other case being symmetric.

We denote by $g_{i}$ the piece of the glove $G$ that falls into the region $T_{i}$. Since $s_{1} \geq 0^{+}, B_{1}$ is above $B_{0}$ and hence the piece of the glove $g_{0}$ must be below $B_{0}$.

This implies that $V_{0}$ is the fundamental domain between the pieces $g_{0}$ and $g_{1}$ and, in general, each $V_{i}$ lies between $g_{i}$ and $g_{i+1}$, in particular including $F_{i}$. This last remark implies that the last digit of the sequence will not change. That is, $k_{n-2}=s_{n-2}$.

Consider $s_{i}$ for $1 \leq i<n-2$. Hence $B_{i}$ lies in $T_{s_{i}}$. By the observations above, either (see Figure 5)

1. $B_{i}$ lies in $V_{s_{i}}$ because the piece $g_{s_{i}}$ lies below $B_{i}$ (case $k_{i}=s_{i}$ ), or
2. $B_{i}$ lies in $V_{s_{i}-1}$ because the piece $g_{s_{i}}$ lies above $B_{i}\left(\right.$ case $\left.k_{i}=s_{i}-1\right)$.


Fig. 5. Example of the two possibilities: the $S$-kneading sequence $02 *$ translating either into 02 * or into $01 *$

It is straightforward to check that the first case occurs if and only if $B_{i+1}$ is above $B_{0}$, i.e., $s_{i+1} \geq 0^{+}$. The second case occurs if and only if $B_{i+1}$ is below $B_{0}$, i.e., $s_{i+1} \leq 0^{-}$.

As an example, consider the $S$-kneading sequence

$$
S=0-200-12300-1200 * \text {. }
$$

After the first step we have

$$
S=0-20^{+} 0^{+}-12330^{+} 0^{+}-120^{+} 0^{+} *,
$$

and after the second step the corresponding $K$-kneading sequence is

$$
K=0-110^{+} 0^{-} 3410^{+} 0^{-} 311 * .
$$

We finally observe that the above 2-step algorithm can also be used in the reverse direction, that is, for a given $K$ with the symbols $0^{+}$and $0^{-}$we obtain, via the inverse algorithm, a unique $S$. The next section will refer to this point when taking into account the admissibility of the given sequence.
2.4. Properties.. Why are we working with two distinct kneading sequences? The answer to this question is based on the fact that the two
sequences have different properties and consequently each is suitable in different circumstances.

More precisely, the $K$-kneading sequences work well when studying the dynamical plane since they are defined using fundamental domains. These domains work for all points of the Julia set and give rise to good symbolic dynamics and consequently to a complete description of the Julia set (see [2]). In contrast, when working in the parameter plane, one can find many different hyperbolic components sharing the same $K$-kneading sequence. For instance, for any $n \in \mathbb{N}$, all hyperbolic components of period $n$ bifurcating from the main cardioid have their $K$-kneading sequence given by $K=0000 \ldots 0$. To fix this uniqueness problem we might consider the symbols $0^{+}$and $0^{-}$as before. But then an admissibility problem arises, without an obvious way to decide if a sequence is admissible or not (except, of course, by going through the inverse algorithm to check if the resulting sequence is possible).

The $S$-kneading sequences do not involve fundamental domains and hence they are not as useful as the $K$-kneading sequences when working in the dynamical plane. However, we prove in the next section that all sequences are admissible; that is, we can find a hyperbolic component $\Omega$ corresponding to any given sequence of integers. Moreover, these sequences give a significant amount of information about the location of the periodic orbit.

We remark that the uniqueness of hyperbolic components having a given $S$-kneading sequence would seem to be a natural result but it is not a straightforward deduction from the construction below.
3. Hyperbolic components. Proof of the main result. Our goal in this section is to construct a parameter value $\lambda$ for which $E_{\lambda}$ has an attracting cycle with any given $S$-kneading sequence. We first consider the special case where the $S$-kneading sequence consists of a single digit; the proof in this case makes use of many of the ideas of the general case, but in a simpler setting.
3.1. The case $0 k *$.. This result follows from the next two propositions. The first proposition can be found in [1], but we give a proof adapted to the general case.

Proposition 3.1. Fix $k \in \mathbb{Z}$. For $a \in \mathbb{R}$, let $\lambda_{a}=a+(2 k+1) \pi i$. Then, for sufficiently large values of $a$, the $\operatorname{map} E_{\lambda_{a}}$ has an attracting cycle of period 3 .

Proof. We assume throughout that $a \geq|2 k+1| \pi$, so that $\left|\operatorname{Arg}\left(\lambda_{a}\right)\right| \leq$ $\pi / 4$, where $\operatorname{Arg}$ denotes the principal branch of the argument. Then $\lambda_{a}=$ $E_{\lambda_{a}}(0)$ lies in the right half plane, but $E_{\lambda_{a}}^{2}(0)=\lambda_{a} \exp \left(\lambda_{a}\right)$ lies in the left half plane since $E_{\lambda_{a}}^{2}(0)=-e^{a} \lambda_{a}$. Choosing $a$ large enough, we may assume that $a<\left|\lambda_{a}\right| \leq a+1$. Since

$$
3 \pi / 4 \leq\left|\operatorname{Arg} E_{\lambda_{a}}^{2}(0)\right| \leq \pi
$$

it follows that

$$
\operatorname{Re} E_{\lambda_{a}}^{2}(0)=\left|\lambda_{a}\right| e^{a} \cos \left(\operatorname{Arg} E_{\lambda_{a}}^{2}(0)\right) \leq-\frac{\left|\lambda_{a}\right|}{\sqrt{2}} e^{a}<-a e^{a} / \sqrt{2}
$$

Let $U_{2}$ be the ball of radius 1 about $E_{\lambda_{a}}^{2}(0)$. The preimage of $U_{2}$ containing $\lambda_{a}$ is an open set $U_{1}$ which is mapped univalently onto $U_{2}$ by $E_{\lambda_{a}}$, and the preimage of $U_{1}$ containing 0 is another open set, say $U_{0}$, which is mapped univalently onto $U_{2}$ by $E_{\lambda_{a}}^{2}$. We claim that there is an attracting cycle of period 3 whose orbit under $E_{\lambda_{a}}$ lies in $U_{0}, U_{1}$, and $U_{2}$. Let $F$ denote the appropriate branch of the inverse of $E_{\lambda_{a}}^{2}$ that takes $U_{2}$ univalently into $U_{0}$. See Figure 6. Define the absolute constant $C_{a}=e^{-a} /\left(4(a+1)^{2}\right)$. By the Koebe 1/4 Theorem, we have

$$
\operatorname{dist}\left(0, \partial U_{0}\right) \geq \frac{1}{4}\left|F^{\prime}\left(E_{\lambda_{a}}^{2}(0)\right)\right|=\frac{1}{4} \cdot\left|\frac{1}{\lambda_{a}}\right| \cdot\left|\frac{1}{\lambda_{a} e^{\lambda_{a}}}\right|=\frac{e^{-a}}{4\left|\lambda_{a}\right|^{2}} \geq C_{a}
$$



Fig. 6. The sets $U_{0}, U_{1}$ and $U_{2}$ in the proof of Proposition 3.1
Now

$$
\left|E_{\lambda_{a}}^{3}(0)\right|=\left|\lambda_{a}\right| \exp \left(\operatorname{Re} E_{\lambda_{a}}^{2}(0)\right) \leq(a+1) \exp \left(-a e^{a} / \sqrt{2}\right) \ll C_{a}
$$

for large $a$. Hence $E_{\lambda_{a}}^{3}(0)$ is contained in $U_{0}$. Moreover, if $w \in U_{2}$, then

$$
\begin{aligned}
\left|E_{\lambda_{a}}(w)-E_{\lambda_{a}}^{3}(0)\right| & \leq \max _{z \in U_{2}}\left|E_{\lambda_{a}}^{\prime}(z)\right| \leq\left|\lambda_{a} \exp \left(\operatorname{Re} E_{\lambda_{a}}^{2}(0)+1\right)\right| \\
& \leq(a+1) e \exp \left(-a e^{a} / \sqrt{2}\right) \ll C_{a}
\end{aligned}
$$

as before. Hence,

$$
\operatorname{dist}\left(0, \partial E_{\lambda_{a}}^{3}\left(U_{0}\right)\right) \leq(a+1)(e+1) \exp \left(-a e^{a} / \sqrt{2}\right) \ll C_{a}
$$

and it follows that $E_{\lambda_{a}}^{3}\left(U_{0}\right)$ is properly contained in $U_{0}$. Thus we have an attracting cycle whose orbit visits $U_{0}, U_{1}$ and $U_{2}$.

Before proceeding, we observe that the above estimates guarantee that the entire half plane $\operatorname{Re} z \leq \operatorname{Re} E_{\lambda_{a}}^{2}(0)+1$ is contained in the basin of the cycle.

We now claim that the $S$-kneading sequence of $\lambda_{a}$ is $0 k *$.
Proposition 3.2. Let $k \in \mathbb{Z}$ and set $\lambda_{a}=a+(2 k+1) \pi i$. Then for sufficiently large $a, E_{\lambda_{a}}$ has an attracting 3 -cycle with $S\left(\lambda_{a}\right)=0 k *$.

Proof. Let $\gamma(t)=t+(2 k+1) \pi i$ for $t \geq a$. Then $E_{\lambda_{a}}(\gamma(t))$ is a straight line which lies to the left of $E_{\lambda_{a}}^{2}(0)$. By the above observation, $E_{\lambda_{a}}(\gamma(t))$ lies in the connected component of the immediate basin of attraction which contains $E_{\lambda_{a}}^{2}(0)$. Hence $\gamma(t)$ lies in the component of the immediate basin which contains $\lambda_{a}$.

Let $S$ be the strip $\{z||\operatorname{Im} z| \leq \pi\}$. There is a preimage of $\gamma(t)$ contained in the interior of $S$, at least for $t$ large. We claim that the entire preimage of $\gamma(t)$ lies in $S$. The preimage of $\gamma(t)$ can never meet the boundary of $S$, for $E_{\lambda_{a}}$ maps the boundary of $S$ into the left half plane, far from $\gamma(t)$. Hence the preimage of $\gamma(t)$ lying in $S$ must be the preimage that contains 0 .

We then consider the set $B$ as above so that $B$ contains $E_{\lambda_{a}}^{3}(0)$. It then follows that $B_{2}$ contains $E_{\lambda_{a}}^{2}(0)$ and $E_{\lambda_{a}}(\gamma(t))$. After taking one more preimage, the big finger $B_{1}$ contains $\lambda_{a}$ and $\gamma(t)$, and its translations contain the half lines $\{t+(2 j+1) \pi \mid t \geq a\}$. Moreover, the finger $B_{0}$ contains 0 and the preimage of $\gamma(t)$ in $S$. It follows then that the fingers are indexed so that $B_{1}=F_{k}$ and hence $S\left(\lambda_{a}\right)=0 k *$.
3.2. The general case.. Now we proceed to the general case. For the remainder of this section we fix a kneading sequence $s=0 s_{1} \ldots s_{n-2} *$. Let $\widehat{s}=\max \left|s_{i}\right|$ and define $M=(2 \widehat{s}+1) \pi$. We assume throughout that $a>M$. Let $H(a)$ denote the closed half strip

$$
H(a)=\{z|\operatorname{Re} z \geq a,|\operatorname{Im}(z)| \leq M\}
$$

We let $L(a)$ denote the left boundary of $H(a)$. We will prove:
Theorem 3.3. For each sufficiently large $a$, there is $\lambda_{a} \in L(a)$ for which $E_{\lambda_{a}}$ has an attracting n-cycle with $S\left(\lambda_{a}\right)=s$.

If we denote the first $n$ points on the orbit of 0 by $w_{i}$, so $w_{0}=0, w_{1}=\lambda_{a}$, $\ldots, w_{n}=E_{\lambda_{a}}^{n}(0)$, as in the previous special case, we will construct $\lambda_{a}$ so that the orbit of 0 under $E_{\lambda_{a}}$ has the following properties:

1. $w_{i} \in H(a)$ for $i=1, \ldots, n-2$ and $\operatorname{Re} w_{i+1} \gg \operatorname{Re} w_{i}$ for $i=0, \ldots, n-3$.
2. $w_{n-1}$ lies in the left half plane and

$$
\left|\operatorname{Re} w_{n-1}\right| \gg \operatorname{Re} w_{n-2}
$$

3. $w_{n}$ lies close to 0 and, as in the period 3 case, there is an attracting cycle of period $n$ lying close to $w_{0}, \ldots, w_{n-1}$.

We will divide the proof into three parts, namely Propositions 3.5-3.7 (Proposition 3.4 is auxiliary). Afterwards we will see how Theorem A (see Section 1) follows.

Let $\nu=\nu(a)=|a+(2 \widehat{s}+1) \pi i|=\max _{z \in L(a)}|z|$, and note that $\nu(a)-a \rightarrow 0$ as $a \rightarrow \infty$.

For $-\widehat{s} \leq i \leq \widehat{s}$, let $H_{i}(a)$ be the substrip of $H(a)$ given by

$$
H_{i}(a)=\{z \in H(a) \mid \operatorname{Re} z \geq a,(2 i-1) \pi \leq \operatorname{Im} z \leq(2 i+1) \pi\}
$$

See Figure 7.


Fig. 7. The sets $H(a), L(a)$ and the substrips $H_{i}(a)$ for the case $\widehat{s}=3$
For $j=1, \ldots, n-2$, define $w_{j}(\lambda)=E_{\lambda}^{j}(0)$. Note that each $w_{j}$ is a function of the parameter $\lambda$ and is analytic. For example, $w_{1}(\lambda)=\lambda$ and $w_{2}(\lambda)=\lambda e^{\lambda}$.

For $j=1, \ldots, n-2$, define

$$
I_{s_{1} \ldots s_{j}}(a)=\left\{\lambda \in L(a) \mid w_{i}(\lambda) \in H_{s_{i}}(a) \text { for } i=1, \ldots, j\right\} .
$$

Note that $I_{s_{1}}(a)=L(a) \cap H_{s_{1}}(a)$ and that the $I_{s_{1} \ldots s_{j}}$ are nested, provided they are nonempty. The following proposition shows that each of the $I_{s_{1} \ldots s_{j}}$ consists of a single vertical segment.

We say that a smooth curve $\mu(t)$ in $H_{s_{i}}(a)$ is a vertical curve if the curve connects the upper and lower boundaries of $H_{s_{i}}(a)$.

Proposition 3.4. For each sufficiently large a and $1 \leq j \leq n-2$, the set $\left\{w_{j}(\lambda) \mid \lambda \in I_{s_{1} \ldots s_{j}}(a)\right\}$ consists of a single vertical curve in $H_{s_{j}}(a)$. Hence $I_{s_{1} \ldots s_{j}}$ is a single vertical segment.

Proof. If $j=1$, there is nothing to prove since $w_{1}(\lambda)=\lambda$ and $I_{s_{1}}(a)=$ $L(a) \cap H_{s_{1}}(a)$. Let $j>1$. We parametrize the segment $I_{s_{1}}(a)$ as $\lambda(t)=$
$a+\left(2 s_{1} \pi+t\right) i$ for $t \in(-\pi, \pi)$ and consider the set

$$
J_{s_{1} s_{2}}(a)=\left\{\lambda \in I_{s_{1}}(a) \mid w_{2}(t) \subset H(a)\right\}
$$

where $w_{2}(t)=w_{2}(\lambda(t))=\lambda(t) e^{\lambda(t)}$. We will show that, given any $\varepsilon>0$ and taking $a$ large enough,
$\left|\operatorname{Arg} w_{2}^{\prime}(t)-\pi / 2\right|<\varepsilon$
for any $t$ such that $\lambda(t) \in J_{s_{1} s_{2}}(a)$. This implies that, when $t$ runs from $-\pi$ to $\pi$, as the curve $w_{2}(t)$ crosses the strip $H(a)$, its tangent vector points upwards and it is almost vertical. It follows that the imaginary part of $w_{2}(t)$ is an increasing function of $t$, and hence the curve crosses the strip only once. We now proceed to show (1).

Set $a_{0}$ large enough so that

$$
|\operatorname{Arg} \lambda(t)|<\frac{\varepsilon}{3(n-3)}=\varepsilon^{\prime}
$$

for all $t \in(-\pi, \pi)$.
The tangent vector to $w_{2}(t)$ is

$$
w_{2}^{\prime}(t)=\lambda^{\prime}(t) e^{\lambda(t)}(1+\lambda(t))=i e^{\lambda(t)}(1+\lambda(t))
$$

and thus

$$
\left|\operatorname{Arg} w_{2}^{\prime}(t)-\pi / 2\right|=\left|\operatorname{Arg} e^{\lambda(t)}+\operatorname{Arg}(1+\lambda(t))\right| \leq\left|\operatorname{Arg} e^{\lambda(t)}\right|+\varepsilon^{\prime}
$$

If $\lambda(t) \in J_{s_{1} s_{2}}$, it is clear that $\left|w_{2}(t)\right|>|\lambda(t)|$. Since both are inside the strip $H(a)$, we have $\varepsilon^{\prime}>\left|\operatorname{Arg} w_{2}(t)\right|=\left|\operatorname{Arg} \lambda(t)+\operatorname{Arg} e^{\lambda(t)}\right|$. It is then easy to see that $\left|\operatorname{Arg} e^{\lambda(t)}\right|<2 \varepsilon^{\prime}$. Plugging this in the expression above, we obtain

$$
\left|\operatorname{Arg} w_{2}^{\prime}(t)-\pi / 2\right|<3 \varepsilon^{\prime}=\frac{\varepsilon}{n-3}<\varepsilon
$$

as required.
We now proceed to investigate $I_{s_{1} s_{2} s_{3}}$ which will illustrate the general case. As above, consider

$$
J_{s_{1} s_{2} s_{3}}(a)=\left\{\lambda \in I_{s_{1} s_{2}}(a) \mid w_{3}(t) \subset H(a)\right\}
$$

where $w_{3}(t)=w_{3}(\lambda(t))=\lambda(t) e^{w_{2}(t)}$ and $w_{2}(t) \in H_{s_{2}}(a)$. For these values of $t$, we will show that

$$
\begin{equation*}
\left|\operatorname{Arg} w_{3}^{\prime}(t)-\pi / 2\right|<\varepsilon \tag{2}
\end{equation*}
$$

Now the tangent vector to $w_{3}(t)$ is

$$
w_{3}^{\prime}(t)=e^{w_{2}(t)}\left(\lambda^{\prime}(t)+\lambda(t) w_{2}^{\prime}(t)\right)=e^{w_{2}(t)}\left(i+\lambda(t) w_{2}^{\prime}(t)\right)
$$

and thus

$$
\operatorname{Arg} w_{3}^{\prime}(t)=\operatorname{Arg} e^{w_{2}(t)}+\operatorname{Arg}\left(i+\lambda(t) w_{2}^{\prime}(t)\right)
$$

We claim that

$$
\pi / 2-4 \varepsilon^{\prime}<\operatorname{Arg}\left(i+\lambda(t) w_{2}^{\prime}(t)\right)<\pi / 2+4 \varepsilon^{\prime}
$$

Indeed, we showed above that

$$
\pi / 2-3 \varepsilon^{\prime}<\operatorname{Arg} w_{2}^{\prime}(t)<\pi / 2+3 \varepsilon^{\prime}
$$

Moreover, since $|\operatorname{Arg} \lambda(t)|<\varepsilon^{\prime}$, we obtain

$$
\pi / 2-4 \varepsilon^{\prime}<\operatorname{Arg}\left(\lambda(t) w_{2}^{\prime}(t)\right)<\pi / 2+4 \varepsilon^{\prime}
$$

Finally, it remains to add the vector $i$ to this expression, which makes the argument even closer to $\pi / 2$.

To finish the proof of (2) observe that, by the same argument as in the first case, $\operatorname{Arg} w_{3}(t)=\operatorname{Arg}\left(\lambda(t) e^{w_{2}(t)}\right)<\varepsilon^{\prime}$ and hence $\left|\operatorname{Arg} e^{w_{2}(t)}\right|<2 \varepsilon^{\prime}$. Putting all this together we have

$$
\pi / 2-6 \varepsilon^{\prime}<\operatorname{Arg} w_{3}^{\prime}(t)<\pi / 2+6 \varepsilon^{\prime}
$$

as we wanted to prove.
It is easy to check that we may iterate this procedure to deduce that, for $j=2, \ldots, n-2$ and for all $t$ such that $\lambda(t) \in J_{s_{1} \ldots s_{j}}(a)$,

$$
\left|\operatorname{Arg} w_{j}^{\prime}(t)-\pi / 2\right|<3(j-1) \varepsilon^{\prime}=(j-1) \frac{\varepsilon}{n-3} \leq \varepsilon
$$

which concludes the proof of the proposition. -
Proposition 3.5. Let $\varepsilon>0$. For each sufficiently large $a$, there is $\lambda_{a} \in$ $L(a)$ satisfying

1. $w_{i}\left(\lambda_{a}\right) \in H_{s_{i}}(a)$ for $i=1, \ldots, n-2$.
2. $\operatorname{Im} w_{n-2}\left(\lambda_{a}\right)=\left(2 s_{n-2}+1\right) \pi$.
3. $E_{a-\varepsilon}^{j-1}(a-\varepsilon) \leq \operatorname{Re} w_{j}\left(\lambda_{a}\right) \leq\left|w_{j}\left(\lambda_{a}\right)\right| \leq E_{a+\varepsilon}^{j-1}(a+\varepsilon)$ for $j=1, \ldots, n-2$, where $E_{b}$ is the real exponential $E_{b}(x)=b e^{x}$.

Proof. By the previous proposition, if $\lambda \in I_{s_{1} \ldots s_{j}}(a)$, then the curve $\lambda \mapsto w_{j}(\lambda)$ is a vertical curve in $H_{s_{j}}(a)$. We will show that, moreover,

$$
E_{a-\varepsilon}^{j-1}(a-\varepsilon) \leq \operatorname{Re} w_{j}(\lambda) \leq E_{a+\varepsilon}^{j-1}(a+\varepsilon)
$$

for each $j$. Then $\lambda_{a}$ will be defined as the upper endpoint of $I_{s_{1} \ldots s_{n-2}}(a)$.
If $\lambda \in V_{s_{1}}(a)$, then $\exp (\lambda)$ lies on a circle of radius $e^{a}$ centered at 0 . Hence $\lambda \mapsto w_{2}(\lambda)=\lambda e^{\lambda}$ is a nearly circular arc contained in the annulus

$$
\begin{equation*}
E_{a}(a) \leq|z| \leq E_{\nu}(\nu) \tag{3}
\end{equation*}
$$

where we recall that $\nu=\max _{z \in L(a)}|z|$. This arc crosses $H_{s_{2}}(a)$ in a single vertical curve $\eta_{2}$, provided $a$ is sufficiently large.

Given $\varepsilon>0$, we claim we may choose $a$ large enough so that if $\lambda \in$ $I_{s_{1} s_{2}}(a)$ then

$$
\begin{equation*}
E_{a-\varepsilon}(a-\varepsilon) \leq \operatorname{Re} w_{2}(\lambda) \leq\left|w_{2}(\lambda)\right| \leq E_{a+\varepsilon}(a+\varepsilon) \tag{4}
\end{equation*}
$$

Indeed, both estimates are deduced from (3). The lower estimate holds since the circle of radius $E_{a}(a)$ meets $H(a)$ in a nearly vertical arc. The upper
estimate follows since $\nu(a)-a \rightarrow 0$ as $a \rightarrow \infty$ and hence we may choose $a$ so that $\nu<a+\varepsilon$.

Now we exponentiate points on $\eta_{2}$. The result is a curve whose endpoints lie in $\mathbb{R}^{-}$. Multiplication of this curve by the appropriate $\lambda \in I_{s_{1} s_{2}}(a)$ expands this curve, but the image must cross $H_{s_{3}}(a)$ in a single vertical curve which we denote by $\eta_{3}$.

As above, we claim that by choosing $a$ large enough we have, for $\lambda \in$ $I_{s_{1} s_{2} s_{3}}(a)$,

$$
\begin{equation*}
E_{a-\varepsilon}^{2}(a-\varepsilon) \leq \operatorname{Re} w_{3}(\lambda) \leq\left|w_{3}(\lambda)\right| \leq E_{a+\varepsilon}^{2}(a+\varepsilon) \tag{5}
\end{equation*}
$$

The upper estimate holds since

$$
\left|w_{3}(\lambda)\right|=|\lambda| \exp \left(\operatorname{Re} w_{2}(\lambda)\right) \leq \nu \exp \left(E_{a+\varepsilon}(a+\varepsilon)\right) \leq E_{a+\varepsilon}^{2}(a+\varepsilon)
$$

To obtain the lower estimate, first set $R_{a, \varepsilon}=a \exp \left(E_{a-\varepsilon}(a-\varepsilon)\right)$ and observe that, by (4),

$$
\left|w_{3}(\lambda)\right|=|\lambda| e^{\operatorname{Re} w_{2}(\lambda)} \geq R_{a, \varepsilon}
$$

By a simple trigonometric argument (see Figure 8) one can see that

$$
\begin{equation*}
\operatorname{Re} w_{3}(\lambda) \geq \sqrt{R_{a, \varepsilon}-M^{2}} \tag{6}
\end{equation*}
$$



Fig. 8. Sketch of the construction above
We then have, on the one hand,

$$
R_{a, \varepsilon}-\sqrt{R_{a, \varepsilon}-M^{2}} \underset{a \rightarrow \infty}{\longrightarrow} 0
$$

and, on the other hand,

$$
R_{a, \varepsilon}-E_{a-\varepsilon}^{2}(a-\varepsilon)=\varepsilon \exp \left(E_{a-\varepsilon}(a-\varepsilon)\right) \underset{a \rightarrow \infty}{\longrightarrow} \infty
$$

Putting everything together, we obtain the lower estimate in (5).
It is now clear that continuing in the same fashion we obtain the required $I_{s_{1} \ldots s_{j}}(a)$. Note that, by construction, if $\lambda$ is the upper endpoint of $I_{s_{1} \ldots s_{j}}(a)$, then $z_{j}(\lambda) \in \partial H_{s_{j}}(a)$. Hence, pick $\lambda$ to be the upper endpoint of $I_{s_{1} \ldots s_{n-2}}(a)$ and then $\operatorname{Im} w_{n-2}\left(\lambda_{a}\right)=\left(2 s_{n-2}+1\right) \pi$.

Proposition 3.6. Choose $\lambda_{a}$ as in Proposition 3.5. Then $E_{\lambda_{a}}$ has an attracting cycle of period $n$.

Proof. By the same arguments as in Proposition 3.5, it is clear that

$$
E_{a-\varepsilon}^{n-2}(a-\varepsilon) \leq\left|w_{n-1}(\lambda)\right| \leq E_{a+\varepsilon}^{n-2}(a+\varepsilon) .
$$

We know that $\operatorname{Im} w_{n-2}\left(\lambda_{a}\right)=\left(2 s_{n-2}+1\right) \pi$, and hence

$$
\operatorname{Re} w_{n-1}\left(\lambda_{a}\right) \leq-E_{a-\varepsilon}^{n-2}(a-\varepsilon) \cos \left(\operatorname{Arg} \lambda_{a}\right)
$$

since $\operatorname{Arg} w_{n-1} \lambda_{a}=\operatorname{Arg} \lambda_{a}+\pi$. Now $\left|\operatorname{Arg} \lambda_{z}\right| \leq \pi / 4$ so that

$$
\operatorname{Re} w_{n-1}\left(\lambda_{a}\right) \leq-\left(E_{a}^{n-2}(a)-1\right) / \sqrt{2} .
$$

Let $B$ be an open ball of radius 1 about $w_{n-1}\left(\lambda_{a}\right)$. The preimages of $B$ containing $w_{j}\left(\lambda_{a}\right)$ for $j=1, \ldots, n-2$ are open sets, and $E_{\lambda_{a}}^{n-1-j}$ maps them univalently onto $B$. Let $U$ be the preimage of $B$ containing 0 . Then $E_{\lambda_{a}}^{n-1}$ maps $U$ univalently onto $B$.

Let $F: B \rightarrow U$ denote the appropriate branch of the inverse of $E_{\lambda_{a}}^{n-1}$ taking $w_{n-1}\left(\lambda_{a}\right)$ to 0 . We have

$$
\begin{aligned}
\left|F^{\prime}\left(w_{n-1}\left(\lambda_{a}\right)\right)\right| & =\left|\frac{1}{\prod_{j=0}^{n-2} E_{\lambda_{a}}^{\prime}\left(w_{j}\left(\lambda_{a}\right)\right)}\right|=\frac{1}{\prod_{j=1}^{n-1}\left|w_{j}\left(\lambda_{a}\right)\right|} \\
& \geq \frac{1}{\prod_{j=1}^{n-1} E_{a+\varepsilon}^{j-1}(a+\varepsilon)}
\end{aligned}
$$

by Proposition 3.5. By the Koebe $1 / 4$ Theorem we have

$$
\operatorname{dist}(0, \partial U) \geq \frac{1}{4}\left|F^{\prime}\left(w_{n-1}\left(\lambda_{a}\right)\right)\right| \geq \frac{1}{4} \frac{1}{\prod_{j=1}^{n-1} E_{a+\varepsilon}^{j-1}(a+\varepsilon)} .
$$

Now consider $w_{n}\left(\lambda_{a}\right)$. We have

$$
\begin{aligned}
\left|w_{n}\left(\lambda_{a}\right)\right| & =\left|E_{\lambda_{a}}\left(w_{n-1}\left(\lambda_{a}\right)\right)\right|=\left|\lambda_{a}\right| \exp \left(\operatorname{Re} w_{n-1}\left(\lambda_{a}\right)\right) \\
& \leq(a+\varepsilon) \exp \left(-\frac{1}{\sqrt{2}} E_{a-\varepsilon}^{n-2}(a-\varepsilon)\right) \ll \frac{1}{4} \frac{1}{\prod_{j=0}^{n-1}\left(E_{a+\varepsilon}^{j-1}(a+\varepsilon)\right)} .
\end{aligned}
$$

The last inequality follows (for $a$ large enough and $\varepsilon$ small enough) since the expression for $\left|E_{\lambda_{a}}\left(w_{n-1}\left(\lambda_{a}\right)\right)\right|$ contains one higher iterate of $E_{a}$. Hence $w_{n}\left(\lambda_{a}\right)$ lies well within $U$. We claim that $E_{\lambda_{a}}(B) \subset U$ as well. Indeed, for $w \in B$, we have

$$
\begin{aligned}
\left|E_{\lambda_{a}}^{\prime}(w)\right| & \leq\left|E_{\lambda_{a}}^{\prime}\left(w_{n-1}\left(\lambda_{a}\right)+1\right)\right|=\left|\lambda_{a}\right| \exp \left(\operatorname{Re} w_{n-1}\left(\lambda_{a}\right)+1\right) \\
& \leq(a+\varepsilon) \exp \left(-\frac{1}{\sqrt{2}} E_{a-\varepsilon}^{n-2}(a-\varepsilon)+1\right) \ll \frac{1}{4} \frac{1}{\prod_{j=1}^{n-1} E_{a+\varepsilon}^{j-1}(a+\varepsilon)}
\end{aligned}
$$

as above. This shows that $E_{\lambda_{a}}(B)$ lies well within $U$ since

$$
\left|E_{\lambda_{a}}(w)-E_{\lambda_{a}}\left(w_{n-1}\left(\lambda_{a}\right)\right)\right| \leq \max _{w \in B}\left|E_{\lambda}^{\prime}(w)\right| .
$$

It follows that $E_{\lambda_{a}}$ has an attracting cycle of period $n$ that lies close to $w_{j}\left(\lambda_{a}\right)$ for $j=0, \ldots, n-1$.

The following proposition completes the proof of Theorem 3.3.
Proposition 3.7. Choose $\lambda_{a}$ as in Proposition 3.5. Then $S\left(\lambda_{a}\right)=$ $0 s_{1} \ldots s_{n-2} *$.

Proof. Let $\gamma(t)=t+\left(2 s_{n-2}+1\right) \pi i$ with $t \geq \operatorname{Re} w_{n-2}\left(\lambda_{a}\right)$ so $w_{n-2}\left(\lambda_{a}\right)$ is the left endpoint of this horizontal line. We claim that $\gamma(t)$ belongs to the basin of attraction of the attracting cycle. Indeed, $E_{\lambda_{a}}(\gamma(t))$ is a straight line lying to the left of $w_{n-1}\left(\lambda_{a}\right)$. Hence $\left|E_{\lambda_{a}}^{2}(\gamma(t))\right| \leq\left|w_{n}\left(\lambda_{a}\right)\right|$ and it follows that this line lies in the immediate basin containing $w_{n-1}\left(\lambda_{a}\right)$.

For any $\varepsilon>0$ we let $\tau=\varepsilon / n$. Then for $a$ sufficiently large we have $\left|\operatorname{Arg} w_{j}\left(\lambda_{a}\right)\right| \leq \tau$ for $j=1, \ldots, n-2$. This follows since $\left|\operatorname{Arg} w_{j}\left(\lambda_{a}\right)\right| \leq$ $|\operatorname{Arg}(a+(2 \widehat{s}+1) \pi i)|$, which may be made arbitrarily small as $a$ increases.

Now let $\mu_{j}(t)$ denote the curve that contains $w_{n-2-j}\left(\lambda_{a}\right)$ and satisfies $E_{\lambda_{a}}^{j}\left(\mu_{j}(t)\right)=\gamma(t)$ for $t \geq \operatorname{Re} w_{n-2}\left(\lambda_{a}\right)$ and $j=1, \ldots, n-2$. So $\mu_{1}(t)$ contains $w_{n-3}\left(\lambda_{a}\right)$ while $\mu_{n-2}(t)$ contains 0 . By construction, each $\mu_{j}$ is in a different component of the immediate basin of the attracting cycle. To prove the result, we will show that $\mu_{j}(t) \subset H_{s_{n-2-j}}(a)$ for each $j \leq n-3$ and $\left|\operatorname{Im} \mu_{n-2}(t)\right|<\pi$.

Consider $\mu_{1}(t)$. We have $E_{\lambda_{a}}\left(\mu_{1}(t)\right)=\gamma(t)$ so that

$$
E_{\lambda_{a}}^{\prime}\left(\mu_{1}(t)\right) \cdot \mu_{1}^{\prime}(t)=\gamma^{\prime}(t)
$$

Therefore

$$
\operatorname{Arg} E_{\lambda_{a}}^{\prime}\left(\mu_{1}(t)\right)+\operatorname{Arg} \mu_{1}^{\prime}(t)=\operatorname{Arg} \gamma^{\prime}(t)=0
$$

and consequently

$$
\left|\operatorname{Arg} \mu_{1}^{\prime}(t)\right|=\left|\operatorname{Arg} E_{\lambda_{a}}^{\prime}\left(\mu_{1}(t)\right)\right|=\left|\operatorname{Arg} E_{\lambda_{a}}\left(\mu_{1}(t)\right)\right|=|\operatorname{Arg} \gamma(t)| \leq \tau
$$

In particular, this implies $\mu_{1}(t)$ lies to the right of its endpoint, $w_{n-3}\left(\lambda_{a}\right)$, for $t>\operatorname{Re} w_{n-2}\left(\lambda_{a}\right)$.

Continuing inductively, we find that

$$
\left|\operatorname{Arg} \mu_{j}^{\prime}(t)\right| \leq \tau j
$$

so that $\left|\operatorname{Arg} \mu_{j}^{\prime}(t)\right| \leq \varepsilon$ for all $j$, and that each $\mu_{j}(t)$ lies to the right of its endpoint, $w_{n-z-j}\left(\lambda_{a}\right)$.

Now suppose that $\operatorname{Im} \mu_{j}\left(t_{0}\right)=(2 k+1) \pi$ for some $k \in \mathbb{Z}$. It follows that $E_{\lambda_{a}}\left(\mu_{j}\left(t_{0}\right)\right)$ lies in the left half plane. But $E_{\lambda_{a}}\left(\mu_{j}(t)\right)=\mu_{j-1}(t)$ if $j>1$ and $E_{\lambda_{a}}\left(\mu_{1}(t)\right)=\gamma(t)$. This contradicts the fact that $\mu_{j-1}\left(t_{0}\right)$ lies to the right of the endpoint of $\mu_{j-1}$. Hence each $\mu_{j}$ must lie in a horizontal strip of width at most $2 \pi$ and contained between the translates of $\gamma(t)$. This implies that $\mu_{j}(t) \subset H_{s_{n-2-j}}(a)$, and the result follows.

This concludes the proof of Theorem 3.3. To complete the proof of Theorem A, observe that the result holds for any $a$ larger than a certain value $a_{0}$. Following the construction, we then see that we have constructed a curve
of $\lambda_{a}$-values, one for each sufficiently large $a \in \mathbb{R}$, having the property that $\operatorname{Re} \lambda_{a}=a$ and $S(\lambda)=s$. Note that $\lambda_{a}$ lies in the intervals $I_{s_{1} \ldots s_{n-2}}(a)$ and, by construction, we have $\operatorname{Im} I_{s_{1} \ldots s_{n-3} \alpha}(a)<\operatorname{Im} I_{s_{1} \ldots s_{n-3} \beta}(a)$ if and only if $\alpha<\beta$. Thus, the hyperbolic components of the same period are ordered lexicographically. The following corollary shows how the components of period $n+1$ may be inserted between the components of period $n$.

Corollary 3.8. Let $\lambda_{a}$ and $\widetilde{\lambda}_{a}$ have kneading sequences $0 s_{1} \ldots s_{n-2}$ * and $0 s_{1} \ldots\left(s_{n-2}+1\right) *$ for a sufficiently large. Then, given any $k \in \mathbb{Z}$, there exists $\lambda_{a}(k)$ with $\operatorname{Re} \lambda_{a}(k)=a$ and $S\left(\lambda_{a}(k)\right)=0 s_{1} \ldots\left(s_{n-2}+1\right) k *$.

Proof. By construction, the $\lambda$-values in the vertical segment in between $\lambda_{a}$ and $\lambda_{a}$ are exactly those belonging to $I_{s_{1} \ldots\left(s_{n-2}+1\right)}(a)$. Hence, if we iterate the process one step further to obtain $\lambda_{a}(k)$ such that $S\left(\lambda_{a}(k)\right)=$ $0 s_{1} \ldots\left(s_{n-2}+1\right) k *$, we must iterate once more for values of $\lambda$ in this segment. Hence each of the $\lambda(k)$ belongs to $I_{s_{1} \ldots\left(s_{n-2}+1\right)}(a)$.

Appendix. Let $E_{\lambda}, \Omega, n, z_{0}(\lambda), \ldots, z_{n-1}(\lambda), A^{*}$ and $A^{*}\left(z_{i}\right)$ be as in Section 2.1. Our goal in this section is to show how the disc $B_{\lambda}$ in the construction of the kneading sequence may be defined for any $\lambda \in \Omega$ so that it varies holomorphically with respect to $\lambda$. Although this is not crucial in this paper, we believe it is interesting in itself.

More precisely, our goal is to prove the following proposition.
Proposition A.1. For any $\lambda \in \Omega$, there exists a topological disc $B_{\lambda}$ such that
(a) $\partial B_{\lambda}$ is a simple closed curve in $\mathbb{C}$;
(b) $0, z_{0} \in B_{\lambda}$;
(c) $\overline{E_{\lambda}^{n}\left(B_{\lambda}\right)} \subset B_{\lambda}$;
(d) $B_{\lambda} \subset A^{*}\left(z_{0}\right)$;
(e) $B_{\lambda}$ depends holomorphically on the parameter $\lambda$. More precisely, the boundary of $B_{\lambda}$ is defined by a map

$$
\gamma: \Omega \times \mathbb{T} \rightarrow \partial\left(B_{\lambda}\right) \subset \mathbb{C}, \quad(\lambda, t) \mapsto \gamma(\lambda, t)
$$

satisfying
(1) for a fixed $t \in \mathbb{T}$, the map $\lambda \mapsto \gamma_{\lambda}(t)=\gamma(\lambda, t)$ from $\Omega$ to $A^{*}\left(z_{0}\right)$ is holomorphic;
(2) for a fixed $\lambda \in \Omega$, the map $t \mapsto \gamma_{\lambda}(t)$ is an injection (and hence $\gamma_{\lambda}$ is a simple closed curve).

Remark A.2. Observe that conditions (1) and (2) imply that the map $\gamma$ (appropriately rewritten after choosing a basepoint $\lambda_{0} \in \Omega$ ) defines a holomorphic motion (see [9]) of $\partial B_{\lambda_{0}}$. Then we can deduce from the $\lambda$-lemma that the map $\gamma$ is jointly continuous.

To prove Proposition A. 1 we shall study how the linearizing coordinates of the attracting cycle behave in its basin of attraction, and use them to define precisely the boundary of the set $B_{\lambda}$. From the construction it will become clear how this curve depends on the parameter.

For any $\lambda \in \Omega$, let $\varrho_{\lambda}=\left(E_{\lambda}^{n}\right)^{\prime}\left(z_{0}(\lambda)\right)$ be the multiplier of the periodic orbit, and let $\phi_{\lambda}$ be the linearizing coordinates defined on a neighborhood $U_{\lambda}$ of $z_{0}$ and conjugating $E_{\lambda}^{n}$ to multiplication by $\varrho_{\lambda}$. That is,

$$
\begin{equation*}
\phi_{\lambda} \circ E_{\lambda}^{n}(z)=\varrho_{\lambda} \cdot \phi_{\lambda}(z) \tag{7}
\end{equation*}
$$

Lemma A.3. The linearizing coordinates $\phi_{\lambda}$ and the neighborhood $U_{\lambda}$ may be chosen such that $0 \in \partial U_{\lambda}, \phi_{\lambda}\left(U_{\lambda}\right)=\mathbb{D}$, and $\phi_{\lambda}(0)=1$.

Proof. Let $\widetilde{\phi}_{\lambda}$ be some linearizing coordinates on a neighborhood $\widetilde{U}_{\lambda}$, i.e., satisfying

$$
\begin{equation*}
\widetilde{\phi}_{\lambda} \circ E_{\lambda}^{n}(z)=\varrho_{\lambda} \cdot \widetilde{\phi}_{\lambda}(z) \tag{8}
\end{equation*}
$$

for $\underset{\sim}{\text { all }} z \in \widetilde{U}_{\lambda}$. We may assume we have restricted $\widetilde{U}_{\lambda}$ so that it is mapped by $\widetilde{\phi}_{\lambda}$ to a round disc centered at 0 .

By construction we know that $0 \in A^{*}\left(z_{0}\right)$. If $0 \in \widetilde{U}_{\lambda}$, we are done by further restricting $\widetilde{U}_{\lambda}$ and composing with a rotation. So we assume this is not the case. Hence there exists $p \in \mathbb{N}$ such that $E_{\lambda}^{n p}(0) \in \widetilde{U}_{\lambda}$. Let $\widetilde{\omega}=\widetilde{\phi}_{\lambda}\left(E_{\lambda}^{n p}(0)\right)$. Then the preimage of the $\operatorname{disc} D(0,|\widetilde{\omega}|)$ under $\widetilde{\phi}_{\lambda}$ is a neighborhood of $z_{0}$ contained inside $\widetilde{U}_{\lambda}$, which we denote by $\widetilde{V}$. Since $\widetilde{V}$ does not contain 0 , we may pull it back by the branch of $\left(E_{\lambda}^{n}\right)^{-1}$ that maps $z_{0}$ to itself, and obtain a new neighborhood of $z_{0}$ that strictly contains $\widetilde{V}$. The map $\widetilde{\phi}_{\lambda}$ can be extended to this new domain by using the functional equation (8) (see Figure 9).

We may repeat this process exactly $p$ times to obtain a nested sequence of (bounded) neighborhoods of $z_{0}$, where the map $\widetilde{\phi}_{\lambda}$ is well defined and whose image is a nested sequence of discs of radii $|\widetilde{\omega}|,\left|\widetilde{\omega} / \varrho_{\lambda}\right|,\left|\widetilde{\omega} / \varrho_{\lambda}^{2}\right|, \ldots,\left|\widetilde{\omega} / \varrho_{\lambda}^{p}\right|$ respectively.


Fig. 9. Sketch of the construction in the proof of Lemma A. 3 for $p=2$. Both maps are univalent on these domains.

We denote the largest neighborhood of $z_{0}$ in the process by $U_{\lambda}$. Observe that, by construction, $U_{\lambda}$ contains 0 in its boundary and it is mapped by $\widetilde{\phi}_{\lambda}$ to a round disc of radius $\left|\widetilde{\omega} / \varrho_{\lambda}^{p}\right|=\left|\widetilde{\phi}_{\lambda}(0)\right|$.

Since the linearizing coordinates are defined uniquely up to multiplication by a nonzero scale factor, the map defined as

$$
\phi_{\lambda}(z)=\widetilde{\phi}_{\lambda}(z) / \widetilde{\phi}_{\lambda}(0)
$$

on the domain $U_{\lambda}$ has the required properties.
Let $V_{0}=U_{\lambda}$. The boundary of $V_{0}$ is a (real-analytic) closed simple curve containing 0 in its boundary. Hence, we may still take one further preimage of $U_{\lambda}$ under $E_{\lambda}^{n}$ (taking the appropriate branch of the inverse), and obtain a finger $V_{1}$ (open, simply connected, unbounded to the right) that strictly contains $V_{0}$ and is mapped one-to-one onto $V_{0}$ by $E_{\lambda}^{n}$. See Figure 10. To see this, one can check that the set of preimages of $V_{0}$ under $E_{\lambda}$ is a collection of disjoint fingers unbounded in the left half plane, which are $2 \pi i$-translations of each other and map univalently onto $V_{0}$. Only one of them contains the point $z_{n-1}$ and hence this finger is contained in $A^{*}\left(z_{n-1}\right)$. We now pull back this finger along the periodic orbit and obtain $V_{1}$. The linearizing map $\phi_{\lambda}$ sends $V_{1}$ onto the disc $D\left(0,\left|\varrho_{\lambda}\right|^{-1}\right)$ univalently.

We are now ready to define the set $B_{\lambda}$ in Proposition A.1. We observe that, by construction, any disc of radius in between 1 and $\left|\varrho_{\lambda}\right|^{-1}$ has a preimage under $\phi_{\lambda}$ which contains 0 , and is mapped under $E_{\lambda}^{n}$ strictly inside itself. In particular, if we take for example

$$
\gamma_{\lambda}(t)=\phi_{\lambda}^{-1}\left(e^{2 \pi i t} \varrho_{\lambda}^{-1 / 2}\right)
$$

for $t \in \mathbb{T}$, and we define $B_{\lambda}$ to be the open set bounded by $\gamma_{\lambda}$, this set is a topological disc in $A^{*}\left(z_{0}\right)$ which contains 0 and $z_{0}$ and maps one-to-one strictly inside itself under $E_{\lambda}^{n}$.


Fig. 10. The unbounded domain $V_{1}$ where $\phi_{\lambda}$ is univalent and the construction of the curve $\gamma_{\lambda}$, boundary of $B_{\lambda}$

It remains to be checked that the curve $\gamma_{\lambda}$ depends holomorphically on the parameter $\lambda$. More precisely we want to show that the map

$$
\gamma: \Omega \times \mathbb{T} \rightarrow \partial\left(B_{\lambda}\right) \subset \mathbb{C}, \quad(\lambda, t) \mapsto \phi_{\lambda}^{-1}\left(e^{2 \pi i t} \varrho_{\lambda}^{-1 / 2}\right),
$$

is well defined and holomorphic in $\lambda$ ．To see this we just need to recall that the multiplier function

$$
\varrho: \Omega \rightarrow \mathbb{D}^{*}, \quad \lambda \mapsto \varrho_{\lambda},
$$

is a universal covering map［8］．Hence no closed loop in $\Omega$ can ever map to a closed loop in $\mathbb{D}$ surrounding the point 0 ．This implies that the principal square root $\varrho_{\lambda}^{1 / 2}$ is well defined and holomorphic in $\lambda$ ．Since $\phi_{\lambda}$ is biholo－ morphic in $z$ and $\lambda$ we conclude that $\gamma$ is holomorphic in $\lambda$ ．

This concludes the proof of Proposition A．1．

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