# INSTABILITY OF RAPIDLY-OSCILLATING PERIODIC SOLUTIONS FOR DISCONTINUOUS DIFFERENTIAL DELAY EQUATIONS 

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Abstract. We study the equation
$(\star) \quad \dot{x}(t)=-h(x(t-1))+f(x(t))$ for $t \geq 0, x_{\left.\right|_{[-1,0]}}=x_{0}$,
where $h$ is an odd function defined by $h(y)$ is equal to $a$ if $0<y<c$, equal to $b$ if $y \geq c, a>b>0$ and $c>0$ and $f$ is an odd $\mathcal{C}^{1}$ function such that $\sup |f(x)|<b$. We first consider the equation $\dot{x}(t)=-h(x(t-1))$, corresponding to $f \equiv 0$. We find the admissible shapes of rapidlyoscillating symmetric periodic solutions and we show that these periodic solutions are all unstable. We then extend these results to our general equation ( $\star$ ).

## 1. Introduction

We study here the scalar delay-differential equation

$$
\begin{equation*}
\dot{x}(t)=-h(x(t-1))+f(x(t)) \quad t \geq 0, \quad x_{\left.\right|_{[-1,0]}}=x_{0} \tag{1}
\end{equation*}
$$

where $x_{0} \in \mathcal{C}([-1,0]), h$ is defined as follows:
$h: \mathbb{R} \rightarrow \mathbb{R}$, odd, $h(y)=\left\{\begin{array}{ll}a & \text { if } 0<y<c \\ b & \text { if } y \geq c\end{array}\right.$ with $a>b>0$ and $c>0$,
$f$ is an odd Lebesgue measurable function such that ess sup $|f(x)|<b$.

[^0]The motivation for considering such systems comes from an automotive control problem, namely the fuel-air ratio regulation problem for spark ignition engines, see [6].

Equation (1) is a particular case of a more general class of first order systems given by

$$
\begin{equation*}
\dot{x}(t)=g(x(t), x(t-1)), \tag{4}
\end{equation*}
$$

which arises in a variety of models in the literature, see Diekmann et al [4].
Say that $g$ satisfies a "negative feedback condition" if

$$
\begin{array}{rll}
g(x, y)<0 & \text { for } & x>0 \text { and } y>0 \\
g(x, y)>0 & \text { for } & x<0 \text { and } y<0,  \tag{5}\\
y g(0, y)<0 & \text { for } & y \neq 0 .
\end{array}
$$

If $g$ satisfies some regularity conditions and in particular the negative feedback condition (5), then Equation (4) has a slowly-oscillating periodic solution, see Mallet-Paret and Nussbaum [12] and Mallet-Paret et al [13]. More general results are proved in [11, 12, 13].

If $g(x, y)=f(x)-\operatorname{sgn} y$, where $f$ is $\mathcal{C}^{1}$ and $\sup |f(x)|<1$, Condition (5) is satisfied. Equation (4) can be rewritten as

$$
\begin{equation*}
\dot{x}(t)=-\operatorname{sgn}(x(t-1))+f(x(t)) . \tag{6}
\end{equation*}
$$

Although $g$ is not regular in this case, a complete study of (6) has been realized by doing some explicit computations.

Define $V(t)$ as the number of zeros of $x$ on $\left[t^{\prime}-1, t^{\prime}\right)$ with $t^{\prime}$ the first time after $t$ such that $x\left(t^{\prime}\right)=0$. Fridman et al [7] have shown that if $V(0)$ is finite then $V(t)$ is even and nonincreasing. For any $n \in \mathbb{N}$, Equation (6) has a unique periodic solution $x$ with $V(t) \equiv 2 n[7]$. The slowest one $(n=0)$ is the only orbitally stable periodic solution, see [7, 2].

In addition, Shustin [17], Nussbaum and Shustin [16] and Akian and Bliman [1] have proved that, for any solution $x, V$ is finite after a finite time (there is no super-high-frequency). For related results, see [8, 9]. Thus, we have a complete description of the asymptotic behavior of (6) in which the slow cycle is the only stable part of the attractor.

Equation ( $1,2,3$ ) is still a particular case of (4) such that $g$ satisfies the negative feedback condition (5). In this paper, we try to generalize the results of $[7,2]$ to our equation $(1,2,3)$. We characterize some of rapidlyoscillating 2-phase periodic solutions and we show that these solutions are all unstable.

## 2. Previous Results

First, we recall some definitions:
Definition 2.1 (Rapidly-oscillating functions). A continuous function $x$ defined on $\left[t_{0},+\infty\right.$ ) is called rapidly-oscillating (with respect to 1 ) if there exist $t, t^{\prime}>t_{0}$, such that $t \neq t^{\prime},\left|t-t^{\prime}\right|<1$ and $x(t)=x\left(t^{\prime}\right)=0$.

Definition 2.2 (2-phase periodic (2PP) functions). A $T$-periodic ( $T>0$ ) continuous function $x$ defined on $\left[t_{0},+\infty\right)$ is called 2-phase periodic if there exists $t \geq t_{0}$ such that $x_{(t, t+T)}$ changes signs (strictly) exactly once.
Definition 2.3 (Symmetric periodic functions). A $T$-periodic ( $T>0$ ) continuous function $x$ defined on $\left[t_{0},+\infty\right)$ is called symmetric if $x\left(t+\frac{T}{2}\right)=-x(t)$ for $t \geq t_{0}$.

The following result proves the existence of rapidly-oscillating periodic solutions of Equation ( $1,2,3$ ):
Theorem 2.4 (Nussbaum [15]). Let $g: \mathbb{R}^{2} \rightarrow \mathbb{R}$ be a Lebesgue measurable function. Assume that for each $\left(x_{0}, y_{0}\right) \in \mathbb{R}^{2}$, there exist $R=R\left(x_{0}, y_{0}\right)$ $>0$ and $C=C\left(x_{0}, y_{0}\right)>0$ such that for every $y \in B_{R}\left(y_{0}\right)=\{y \in \mathbb{R}$ : $\left.\left|y-y_{0}\right|<R\right\}$, the map $x \mapsto g(x, y)$ is Lipschitzian with Lipschitz constant $C$ on $B_{R}\left(x_{0}\right)$. Assume that there are positive constants $m$ and $M$ such that $-M \leq g(x, y) \leq-m$ for all $(x, y) \in \mathbb{R}^{2}$ with $y>0$ and $m \leq g(x, y) \leq M$ for all $(x, y) \in \mathbb{R}^{2}$ with $y<0$. Then, for every integer $n \geq 0$, there exists a 2phase periodic Lipschitzian function $x_{n}$ which satisfies Equation (4) almost everywhere in $\mathbb{R}$ and has the property that, if $x_{n}(\tau)=0$ for some $\tau$ then $\operatorname{card}\left\{t \in(\tau-1, \tau): x_{n}(t)=0\right\}=2 n$.

Example. Let us consider Equation (1) with $f \equiv 0$ and $h$ given by (2). We seek some periodic solutions. If $x(t)$ denotes a 2 -phase periodic solution of (1), then either (a) $|x(t)| \leq c$ for all t , or (b) $|x(t)|>c$ for some $t$. If case (a) holds, then $h(x(t-1))=a \operatorname{sgn}(x(t-1))$ almost everywhere. Hence the results on the periodic solutions of (6) and their stability can be applied. We give an explicit form of the family of periodic solutions $x_{n}$ of Equation (1) such that $\left\|x_{n}\right\|_{\infty} \leq c$ : for $t \in\left[0, \frac{T_{n}}{2}\right], x_{n}(t)=-\frac{a T_{n}}{4}+a t$, and for $t \in\left[\frac{T_{n}}{2}, T_{n}\right]$, $x_{n}(t)=\frac{a T_{n}}{4}-a\left(t-\frac{T_{n}}{2}\right)$ and the period $T_{n}$ is equal to $\frac{4}{4 n+1}$. For a given $n$, $x_{n}(t)$ solves (1) if and only if $\frac{a}{4 n+1} \leq c$.

Ivanov and Losson [10] claim that first order differential delay equations with negative feedback can possess asymptotically stable rapidly oscillating solutions. They present an analytically tractable example for Equation (1) with $h$ piecewise constant.

Here, we are interested in symmetric 2-phase periodic solutions of (1) which exceed c in magnitude (case(b) in the example). There are two changes of slopes in each monotone phase of $x$. We start out with $f \equiv 0$. The evolution of system (1) around each periodic solution can be described with a linear discrete dynamical system. We shall show that the fixed point of the discrete system associated to our problem is unstable. Then, we extend these results to the class of functions $f$ verifying (3) for which the discrete system is nonlinear.

$$
\text { 3. STUDY OF } \dot{x}(t)=-h(x(t-1))
$$

Let $x$ be a solution of (1). We define the set $Z$ of points where $x$ vanishes and changes sign, and the cardinality $V$ of the number of zeros with changes of sign: $Z=\left\{t \geq-1: x(t)=0\right.$ and $\forall \epsilon>0 \exists t^{\prime} \in(t, t+\epsilon], t^{\prime \prime} \in[t-$ $\left.\epsilon, t) / x\left(t^{\prime}\right) x\left(t^{\prime \prime}\right)<0\right\}, V(t)=\operatorname{card}\left(Z \cap\left[t^{\prime}-1, t^{\prime}\right)\right), \quad$ where $t^{\prime}=\inf \{s \geq$ $t, s \in Z\}$. We write $V(t)=\infty$ if $Z \cap\left[t^{\prime}-1, t^{\prime}\right)$ is not finite.

Theorem 3.1. Let us consider the following equation

$$
\begin{equation*}
\dot{x}(t)=-h(x(t-1)) \quad t \geq 0, \quad x_{[-1,0]}=x_{0}, \tag{7}
\end{equation*}
$$

where $h$ is given by (2), and $x_{0} \in \mathcal{C}([-1,0])$ has finitely many zeros.
For any $x_{0}$ with a finite number of zeros, there exists a unique function $x \in \mathcal{C}([-1,+\infty)$ ), absolutely continuous on $[0,+\infty)$, satisfying Equation (7) almost everywhere. For such a solution, $V(t)$ is nonincreasing and even-valued. For any periodic solution of (7), $V(t)$ is constant. Rapidlyoscillating symmetric 2-phase periodic solutions $x$, such that $\|x\|_{\infty}>c$, are characterized as in Cases $\# 2, \# 2^{\prime}, \# 3, \# 3^{\prime}, \# 4, \# 4^{\prime}, \# 5$ and $\# 6$, defined below. These periodic solutions are all unstable.
Proof. Existence and uniqueness of a solution of Equation (7) can be proved by induction. Since $x_{0}$ has a finite number of zeros, $-h(x(t-1))$ is a piecewise constant function on $[0,1]$, then a solution of ( 7 ) is easily determined.

The rest of the proof is done in the following subsections.
3.1. Evolution of the number of zeros. The integer-valued Lyapunov function $V(t)$ has already been exploited for smooth function $h$ by J. MalletParet, see [11]. In [11, Theorem A,(i)], it has been proved that $V$ is nonincreasing and even-valued for any solution of (4) if $g$ satisfies the negative feedback condition (5). Indeed, the proof can be applied for non smooth $g$. Therefore, for any solution, $V$ is constant after a finite time. Moreover, for any periodic solution, $V$ is constant. This solution is rapidly oscillating if and only if $V>0$.

Shape \#1


Shape \#3

c


Shape \#5


Symmetric version: shape $\# 1^{\prime}$


Symmetric version: shape $\# 2^{\prime}$


Symmetric version: shape $\# 3^{\prime}$


Since we want to prove a result of local instability, we only consider perturbations of periodic solutions for which $V$ does not change. From now on, our argument reduces to the study of discrete systems describing the evolution on each level $V(t) \equiv 2 n, n \in \mathbb{N} \backslash\{0\}$.
3.2. Shapes of rapidly-oscillating 2PP solutions. The reader can check that any 2 PP solution $x$ of (7) increases strictly from its global minimum to its global maximum, decreases from its global maximum to its global minimum and then repeats. By exploiting the oddness of $h$, one obtains for each $n \geq 0$ a nonconstant periodic solution $x_{n}(t)$ of $(7)$ as in Theorem 2.4 such that $x_{n}$ has minimal period $T_{n}$ and $x_{n}$ is symmetric, i.e., $x_{n}\left(t+\frac{T_{n}}{2}\right)=$ $-x_{n}(t)$ for all $t$. Consequently, $T_{n}$ verifies $\frac{1}{n+1 / 2}<T_{n}<\frac{1}{n}$. We shall restrict our attention to such symmetric rapidly-oscillating (thus $n \geq 1$ ) periodic solutions of $(7)$ such that $\left\|x_{n}\right\|_{\infty}>c$. In the following we draw all a priori possible shapes of these solutions from their minimum to their maximum (so in time interval of length $\frac{T_{n}}{2}$ ). The key point is to determine where changes of slope are situated with respect to the straight lines $x=-c, x=0$ and $x=c$. We only consider "non degenerate" solutions, that is solutions such that the changes of slope do not intersect the straight lines $x=-c, x=0$ or $x=c$.

We show that Shapes $\# 1$ and $\# 1^{\prime}$ are impossible. The existence of a symmetric 2 PP solution of Equation (7) with a given other shape (shapes $\# 2$ to $\# 6$ or $\# 2^{\prime}$ to $\# 4^{\prime}$ ) depends on the constants $a, b, c$ and on $n$.
3.3. Evolution on level $V(t) \equiv 2 n, n \in \mathbb{N} \backslash\{0\}$. Let $x$ be a solution of (7) such that $V(t)=2 n \forall t \geq 0$ with $n \in \mathbb{N} \backslash\{0\}$. Let us consider an interval $I=[t-1, t)$ with $x(t)=0$. As $V(t)=2 n$, the interval $I$ "splits" in $2 n+1$ intervals $\left(t-1, t_{1}\right),\left(t_{1}, t_{2}\right), \ldots,\left(t_{2 n}, t\right)$ and on each of these intervals, $x$ is of constant sign. First of all, we consider $\left(t-1, t_{1}\right)$ on which $x$ is positive for instance and $x(t-1) \neq 0$ and the interval $\left(t_{1}, t_{2}\right)$ on which $x$ is negative. We divide $\left(t-1, t_{1}\right)$ into three parts (possibly empty) : the first part is the time interval, with length $a_{1}^{\prime}$ on which $x$ is increasing and between 0 and $c$, the second part is the time interval, with length $a_{2}^{\prime}$ on which $x$ is above $c$ and finally the third one is the time interval, with length $a_{3}^{\prime}$, on which $x$ is decreasing from $c$ to 0 . Concerning $\left(t_{1}, t_{2}\right)$, we divide it into three parts: the first part is the time interval, with length $a_{1}$, on which $x(t)$ is decreasing from 0 to $-c$, the second part is the time interval, with length $a_{2}$, on which $x<-c$ and the third part is the time interval, with length $a_{3}$, on which $x$ is increasing from $-c$ to 0 . We define in a similar way $a_{4}, \ldots, a_{6 n}$. From the
definition of $h$, we see that the value of $\left.h(x())\right|_{I$.$} is uniquely determined by$ the value of $a=\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}, a_{1}, \ldots, a_{6 n}\right)$.

Let us define the simplex $\sum_{6 n+3}=\left\{\bar{a}=\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}, a_{1}, a_{2}, a_{3}, \ldots, a_{6 n-2}\right.\right.$, $\left.a_{6 n-1}, a_{6 n}\right) \in \mathbb{R}^{6 n+3}: a_{i}^{\prime} \geq 0$ for $1 \leq i \leq 3, \sum_{i=1}^{3} a_{i}^{\prime}>0, a_{i}>0$ for $1 \leq i \leq$ $6 n$ and $\left.\sum_{i=1}^{3} a_{i}^{\prime}+\sum_{i=1}^{6 n} a_{i}=1\right\}$ and let $\Phi_{n}$ be the flow on $\Sigma_{6 n+3}$ associated with the evolution of ( 7 ) on the level $V \equiv 2 n$, corresponding to the map $x_{\mid t-1, t)} \mapsto x_{\left[t^{\prime}-1, t^{\prime}\right)}$, where $t^{\prime}$ is the first time $t^{\prime}>t$ such that $x\left(t^{\prime}\right)=0$. We have,

$$
\begin{align*}
& \Phi_{n}: \Sigma_{6 n+3} \longrightarrow \mathbb{R}^{6 n+3}, \\
& \bar{a}=\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}, a_{1}, \ldots, a_{6 n}\right) \longmapsto \Phi_{n}(\bar{a})=\left(a_{1}^{\prime \prime}, a_{2}^{\prime \prime}, a_{3}^{\prime \prime}, a_{4}, \ldots, a_{6 n+3}\right) . \tag{8}
\end{align*}
$$

In the following, we compute $\Phi_{n}(\bar{a})$, when $\bar{a}$ is near $\bar{a}^{\star}$, where $\bar{a}^{\star}$ corresponds to a "nondegenerate" symmetric periodic orbit of Equation (7). Thus, we can suppose that the shape of the solution $x$ corresponding to $\bar{a}$ does not change with time and that $\bar{a}$ is such that $\Phi_{n}(\bar{a}) \in \Sigma_{6 n+3}$.

If for instance $x$ is negative on $\left(t, t^{\prime}\right)$ and if $-\underline{x}$ denotes the minimum of $x$ on $\left[t, t^{\prime}\right]$, then

$$
\begin{equation*}
\underline{x}=a a_{1}^{\prime}+b a_{2}^{\prime}+a a_{3}^{\prime}=a\left(a_{1}-a_{1}^{\prime \prime}\right)+b\left(a_{2}-a_{2}^{\prime \prime}\right)+a\left(a_{3}-a_{3}^{\prime \prime}\right) . \tag{9}
\end{equation*}
$$

With the additional conditions $0 \leq a_{i}^{\prime \prime} \leq a_{i}$ for $1 \leq i \leq 3$ and ( $a_{i}^{\prime \prime}>0 \Rightarrow$ $a_{i+1}^{\prime \prime}=a_{i+1}$ ) for $i=1,2$, one determines $a_{1}^{\prime \prime}, a_{2}^{\prime \prime}$ and $a_{3}^{\prime \prime}$. The values of $a_{6 n+1}$, $a_{6 n+2}$ and $a_{6 n+3}$ can then be computed. Their formula depends on the chosen shape. Indeed, for each shape, we shall obtain an affine function $\Phi_{n}(\bar{a})$ of $\bar{a}$. For non degenerate cases, this function has at most one fixed point $\bar{a}$, which corresponds to the unique (if it exists) symmetric periodic solution of Equation (7) with the corresponding shape. The fixed point $\bar{a}$ has to verify: $a_{1}^{\prime \prime}=a_{1}^{\prime}, a_{2}^{\prime \prime}=a_{2}^{\prime}, a_{3}^{\prime \prime}=a_{3}^{\prime}, a_{i}=a_{i+3}$ for $1 \leq i \leq 6 n$, thus $a_{i}=a_{6 n+i}$ for $1 \leq$ $i \leq 3$. If $T_{n}$ is the minimal period of the periodic solution, then $a_{1}+a_{2}+a_{3}=$ $\frac{T_{n}}{2}$ and $a_{1}^{\prime}+a_{2}^{\prime}+a_{3}^{\prime}=1-n T_{n}>0$. Thus, $\frac{1}{n+1 / 2}<T_{n}<\frac{1}{n}$. Now, we determine more precisely the value of $\Phi_{n}$ corresponding to each shape. We compute the corresponding fixed point and if it exists, we prove that it is unstable for the dynamic $\Phi_{n}$. This will imply the instability of the corresponding symmetric 2PP solution of Equation (7). A sufficient condition is that the tangent map of $\Phi_{n}$ at the fixed point has an eigenvalue of modulus strictly greater than 1 , or that the tangent map of $\Phi_{n}^{-1}$ at the fixed point has an eigenvalue of modulus strictly less than 1 . We shall thus compute the matrix associated to the tangent map of $\Phi_{n}$ and the characteristic polynomial of this matrix. Then, we prove that this polynomial has at least one root with modulus strictly greater than 1 , using homotopy techniques. Since these
characteristic polynomials are essentially of three types, we gather important results concerning their roots into Section 3.5.

### 3.4. Study of $\Phi_{n}$ for all cases of shape.

3.4.1. Cases $\# 1$ and $\# 1^{\prime}$.

Lemma 3.2. $\Phi_{n}$ has no fixed point or equivalently, there is no symmetric 2PP solution with shape \#1 or \#1' .
Proof. For Case \#1, $a_{6 n+2} \geq a_{2}^{\prime}+a_{3}^{\prime}$ and $a_{2}^{\prime \prime}=a_{2}$. Since any fixed point satisfies $a_{6 n+2}=a_{2}$ and $a_{3}^{\prime}>0$, we obtain a contradiction. So $\Phi_{n}$ has no fixed point. For Case $\# 1^{\prime}$, we get $a_{6 n+2} \geq a_{1}+a_{2}$ and $a_{1}>0$, thus $\Phi_{n}$ has no fixed point.

### 3.4.2. Case \#2.

Proposition 3.3. Any fixed point of $\Phi_{n}$ with shape $\# 2$ is unstable.
Let us first compute the flow $\Phi_{n}$ near a fixed point with shape $\# 2$. It is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=0$ and

$$
\begin{aligned}
a_{2}^{\prime \prime} & =-a_{2}^{\prime}-\frac{a}{b} a_{3}^{\prime}+\frac{a}{b} a_{1}+a_{2}, \quad a_{3}^{\prime \prime}=a_{3}, \\
a_{6 n+1} & =\frac{c}{b}, \quad a_{6 n+2}=a_{2}^{\prime}\left(1+\frac{b}{a}\right)+2 a_{3}^{\prime}-\frac{c}{b}-\frac{c}{a}, \\
a_{6 n+3} & =a_{1}\left(1-\frac{a}{b}\right)+a_{2}^{\prime}\left(1-\frac{b}{a}\right)+a_{3}^{\prime}\left(-1+\frac{a}{b}\right)+\frac{c}{a} .
\end{aligned}
$$

The fixed point is defined by the following equations:

$$
\begin{aligned}
& a_{2}^{\prime}=-a_{2}^{\prime}-\frac{a}{b} a_{3}^{\prime}+\frac{a}{b} a_{1}+a_{2}, \quad a_{3}^{\prime}=a_{3}, \quad a_{1}=\frac{c}{b}, \\
& a_{2}=a_{2}^{\prime}\left(1+\frac{b}{a}\right)+2 a_{3}^{\prime}-\frac{c}{b}-\frac{c}{a}, \quad 1=2 n\left(a_{1}+a_{2}+a_{3}\right)+a_{2}^{\prime}+a_{3}^{\prime},
\end{aligned}
$$

which solution is given by $a_{2}^{\prime}=\frac{1}{d}\left((2 n+1) \frac{c}{b}\left(-4+3 \frac{a}{b}-\frac{b}{a}\right)+2-\frac{a}{b}+\frac{c}{b}\left(3-2 \frac{a}{b}\right)\right)$, $a_{3}^{\prime}=\frac{1}{d}\left((2 n+1) \frac{c}{b}\left(3 \frac{b}{a}-\frac{a}{b}\right)+1-\frac{b}{a}+\frac{c}{b}\left(1-2 \frac{b}{a}\right)\right), a_{1}=\frac{c}{b}, a_{2}=\frac{1}{d}\left(2(2 n+1) \frac{c}{b}\left(\frac{a}{b}-\right.\right.$ $\left.\left.\frac{b}{a}-2\right)+3-\frac{a}{b}+2 \frac{c}{b}\left(2-\frac{a}{b}\right)\right), a_{3}=a_{3}^{\prime}, d=(2 n+1)\left(4-\frac{a}{b}-\frac{b}{a}\right)-1$. We suppose that $d \neq 0$, so that the fixed point exists and is unique. To be compatible with the shape $\# 2$, the fixed point $\bar{a}$ has to satisfy $a_{2}^{\prime}>\frac{c}{b}, a_{3}^{\prime}>0$ and $0<a_{2}-a_{2}^{\prime}<\frac{c}{b}$, which are equivalent to the conditions $\frac{a}{b}>2$ with $(2 n+1)\left(4-\frac{a}{b}-\frac{b}{a}\right)<1$, $2 n+1>\left(\frac{b}{c}\left(1-\frac{b}{a}\right)+1-2 \frac{b}{a}\right) /\left(\frac{a}{b}-3 \frac{b}{a}\right)$ and $\left(\frac{b}{c}+1\right) /\left(\frac{a}{b}+\frac{b}{a}\right)<2 n+1<\frac{1}{2}\left(\frac{b}{2 c}+1\right)$, or $(2 n+1)\left(4-\frac{a}{b}-\frac{b}{a}\right)>1,2 n+1<\left(\frac{b}{c}\left(1-\frac{b}{a}\right)+1-2 \frac{b}{a}\right) /\left(\frac{a}{b}-3 \frac{b}{a}\right)$ and $\left(\frac{b}{c}+1\right) /\left(\frac{a}{b}+\frac{b}{a}\right)>2 n+1>\frac{1}{2}\left(\frac{b}{2 c}+1\right)$. These conditions can be satisfied for some constants $a, b, c$ and some $n \geq 1$.

Let us now compute the tangent map of $\Phi_{n}$. Any sufficiently small perturbation $\bar{a}$ of the fixed point (or any fixed point) of $\Phi_{n}$ satisfies

$$
\begin{equation*}
a_{1}^{\prime}=0 \quad \text { and } \quad a_{3 k-2}=\frac{c}{b}, \text { for } k=1, \ldots, 2 n . \tag{10}
\end{equation*}
$$

We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+2}$ of $\Sigma_{6 n+3}$ defined by (10) and we replace $\Phi_{n}$ by the map $\widetilde{\Phi}_{n}:\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{2}, a_{3}, a_{5}, \ldots, a_{6 n}\right) \mapsto\left(a_{2}^{\prime \prime}, a_{3}^{\prime \prime}, a_{5}, a_{6}, \ldots\right.$, $\left.a_{6 n+3}\right)$. Let us denote by $M \in \mathbb{R}^{(4 n+2) \times(4 n+2)}$, the matrix associated to the tangent map of $\widetilde{\Phi}_{n}$, its transposed matrix $M^{T}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
-1 & 0 & \cdots & \cdots & \cdots & 1+\frac{b}{a} & 1-\frac{b}{a} \\
-\frac{a}{b} & 0 & \cdots & \cdots & \cdots & 2 & -1+\frac{a}{b} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \cdots & \cdots & \cdots & 0 \\
\vdots & 0 & \ddots & \ddots & \vdots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right)
$$

The following result will be useful to compute characteristic polynomials of tangent map matrices of $\Phi_{n}$.
Lemma 3.4. Let $G$ and $T$ be $k \times k$ real matrices, with $k$ a positive integer. For any positive integer $N$, denote $I_{N}$ the identity $N \times N$ matrix, and for any real e, denote $\mathcal{M}_{e, N}=\mathcal{G}_{N}+e \mathcal{T}_{N}$ with

$$
\mathcal{G}_{N}=\left(\begin{array}{c|c}
G \mid & 0 \\
0 \\
\hline I_{N-k} & 0
\end{array}\right), \quad \mathcal{T}_{N}=\left(\begin{array}{c|c}
0 & T \\
\hline 0 & 0
\end{array}\right)
$$

Then, for any positive integer $m$ such that $N \geq k(m+1)$, the characteristic polynomial of $\mathcal{M}_{e, N}$ satisfies:

$$
\operatorname{det}\left(\mathcal{M}_{e, N}-X I_{N}\right)=(-X)^{k m} \operatorname{det}\left(\mathcal{M}_{\frac{e}{X^{m}}, N-k m}-X I_{N-k m}\right)
$$

Proof. If $N-k \geq k$, we have $\mathcal{G}_{N}=\left(\begin{array}{c|c|c}\mathcal{G}_{N-k} & 0 \\ \hline 0 & I_{k} & 0\end{array}\right)$, so

$$
\begin{aligned}
& \operatorname{det}\left(\mathcal{M}_{e, N}-X I_{N}\right)=\operatorname{det}\left(\right), \\
& =\operatorname{det}\left(-X I_{k}\right) \times \operatorname{det}\left(\left(\mathcal{G}_{N-k}-X I_{N-k}\right)+\frac{1}{X}\left(\frac{e T}{0}\right)\left(\begin{array}{lll}
0 & I_{k}
\end{array}\right)\right),
\end{aligned}
$$

$$
\begin{aligned}
& =(-X)^{k} \times \operatorname{det}\left(\left(\mathcal{G}_{N-k}-X I_{N-k}\right)+\frac{1}{X}\left(\begin{array}{c|c}
0 & e T \\
\hline 0 & 0
\end{array}\right)\right), \\
& =(-X)^{k} \times \operatorname{det}\left(\mathcal{M}_{\frac{e}{X}, N-k}-X I_{N-k}\right) .
\end{aligned}
$$

One concludes by induction on $m$.
Using Lemma 3.4 with $N=4 n+2, k=2$ and $m=2 n$, we obtain:
Lemma 3.5. The characteristic polynomial of $M$ is equal to:

$$
R(X)=X^{4 n+2}+X^{4 n+1}-(2+B) X^{2 n+1}+B,
$$

where $B=-2+\frac{a}{b}+\frac{b}{a}$ is such that $B>0$.
Proof of Proposition 3.3. The fixed point $\bar{a}$ of $\Phi_{n}$ found above is unstable if the tangent map of $\Phi_{n}$ at $\bar{a}$ has an eigenvalue with modulus strictly greater than 1 , or equivalently if the characteristic polynomial $R$ of $M$ has one root $z$, such that $|z|>1$.
a) If $B>1$, then the product of all roots of $R$ is greater than 1 , so $R$ has clearly one root $z$, such that $|z|>1$.
b) If $0<B \leq 1$, we use Theorem 3.19 below, which implies that, for $n \geq 1$, there exists at least one root $z$ such that $|z|>1$.

### 3.4.3. Case \#2'.

Proposition 3.6. Any fixed point of $\Phi_{n}$ with shape $\# 2^{\prime}$ is unstable.
The flow $\Phi_{n}$ near a fixed point with shape $\# 2^{\prime}$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=0$ and

$$
\begin{aligned}
& a_{2}^{\prime \prime}=-a_{2}^{\prime}-\frac{a}{b} a_{3}^{\prime}+\frac{a}{b} a_{1}+a_{2}, \quad a_{3}^{\prime \prime}=a_{3}, \quad a_{6 n+1}=a_{2}^{\prime}\left(1-\frac{b}{a}\right)+\frac{c}{a}, \\
& a_{6 n+2}=a_{2}^{\prime}\left(1+\frac{b}{a}\right)+a_{3}^{\prime}\left(1+\frac{a}{b}\right)+a_{1}\left(1-\frac{a}{b}\right)-\frac{c}{b}-\frac{c}{a}, \quad a_{6 n+3}=\frac{c}{b} .
\end{aligned}
$$

The fixed point is defined by the following equations:

$$
\begin{aligned}
& a_{2}^{\prime}=\frac{1}{d}\left(2 n \frac{c}{b}\left(\frac{b}{a}+\frac{a}{b}\right)+\frac{c}{b}-1\right) \quad a_{3}^{\prime}=\frac{c}{b}, \quad a_{1}=\left(1-\frac{b}{a}\right) a_{2}^{\prime}+\frac{c}{a}, \\
& a_{2}=\left(3-\frac{a}{b}\right) a_{2}^{\prime}+\frac{c}{b}\left(\frac{a}{b}-1\right), \quad a_{3}=\frac{c}{b}, \quad d=2 n\left(\frac{a}{b}+\frac{b}{a}-4\right)-1 .
\end{aligned}
$$

We suppose that $d \neq 0$, so that the fixed point exists and is unique. To be compatible with the shape $\# 2^{\prime}$, the fixed point $\bar{a}$ has to satisfy $0<a_{2}^{\prime}<\frac{c}{b}$, $a_{2}-a_{2}^{\prime}>\frac{c}{b}$ and $a_{1}>0$ which are equivalent to the conditions: $\frac{a}{b}>2$ with $2 n\left(\frac{a}{b}+\frac{b}{a}-4\right)>1$ and $\left(\frac{b}{c}-1\right) /\left(\frac{a}{b}+\frac{b}{a}\right)<2 n<\frac{1}{2}\left(\frac{b}{2 c}-1\right)$, or $2 n\left(\frac{a}{b}+\frac{b}{a}-4\right)<1$
and $\left(\frac{b}{c}-1\right) /\left(\frac{a}{b}+\frac{b}{a}\right)>2 n>\frac{1}{2}\left(\frac{b}{2 c}-1\right)$. These conditions can be satisfied for some constants $a, b, c$ and some $n \geq 1$.

Any sufficiently small perturbation $\bar{a}$ of a fixed point satisfies

$$
\begin{equation*}
a_{1}^{\prime}=0 \quad \text { and } \quad a_{3}^{\prime}=a_{3 k}=\frac{c}{b}, \text { for } k=1, \ldots, 2 n \tag{11}
\end{equation*}
$$

We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+1}$ of $\Sigma_{6 n+3}$ defined by (11) and we replace $\Phi_{n}$ by the map $\widetilde{\Phi}_{n}:\left(a_{2}^{\prime}, a_{1}, a_{2}, a_{4}, \ldots, a_{6 n-1}\right) \longmapsto\left(a_{2}^{\prime \prime}, a_{4}, a_{5}, \ldots, a_{6 n+2}\right)$.

The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+1) \times(4 n+1)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
-1 & 0 & \cdots & \cdots & \cdots & 1-\frac{b}{a} & 1+\frac{b}{a} \\
\frac{a}{b} & 0 & \cdots & \cdots & \cdots & 0 & 1-\frac{a}{b} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \cdots & \cdots & \cdots & 0 \\
\vdots & 0 & \ddots & \ddots & \vdots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right) .
$$

Using Lemma 3.4 with $N=4 n+1, k=2$ and $m=2 n-1$, we obtain:
Lemma 3.7. The characteristic polynomial of $M$ is equal to:

$$
R(X)=-\left(X^{4 n+1}+X^{4 n}-(2+B) X^{2 n}+B\right)
$$

where $B=-2+\frac{a}{b}+\frac{b}{a}$ is such that $B>0$.
The proof of Proposition 3.6 follows by the same arguments as for Proposition 3.3, using Theorem 3.20 instead of Theorem 3.19.

### 3.4.4. Case \#3.

Proposition 3.8. Any fixed point of $\Phi_{n}$ with shape $\# 3$ is unstable.
The flow $\Phi_{n}$ near a fixed point with shape $\# 3$ is given by (8) and

$$
\begin{aligned}
& a_{1}^{\prime \prime}=-a_{1}^{\prime}-\frac{b}{a} a_{2}^{\prime}+a_{1}-a_{3}^{\prime}, \quad a_{2}^{\prime \prime}=a_{2}, \quad a_{3}^{\prime \prime}=a_{3}, \\
& a_{6 n+1}=a_{2}^{\prime}\left(1-\frac{b}{a}\right)+\frac{c}{a}, \quad a_{6 n+2}=2\left(a_{1}^{\prime}+\frac{b}{a} a_{2}^{\prime}+a_{3}^{\prime}-\frac{c}{a}\right), \quad a_{6 n+3}=\frac{c}{a} .
\end{aligned}
$$

The fixed point is unique and defined by the following equations:

$$
a_{1}^{\prime}=\left(\frac{1}{2}-\frac{b}{a}\right) a_{2}, \quad a_{2}^{\prime}=a_{2}, \quad a_{3}^{\prime}=\frac{c}{a},
$$

$$
a_{1}=\left(1-\frac{b}{a}\right) a_{2}+\frac{c}{a}, \quad a_{2}=\frac{1-(4 n+1) \frac{c}{a}}{(2 n+1)\left(2-\frac{b}{a}\right)-\frac{1}{2}}, \quad a_{3}=\frac{c}{a}
$$

To be compatible with shape $\# 3$, the fixed point $\bar{a}$ has to satisfy $a a_{1}^{\prime}+b a_{2}^{\prime}<$ $c<a a_{1}^{\prime}+b a_{2}^{\prime}+a a_{3}^{\prime}, a_{1}^{\prime}>0$ and $a_{2}>0$ which are equivalent to $\frac{a}{b}>2$ and $\left(\frac{a}{2 c}+1\right) /\left(3-\frac{b}{a}\right)<2 n+1<\frac{1}{2}\left(\frac{a}{c}+1\right)$. Any sufficiently small perturbation $\bar{a}$ of the fixed point satisfies $a_{3}^{\prime}=a_{3 k}=\frac{c}{a}$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+1}$ of $\Sigma_{6 n+3}$ defined by these conditions and we replace $\Phi_{n}$ by the map

$$
\begin{equation*}
\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{1}, a_{2}, a_{4}, \ldots, a_{6 n-1}\right) \quad \stackrel{\widetilde{\Phi}_{n}}{\longmapsto}\left(a_{1}^{\prime \prime}, a_{2}^{\prime \prime}, a_{4}, \ldots, a_{6 n+2}\right) \tag{12}
\end{equation*}
$$

The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+2) \times(4 n+2)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
-1 & 0 & \cdots & \cdots & \cdots & 0 & 2 \\
-\frac{b}{a} & 0 & \cdots & \cdots & 0 & 1-\frac{b}{a} & 2 \frac{b}{a} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \cdots & \cdots & \cdots & 0 \\
\vdots & 0 & \ddots & \ddots & \vdots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right)
$$

Using Lemma 3.4 with $N=4 n+2, k=2$ and $m=2 n$, we obtain:
Lemma 3.9. The characteristic polynomial of $M$ is equal to:

$$
R(X)=X^{4 n+2}+X^{4 n+1}-(2+B) X^{2 n+1}+B
$$

with $B=2 \frac{b}{a}-2<0$.
Proof of Proposition 3.8. The proof of Proposition 3.8 follows by the same arguments as for Proposition 3.3, using Theorem 3.19. Indeed, since we know that $a>2 b$, then $|B|>1$ and $R$ has clearly one root $z$ such that $|z|>1$. This shows that the fixed point of $\Phi_{n}$ corresponding to shape $\# 3$ is unstable.

### 3.4.5. Case $\# 3^{\prime}$.

Proposition 3.10. Any fixed point of $\Phi_{n}$ with shape $\# 3^{\prime}$ is unstable.
The flow $\Phi_{n}$ near a fixed point with shape $\# 3^{\prime}$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=a_{2}^{\prime}=a_{2}^{\prime \prime}=0$ and

$$
a_{3}^{\prime \prime}=-a_{3}^{\prime}+a_{1}+\frac{b}{a} a_{2}+a_{3}, \quad a_{6 n+1}=\frac{c}{a}
$$

$$
a_{6 n+2}=2\left(a_{3}^{\prime}-\frac{c}{a}\right), \quad a_{6 n+3}=a_{2}\left(1-\frac{b}{a}\right)+\frac{c}{a}
$$

The fixed point is unique and defined by the following equations:

$$
a_{3}^{\prime}=\frac{1}{2} a_{2}+\frac{c}{a}, \quad a_{1}=\frac{c}{a}, \quad a_{2}=\frac{1-(4 n+1) \frac{c}{a}}{2 n\left(2-\frac{b}{a}\right)+\frac{1}{2}}, \quad a_{3}=a_{2}\left(1-\frac{b}{a}\right)+\frac{c}{a}
$$

To be compatible with the shape $\# 3^{\prime}$, the fixed point $\bar{a}$ has to satisfy $a_{3}^{\prime} a>$ $c>a_{3}^{\prime} a-a_{1} a, a_{3}^{\prime}<a_{3}$ and $a_{2}>0$ which are equivalent to $\frac{a}{b}>2$ and $\left(\frac{a}{2 c}-1\right) /\left(3-\frac{b}{a}\right)<2 n<\frac{1}{2}\left(\frac{a}{c}-1\right)$. Any sufficiently small perturbation $\bar{a}$ of the fixed point satisfies $a_{1}^{\prime}=a_{2}^{\prime}=0$ and $a_{3 k-2}=\frac{c}{a}$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+1}$ of $\Sigma_{6 n+3}$ defined by these conditions and we replace $\Phi_{n}$ by the map

$$
\begin{equation*}
\left(a_{3}^{\prime}, a_{2}, a_{3}, a_{5}, \ldots, a_{6 n}\right) \stackrel{\widetilde{\Phi}_{n}}{\longleftrightarrow}\left(a_{3}^{\prime \prime}, a_{5}, a_{6}, \ldots, a_{6 n+3}\right) . \tag{13}
\end{equation*}
$$

The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+1) \times(4 n+1)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
-1 & 0 & \cdots & \cdots & \cdots & 2 & 0 \\
\frac{b}{a} & 0 & \cdots & \cdots & \cdots & 0 & 1-\frac{b}{a} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \cdots & \cdots & \cdots & 0 \\
\vdots & 0 & \ddots & \ddots & \vdots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right) .
$$

Using Lemma 3.4 with $N=4 n+1, k=2$ and $m=2 n-1$, we obtain:
Lemma 3.11. The characteristic polynomial of $M$ is equal to:

$$
R(X)=-\left(X^{4 n+1}+X^{4 n}-(2+B) X^{2 n}+B\right),
$$

with $B=2 \frac{b}{a}-2<0$.
The proof of Proposition 3.10 may follow by the same arguments as for Proposition 3.3, using Theorem 3.20, or by that of Proposition $3.8(a>2 b$ implies $|B|>1)$.
3.4.6. Case \#4.

Proposition 3.12. Any fixed point of $\Phi_{n}$ with shape $\# 4$ is unstable.

The flow $\Phi_{n}$ near a fixed point with shape $\# 4$ is given by (8) and

$$
\begin{aligned}
& a_{1}^{\prime \prime}=-a_{1}^{\prime}-\frac{b}{a} a_{2}^{\prime}-a_{3}^{\prime}+a_{1}, \quad a_{2}^{\prime \prime}=a_{2}, \quad a_{3}^{\prime \prime}=a_{3}, \quad a_{6 n+1}=a_{1}^{\prime}\left(1-\frac{a}{b}\right)+\frac{c}{b}, \\
& a_{6 n+2}=a_{2}^{\prime}\left(1+\frac{b}{a}\right)+a_{1}^{\prime}\left(1+\frac{a}{b}\right)+2 a_{3}^{\prime}-\frac{c}{a}-\frac{c}{b}, \quad a_{6 n+3}=\frac{c}{a}
\end{aligned}
$$

The fixed point is unique and defined by the following equations :

$$
\begin{aligned}
& a_{1}^{\prime}=\frac{1}{d}\left((2 n+1) \frac{c}{a}\left(\frac{a}{b}+\frac{b}{a}\right)-\frac{b}{a}-\frac{c}{a}\right), \quad a_{2}^{\prime}=a_{2}, \quad a_{3}^{\prime}=\frac{c}{a}, \\
& a_{1}=\frac{1}{d}\left((2 n+1) \frac{c}{a}\left(3 \frac{a}{b}+\frac{b}{a}-2\right)-\frac{b}{a}-\frac{c}{a}+1\right), \\
& a_{2}=\frac{1}{d}\left(-4(2 n+1) \frac{c}{b}+1+\frac{a}{b}+2 \frac{c}{b}\right), \\
& a_{3}=\frac{c}{a}, \quad d=(2 n+1)\left(2+\frac{a}{b}-\frac{b}{a}\right)-1>0 .
\end{aligned}
$$

To be compatible with the shape \#4, the fixed point $\bar{a}$ has to satisfy $0<$ $a_{1}^{\prime}<\frac{c}{a}<a_{1}^{\prime}+a_{2}^{\prime} \frac{b}{a}$ and $a_{1}>0$, which are equivalent to the conditions: $\left(\frac{b}{c}+1\right) /\left(\frac{a}{b}+\frac{b}{a}\right)<2 n+1<\left(\frac{a}{2 c}+1\right) /\left(3-\frac{b}{a}\right)$. These conditions can be satisfied for some constants $a, b, c$ and some $n \geq 1$, but they imply $\frac{a}{b}>2$.

Since any sufficiently small perturbation $\bar{a}$ of the fixed point satisfies $a_{3}^{\prime}=$ $a_{3 k}=\frac{c}{a}$ for $k=1, \ldots, 2 n$, we use the notations $\Sigma_{4 n+2}$ and $\widetilde{\Phi}_{n}$ defined for Case $\# 3$, see (12). The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+2) \times(4 n+2)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$, is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
-1 & 0 & \cdots & \cdots & 0 & 1-\frac{a}{b} & 1+\frac{a}{b} \\
-\frac{b}{a} & 0 & \cdots & \cdots & \cdots & 0 & 1+\frac{b}{a} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \cdots & \cdots & \cdots & 0 \\
\vdots & 0 & \ddots & \ddots & \vdots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right) .
$$

Using Lemma 3.4 with $N=4 n+2, k=2$ and $m=2 n$, we obtain:
Lemma 3.13. The characteristic polynomial of $M$ is equal to:

$$
R(X)=X^{4 n+2}+X^{4 n+1}-(2+B) X^{2 n+1}+B, \quad \text { with } B=\frac{b}{a}-\frac{a}{b}<0 .
$$

The proof of Proposition 3.12 follows by the same arguments as for Proposition 3.3, using Theorem 3.19 or by that of Proposition 3.8 ( $a>2 b$ implies $|B|>1$ ).

### 3.4.7. Case $\# 4^{\prime}$.

Proposition 3.14. Any fixed point of $\Phi_{n}$ with shape $\# 4^{\prime}$ is unstable.
The flow $\Phi_{n}$ near a fixed point with shape $\# 4^{\prime}$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=a_{2}^{\prime}=a_{2}^{\prime \prime}=0$ and

$$
\begin{aligned}
& a_{3}^{\prime \prime}=-a_{3}^{\prime}+a_{1}+\frac{b}{a} a_{2}+a_{3}, \quad a_{6 n+1}=\frac{c}{a}, \\
& a_{6 n+2}=a_{3}^{\prime}\left(1+\frac{a}{b}\right)+a_{1}\left(1-\frac{a}{b}\right)-\frac{c}{a}-\frac{c}{b}, \\
& a_{6 n+3}=a_{3}^{\prime}\left(1-\frac{a}{b}\right)-a_{1}\left(1-\frac{a}{b}\right)+a_{2}\left(1-\frac{b}{a}\right)+\frac{c}{b} .
\end{aligned}
$$

The fixed point is unique and defined by the following equations:

$$
\begin{aligned}
& a_{3}^{\prime}=\frac{1}{d}\left(1-4 n \frac{c}{a}\left(1-\frac{a}{b}\right)\right), \quad a_{1}=\frac{c}{a}, \quad a_{2}=\frac{1}{d}\left(1+\frac{a}{b}-2 \frac{c}{b}-8 n \frac{c}{b}\right) \\
& a_{3}=\frac{1}{d}\left(1+\frac{c}{a}-\frac{b}{a}+2 n \frac{c}{a}\left(-2+3 \frac{a}{b}+\frac{b}{a}\right)\right), \quad d=2 n\left(2+\frac{a}{b}-\frac{b}{a}\right)+1
\end{aligned}
$$

To be compatible with the shape \#4', the fixed point $\bar{a}$ has to satisfy $a_{3}^{\prime} a-$ $a_{1} a>c>a_{3} a-a_{3}^{\prime} a>0$ and $a_{2}>0$ which are equivalent to the conditions: $\left(\frac{b}{c}-1\right) /\left(\frac{a}{b}+\frac{b}{a}\right)<2 n<\left(\frac{a}{2 c}-1\right) /\left(3-\frac{b}{a}\right)$. These conditions can be satisfied for some constants $a, b, c$ and some $n \geq 1$, but they imply $\frac{a}{b}>2$.

Since any sufficiently small perturbation $\bar{a}$ of the fixed point satisfies $a_{1}^{\prime}=$ $a_{2}^{\prime}=0$ and $a_{3 k-2}=\frac{c}{a}$ for $k=1, \ldots, 2 n$, we use the notations $\Sigma_{4 n+1}$ and $\widetilde{\Phi}_{n}$ defined for Case $\# 3^{\prime}$, see (13). The transposed $M^{T}$ of the matrix $M \in$ $\mathbb{R}^{(4 n+1) \times(4 n+1)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$, is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
-1 & 0 & \cdots & \cdots & \cdots & 1+\frac{a}{b} & 1-\frac{a}{b} \\
\frac{b}{a} & 0 & \cdots & \cdots & \cdots & 0 & 1-\frac{b}{a} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \cdots & \cdots & \cdots & 0 \\
\vdots & 0 & \ddots & \ddots & \vdots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right) .
$$

Using Lemma 3.4 with $N=4 n+1, k=2$ and $m=2 n-1$, we obtain:
Lemma 3.15. The characteristic polynomial of $M$ is equal to:

$$
R(X)=-\left(X^{4 n+1}+X^{4 n}-(2+B) X^{2 n}+B\right), \quad \text { with } B=\frac{b}{a}-\frac{a}{b}<0 .
$$

The proof of Proposition 3.14 may follow by the same arguments as for Proposition 3.3, using Theorem 3.20, or by that of Proposition $3.8(a>2 b$ implies $|B|>1)$.

### 3.4.8. Case \#5.

Proposition 3.16. Any fixed point of $\Phi_{n}$ with shape $\# 5$ is unstable.
The flow $\Phi_{n}$ near a fixed point with shape $\# 5$ is given by (8), with $a_{1}^{\prime}=$ $a_{1}^{\prime \prime}=0$ and
$a_{2}^{\prime \prime}=-a_{2}^{\prime}-\frac{a}{b} a_{3}^{\prime}+\frac{a}{b} a_{1}+a_{2}, \quad a_{3}^{\prime \prime}=a_{3}, \quad a_{6 n+1}=a_{2}^{\prime}\left(1-\frac{b}{a}\right)+\frac{c}{a}$,
$a_{6 n+2}=2\left(\frac{b}{a} a_{2}^{\prime}+a_{3}^{\prime}-\frac{c}{a}\right), a_{6 n+3}=\frac{c}{a}+\left(1-\frac{a}{b}\right) a_{1}+\left(1-\frac{b}{a}\right) a_{2}^{\prime}+\left(\frac{a}{b}-1\right) a_{3}^{\prime}$.
If $\frac{a}{b} \neq 2$, the fixed point is unique and verifies:
$a_{2}^{\prime}=\frac{a_{2}}{2}, \quad a_{3}^{\prime}=a_{3}, \quad a_{1}=a_{3}, \quad a_{2}=\frac{1-\frac{c}{a}(4 n+1)}{\left(1-\frac{b}{2 a}\right)(4 n+1)}, \quad a_{3}=\frac{c}{a}+\frac{a_{2}}{2}\left(1-\frac{b}{a}\right)$.
To be compatible with the shape $\# 5$, the fixed point $\bar{a}$ has to satisfy $0<$ $a_{2}^{\prime} b<c<a_{2}^{\prime} b+a_{3}^{\prime} a, 0<\left(a_{2}-a_{2}^{\prime}\right) b<c$ and $a_{3}>0$ which are equivalent to the conditions: $\frac{b}{2 c}<4 n+1<\frac{a}{c}$. Since any sufficiently small perturbation $\bar{a}$ of a fixed point satisfies $a_{1}^{\prime}=0$, we restrict $\Phi_{n}$ to the subspace $\Sigma_{6 n+2}$ of $\Sigma_{6 n+3}$ defined by this condition and we replace $\Phi_{n}$ by the map

$$
\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}, a_{2}, \ldots, a_{6 n}\right) \stackrel{\widetilde{\Phi}_{n}}{\longmapsto}\left(a_{2}^{\prime \prime}, a_{3}^{\prime \prime}, a_{4}, a_{5}, \ldots, a_{6 n+3}\right) .
$$

The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(6 n+2) \times(6 n+2)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
-1 & 0 & \cdots & 0 & 1-\frac{b}{a} & 2 \frac{b}{a} & 1-\frac{b}{a} \\
-\frac{a}{b} & 0 & \cdots & \cdots & 0 & 2 & \frac{a}{b}-1 \\
\frac{a}{b} & 0 & \cdots & \cdots & \cdots & 0 & 1-\frac{a}{b} \\
1 & 0 & \cdots & \cdots & \cdots & \vdots & 0 \\
0 & \ddots & 0 & \vdots & \vdots & \vdots & 0 \\
\vdots & \ddots & \ddots & \ddots & 0 & \vdots & \vdots \\
0 & \cdots & 0 & 1 & 0 & \cdots & 0
\end{array}\right) .
$$

Using Lemma 3.4 with $N=6 n+2, k=3$ and $m=2 n-1$, we get:
Lemma 3.17. The characteristic polynomial of $M$ is equal to:

$$
R(X)=X^{6 n+2}+X^{6 n+1}-(2+B) X^{4 n+1}+B,
$$

where $B=2\left(\frac{a}{b}+\frac{b}{a}-2\right)$ is such that $B>0$.
The proof of Proposition 3.16 follows by the same arguments as for Proposition 3.3, using Theorem 3.21 instead of Theorem 3.19.

### 3.4.9. Case \#6.

Proposition 3.18. Any fixed point of $\Phi_{n}$ with shape $\# 6$ is unstable.
The flow $\Phi_{n}$ near a fixed point with shape $\# 6$ is given by (8), with $a_{1}^{\prime}=$ $a_{1}^{\prime \prime}=0$ and

$$
\begin{aligned}
& a_{2}^{\prime \prime}=-a_{2}^{\prime}-a_{3}^{\prime} \frac{a}{b}+a_{1} \frac{a}{b}+a_{2}, \quad a_{3}^{\prime \prime}=a_{3}, \\
& a_{6 n+2}=2 a_{2}^{\prime}+a_{3}^{\prime}\left(1+\frac{a}{b}\right)+a_{1}\left(1-\frac{a}{b}\right)-2 \frac{c}{b}, \quad a_{6 n+1}=\frac{c}{b}, \quad a_{6 n+3}=\frac{c}{b} .
\end{aligned}
$$

The fixed point is unique and defined by:

$$
a_{2}^{\prime}=\frac{1}{2} a_{2}, \quad a_{3}^{\prime}=\frac{c}{b}, \quad a_{1}=\frac{c}{b}, \quad a_{2}=\frac{2}{4 n+1}-\frac{2 c}{b}, \quad a_{3}=\frac{c}{b}
$$

To be compatible with the shape $\# 6$, the fixed point $\bar{a}$ has to satisfy $c<$ $a_{2}^{\prime} b<a_{2} b-c$ which is equivalent to the condition: $\frac{b}{2 c}>4 n+1$. Since any sufficiently small perturbation $\bar{a}$ of the fixed point satisfies $a_{1}^{\prime}=0$ and $a_{3}^{\prime}=a_{3 k-2}=a_{3 k}=\frac{c}{a}$ for $k=1, \ldots, 2 n$, we restrict $\Phi_{n}$ to the subspace $\Sigma_{2 n+1}$ of $\Sigma_{6 n+3}$ defined by these conditions and we replace $\Phi_{n}$ by the map $\widetilde{\Phi}_{n}:\left(a_{2}^{\prime}, a_{2}, a_{5}, \ldots, a_{6 n-1}\right) \longmapsto\left(a_{2}^{\prime \prime}, a_{5}, \ldots, a_{6 n+2}\right)$.

The flow $\widetilde{\Phi}_{n}$ has an inverse $\widetilde{\Phi}_{n}^{-1}$ given by:

$$
\left(a_{2}^{\prime \prime}, a_{5}, a_{8}, \ldots, a_{6 n+2}\right) \stackrel{\widetilde{\Phi}_{n}^{-1}}{\longrightarrow}\left(a_{2}^{\prime}=\frac{a_{6 n+2}}{2}, a_{2}=a_{2}^{\prime \prime}+\frac{a_{6 n+2}}{2}, a_{5}, \ldots, a_{6 n-1}\right) .
$$

If $M \in \mathbb{R}^{(2 n+1) \times(2 n+1)}$ denotes the matrix associated to the tangent map of $\widetilde{\Phi}$, then $M^{-1}$ is that associated to $\widetilde{\Phi}_{n}^{-1}$. It is equal to:

$$
M^{-1}=\left(\begin{array}{ccccccc}
0 & \cdots & \cdots & \cdots & \cdots & 0 & \frac{1}{2} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & \frac{1}{2} \\
0 & 1 & \ddots & \ldots & \ldots & \ldots & 0 \\
\vdots & \ddots & 1 & 0 & \cdots & \cdots & 0 \\
\vdots & \vdots & 0 & \ddots & \ddots & \vdots & 0 \\
\vdots & \vdots & \vdots & \ddots & 1 & 0 & \vdots \\
0 & \cdots & 0 & \cdots & 0 & 1 & 0
\end{array}\right) .
$$

The matrix $M^{-1}$ has nonnegative coefficients, and the sum of the columns are 1 . Thus, $z=1$ is an eigenvalue associated to an eigenvector with positive coefficients. According to Perron-Frobenius'theorem, see for instance [3], the other roots are of modulus strictly less than 1 . Therefore, any fixed point of $\widetilde{\Phi}_{n}^{-1}$ is stable for the flow $\widetilde{\Phi}_{n}^{-1}$ and unstable for $\widetilde{\Phi}_{n}$. The argument here is related to that given in [7, 2] for the instability of rapidly-oscillating 2 PP solutions of a non linear equation (1) with $h(y)=\operatorname{sgn} y$ (that is $a=b=1$ ). One can also prove, by the same technique as in the previous sections, that the characteristic polynomial of the matrix of the tangent map of $\widetilde{\Phi}_{n}$ is equal to $R(X)=-\left(X^{2 n+1}+X^{2 n}-2\right)$. This polynomial corresponds to the limit case $a=b$ or equivalently $B=0$ in any of the previous polynomials, up to a product by a power of $X$. By Lemma $3.23, R$ has $2 n$ roots with $|z|>1$, and the other root is 1 .
3.5. Roots of the characteristic polynomials. The following results are proved using homotopy techniques.
Theorem 3.19. Assume that $-1 \leq B \leq 1$. The equation

$$
X^{4 n+2}+X^{4 n+1}-(2+B) X^{2 n+1}+B=0
$$

has precisely $2 n$ solutions $z$ (counting algebraic multiplicity) with $|z|>1$.
Theorem 3.20. Assume that $-1 \leq B \leq 1$. The equation

$$
\begin{equation*}
X^{4 n+1}+X^{4 n}-(2+B) X^{2 n}+B=0 \tag{14}
\end{equation*}
$$

has precisely $2 n$ solutions $z$ (counting algebraic multiplicity) with $|z|>1$.
Theorem 3.21. Assume that $0 \leq B \leq 1$. The equation

$$
R(X)=X^{6 n+2}+X^{6 n+1}-(2+B) X^{4 n+1}+B=0
$$

has precisely $2 n$ solutions $z$ (counting algebraic multiplicity) with $|z|>1$.
Let us denote by $R$ the polynomial of Theorem 3.19 and consider the homotopy

$$
\begin{equation*}
f_{t}(z)=Q(z)-t P(z) \tag{15}
\end{equation*}
$$

with

$$
\begin{align*}
P(X) & =X^{2 n}+\cdots+1=\frac{X^{2 n+1}-1}{X-1}  \tag{16}\\
Q(X) & =X^{4 n+1}+2 X^{4 n}+\cdots+2 X^{2 n+1}=X^{2 n+1} \frac{S(X)}{X-1}  \tag{17}\\
S(X) & =X^{2 n+1}+X^{2 n}-2 \tag{18}
\end{align*}
$$

We have $R(X)=(X-1) f_{B}(X)$. So if we denote $D(0,1)=\{z \in \mathbb{C},|z|<1\}$ and $\bar{D}(0,1)^{c}=\{z \in \mathbb{C},|z|>1\}$ the complement of its closure, the roots of $R$ in $D(0,1)$ or $\bar{D}(0,1)^{c}$ are exactly that of $f_{B}$ (with same multiplicity).
Lemma 3.22. Consider $n \geq 1$ and $f_{t}$ defined by ( $15,16,17,18$ ). If $|z|=1$ and $-1<t \leq 1$, then $f_{t}(z) \neq 0$. Moreover, $f_{-1}(z)=0 \Leftrightarrow z=-1$ and $z=-1$ is a simple root of $f_{-1}$.
Proof. Consider $z$ such that $|z|=1$ and put $w=z^{2 n+1}$. If $f_{t}(z)=0$, we get

$$
w^{2}+\frac{1}{z} w^{2}-(2+t) w+t=f_{t}(z)(z-1)=0
$$

or

$$
\begin{equation*}
\frac{1}{z}=\frac{(2+t) w-t-w^{2}}{w^{2}} \tag{19}
\end{equation*}
$$

Taking modulus on both sides of (19) gives $1=\left|w^{2}-(2+t) w+t\right|^{2}$. Writing $w=e^{i \theta}$ and simplifying gives

$$
\begin{equation*}
(\cos \theta-1)\left(4 t \cos \theta-2 t-2 t^{2}-4\right)=0 . \tag{20}
\end{equation*}
$$

Thus, for $t \neq 0$, we either have $\cos \theta=1$ or $\cos \theta=\frac{1}{t}+\frac{t}{2}+\frac{1}{2}=g(t)$. Note that $g(1)=2>1, g(-1)=-1$ and $g^{\prime}(t)=-\frac{1}{t^{2}}+\frac{1}{2}<0$ for $-1 \leq t \leq 1$, thus $g(t)>1$ for $0<t \leq 1$ and $g(t)<-1$ for $-1<t<0$. This implies that the equality $\cos \theta=g(t)$ is impossible for $-1<t \leq 1, t \neq 0$. Then, for $-1<t \leq 1, \cos \theta=1$, so $w=1$. But if $w=z^{2 n+1}=1$ and $f_{t}(z)=0$, we find by (19), $z=1$, which is impossible since $f_{t}(1)=4 n+1-t(2 n+1)>0$ for $n \geq 1$ and $t \leq 1$.

Suppose now that $t=-1$. Using (20), we obtain the two possibilities $\cos \theta=1$ or $\cos \theta=-1$. the first one is impossible as before, whereas the second one implies $z=-1$, by using (19). Conversely, we have $f_{-1}(-1)=0$ and $f_{-1}^{\prime}(-1)=n+1>0$, so that $z=-1$ is a simple root of $f_{-1}$.
Lemma 3.23. The equation $S(z)=z^{2 n+1}+z^{2 n}-2=0$ has no solutions $z$ with $|z|<1,2 n$ solutions $z$ with $|z|>1$ (counting algebraic multiplicity) and 1 is the only root with $|z|=1$.
Proof. For $|z| \leq 1$, we have $\left|z^{2 n+1}+z^{2 n}\right| \leq 2$, with equality if and only if $|z|=1$ and $|z+1|=2$ or equivalently $z=1$. So $z=1$ is the only root of $S$ such that $|z| \leq 1$. It satisfies $S^{\prime}(1)=4 n+1>0$ so that 1 is a simple root. Since $S$ is a polynomial with degree $2 n+1$, this implies that $S$ has exactly $2 n$ roots such that $|z|>1$ (counting multiplicity).
Proof of Theorem 3.19 for $B \neq-1$. The function $f_{t}(z)$ defined in (15), (16), (17), (18) is continuous with respect to $t$ and polynomial (thus analytic)
with respect to $z$. By Lemma $3.23, S$ has no roots in $D(0,1)$. So the only root of $f_{0}=Q$ in $D(0,1)$ is 0 with multiplicity $2 n+1$. This shows that $\operatorname{deg}\left(f_{0}, D(0,1), 0\right)=2 n+1$. From Lemma 3.22, $f_{t}$ has no roots on $\partial D(0,1)$ if $-1<t \leq 1$. Applying Rouché's theorem to the homotopy $f_{t}$ on $D(0,1)$, we get

$$
\operatorname{deg}\left(f_{t}, D(0,1), 0\right)=\operatorname{deg}\left(f_{0}, D(0,1), 0\right)=2 n+1 \quad \text { for }-1<t \leq 1
$$

This implies that the equation $f_{t}(z)=0$ has (counting algebraic multiplicity) $2 n+1$ solutions $z$ in $D(0,1)$. Since $f_{t}$ is a polynomial with degree $4 n+1$, $f_{t}(z)=0$ has $2 n$ solutions in $\bar{D}(0,1)^{c}$.
Proof of Theorem 3.19 for $B=-1$. By Lemma 3.22, $f_{t}$ has no roots on $\partial D(0,1)$ when $-1<t \leq 1$ and -1 is the only root of $f_{-1}$ on $\partial D(0,1)$. Since $f_{-1}^{\prime}(-1) \neq 0$, the implicit function theorem shows that there exists $\delta>0$, $\rho>0$ and a $\mathcal{C}^{1}$ function $Z:(-1-\delta,-1+\delta) \rightarrow \mathbb{C}$, such that $f_{t}(z)=0$ and $|z+1| \leq \rho$ if and only if $z=Z(t)$. So $Z(-1)=-1$ and

$$
Z^{\prime}(-1)=-\frac{\left.\frac{\partial f_{t}(z)}{\partial t}\right|_{t=-1, z=-1}}{\left.\frac{\partial f_{t}(z)}{\partial z}\right|_{t=-1, z=-1}}=\frac{P(-1)}{f_{-1}^{\prime}(-1)}=\frac{1}{n+1}>0
$$

which implies $|Z(t)|<1$ for $-1<t<-1+\delta$ (possibly by shrinking $\delta$ ). Therefore, when $-1<t<-1+\delta, f_{t}(z)=0$ and $|z+1| \leq \rho$ imply $|z|<1$. Denote $G_{\rho}=\{z \in \mathbb{C}$ such that $|z|<1$ or $|z+1|<\rho\}$. The properties of $Z$ imply that for $-1 \leq t<-1+\delta, f_{t}$ has no roots on $\partial G_{\rho}$, and that for $-1<t<-1+\delta$, the number of roots of $f_{t}$ in $G_{\rho}$ is the same as in $D(0,1)$, that is $2 n+1$. Applying Rouché's theorem to the homotopy $f_{t}$ on $G_{\rho}$, we get

$$
\operatorname{deg}\left(f_{-1}, G_{\rho}, 0\right)=\operatorname{deg}\left(f_{t}, G_{\rho}, 0\right)=2 n+1 \quad \text { for }-1<t<-1+\delta
$$

so that $f_{-1}$ has $2 n+1$ roots in $G_{\rho}$ and $2 n$ roots outside the closure of $G_{\rho}$ ( $f_{t}$ is a polynomial with degree $4 n+1$ ). Since $f_{-1}$ has no roots such that $0<|z+1| \leq \rho$, it has exactly $2 n$ roots in $\bar{D}(0,1)^{c}$.

Theorem 3.20 is proved by the same technique as for the case $B \neq 1$ of Theorem 3.19, changing $P$ and $Q$ into:

$$
\begin{align*}
& P(X)=X^{2 n-1}+\cdots+1=\frac{X^{2 n}-1}{X-1}  \tag{21}\\
& Q(X)=X^{4 n}+2 X^{4 n-1}+\cdots+2 X^{2 n}=X^{2 n} \frac{S(X)}{X-1} . \tag{22}
\end{align*}
$$

Indeed, by the same technique as for Lemma 3.22, one proves:

Lemma 3.24. Consider $n \geq 1$ and $f_{t}$ defined by (15), (21), (22), (18). If $|z|=1$ and $-1 \leq t \leq 1$, then $f_{t}(z) \neq 0$.
Proof. Consider $z$ such that $|z|=1$ and $f_{t}(z)=0$, and put $w=z^{2 n}$. We obtain

$$
z=\frac{(2+t) w-t-w^{2}}{w^{2}}
$$

which leads, as in Lemma 3.22, to the only possibilities $w=z=1$, or $w=z=-1$ and $t=-1$. Since $f_{t}(1)=4 n+1-2 n t>0$ and $f_{-1}(-1) \neq 0$ (indeed $w=z^{2 n}=1$ for $z=-1$ ), we obtain the result of the lemma.

The arguments for the proof of Theorem 3.21 are essentially the same as for Theorem 3.19, changing $P$ and $Q$ into:

$$
\begin{align*}
& P(X)=X^{4 n}+\cdots+1=\frac{X^{4 n+1}-1}{X-1}  \tag{23}\\
& Q(X)=X^{6 n+1}+2 X^{6 n}+\cdots+2 X^{4 n+1}=X^{4 n+1} \frac{S(X)}{X-1} . \tag{24}
\end{align*}
$$

Lemma 3.25. Consider $n \geq 1$, and $P$ and $Q$ defined by (23), (24), (18). If $|z|=1$ and $0 \leq t<1$, then $t|P(z)|<|Q(z)|$.
Proof. Indeed

$$
\begin{equation*}
\left(|Q(z)|^{2}-|P(z)|^{2}\right)|z-1|^{2}=\left|z^{2 n}-1\right|^{2}\left|z^{2 n+1}-1\right|^{2} \geq 0 \quad \text { if }|z|=1 \tag{25}
\end{equation*}
$$

and $Q(1)=4 n+1=P(1)$, so that

$$
\begin{equation*}
|P(z)| \leq|Q(z)| \quad \text { for }|z|=1 \tag{26}
\end{equation*}
$$

Moreover, $Q(z)=0$ and $|z|=1$ imply $S(z)=0$. But by Lemma 3.23, the only root of $S$ with $|z|=1$ is $z=1$ for which $Q(1) \neq 0$. Thus, $Q(z) \neq 0$ for all $z$ such that $|z|=1$, which with (26) leads to the result of the lemma.
Proof of Theorem $\mathbf{3 . 2 1}$ for $B \neq 1$. Let $f_{t}$ be defined by $(15,23,24,18)$. From Lemma 3.25, one deduces by Rouché's theorem that $f_{B}$ and $f_{0}$ have same number of roots in $D(0,1)$ (counting multiplicity). By Lemma 3.23, 0 is the only root of $f_{0}=Q$ in $D(0,1)$ and it has multiplicity $4 n+1$. Since $f_{B}$ is a polynomial with degree $6 n+1$ and has no roots on $\partial D(0,1)$ (due to Lemma 3.25), $f_{B}$ has $2 n$ roots in $\bar{D}(0,1)^{c}$. So does $R=(X-1) f_{B}$.
Proof of Theorem $\mathbf{3 . 2 1}$ for $B=1$. We use the same technique as for the case $B=-1$ of Theorem 3.19. By Lemma 3.25, $f_{t}$ has no roots on $\partial D(0,1)$ for $0 \leq t<1$. For $t=1, f_{1}(z)=0$ implies, by $(25), z^{2 n}=1$ or $z^{2 n+1}=1$. Computing $f_{1}(z)$ for such a $z$, one obtains a contradiction, except for $z=1$ which is indeed the only root of $f_{1}$ on $\partial D(0,1)$.

Since $f_{1}^{\prime}(1)=12 n^{2}+6 n+1>0$, the implicit function theorem shows that there exists $\delta>0, \rho>0$ and a $\mathcal{C}^{1}$ function $Z:(1-\delta, 1+\delta) \rightarrow \mathbb{C}$, such that $f_{t}(z)=0$ and $|z-1| \leq \rho$ if and only if $z=Z(t)$. So $Z(1)=1$ and

$$
Z^{\prime}(1)=-\frac{\left.\frac{\partial f_{t}(z)}{\partial t}\right|_{t=1, z=1}}{\left.\frac{\partial f_{t}(z)}{\partial z}\right|_{t=1, z=1}}=\frac{P(1)}{f_{1}^{\prime}(1)}>0
$$

which implies $|Z(t)|<1$ for $1-\delta<t<1$ (possibly by shrinking $\delta$ ). Denote $G_{\rho}=\{z \in \mathbb{C}$ such that $|z|<1$ or $|z-1|<\rho\}$. We obtain that for $1-\delta<$ $t \leq 1, f_{t}$ has no roots on $\partial G_{\rho}$, and that for $1-\delta<t<1$, the number of roots of $f_{t}$ in $G_{\rho}$ is the same as in $D(0,1)$, that is $4 n+1$. So by Rouché's theorem, the number of roots of $f_{1}$ in $G_{\rho}$ is also $4 n+1$, and then $f_{1}$ has $2 n$ roots outside the closure of $G_{\rho}$. Since $f_{1}$ has no roots such that $0<|z-1| \leq \rho$, $f_{1}$ has exactly $2 n$ roots in $\bar{D}(0,1)^{c}$.

$$
\text { 4. STUDY OF } \dot{x}(t)=-h(x(t-1))+f(x(t)) \text {. }
$$

The following result generalizes [1, Theorem 1].
Theorem 4.1. Let $f$ be a Lebesgue measurable function such that

$$
\text { ess sup }|f(x)|<b
$$

and let $h$ be given by (2). Then, for any $x_{0} \in \mathcal{C}([-1,0])$, such that ${ }^{1}$ meas $\{t \in$ $\left.[-1,0]: x_{0}(t)=0\right\}=0$, there exists at least one function $x \in \mathcal{C}([-1,+\infty))$, absolutely continuous on $[0,+\infty)$, such that the composition $f \circ x$ is Lebesgue measurable and satisfying (1) almost everywhere. For any Lebesgue measurable function $\tilde{f}$ such that $\tilde{f}(x)=f(x)$ a.e. in $\mathbb{R}, x$ satisfies $\tilde{f}(x(t))=f(x(t))$ a.e. By definition, such a function $x$ is called a solution of Equation (1) on $[0,+\infty)$.
Sketch of Proof. We proceed by the same way as in [1]. We prove the existence of a solution of Equation (1) on $[0,1]$. The global existence is obtained by induction.

Let us denote $g(t)=-h(x(t-1))$ for $t \in[0,1]$. Since $g$ is measurable on $[0,1]$ and takes its values in $\{-a,-b, a, b\}$ a.e., one can construct a sequence $\left(g_{n}\right)_{n \geq 1}$ of piecewise constant functions $[0,1] \rightarrow\{-a,-b, b, a\}$ such that $\| g_{n}-$ $g \|_{L^{1}}<\frac{C}{n}$, with $C>0$. Indeed, let $\alpha_{i}, 1 \leq i \leq 4$, be the elements of $\{-a,-b, b, a\}$. Define $A_{i}=g^{-1}\left(\alpha_{i}\right)$ for $1 \leq i \leq 4$. Since $A_{1}$ is a measurable

[^1]set, there exists a set $O_{1}^{n}$ such that $\lambda\left(O_{1}^{n} \Delta A_{1}\right)<\frac{1}{n}$, where $\lambda$ is the Lebesgue measure, and $O_{1}^{n}$ is a finite union of open intervals. We set $A_{2}^{\prime}=\left(A_{1} \cup\right.$ $\left.A_{2}\right) \backslash O_{1}^{n}$ and for $i=3,4, A_{i}^{\prime}=A_{i} \backslash O_{1}^{n}$. We have $\lambda\left(A_{i}^{\prime} \Delta A_{i}\right)<\frac{1}{n}$ for $i=2,3,4$. In the same way, we modify successively $A_{i}^{\prime}, i=2,3,4$ to find the associated sets $O_{i}^{n}, i=2,3,4$, that are all finite unions of intervals and such that $\lambda\left(O_{i}^{n} \Delta A_{i}\right)<\frac{3}{n}$. Then, $g_{n}=\sum_{i=1}^{4} \alpha_{i} 1_{O_{i}^{n}}$ is a piecewise constant function such that $\left\|g_{n}-g\right\|_{L^{1}}<\frac{3 \sum_{i=1}^{4}\left|\alpha_{i}\right|}{n}$.

For any $n \geq 1$, the following equation

$$
\dot{x}(t)=g_{n}(t)+f(x(t)), \quad t \in[0,1]
$$

admits a unique solution $x_{n}$ which is continuous and piecewise $\mathcal{C}^{1}$ in $[0,1]$. To prove this result, let us define the increasing functions

$$
\begin{equation*}
\mathcal{F}_{d}(e)=\int_{0}^{e} \frac{\operatorname{sgn} d}{d+f(x)} d x, \quad e \in \mathbb{R}, d \in\{-a,-b, b, a\} \tag{27}
\end{equation*}
$$

The functions $\mathcal{F}_{d}$ are Lipschitz continuous and such that $\mathcal{F}_{d}^{\prime}(x)=\frac{\operatorname{sgn} d}{d+f(x)}$ a.e. Moreover, for $d \in\{-a,-b, b, a\}$, we get
$0<\frac{1}{a+\operatorname{esssup}_{x \in \mathbb{R}}|f(x)|} \leq \mathcal{F}_{d}^{\prime}(e) \leq \frac{1}{b-\operatorname{ess} \sup }{ }_{x \in \mathbb{R}}|f(x)|, \quad$ for a.e. $e \in \mathbb{R}$.
Hence $\mathcal{F}_{d}$ is invertible with Lipschitz continuous inverse. We deduce, by the same way as in [1] that on any interval $\left[t_{1}, t_{2}\right]$ on which $g_{n}(t) \equiv d, x_{n}$ is given by

$$
\begin{equation*}
x_{n}(t)=\mathcal{F}_{d}^{-1}\left(\mathcal{F}_{d}\left(x_{n}\left(t_{1}\right)\right)+\left(t-t_{1}\right) \operatorname{sgn} d\right) \tag{28}
\end{equation*}
$$

The rest of the proof is the same as for [1, Theorem 1]: there exists a limit point $x$ of the sequence $x_{n}$ in $\mathcal{C}([0,1])$, this limit point is such that $f \circ x$ is measurable, $x$ is solution of (1) and $f \circ x=\tilde{f} \circ x$ a.e. when $f=\tilde{f}$ a.e.
Theorem 4.2. Let us consider Equation (1), where $h$ is given by (2), $f$ satisfies (3) and $x_{0} \in \mathcal{C}([-1,0])$ has finitely many zeros. For any $x_{0}$ with a finite number of zeros, there exists a unique solution $x$ of (1). For such a solution, $V(t)$ is nonincreasing and even-valued. For any periodic solution of (1), $V(t)$ is constant. Rapidly-oscillating symmetric 2 -phase periodic solutions $x$, such that $\|x\|_{\infty}>c$, have qualitatively the same pictures as Shapes $\# 2, \# 2^{\prime}, \# 3, \# 3^{\prime}, \# 4, \# 4^{\prime}, \# 5$ and $\# 6$, defined above except that straight lines are replaced by curves. These periodic solutions are all unstable.

Proof. Existence and uniqueness of a solution of (1) follows by Theorem 4.1 and its proof, see (28). For the evolution of the number $V$ of zeros, see Section 3.1.

For the rest of the proof, we use the notations $\bar{a}, \Sigma_{6 n+3}$ and $\Phi_{n}$ of Section 3.3. The properties stated in Section 3.3 still hold here, except that the flow $\Phi_{n}$ near a fixed point is no more affine, and that $\left(a_{1}^{\prime \prime}, a_{2}^{\prime \prime}, a_{3}^{\prime \prime}\right)$ is determined by the equation:

$$
\begin{aligned}
& \mathcal{F}_{-a}^{-1}\left(\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\left(\mathcal{F}_{-b} \mathcal{F}_{-a}^{-1}\left(-a_{1}^{\prime}\right)-a_{2}^{\prime}\right)-a_{3}^{\prime}\right)=-\underline{x} \\
& \quad=\mathcal{F}_{a}^{-1}\left(\mathcal{F}_{a} \mathcal{F}_{b}^{-1}\left(\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\left(a_{3}^{\prime \prime}-a_{3}\right)+a_{2}^{\prime \prime}-a_{2}\right)+a_{1}^{\prime \prime}-a_{1}\right)
\end{aligned}
$$

which replaces (9), and where the functions $\mathcal{F}_{d}$, for $d \in\{a, b,-a,-b\}$, are defined in (27). Since $f$ is odd, we have $\mathcal{F}_{d}(e)=-\mathcal{F}_{-d}(-e)$. Thus, the flow $\Phi_{n}$ is independent of the sign of $x$ on $\left(t, t^{\prime}\right)$, where $t$ and $t^{\prime}$ are defined in Section 3.3.

In Section 4.1 below, we determine as in Section 3.4, the flow $\Phi_{n}$ corresponding to each shape (using the functions $\mathcal{F}_{d}$ ) and prove that any fixed point (if it exists) is unstable. We do not give the equations of the fixed point, since they cannot be solved explicitly as before. We only determine $\Phi_{n}$, the matrix $M$ associated to the tangent map of its restriction $\widetilde{\Phi}_{n}$ at some fixed point and the characteristic polynomial $R$ of $M$. The restriction $\widetilde{\Phi}_{n}$ is defined in the same way as in the linear case and takes its values in a similar subspace of $\Sigma_{6 n+3}$. As for the linear case, we shall prove that $R$ has at least one root with modulus strictly greater than 1 , by using homotopy techniques. We gather all results concerning roots of characteristic polynomials into Section 4.2.

### 4.1. Study of $\Phi_{n}$ for all cases of generalized shape.

4.1.1. Generalized cases $\# 1$ and $\# 1^{\prime}$.

Lemma 4.3. Generalized shapes $\# 1$ and $\# 1^{\prime}$ are impossible.
Proof. We proceed as for Lemma 3.2 to show that $\Phi_{n}$ has no fixed point with shape $\# 1$ or $\# 1^{\prime}$.

### 4.1.2. Generalized case \#2.

Proposition 4.4. Any fixed point of $\Phi_{n}$ with generalized shape $\# 2$ is unstable.

The flow $\Phi_{n}$ near a fixed point with generalized shape $\# 2$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=0$ and

$$
\begin{aligned}
& a_{2}^{\prime \prime}=a_{2}+\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right), \quad a_{3}^{\prime \prime}=a_{3}, \quad a_{6 n+1}=-\mathcal{F}_{-b}(-c), \\
& a_{6 n+2}=a_{2}^{\prime}+a_{3}^{\prime}+\mathcal{F}_{-b}(-c)+\mathcal{F}_{a}(-c)-\mathcal{B}\left(a_{2}^{\prime}, a_{3}^{\prime}\right), \\
& a_{6 n+3}=a_{1}-\mathcal{F}_{a}(-c)+\mathcal{B}\left(a_{2}^{\prime}, a_{3}^{\prime}\right)-\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right),
\end{aligned}
$$

where

$$
\begin{aligned}
& \mathcal{B}\left(a_{2}^{\prime}, a_{3}^{\prime}\right)=\mathcal{F}_{a} \mathcal{F}_{-a}^{-1}\left(-a_{3}^{\prime}+\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\left(-a_{2}^{\prime}\right)\right) \\
& \mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right)=\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\left(a_{1}+\mathcal{B}\left(a_{2}^{\prime}, a_{3}^{\prime}\right)\right)
\end{aligned}
$$

Any point $\bar{a}$ in the image of the $2 n$-th iterate $\left(\Phi_{n}\right)^{2 n}$ of $\Phi_{n}$, and so any sufficiently small perturbation $\bar{a}$ of a fixed point of $\Phi_{n}$, satisfies $a_{1}^{\prime}=0$ and $a_{3 k-2}=-\mathcal{F}_{-b}(-c)$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+2}$ of $\Sigma_{6 n+3}$ defined by these conditions and use the notation $\widetilde{\Phi}_{n}$ defined for Case $\# 2$, see Section 3.4.2. The transposed $M^{T}$ of the matrix $M \in$ $\mathbb{R}^{(4 n+2) \times(4 n+2)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ at any point $\bar{a} \in \Sigma_{4 n+2}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
\mathcal{A}_{a_{2}^{\prime}}^{\prime} & 0 & \cdots & \cdots & 0 & 1-\mathcal{B}_{a_{2}^{\prime}}^{\prime} & \mathcal{B}_{a_{2}^{\prime}}^{\prime}-\mathcal{A}_{a_{2}^{\prime}}^{\prime} \\
\mathcal{A}_{a_{3}^{\prime}}^{\prime} & 0 & \cdots & \cdots & \cdots & 1-\mathcal{B}_{a_{3}^{\prime}}^{\prime} & \mathcal{B}_{a_{3}^{\prime}}^{\prime}-\mathcal{A}_{a_{3}^{\prime}}^{\prime} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \vdots & \vdots & \vdots & 0 \\
\vdots & 0 & \ddots & \vdots & \ddots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right)
$$

where for any function $f, f_{x}^{\prime}$ denotes the partial derivative of $f$ with respect to $x$.

Lemma 4.5. The characteristic polynomial of $M$ is equal to:

$$
R(X)=X^{4 n+2}+A X^{4 n+1}-(1+A+B) X^{2 n+1}+B
$$

where $A=-\mathcal{A}_{a_{2}^{\prime}}^{\prime}>0$ and $B=\left(\mathcal{A}_{a_{1}}^{\prime}-1\right)\left(\mathcal{B}_{a_{2}^{\prime}}^{\prime}-\mathcal{B}_{a_{3}^{\prime}}^{\prime}\right)>0$.
Proof. We use the same arguments as for Lemma 3.5 and the property that $\mathcal{A}_{a_{i}^{\prime}}^{\prime}=\mathcal{A}_{a_{1}}^{\prime} \mathcal{B}_{a_{i}^{\prime}}^{\prime}$ for $i=2,3$. In particular, $B=\left(\mathcal{A}_{a_{1}}^{\prime}-1\right)\left(\mathcal{B}_{a_{2}^{\prime}}^{\prime}-\mathcal{B}_{a_{3}^{\prime}}^{\prime}\right)=$ $\left(\left(\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\right)^{\prime}(x)-1\right) \mathcal{B}_{a_{3}^{\prime}}^{\prime}\left(\left(\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\right)^{\prime}\left(-a_{2}^{\prime}\right)-1\right)$ for some $x \in \mathbb{R}$. Since $\mathcal{B}_{a_{3}^{\prime}}^{\prime}<0$ and

$$
\begin{gathered}
\left(\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\right)^{\prime}(x)=\frac{\mathcal{F}_{b}^{\prime} \mathcal{F}_{a}^{-1}(x)}{\mathcal{F}_{a}^{\prime} \mathcal{F}_{a}^{-1}(x)}=\frac{a+f\left(\mathcal{F}_{a}^{-1}(x)\right)}{b+f\left(\mathcal{F}_{a}^{-1}(x)\right)}>1 \\
\left(\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\right)^{\prime}(x)=\frac{b-f\left(\mathcal{F}_{-b}^{-1}(x)\right)}{a-f\left(\mathcal{F}_{-b}^{-1}(x)\right)}<1
\end{gathered}
$$

for any $x \in \mathbb{R}$ when $a>b$, we obtain $B>0$. Indeed, as in Lemma 3.5, we also obtain that $B>0$ when $a<b$ (and ess sup $\left.{ }_{x \in \mathbb{R}}|f(x)|<a\right)$.

The proof of Proposition 4.4 follows by the same arguments as for Proposition 3.3 , using Theorem 4.20 below instead of Theorem 3.19.

### 4.1.3. Generalized case $\# 2^{\prime}$.

Proposition 4.6. Any fixed point of $\Phi_{n}$ with generalized shape $\# 2^{\prime}$ is unstable.

The flow $\Phi_{n}$ near a fixed point with generalized shape $\# 2^{\prime}$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=0$ and

$$
\begin{aligned}
& a_{2}^{\prime \prime}=a_{2}+\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right), \quad a_{3}^{\prime \prime}=a_{3}, \quad a_{6 n+1}=a_{2}^{\prime}+\mathcal{B}\left(a_{2}^{\prime}\right)-\mathcal{F}_{-a}(-c) \\
& a_{6 n+2}=a_{3}^{\prime}+a_{1}-\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right)-\mathcal{B}\left(a_{2}^{\prime}\right)+\mathcal{F}_{-a}(-c)+\mathcal{F}_{b}(-c) \\
& a_{6 n+3}=-\mathcal{F}_{b}(-c)
\end{aligned}
$$

where

$$
\mathcal{B}\left(a_{2}^{\prime}\right)=\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\left(-a_{2}^{\prime}\right), \mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right)=\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\left(a_{1}+\mathcal{F}_{a} \mathcal{F}_{-a}^{-1}\left(-a_{3}^{\prime}+\mathcal{B}\left(a_{2}^{\prime}\right)\right)\right)
$$

Any point $\bar{a}$ in the image of $\left(\Phi_{n}\right)^{2 n+1}$, and so any sufficiently small perturbation $\bar{a}$ of a fixed point of $\Phi_{n}$, satisfies $a_{1}^{\prime}=0$ and $a_{3}^{\prime}=a_{3 k}=-\mathcal{F}_{b}(-c)$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+1}$ of $\Sigma_{6 n+3}$ defined by these conditions and use the notation $\widetilde{\Phi}_{n}$ defined for Case $\# 2^{\prime}$, see Section 3.4.3. The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+1) \times(4 n+1)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ at any point $\bar{a} \in \Sigma_{4 n+1}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
\mathcal{A}_{a_{2}^{\prime}}^{\prime} & 0 & \cdots & \cdots & 0 & 1+\mathcal{B}_{a_{2}^{\prime}}^{\prime} & -\mathcal{B}_{a_{2}^{\prime}}^{\prime}-\mathcal{A}_{a_{2}^{\prime}}^{\prime} \\
\mathcal{A}_{a_{1}}^{\prime} & 0 & \cdots & \cdots & \cdots & 0 & 1-\mathcal{A}_{a_{1}}^{\prime} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \vdots & \vdots & \vdots & 0 \\
\vdots & 0 & \ddots & \vdots & \ddots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right)
$$

Lemma 4.7. The characteristic polynomial of $M$ is equal to:

$$
R(X)=-\left(X^{4 n+1}+A X^{4 n}-(1+A+B) X^{2 n}+B\right)
$$

where $A=-\mathcal{A}_{a_{2}^{\prime}}^{\prime}>0$ and $B=\left(\mathcal{A}_{a_{1}}^{\prime}-1\right)\left(1+\mathcal{B}_{a_{2}^{\prime}}^{\prime}\right)>0$.
Here also we obtain $B>0$ for any $a \neq b$, with ess sup $|f|<\min (a, b)$.
The proof of Proposition 4.6 follows by the same arguments as for Proposition 3.3 , using Theorem 4.21 below.

### 4.1.4. Generalized case \#3.

Proposition 4.8. Any fixed point of $\Phi_{n}$ with generalized shape $\# 3$ is unstable.

The flow $\Phi_{n}$ near a fixed point with generalized shape $\# 3$ is given by (8) and

$$
\begin{aligned}
& a_{1}^{\prime \prime}=a_{1}+\mathcal{A}\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}\right), \quad a_{2}^{\prime \prime}=a_{2}, \quad a_{3}^{\prime \prime}=a_{3}, \\
& a_{6 n+1}=a_{1}^{\prime}+a_{2}^{\prime}+\mathcal{B}\left(a_{1}^{\prime}, a_{2}^{\prime}\right)-\mathcal{F}_{-a}(-c), \\
& a_{6 n+2}=a_{3}^{\prime}+\mathcal{F}_{-a}(-c)-\mathcal{B}\left(a_{1}^{\prime}, a_{2}^{\prime}\right)-\mathcal{A}\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}\right)+\mathcal{F}_{a}(-c), \\
& a_{6 n+3}=-\mathcal{F}_{a}(-c),
\end{aligned}
$$

where

$$
\begin{aligned}
\mathcal{B}\left(a_{1}^{\prime}, a_{2}^{\prime}\right) & =\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\left(-a_{2}^{\prime}+\mathcal{F}_{-b} \mathcal{F}_{-a}^{-1}\left(-a_{1}^{\prime}\right)\right), \\
\mathcal{A}\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}\right) & =\mathcal{F}_{a} \mathcal{F}_{-a}^{-1}\left(-a_{3}^{\prime}+\mathcal{B}\left(a_{1}^{\prime}, a_{2}^{\prime}\right)\right) .
\end{aligned}
$$

Any sufficiently small perturbation $\bar{a}$ of a fixed point of $\Phi_{n}$ satisfies $a_{3}^{\prime}=$ $a_{3 k}=-\mathcal{F}_{a}(-c)$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+2}$ of $\Sigma_{6 n+3}$ defined by these conditions and use the notation $\widetilde{\Phi}_{n}$ defined for Case $\# 3$, see (12). The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+2) \times(4 n+2)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ at any point $\bar{a} \in \Sigma_{4 n+2}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
\mathcal{A}_{a_{1}^{\prime}}^{\prime} & 0 & \cdots & \cdots & 0 & 1+\mathcal{B}_{a_{1}^{\prime}}^{\prime} & -\mathcal{B}_{a_{1}^{\prime}}^{\prime}-\mathcal{A}_{a_{1}^{\prime}}^{\prime} \\
\mathcal{A}_{a_{2}^{\prime}}^{\prime} & 0 & \cdots & \cdots & \cdots & 1+\mathcal{B}_{a_{2}^{\prime}}^{\prime} & -\mathcal{B}_{a_{2}^{\prime}}^{\prime}-\mathcal{A}_{a_{2}^{\prime}}^{\prime} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \vdots & \vdots & \vdots & 0 \\
\vdots & 0 & \ddots & \vdots & \ddots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right) .
$$

Lemma 4.9. The characteristic polynomial of $M$ is equal to:

$$
R(X)=X^{4 n+2}+A X^{4 n+1}-(1+A+B) X^{2 n+1}+B,
$$

with $A=-\mathcal{A}_{a_{1}^{\prime}}^{\prime}>0$ and $B=\left(1-\mathcal{A}_{a_{3}^{\prime}}^{\prime}\right)\left(\mathcal{B}_{a_{1}^{\prime}}^{\prime}-\mathcal{B}_{a_{2}^{\prime}}^{\prime}\right)>0$.
Here we obtain $B>0$ for $a>b$ and $B<0$ for $a<b$ (with ess sup $|f|<$ $\min (a, b))$.

The proof of Proposition 4.8 follows by the same arguments as for Proposition 4.4, that are the arguments of Proposition 3.3, using Theorem 4.20 below.

### 4.1.5. Generalized case $\# 3^{\prime}$.

Proposition 4.10. Any fixed point of $\Phi_{n}$ with generalized shape $\# 3^{\prime}$ is unstable.

The flow $\Phi_{n}$ near a fixed point with generalized shape $\# 3^{\prime}$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=a_{2}^{\prime}=a_{2}^{\prime \prime}=0$ and

$$
\begin{aligned}
& a_{3}^{\prime \prime}=a_{3}+\mathcal{A}\left(a_{3}^{\prime}, a_{1}, a_{2}\right), \quad a_{6 n+1}=-\mathcal{F}_{-a}(-c), \\
& a_{6 n+2}=a_{3}^{\prime}-\mathcal{B}\left(a_{3}^{\prime}\right)+\mathcal{F}_{-a}(-c)+\mathcal{F}_{a}(-c), \\
& a_{6 n+3}=a_{1}+a_{2}-\mathcal{A}\left(a_{3}^{\prime}, a_{1}, a_{2}\right)+\mathcal{B}\left(a_{3}^{\prime}\right)-\mathcal{F}_{a}(-c),
\end{aligned}
$$

where

$$
\mathcal{B}\left(a_{3}^{\prime}\right)=\mathcal{F}_{a} \mathcal{F}_{-a}^{-1}\left(-a_{3}^{\prime}\right), \quad \mathcal{A}\left(a_{3}^{\prime}, a_{1}, a_{2}\right)=\mathcal{F}_{a} \mathcal{F}_{b}^{-1}\left(a_{2}+\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\left(a_{1}+\mathcal{B}\left(a_{3}^{\prime}\right)\right)\right) .
$$

Any sufficiently small perturbation $\bar{a}$ of a fixed point of $\Phi_{n}$ satisfies $a_{1}^{\prime}=$ $a_{2}^{\prime}=0$ and $a_{3 k-2}=-\mathcal{F}_{-a}(-c)$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+1}$ of $\Sigma_{6 n+3}$ defined by these conditions and use the notation $\widetilde{\Phi}_{n}$ defined for Case $\# 3^{\prime}$, see (13). The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+1) \times(4 n+1)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ at any point $\bar{a} \in \Sigma_{4 n+1}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
\mathcal{A}_{a_{3}^{\prime}}^{\prime} & 0 & \cdots & \cdots & 0 & 1-\mathcal{B}_{a_{3}^{\prime}}^{\prime} & \mathcal{B}_{a_{3}^{\prime}}^{\prime}-\mathcal{A}_{a_{3}^{\prime}}^{\prime} \\
\mathcal{A}_{a_{2}}^{\prime} & 0 & \cdots & \cdots & \cdots & 0 & 1-\mathcal{A}_{a_{2}}^{\prime} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \vdots & \vdots & \vdots & 0 \\
\vdots & 0 & \ddots & \vdots & \ddots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right)
$$

Lemma 4.11. The characteristic polynomial of $M$ is equal to:

$$
R(X)=-\left(X^{4 n+1}+A X^{4 n}-(1+A+B) X^{2 n}+B\right)
$$

where $A=-\mathcal{A}_{a_{3}^{\prime}}^{\prime}>0$ and $B=\left(\mathcal{A}_{a_{2}}^{\prime}-1\right)\left(1-\mathcal{B}_{a_{3}^{\prime}}^{\prime}\right)<0$.
Indeed, $B$ has same sign as $b-a$ (when ess sup $|f|<\min (a, b))$.
The proof of Proposition 4.10 follows by the same arguments as for Proposition 4.6, that are the arguments of Proposition 3.3 using Theorem 4.21 below.

### 4.1.6. Generalized case \#4.

Proposition 4.12. Any fixed point of $\Phi_{n}$ with generalized shape $\# 4$ is unstable.

The flow $\Phi_{n}$ near a fixed point with generalized shape \#4 is given by (8) and
$a_{1}^{\prime \prime}=a_{1}+\mathcal{A}\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}\right), a_{2}^{\prime \prime}=a_{2}, a_{3}^{\prime \prime}=a_{3}, \quad a_{6 n+1}=a_{1}^{\prime}+\mathcal{B}\left(a_{1}^{\prime}\right)-\mathcal{F}_{-b}(-c)$,
$a_{6 n+2}=a_{2}^{\prime}+a_{3}^{\prime}+\mathcal{F}_{-b}(-c)-\mathcal{B}\left(a_{1}^{\prime}\right)+\mathcal{F}_{a}(-c)-\mathcal{A}\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}\right)$,
$a_{6 n+3}=-\mathcal{F}_{a}(-c)$,
where
$\mathcal{B}\left(a_{1}^{\prime}\right)=\mathcal{F}_{-b} \mathcal{F}_{-a}^{-1}\left(-a_{1}^{\prime}\right), \mathcal{A}\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}\right)=\mathcal{F}_{a} \mathcal{F}_{-a}^{-1}\left(-a_{3}^{\prime}+\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\left(-a_{2}^{\prime}+\mathcal{B}\left(a_{1}^{\prime}\right)\right)\right)$.
Any sufficiently small perturbation $\bar{a}$ of a fixed point of $\Phi_{n}$ satisfies $a_{3}^{\prime}=$ $a_{3 k}=-\mathcal{F}_{a}(-c)$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+2}$ of $\Sigma_{6 n+3}$ defined by these conditions and use the notation $\widetilde{\Phi}_{n}$ defined for Case $\# 3$, see (12). The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+2) \times(4 n+2)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ at any point $\bar{a} \in \Sigma_{4 n+2}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
\mathcal{A}_{a_{1}^{\prime}}^{\prime} & 0 & \cdots & \cdots & 0 & 1+\mathcal{B}_{a_{1}^{\prime}}^{\prime} & -\mathcal{B}_{a_{1}^{\prime}}^{\prime}-\mathcal{A}_{a_{1}^{\prime}}^{\prime} \\
\mathcal{A}_{a_{2}^{\prime}}^{\prime} & 0 & \cdots & \cdots & \cdots & 0 & 1-\mathcal{A}_{a_{2}^{\prime}}^{\prime} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \vdots & \vdots & \vdots & 0 \\
\vdots & 0 & \ddots & \vdots & \ddots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right) .
$$

Lemma 4.13. The characteristic polynomial of $M$ is equal to:

$$
R(X)=X^{4 n+2}+A X^{4 n+1}-(1+A+B) X^{2 n+1}+B
$$

with $A=-\mathcal{A}_{a_{1}^{\prime}}^{\prime}>0$ and $B=\left(1+\mathcal{B}_{a_{1}^{\prime}}^{\prime}\right)\left(1-\mathcal{A}_{a_{2}^{\prime}}^{\prime}\right)<0$.
Indeed, $B$ has same sign as $b-a$ (when ess sup $|f|<\min (a, b)$ ).
The proof of Proposition 4.12 follows by the same arguments as for Proposition 4.4.

### 4.1.7. Generalized case \#4'.

Proposition 4.14. Any fixed point of $\Phi_{n}$ with generalized shape $\# 4^{\prime}$ is unstable.

The flow $\Phi_{n}$ near a fixed point with generalized shape $\# 4^{\prime}$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=a_{2}^{\prime}=a_{2}^{\prime \prime}=0$ and

$$
\begin{aligned}
& a_{3}^{\prime \prime}=a_{3}+\mathcal{A}\left(a_{3}^{\prime}, a_{1}, a_{2}\right), \quad a_{6 n+1}=-\mathcal{F}_{-a}(-c), \\
& a_{6 n+2}=a_{3}^{\prime}+a_{1}-\mathcal{B}\left(a_{3}^{\prime}, a_{1}\right)+\mathcal{F}_{-a}(-c)+\mathcal{F}_{b}(-c), \\
& a_{6 n+3}=a_{2}-\mathcal{A}\left(a_{3}^{\prime}, a_{1}, a_{2}\right)+\mathcal{B}\left(a_{3}^{\prime}, a_{1}\right)-\mathcal{F}_{b}(-c),
\end{aligned}
$$

where
$\mathcal{B}\left(a_{3}^{\prime}, a_{1}\right)=\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\left(a_{1}+\mathcal{F}_{a} \mathcal{F}_{-a}^{-1}\left(-a_{3}^{\prime}\right)\right), \mathcal{A}\left(a_{3}^{\prime}, a_{1}, a_{2}\right)=\mathcal{F}_{a} \mathcal{F}_{b}^{-1}\left(a_{2}+\mathcal{B}\left(a_{3}^{\prime}, a_{1}\right)\right)$.
Any sufficiently small perturbation $\bar{a}$ of a fixed point of $\Phi_{n}$ satisfies $a_{1}^{\prime}=$ $a_{2}^{\prime}=0$ and $a_{3 k-2}=-\mathcal{F}_{-a}(-c)$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{4 n+1}$ of $\Sigma_{6 n+3}$ defined by these conditions and use the notation $\widetilde{\Phi}_{n}$ defined for Case $\# 3^{\prime}$, see (13). The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(4 n+1) \times(4 n+1)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ at any point $\bar{a} \in \Sigma_{4 n+1}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
\mathcal{A}_{a_{3}^{\prime}}^{\prime} & 0 & \cdots & \cdots & 0 & 1-\mathcal{B}_{a_{3}^{\prime}}^{\prime} & \mathcal{B}_{a_{3}^{\prime}}^{\prime}-\mathcal{A}_{a_{3}^{\prime}}^{\prime} \\
\mathcal{A}_{a_{2}}^{\prime} & 0 & \cdots & \cdots & \cdots & 0 & 1-\mathcal{A}_{a_{2}}^{\prime} \\
1 & 0 & \cdots & \cdots & \cdots & 0 & 0 \\
0 & 1 & 0 & \vdots & \vdots & \vdots & 0 \\
\vdots & 0 & \ddots & \vdots & \ddots & \vdots & \vdots \\
\vdots & \vdots & \ddots & \ddots & 0 & \vdots & 0 \\
0 & \cdots & \cdots & 0 & 1 & 0 & 0
\end{array}\right)
$$

Lemma 4.15. The characteristic polynomial of $M$ is equal to:

$$
R(X)=-\left(X^{4 n+1}+A X^{4 n}-(1+A+B) X^{2 n}+B\right),
$$

where $A=-\mathcal{A}_{a_{3}^{\prime}}^{\prime}>0$ and $B=\left(\mathcal{A}_{a_{2}}^{\prime}-1\right)\left(1-\mathcal{B}_{a_{3}^{\prime}}^{\prime}\right)<0$.
Indeed, $B$ has same sign as $b-a$ (when ess sup $|f|<\min (a, b))$.
The proof of Proposition 4.14 follows by the same arguments as for Proposition 4.6.
4.1.8. Generalized case \#5.

Proposition 4.16. Any fixed point of $\Phi_{n}$ with generalized shape $\# 5$ is unstable.

The flow $\Phi_{n}$ near a fixed point with generalized shape $\# 5$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=0$ and

$$
\begin{aligned}
& a_{2}^{\prime \prime}=a_{2}+\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right), \quad a_{3}^{\prime \prime}=a_{3}, \quad a_{6 n+1}=a_{2}^{\prime}+\mathcal{B}\left(a_{2}^{\prime}\right)-\mathcal{F}_{-a}(-c), \\
& a_{6 n+2}=a_{3}^{\prime}-\mathcal{B}\left(a_{2}^{\prime}\right)-\mathcal{C}\left(a_{2}^{\prime}, a_{3}^{\prime}\right)+\mathcal{F}_{-a}(-c)+\mathcal{F}_{a}(-c), \\
& a_{6 n+3}=a_{1}+\mathcal{C}\left(a_{2}^{\prime}, a_{3}^{\prime}\right)-\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right)-\mathcal{F}_{a}(-c) .
\end{aligned}
$$

where

$$
\begin{aligned}
& \mathcal{B}\left(a_{2}^{\prime}\right)=\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\left(-a_{2}^{\prime}\right), \quad \mathcal{C}\left(a_{2}^{\prime}, a_{3}^{\prime}\right)=\mathcal{F}_{a} \mathcal{F}_{-a}^{-1}\left(-a_{3}^{\prime}+\mathcal{B}\left(a_{2}^{\prime}\right)\right), \\
& \mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right)=\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\left(a_{1}+\mathcal{C}\left(a_{2}^{\prime}, a_{3}^{\prime}\right)\right) .
\end{aligned}
$$

Any sufficiently small perturbation $\bar{a}$ of a fixed point of $\Phi_{n}$ satisfies $a_{1}^{\prime}=0$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{6 n+2}$ of $\Sigma_{6 n+3}$ defined by this condition and use the notation $\widetilde{\Phi}_{n}$ defined for Case $\# 5$. The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(6 n+2) \times(6 n+2)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ at any point $\bar{a} \in \Sigma_{6 n+2}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
\mathcal{A}_{a_{2}^{\prime}}^{\prime} & 0 & \cdots & 0 & 1+\mathcal{B}_{a_{2}^{\prime}}^{\prime} & -\mathcal{B}_{a_{2}^{\prime}}^{\prime}-\mathcal{C}_{a_{2}^{\prime}}^{\prime} & \mathcal{C}_{a_{2}^{\prime}}^{\prime}-\mathcal{A}_{a_{2}^{\prime}}^{\prime} \\
\mathcal{A}_{a_{3}^{\prime}}^{\prime} & 0 & \ldots & \ldots & 0 & 1-\mathcal{C}_{a_{3}^{\prime}}^{\prime} & \mathcal{C}_{a_{3}^{\prime}}^{\prime}-\mathcal{A}_{a_{3}^{\prime}}^{\prime} \\
\mathcal{A}_{a_{1}}^{\prime} & 0 & \ldots & \ldots & \ldots & 0 & 1-\mathcal{A}_{a_{1}}^{\prime} \\
1 & 0 & \ldots & \cdots & \ldots & \cdots & 0 \\
0 & \ddots & 0 & \vdots & \vdots & \vdots & 0 \\
\vdots & \ddots & \ddots & \ddots & 0 & \vdots & \vdots \\
0 & \cdots & 0 & 1 & 0 & \cdots & 0
\end{array}\right) .
$$

Lemma 4.17. The characteristic polynomial of $M$ is equal to:

$$
R(X)=X^{6 n+2}+A X^{6 n+1}-(1+A+B) X^{4 n+1}+B
$$

with $A=-\mathcal{A}_{a_{2}^{\prime}}^{\prime}>0$ and $B=\left(1+\mathcal{B}_{a_{2}^{\prime}}^{\prime}\right)\left(1-\mathcal{C}_{a_{3}^{\prime}}^{\prime}\right)\left(\mathcal{A}_{a_{1}}^{\prime}-1\right)>0$.
Here $B>0$ for any $a \neq b$ (when ess sup $|f|<\min (a, b)$ ).
The proof of Proposition 4.16 follows by the same arguments as for Proposition 3.3, using Theorem 4.22 below instead of Theorem 3.19.

### 4.1.9. Generalized case \#6.

Proposition 4.18. Any fixed point of $\Phi_{n}$ with generalized shape $\# 6$ is unstable.

The flow $\Phi_{n}$ near a fixed point with generalized shape $\# 6$ is given by (8), with $a_{1}^{\prime}=a_{1}^{\prime \prime}=0$ and
$a_{2}^{\prime \prime}=a_{2}+\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right), \quad a_{3}^{\prime \prime}=a_{3}, \quad a_{6 n+1}=-\mathcal{F}_{-b}(-c)$,
$a_{6 n+2}=a_{2}^{\prime}+a^{\prime}-3+a_{1}-\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right)+\mathcal{F}_{-b}(-c)+\mathcal{F}_{b}(-c)=a_{2}^{\prime}+\mathcal{A}\left(a_{2}^{\prime}\right)$,
$a_{6 n+3}=-\mathcal{F}_{b}(-c)$.
where

$$
\mathcal{A}\left(a_{2}^{\prime}, a_{3}^{\prime}, a_{1}\right)=\mathcal{F}_{b} \mathcal{F}_{a}^{-1}\left(a_{1}+\mathcal{F}_{a} \mathcal{F}_{-a}^{-1}\left(-a_{3}^{\prime}+\mathcal{F}_{-a} \mathcal{F}_{-b}^{-1}\left(-a_{2}^{\prime}\right)\right)\right) .
$$

Any sufficiently small perturbation $\bar{a}$ of a fixed point of $\Phi_{n}$ satisfies $a_{3}^{\prime}=$ $a_{3 k}=\mathcal{F}_{-b}(c)$ and $a_{3 k-2}=\mathcal{F}_{b}(c)$ for $k=1, \ldots, 2 n$. We thus restrict $\Phi_{n}$ to the subspace $\Sigma_{2 n+1}$ of $\Sigma_{6 n+3}$ defined by these conditions and use the notation $\widetilde{\Phi}_{n}$ defined for Case $\# 6$. The transposed $M^{T}$ of the matrix $M \in \mathbb{R}^{(2 n+1) \times(2 n+1)}$ associated to the tangent map of $\widetilde{\Phi}_{n}$ at any point $\bar{a} \in \Sigma_{2 n+1}$ is given by:

$$
M^{T}=\left(\begin{array}{ccccccc}
\mathcal{A}_{a_{2}^{\prime}}^{\prime} & 0 & \cdots & \cdots & \cdots & 0 & 1-\mathcal{A}_{a_{2}^{\prime}}^{\prime} \\
1 & 0 & \cdots & \cdots & \cdots & \cdots & 0 \\
0 & 1 & \ddots & \cdots & \cdots & \ldots & 0 \\
\vdots & \ddots & 1 & 0 & \cdots & \cdots & 0 \\
\vdots & \vdots & 0 & \ddots & \ddots & \vdots & 0 \\
\vdots & \vdots & \vdots & \ddots & 1 & 0 & \vdots \\
0 & \cdots & 0 & \cdots & 0 & 1 & 0
\end{array}\right) .
$$

Lemma 4.19. The characteristic polynomial of $M$ is equal to:

$$
R(X)=-\left(X^{2 n+1}+A X^{2 n}-(1+A)\right),
$$

with $A=-\mathcal{A}_{a_{2}^{\prime}}^{\prime}>0$.
The proof of Proposition 4.18 may follow by the same arguments as for Proposition 4.4, that is the arguments of Proposition 3.3, using Theorem 4.20 with $B=0$ (that can be proved directly in a similar way as in Lemma 3.23). Indeed, since $A=-\mathcal{A}_{a_{2}^{\prime}}^{\prime}>0$, then $1+A>1$, and $R$ has clearly one root $z$ such that $|z|>1$. This shows that the fixed point of $\Phi_{n}$ corresponding to the generalized shape \#6 is unstable.
4.2. Roots of the characteristic polynomials. The following results are proved using homotopy techniques.

Theorem 4.20. Assume that $-1 \leq B \leq 1$ and $A>0$. If $B \neq 1$ or $A<2+\frac{1}{n}$, then the equation

$$
X^{4 n+2}+A X^{4 n+1}-(1+A+B) X^{2 n+1}+B=0
$$

has precisely $2 n$ solutions $z$ (counting algebraic multiplicity) such that $|z|>$ 1. Otherwise it has $2 n-1$ such solutions.

Theorem 4.21. If $-1 \leq B \leq 1$ and $A>0$, the equation

$$
X^{4 n+1}+A X^{4 n}-(1+A+B) X^{2 n}+B=0
$$

has precisely $2 n$ solutions $z$ (counting algebraic multiplicity) such that $|z|>$ 1.

Theorem 4.22. Assume that $0 \leq B \leq 1$ and $A>0$. If $A \geq 1-\frac{4 n+1}{2 n}(1-B)$, the equation

$$
X^{6 n+2}+A X^{6 n+1}-(1+A+B) X^{4 n+1}+B=0
$$

has precisely $2 n$ solutions $z$ (counting algebraic multiplicity) such that $|z|>$ 1. Otherwise it has $2 n+1$ such solutions.

Let us denote by $R$ the polynomial of Theorem 4.20 and consider the homotopy

$$
\begin{equation*}
g_{s}(z)=f_{B}(z)+(s-1) P(z), \tag{29}
\end{equation*}
$$

with

$$
\begin{align*}
P(X) & =X^{4 n}+\cdots+X^{2 n+1}=\frac{X^{4 n+1}-X^{2 n+1}}{X-1}  \tag{30}\\
f_{B}(X) & =X^{4 n+1}+2 X^{4 n}+\cdots+2 X^{2 n+1}-B X^{2 n}+\cdots-B \\
& =\frac{X^{4 n+2}+X^{4 n+1}-(2+B) X^{2 n+1}+B}{X-1} \tag{31}
\end{align*}
$$

We have $R(X)=(X-1) g_{A}(X)$ and $(X-1) f_{B}(X)$ is the polynomial of Theorem 3.19.

Using the same technique as for Lemma 3.22, we obtain:
Lemma 4.23. Consider $n \geq 1$ and $g_{s}$ defined by (29,30,31). If $|z|=1$, $s>0$ and $-1<B \leq 1$, then $g_{s}(z) \neq 0$. If $|z|=1, s>0$ and $B=-1$, then $g_{s}(z)=0 \Leftrightarrow z=-1$.

Proof. Consider $z$ such that $|z|=1$ and $g_{s}(z)=0$, and put $w=z^{2 n+1}$. We get

$$
\begin{equation*}
\frac{s}{z} w^{2}=-w^{2}+(1+s+B) w-B \tag{32}
\end{equation*}
$$

Taking modulus on both sides gives $s^{2}=\left|-w^{2}+(1+s+B) w-B\right|^{2}$, writing $w=e^{i \theta}$ gives $(1-\cos \theta)\left(1+B^{2}+(1+B) s-2 B \cos \theta\right)=0$. Thus, we have either $\cos \theta=1$ or $B \neq-1$ and $s=-\frac{\left(1+B^{2}-2 B \cos \theta\right)}{1+B} \leq 0$ which is impossible, or $B=-1$ and $\cos \theta=-1$. The condition $\cos \theta=1$ implies $w=1$. Plugging $w=1$ in (32), we obtain $z=1(s>0)$, which is impossible since $g_{s}(1)=(2 n+1)(1-B)+2 n s>0$. So the only possibility is $B=-1$ and $\cos \theta=-1$, that is $w=-1$. With (32), this implies $z=-1$, which is a root of $g_{s}$ when $B=-1$.

Proof of Theorem 4.20 for $B \neq-1$. The function $g_{s}(z)$ defined in (29), $(30),(31)$ is continuous with respect to $s$ and polynomial with constant degree $4 n+1$ with respect to $z$. From Lemma $4.23, g_{s}$ has no roots on $\partial D(0,1)$ for $s>0$. So by Rouché's theorem, the number of roots of $g_{s}$ in $D(0,1)$ (resp. $\left.\bar{D}(0,1)^{c}\right)$ is constant with respect to $s$. According to Theorem $3.19, g_{1}(z)=0$ has precisely $2 n$ solutions in $\bar{D}(0,1)^{c}$, so does $g_{s}(z)=0$ for $s>0$.

Proof of Theorem 4.20 for $B=-1$. When $B$ is fixed equal to $-1, z=$ -1 is a root of $g_{s}$ for all $s>0$ (see Lemma 4.23), so $g_{s}(z)=(z+1) \tilde{g}_{s}(z)$ with

$$
\begin{equation*}
\tilde{g}_{s}(z)=\tilde{f}_{B}(z)+(s-1) \widetilde{P}(z) \tag{33}
\end{equation*}
$$

and

$$
\begin{align*}
\widetilde{P}(X) & =\frac{P(X)}{X+1}=X^{4 n-1}+X^{4 n-3}+\cdots+X^{2 n+3}+X^{2 n+1}  \tag{34}\\
\tilde{f}_{B}(X) & =\frac{f_{B}(X)}{X+1}  \tag{35}\\
& =X^{4 n}+X^{4 n-1}+\cdots+X^{2 n+1}+X^{2 n}+X^{2 n-2}+\cdots+X^{2}+1
\end{align*}
$$

We have $\tilde{g}_{s}(-1)=g_{s}^{\prime}(-1)=2 n+1-s n \neq 0$ for $s \neq \bar{s}=2+\frac{1}{n}$. So, with Lemma 4.23, we obtain that, for $s \neq \bar{s}, \tilde{g}_{s}$ has no roots on $\partial D(0,1)$. It follows, by Rouché's theorem, that the number of roots of $\tilde{g}_{s}$ (and thus of $\left.g_{s}\right)$ in $D(0,1)$ (resp. $\left.\bar{D}(0,1)^{c}\right)$ is constant with respect to $s$ in each of the intervals $(0, \bar{s})$ and $(\bar{s},+\infty)$. According to Theorem 3.19, $g_{1}$ has precisely $2 n$ roots in $\bar{D}(0,1)^{c}$, so does $g_{s}$ for $0<s<\bar{s}$.

For $s=\bar{s}$, we have $\tilde{g}_{\bar{s}}^{\prime}(-1)=\frac{g_{\bar{s}}^{\prime \prime}(-1)}{2}=n(2 n+1)>0$. The implicit function theorem shows that there exists $\delta>0, \rho>0$ and a $\mathcal{C}^{1}$ function
$Z:(\bar{s}-\delta, \bar{s}+\delta) \rightarrow \mathbb{C}$, such that $\tilde{g}_{s}(z)=0$ and $|z+1| \leq \rho$ if and only if $z=Z(s)$. So $Z(\bar{s})=-1$ and

$$
Z^{\prime}(\bar{s})=-\frac{\left.\frac{\partial \tilde{g}_{s}(z)}{\partial s}\right|_{s=\bar{s}, z=-1}}{\left.\frac{\left.\partial \tilde{g}_{s} z\right)}{\partial z}\right|_{s=\bar{s}, z=-1}}=-\frac{\widetilde{P}(-1)}{\tilde{g}_{\bar{s}}^{\prime}(-1)}=\frac{1}{2 n+1}>0,
$$

which implies $|Z(s)|<1$ for $\bar{s}<s<\bar{s}+\delta$ and $|Z(s)|>1$ for $\bar{s}-\delta<s<\bar{s}$ (possibly by shrinking $\delta$ ). Since the root $Z(s)$ of $\tilde{g}_{s}$ goes from $\bar{D}(0,1)^{c}$ to $D(0,1)$ when $s$ increases around $\bar{s}$, this implies that
$\operatorname{deg}\left(\tilde{g}_{s}, D(0,1), 0\right)=1+\operatorname{deg}\left(\tilde{g}_{s^{\prime}}, D(0,1), 0\right) \quad$ for $\bar{s}-\delta<s^{\prime}<\bar{s}<s<\bar{s}+\delta$
(one may use the set $G_{\rho}=\{z \in \mathbb{C}$ such that $|z|<1$ or $|z+1|<\rho\}$ to prove this result). Since $\tilde{g}_{A}$ has $2 n$ roots in $\bar{D}(0,1)^{c}$ when $0<A<\bar{s}$, it follows that $\tilde{g}_{A}$ has $2 n-1$ roots when $\bar{s}<A$ and also when $A=\bar{s}$. So does $g_{A}$ and $R=(X-1) g_{A}(X)$.

Theorem 4.21 is proved by the same technique as the case $B \neq 1$ of Theorem 4.20, changing $P$ and $Q$ into

$$
\begin{align*}
P(X) & =\frac{X^{4 n}-X^{2 n}}{X-1}  \tag{36}\\
f_{B}(X) & =\frac{X^{4 n+1}+X^{4 n}-(2+B) X^{2 n}+B}{X-1} \tag{37}
\end{align*}
$$

Indeed, using the same technique as for Lemma 3.24, one gets:
Lemma 4.24. Consider $n \geq 1$ and $g_{s}$ defined by $(29,36,37)$. If $|z|=1$, $s>0$ and $-1 \leq B \leq 1$, then $g_{s}(z) \neq 0$.
Proof. Consider $z$ such that $|z|=1$ and $g_{s}(z)=0$, and put $w=z^{2 n}$. We obtain

$$
\begin{equation*}
\frac{z}{s} w^{2}=-w^{2}+\left(1+\frac{1}{s}+\frac{B}{s}\right) w-\frac{B}{s} \tag{38}
\end{equation*}
$$

which leads, as (32), to the only possibilities $w=z=1$ or $\frac{B}{s}=-1$ and $w=z=-1$. This last case is impossible since $w=z^{2 n}$. The first one also since $g_{s}(1)=1+2 n(1-B+s)>0$.

Theorem 4.22 is proved by the same technique as Theorem 4.20, changing $P$ and $Q$ into

$$
\begin{align*}
P(X) & =\frac{X^{6 n+1}-X^{4 n+1}}{X-1}  \tag{39}\\
f_{B}(X) & =\frac{X^{6 n+2}+X^{6 n+1}-(2+B) X^{4 n+1}+B}{X-1} \tag{40}
\end{align*}
$$

Lemma 4.25. Consider $n \geq 1$ and $g_{s}$ defined by $(29,39,40)$. If $|z|=1$, $s>0$ and $0 \leq B \leq 1$, then $g_{s}(z)=0 \Rightarrow z=1$. Moreover, $z=1$ is a root of $g_{s}$ if and only if $s=1-\frac{4 n+1}{2 n}(1-B)$. In this case, it is a simple root and $g_{s}^{\prime}(1)>0$.

Proof. Consider $z$ such that $|z|=1$ and $g_{s}(z)=0$, and put $w=z^{2 n+1}$, $v=z^{2 n}$. We get

$$
\begin{equation*}
w+s v-(1+s+B)+B(w v)^{-1}=0 \tag{41}
\end{equation*}
$$

Set $w=e^{i \theta}, v=e^{i \theta^{\prime}}$. Multiplying (41) by $B v-1$, and taking the real part, we obtain
$\left(1-B^{2}\right)(1-\cos \theta)+\left((1-s)(s+B)+(s-B)^{2}\right)\left(1-\cos \theta^{\prime}\right)+2 B s\left(1-\cos \theta^{\prime}\right)^{2}=0$.
Suppose first that $0<s \leq 1$ and $0 \leq B \leq 1$. Since all the terms in the previous sum are non negative, and the sum is zero, all the terms are zero. This implies, in particular, $\cos \theta^{\prime}=1$ or $v=1$. Plugging $v=1$ in (41) and taking the real part, we obtain $(1+B)(\cos \theta-1)=0$, so $\cos \theta=1$ or $w=1$.

Equation (41) can be rewritten

$$
v+\frac{1}{s} w-\left(1+\frac{1}{s}+\frac{B}{s}\right)+\frac{B}{s}(w v)^{-1}=0
$$

which is nothing but (41) with $w, v, s, B$ replaced respectively by $v, w, \frac{1}{s}, \frac{B}{s}$. Hence, the same conclusion $v=w=1$ is obtained when $0<\frac{1}{s} \leq 1$ and $0 \leq \frac{B}{s} \leq 1$, and in particular when $0 \leq B \leq 1 \leq s$. Therefore, one deduces $z=\frac{\dot{w}}{v}=1$ for any $B$ and $s$ such that $0 \leq B \leq 1$ and $s>0$.

The last points of the lemma are obtained by simple computations: $g_{s}(1)=$ $2 n(s-1)+(4 n+1)(1-B)$ for all $s$ and $2 g_{s}^{\prime}(1)=\left(g_{s}(z)(z-1)\right)^{\prime \prime}(1)=$ $2 n+1+B(4 n+1)(6 n+1)>0$ for $s=1-\frac{4 n+1}{2 n}(1-B)$.

Proof of Theorem 4.22. Let $g_{s}$ be defined by $(29,39,40)$ with $B$ fixed such that $0 \leq B \leq 1$ and denote $\bar{s}=\max \left(0,1-\frac{4 n+1}{2 n}(1-B)\right)$. From Lemma 4.25, $g_{s}$ has no roots on $\partial D(0,1)$ for $s \neq \bar{s}$. It follows, by Rouché's theorem, that the number of roots of $g_{s}$ in $D(0,1)$ (resp. $\left.\bar{D}(0,1)^{c}\right)$ is constant with respect to $s$ in each of the intervals $(0, \bar{s})$ and $(\bar{s},+\infty)$.

Suppose first that $\bar{s}=0$. According to Theorem 3.21, $g_{1}$ has precisely $2 n$ roots in $\bar{D}(0,1)^{c}$, so does $g_{s}$ for all $s>0$.

Suppose now that $\bar{s}>0$. By Lemma 4.25, 1 is a simple root of $g_{\bar{s}}$, so one can apply the same technique as for the case $B=-1$ of Theorem 4.20. By the implicit function theorem, there exists $\delta>0, \rho>0$ and a $\mathcal{C}^{1}$ function $Z:(\bar{s}-\delta, \bar{s}+\delta) \rightarrow \mathbb{C}$, such that $g_{s}(z)=0$ and $|z-1| \leq \rho$ if and only if
$z=Z(s)$. Then, $Z(\bar{s})=1$ and

$$
Z^{\prime}(\bar{s})=-\frac{\left.\frac{\partial g_{s}(z)}{\partial s}\right|_{s=\bar{s}, z=1}}{\left.\frac{\partial g_{s}(z)}{\partial z}\right|_{s=\bar{s}, z=1}}=-\frac{P(1)}{g_{\bar{s}}^{\prime}(1)}=-\frac{2 n}{g_{\bar{s}}^{\prime}(1)}<0
$$

which implies that the root $Z(s)$ of $g_{s}$ goes from $D(0,1)$ to $\bar{D}(0,1)^{c}$ when $s$ decreases around $\bar{s}$. It follows that the number of roots of $g_{s}$ in $\bar{D}(0,1)^{c}$ is constant for $s \in[\bar{s},+\infty$ ), and thus equal to $2 n$ by Theorem 3.21 (using $\bar{s} \leq 1$ ), and that it is constant equal to $2 n+1$ for $s \in(0, \bar{s})$. The same result holds for $R=(X-1) g_{A}(X)$ and $s=A$.

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[^1]:    ${ }^{1}$ One can eliminate this condition by replacing the value of $h$ in $0(h(0)=0$ here due to the oddness of $h$ ) by any value in $\{-a,-b, a, b\}$, or by replacing $h(x(t-1))$ by $h(t, x(t-1))$, where $t \in(0,+\infty) \mapsto h(t,$.$) is Lebesgue measurable and (h(t, x) \in\{-a,-b, a, b\} \forall x \in \mathbb{R})$, for a.e. $t \in(0,+\infty)$, see [1].

