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# Global stability for the 2-dimensional logistic map 

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#### Abstract

For the delayed logistic equation $x_{n+1}=a x_{n}\left(1-x_{n-1}\right)$ it is well known that the nontrivial fixed point is locally stable for $1<a \leq 2$, and unstable for $a>2$. We prove that for $1<a \leq 2$ the fixed point is globally stable; in the sense that it is locally stable and attracts all points of $S$, where $S$ contains those $\left(x_{0}, x_{1}\right) \in \mathbb{R}_{+}^{2}$, for which the sequence $\left\{x_{n}\right\} \subseteq \mathbb{R}_{+}$. The proof is a combination of analytical and reliable numerical methods.


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## 1. Introduction

One of the most studied nonlinear maps is the logistic map

$$
[0,1] \ni x \mapsto a x(1-x) \in \mathbb{R},
$$

with parameter $a>0$. For $0<a \leq 1$, it is well known (see eg [8]) that $x=0$ is the unique fixed point in $[0,1]$, and it is globally stable (ie stable and attracts all points in $[0,1]$ ). For $1<a \leq 3$, there is a nontrivial fixed point $x_{*}=1-(1 / a)$ which is stable and attracts all points in ( 0,1 ). At $a=3$ a period doubling (flip) bifurcation takes place, and the fixed point $x_{*}$ becomes unstable for $a>3$. As $a$ increases, there is a sequence of bifurcation points, and for some larger value of $a$, chaotic behaviour can be shown.

In 1968, Maynard Smith [13] considered the 'delayed' version

$$
x_{k+1}=a x_{k}\left(1-x_{k-1}\right)
$$

of the logistic difference equation. This is natural in the context of population models: the size of the subsequent generation of the population depends not only on the size in the previous year, but also on the size of the two-year earlier population.

Introducing $y_{k}=x_{k+1}$, the second-order difference equation is equivalent to

$$
\left(x_{n+1}, y_{n+1}\right)=F_{a}\left(x_{n}, y_{n}\right)
$$

with

$$
\begin{equation*}
F_{a}: \mathbb{R}^{2} \ni(x, y) \mapsto(y, a y(1-x)) \in \mathbb{R}^{2} . \tag{1}
\end{equation*}
$$

[^0]We study the map $F_{a}$ for those $(x, y) \in \mathbb{R}_{+}^{2}=[0, \infty) \times[0, \infty)$ for which all iterates of $F_{a}$ remain in $\mathbb{R}_{+}^{2}$, ie $F_{a}^{k}(x, y) \in \mathbb{R}_{+}^{2}$, for every $k \in \mathbb{N}$. Here $F_{a}^{k}$ denotes the $k$-fold iteration of $F_{a}$, ie $F_{a}^{0}=i d, F_{a}^{k}=F_{a}\left(F_{a}^{k-1}\right), k \in \mathbb{N}$. As we will see, for $0<a \leq 2$ the set

$$
S_{0}=S_{0}(a)=\left\{(x, y) \in \mathbb{R}^{2}: 0 \leq x \leq 1 ; \quad 0 \leq y \leq 1 ; \quad a y(1-x) \leq 1\right\}
$$

is positively invariant under $F_{a}$, that is, $F_{a}\left(S_{0}\right) \subseteq S_{0}$.
For $0<a \leq 1$, we have $S_{0}=[0,1] \times[0,1]$, the only fixed point in $S_{0}$ is $(0,0)$, which is locally stable and $F_{a}^{k}(x, y) \rightarrow(0,0)$ as $k \rightarrow \infty$, for every $(x, y) \in S_{0}$. For $a>1$, the nontrivial fixed point $(A, A) \in S_{0}$ with $A=1-(1 / a)$ appears, which is locally asymptotically stable for $a \in(1,2)$, and it is unstable for $a>2$. A Neimark-Sacker bifurcation takes place at $a=2$ [see eg in 9 Example 4.3] and for $a>2, a$ close to 2 , there is a stable invariant curve. As we increase $a$, the size of the invariant curve is getting larger; at about $a=2.27$, the curve touches the $x$-axis, and complicated dynamics occurs. For profound numerical studies, see [1, 7].

The aim of this paper is to show that, for $1<a \leq 2$, the nontrivial fixed point $(A, A)$ is globally stable in the sense that $(A, A)$ is locally stable, and for each $(x, y)$ in

$$
S=S(a)=\left\{(x, y) \in \mathbb{R}^{2}: 0 \leq x<1 ; \quad 0<y<1 ; \quad \text { ay }(1-x)<1\right\}
$$

$F_{a}^{k}(x, y) \rightarrow(A, A)$ as $k \rightarrow \infty$. Consequently, the local stability of $(A, A)$ implies its global stability. For similar results on the global stability of other delayed difference equations, the reader is referred for example to [3-5, 10, 11].

We emphasize that we prove the stability even in the critical parameter value $a=2$. However, we do not consider the case $a>2$. According to numerical studies [ 1,7 see], the invariant curve is globally stable for parameter $a>2$ close to the critical value 2 .

For $a \leq 2$ the proof of the global stability is a combination of analytical and computeraided tools. It is based on the method in $[3,4]$. We elaborate the analytical part such that it can be easily applied to similar models. Furthermore, a quite important aim is to have a clear picture of the method in order to be able to prove similar results for higher dimensional models, for example, the 3-dimensional logistic map $x_{n+1}=a x_{n}\left(1-x_{n-2}\right)$, where further difficulties arise.

With analytical tools we construct an attracting neighbourhood $\mathcal{N}$ around the nontrivial fixed point $(A, A)$. Then we show that every $\left(x_{0}, y_{0}\right) \in S \backslash \mathcal{N}$ will eventually step into $\mathcal{N}$, that is, there exist an $n \in \mathbb{N}$, such that $\left(x_{n}, y_{n}\right) \in \mathcal{N}$, where $\left(x_{n}, y_{n}\right)=F_{a}^{n}\left(x_{0}, y_{0}\right)$. So these points are also in the region of attraction of the fixed point $(A, A)$. We use a computer, applying reliable numerical methods, to show the second step. In this context, reliable means, we use interval arithmetic tools to control every occurring numerical error; consequently, the method is suitable to prove mathematical statements [see eg 14].

In Section 3, for smaller parameter values $a$, ie for $a \in(1.5,1.95$ ] we use the linearized map to construct the attracting neighbourhood $\mathcal{N}$. However, as we will see it later, this neighbourhood shrinks to the fixed point as $a$ tends to critical value 2 . Therefore, for parameter values $a$ close to 2 this neighbourhood is not big enough for computer use in the second part of the method. Thus we need another approach to construct an attracting neighbourhood $\mathcal{N}$. In Section 4, for these parameter values a close to 2, we use the normal form of the Neimark-Sacker bifurcation. More precisely, with smooth and invertible maps,
we transform map (1) into its normal form; hereby we obtain an attracting neighbourhood $\mathcal{N}$ around the fixed point $(A, A)$, whose size is independent of the parameter $a \in[1.95,2]$.

Since we need the size of the constructed neighbourhood $\mathcal{N}$ for computer use, it is not enough to determine only the lower order terms during the normal form transformation, like we would do in a regular bifurcation analysis. These lower order terms only assure the existence of such a sufficiently small neighbourhood, whose size is not explicitly determined by them. Therefore, it is essential during the transformation to trace the higher order terms and to estimate them as well as possible, in order to obtain a sufficiently big neighbourhood $\mathcal{N}$, since the computer method is more and more compute-intensive and time-consuming, as we get closer to the fixed point.

In Section 5, we consider those points, which lie outside the attracting neighbourhood, ie the points of $S \backslash \mathcal{N}$. We cover $S$ with finitely many small squares. Considering these squares as vertices of a graph, we introduce a directed graph, which, to a certain extent, describes the behaviour of map (1) on these squares. Therefore we convert the problem of examining infinitely many points into a finite graph problem, which can be handled by computer. To construct the edges of the graph we use reliable numerical methods in order to handle the rounding errors of the computer. We show with the help of this graph that every point from $S$ enters the neighbourhood $\mathcal{N}$ constructed before. With this we will prove our main result:

Theorem 1.1: For all $a \in(1,2]$ the fixed point $(A, A)$ is locally asymptotically stable, and $\lim _{n \rightarrow \infty} F_{a}^{n}(x, y)=(A, A)$ for every $(x, y) \in S$, where $A=1-(1 / a)$.

## 2. Preliminaries

Throughout the article $\mathbb{R}_{+}, \mathbb{N}$ and $\mathbb{N}_{0}$ denote the nonnegative real numbers, the positive integers and nonnegative integers, respectively.

In this section, we study the dynamics of map (1) for $a>0$ in the positive quadrant. Introduce the following sets (see Figure 1):

$$
\begin{aligned}
S= & \left\{(x, y) \in \mathbb{R}^{2}: 0 \leq x<1, \quad 0<y<1, \quad \text { ay }(1-x)<1\right\}, \\
T_{0}= & \left\{(x, 0) \in \mathbb{R}^{2}: x \geq 0\right\} \cup\left\{(1, y) \in \mathbb{R}^{2}: y>0\right\} \cup\left\{(x, 1) \in \mathbb{R}^{2}: 0 \leq x<1\right\} \\
& \cup\left\{(x, y) \in \mathbb{R}^{2}: \text { ay }(1-x)=1, \quad 0 \leq x \leq A\right\}, \\
T_{1}= & \left\{(x, y) \in \mathbb{R}^{2}: x \geq 0, \quad 0<y<1, \quad \text { ay }(1-x)>1\right\}, \\
T_{2}= & \left\{(x, y) \in \mathbb{R}^{2}: 0 \leq x<1, \quad y>1\right\}, \\
T_{3}= & \left\{(x, y) \in \mathbb{R}^{2}: x>1, \quad y>0\right\} .
\end{aligned}
$$

Clearly, $\mathbb{R}_{+}^{2}=S \cup T_{0} \cup T_{1} \cup T_{2} \cup T_{3}$; furthermore, $S=[0,1) \times(0,1)$ and $T_{1}=\emptyset$ for $a \in$ ( 0,1 ].

Proposition 2.1: For all $a>0$, we have

$$
F_{a}^{4}\left(T_{0}\right)=\{(0,0)\}, \quad F_{a}\left(T_{1}\right) \subseteq T_{2}, \quad F_{a}\left(T_{2}\right) \subseteq T_{3}, \quad F_{a}\left(T_{3}\right) \cap \mathbb{R}_{+}^{2}=\emptyset
$$

and furthermore, if $a \in(0,2]$ then $F_{a}(S) \subseteq S$.


Figure 1. The subdivision of the positive quadrant.

Proof: From the definition of $F_{a}$ it is obvious that $F_{a}^{4}\left(T_{0}\right)=\{(0,0)\}$. It is also straightforward to check the relations $F_{a}\left(T_{1}\right) \subseteq T_{2}, F_{a}\left(T_{2}\right) \subseteq T_{3}$ and $F_{a}\left(T_{3}\right) \cap \mathbb{R}_{+}^{2}=\emptyset$.

If $a \in(0,2]$ and $(x, y) \in S$, then $0<y<1,0<a y(1-x)<1$ and

$$
a^{2} y(1-x)(1-y) \leq 4(1-x) \max _{0 \leq y \leq 1} y(1-y) \leq 1,
$$

therefore $F_{a}(S) \subseteq S$.
Consequently, in the rest of the paper we can assume $(x, y) \in S$. For small $a$, the dynamics in $S$ is quite simple. The following statement easily follows from the monotonicity of $\left\{x_{n}\right\}_{n=0}^{\infty}$ for $0<a \leq 1$.

Proposition 2.2: If $a \in(0,1]$, then $\lim _{k \rightarrow \infty} F_{a}^{k}(x, y) \rightarrow(0,0)$ for all $(x, y) \in[0,1]^{2}$.
For $a \in(1,2]$, we divide the positive quadrant into four subsets with lines $x=A, y=A$, and introduce the following sets:

$$
\begin{array}{ll}
S_{1}=\{(x, y) \in S: x \leq A, & y<A\}, \\
S_{2}=\{(x, y) \in S: x<A, & A \leq y\}, \\
S_{3}=\{(x, y) \in S: A \leq x, & A<y\}, \\
S_{4}=\{(x, y) \in S: A<x, & y \leq A\} .
\end{array}
$$

Clearly $S=\bigcup_{i=1}^{4} S_{i} \cup\{(A, A)\}$. Introduce the notation $t_{n}=\left(x_{n}, y_{n}\right)=F_{a}^{n}\left(x_{0}, y_{0}\right)$.
Proposition 2.3: For every $t_{0}=\left(x_{0}, y_{0}\right) \in S$ and $a \in(1,2]$ the sequence $\left\{t_{n}\right\}_{n=0}^{\infty}$, defined by $t_{n+1}=\left(x_{n+1}, y_{n+1}\right)=F_{a}\left(x_{n}, y_{n}\right)$, fulfils one of the following cases:

- $\lim _{n \rightarrow \infty} t_{n}=(A, A)$,


Figure 2. The dynamics in $S$.

- the sequence $\left\{t_{n}\right\}_{n=0}^{\infty}$ goes around the fixed point along the cycle $S_{1} \rightarrow S_{2} \rightarrow S_{3} \rightarrow S_{4} \rightarrow$ $S_{1}$; and in the course of one cycle, there can be more than one elements of the sequence in both $S_{1}$ and $S_{3}$, but the number of these elements are finite.

Proof: The transitions between the aforementioned subsets are the following (see Figure 2):

- For $t_{0} \in S_{1}$ we obtain $x_{1}<A$, therefore $S_{1} \rightarrow\left\{S_{1}, S_{2}\right\}$. That is, $t_{1} \in S_{1}$ or $t_{1} \in S_{2}$.
- For $t_{0} \in S_{2}: x_{1} \geq A$ and $y_{1}=a y_{0}\left(1-x_{0}\right)>a A(1-A)=A$, so $S_{2} \rightarrow S_{3}$.
- For $t_{0} \in S_{3}: x_{1}>A$, so $S_{3} \rightarrow\left\{S_{3}, S_{4}\right\}$.
- For $t_{0} \in S_{4}: x_{1} \leq A$ and $y_{1}=a y_{0}\left(1-x_{0}\right)>a A(1-A)=A$, so $S_{4} \rightarrow S_{1}$.

We obtain there is a cycle $S_{1} \rightarrow S_{2} \rightarrow S_{3} \rightarrow S_{4} \rightarrow S_{1}$. But during a cycle the points of the sequence can spend more time in $S_{1}$ or $S_{3}$, possibly, the sequence can stay in $S_{1}$ or $S_{3}$ forever. We only need to show that, if a sequence gets stuck in $S_{1}$ or $S_{3}$, then it converges to the fixed point $(A, A)$.

Notice that $y_{n+1} \geq y_{n}$ as long as $x_{n} \leq A$, and similarly, $x_{n} \geq A$ implies $y_{n+1} \leq y_{n}$. According to this, as long as $t_{n} \in S_{1} \cup S_{2}$, the sequence $\left\{y_{n}\right\}=\left\{x_{n+1}\right\}$ increases, until the sequence $\left\{t_{n}\right\}$ steps into $S_{3}$. Similarly, as long as $t_{n} \in S_{3} \cup S_{4}$, the sequence $\left\{y_{n}\right\}=\left\{x_{n+1}\right\}$ decreases, until the sequence $\left\{t_{n}\right\}$ steps into $S_{1}$. Consequently, if a sequence $\left\{t_{n}\right\}$ stays in $S_{1}$ for all large $n$, we gain a monotonically increasing, bounded sequence $\left\{x_{n}\right\}$, which converges to $B \leq A$. Taking the limit of both sides of $x_{n+1}=a x_{n}\left(1-x_{n-1}\right)$, we obtain $B=A$, and consequently $t_{0}$ is in the region of attraction of $(A, A)$. Similarly, if a sequence gets stuck in $S_{3}$, it also converges to $(A, A)$.

Now we assume $1<a \leq \frac{3}{2}$ and show that for every $t_{0} \in S$ the sequence $\left\{t_{n}\right\}$ converges to the nontrivial fixed point $(A, A)$. Combining this fact with the local asymptotic stability of the fixed point (see at the beginning of the following section), Theorem 1.1 is proved for these parameter values.

Proposition 2.4: If $a \in\left(1, \frac{3}{2}\right]$, then $\lim _{n \rightarrow \infty} F_{a}^{n}\left(x_{0}, y_{0}\right)=(A, A)$ for every $\left(x_{0}, y_{0}\right) \in S$.
Proof: It is clear from Proposition 2.3, we only need to consider the case when the sequence $\left\{t_{n}\right\}=\left\{\left(x_{n}, y_{n}\right)\right\}$ goes around the fixed point, not getting stuck in $S_{1}$ or $S_{3}$. It means that there exist subsequences $n_{k}$ and $m_{k}$, such that

$$
t_{n_{k}}, t_{n_{k}+1}, \ldots, t_{m_{k}-2} \in S_{1} ; \quad t_{m_{k}-1} \in S_{2} ; \quad t_{m_{k}}, t_{m_{k}+1}, \ldots, t_{n_{k+1}-2} \in S_{3} ; \quad t_{n_{k+1}-1} \in S_{4}
$$

for all $k \geq 0$. Clearly $n_{k}+2 \leq m_{k} \leq n_{k+1}-2$ also holds. Without loss of generality, we can assume $t_{0} \in S_{1}$.

Now, consider the sequence $\left\{s_{n}\right\}$, where $s_{0}=0$ and $s_{n}=f^{n}\left(s_{0}\right)$ with the function $f(x)=$ $a(a-1)(1-x)^{2}$. Denote by $\left\{h_{n}\right\}$ and $\left\{g_{n}\right\}$ the even and odd indexed subsequences of $\left\{s_{n}\right\}$, ie $h_{n}=s_{2 n}$ and $g_{n}=s_{2 n+1}$. Furthermore introduce the following subsets of $S$ :

$$
\mathcal{H}_{k}=\left\{(x, y) \in S: h_{k} \leq x, y\right\} \text { and } \mathcal{G}_{k}=\left\{(x, y) \in S: x, y \leq g_{k}\right\} .
$$

Clearly, $\mathcal{H}_{0}=S$.
It is easy to see that, if $t_{n_{k}} \in \mathcal{H}_{k}$ then $t_{n_{k}}, t_{n_{k}+1}, \ldots, t_{m_{k}} \in \mathcal{H}_{k}$ and because of the inequality

$$
y_{m_{k}}=a y_{m_{k}-1}\left(1-x_{m_{k}-1}\right)=a^{2} y_{m_{k}-2}\left(1-x_{m_{k}-2}\right)\left(1-x_{m_{k}-1}\right) \leq a^{2} A\left(1-h_{k}\right)^{2}=g_{k}
$$

$t_{n_{k}}, t_{n_{k}+1}, \ldots, t_{m_{k}} \in \mathcal{G}_{k}$ also holds. Similarly, $t_{m_{k}}, t_{m_{k}+1}, \ldots, t_{n_{k+1}} \in \mathcal{G}_{k} \cap \mathcal{H}_{k+1}$ provided that $t_{m_{k}} \in \mathcal{G}_{k}$. It follows from the construction that $\mathcal{H}_{k+1} \subseteq \mathcal{H}_{k}$ and $\mathcal{G}_{k+1} \subseteq \mathcal{G}_{k}$. Consequently $\left\{h_{n}\right\}$ is increasing and bounded above by $A$, so $\lim _{n \rightarrow \infty} h_{n}=h_{\infty} \leq A$. Similarly $\left\{g_{n}\right\}$ is decreasing and bounded below by $A$, so $\lim _{n \rightarrow \infty} g_{n}=g_{\infty} \geq A$. Therefore, we only need to show that $h_{\infty}=g_{\infty}=A$.

It is clear that $h_{\infty}$ and $g_{\infty}$ need to be fixed points of $f_{2}(x)=f(f(x))$. Observe that $f_{2}(A)=A, f_{2}^{\prime}(A)=4(a-1)^{2}$. Consequently $0<f_{2}^{\prime}(A) \leq 1$ for $1<a \leq \frac{3}{2}$. Furthermore, for $0<x<A$

$$
f_{2}^{\prime \prime}(x)=4 a^{2}(1-a)^{2}(3 f(x)-1)>0
$$

since $f(x)>A \geq \frac{1}{3}$. We can conclude that $A$ is the only solution of $f_{2}(x)=x$ in the interval $(0, A]$, so $\lim _{n \rightarrow \infty} h_{n}=A$. From the definition of $g_{n}$, it is clear, $\lim _{n \rightarrow \infty} g_{n}=A$, too.

In the rest of the paper, we assume $a \in\left(\frac{3}{2}, 2\right]$. For these parameter values, the above argument does not guarantee convergence for every $t \in S$, but we show it is enough to consider a subset of $S$ later on.

Proposition 2.5: For every $a \in\left(\frac{3}{2}, 2\right]$ the set

$$
\tilde{S}=\{(x, y) \in S: x, y \in[0.072,0.8]\}
$$

is invariant, ie $F(\tilde{S}) \subseteq \tilde{S}$. Furthermore, for every $t \in S$, there exists $N=N(t)$, such that for every $n>N, F^{n}(t) \in \tilde{S}$.

Proof: Using the arguments and notations of the previous proposition, we can assume $t_{0} \in$ $S_{1}$ and the sequence $\left\{t_{n}\right\}=\left\{\left(x_{n}, y_{n}\right)\right\}$ goes around the fixed point. We need to show that $y_{m_{0}} \leq \frac{4}{5}$ and $y_{n_{1}} \geq 0.072$. Since $y_{1}=a y_{0}\left(1-x_{0}\right) \leq a y_{0}=a x_{1}$ follows from $\left(x_{0}, y_{0}\right) \in S$, we can also assume that $y \leq a x$ for every $t \in S$.

For $y_{m_{0}}$ we have to find the maximum of $a^{2} y(1-y)(1-x)$ assuming $(x, y) \in S_{1}$ and $y \leq a x$. Since $y(1-y)$ is increasing on $[0, A]$ we are looking for the maximum of $f(x)=$ $a^{3} x(1-a x)(1-x)$ on $\left[0, \frac{A}{a}\right]$ and $g(x)=(a-1)(1-x)$ on $\left[\frac{A}{a}, A\right]$. The maxima of $g(x)$ and $f(x)$ are $\frac{3}{4}$ and $\frac{4}{3 \sqrt{3}}$, respectively, so $y_{m_{0}} \leq \frac{4}{5}$.

Similarly, for every $a \in[1.5,2]$ we are looking for the minimum of $a^{2} y(1-y)(1-x)$ on $S_{3}$, assuming $x, y \leq \frac{4}{5}$. It is easy to see that this is 0.072 .

We apply this proposition in the computer-assisted part of the proof, since it is useful to exclude a small neighbourhood of the trivial fixed point $(0,0)$, as we see it later. For more general results on absorbing sets like $\tilde{S}$, the reader is referred to [6].

## 3. Attracting neighbourhood with linearization

In this section using the linearization of map (1), for a fixed parameter $a \in\left(\frac{3}{2}, 2\right)$, we give a neighbourhood $\mathcal{N}(a)$ around $(A, A)$, which is inside the region of attraction of this fixed point, ie $\lim _{n \rightarrow \infty} F_{a}^{n}\left(x_{0}, y_{0}\right)=(A, A)$ for every $\left(x_{0}, y_{0}\right) \in \mathcal{N}(a)$.

Introducing the new variables $u=x-A$ and $v=y-A$, map (1) can be written in the following form:

$$
\begin{equation*}
\binom{u}{v} \mapsto J(a)\binom{u}{v}+f_{a}(u, v), \tag{2}
\end{equation*}
$$

where

$$
J(a)=\left(\begin{array}{cc}
0 & 1 \\
1-a & 1
\end{array}\right), \quad f_{a}(u, v)=\binom{0}{-a u v}
$$

For $a \in\left(\frac{3}{2} ; 2\right]$ the eigenvalues of $J(a)$ are $\lambda:=\lambda_{1}(a)=\overline{\lambda_{2}(a)}=(1+i \sqrt{4 a-5}) / 2$ and the corresponding eigenvectors are $q_{1,2}(a)=\left(1, \lambda_{1,2}(a)\right)$. It is easy to see, that $\left|\lambda_{i}(a)\right|<1$ for $a \in(1,2), i=1,2,\left|\lambda_{i}(2)\right|=1$ and $\left|\lambda_{i}(a)\right|>1$ for $a>2$, where $i=1,2$. Introduce the notation $q=q(a)=q_{1}(a)$ and denote by $p=p(a)$ the eigenvector of the transposed matrix $J^{\mathrm{T}}(a)$ corresponding to $\overline{\lambda(a)}$, normalized to $\langle p, q\rangle=1$, where $\left\langle\left(a_{1}, a_{2}\right),\left(b_{1}, b_{2}\right)\right\rangle=$ $\sum_{i=1}^{2} \bar{a}_{i} b_{i},\left(a_{1}, a_{2}\right),\left(b_{1}, b_{2}\right) \in \mathbb{C}^{2}$. We obtain $\bar{p}=d(\lambda-1,1)$, where $d=d(a)=(2 \lambda(a)-$ $1)^{-1}$.

Introduce the vector $U=(u, v)^{\mathrm{T}}$ and the complex variable $z=\langle p, U\rangle$. The variable $U$ can also be expressed by $z$ :

$$
\binom{u}{v}=U=q(a) z+\overline{q(a)} \bar{z}=\binom{z+\bar{z}}{\lambda(a) z+\overline{\lambda(a)} \bar{z}} .
$$

Moreover, map (2) can be written in the following form:

$$
z \mapsto\left\langle p(a), J(a) U+f_{a}(U)\right\rangle=\lambda(a) z+d(a) g_{a}(z)
$$

where $g_{a}: \mathbb{C} \rightarrow \mathbb{R}$ is the following

$$
g_{a}(z)=g_{a}(z, \bar{z})=-a(z+\bar{z})(\lambda(a) z+\overline{\lambda(a)} \bar{z})
$$

At first we use the map

$$
\begin{equation*}
z \mapsto \lambda(a) z+d(a) g_{a}(z):=G(z) \tag{3}
\end{equation*}
$$

without further transformation to construct $\mathcal{N}(a)$.
Proposition 3.1: For every $a \in\left(\frac{3}{2}, 2\right)$ define $\varepsilon(a)$ by

$$
\varepsilon(a)=\frac{\sqrt{4 a-5}(1-\sqrt{a-1})}{a(2 \sqrt{a-1}+1)} \cdot \frac{\sqrt{4 a-5}}{\sqrt{a+1}}
$$

Then the set

$$
\mathcal{N}(a)=\{(x, y) \in S:|x-A|,|y-A| \leq \varepsilon(a)\}
$$

is in the region of attraction of the fixed point $(A, A)$ of $F_{a}$.
Proof: At first we show, there exists a $\zeta_{0}>0$, such that $\left|G_{a}(z)\right|<|z|$ for every $0<|z|<\zeta_{0}$. If such a $\zeta_{0}$ exists, it is clear that the open ball $B_{\zeta_{0}}^{\circ}$ around the origin is invariant and we show that every point of $B_{\zeta_{0}}^{\circ}$ tends to the origin. Let $z_{0}$ be an arbitrary point from $B_{\zeta_{0}}^{\circ}$ and consider the nonnegative, strictly decreasing sequence $\left\{\left|z_{n}\right|\right\}_{n=0}^{\infty}$, where $z_{n+1}=G_{a}\left(z_{n}\right)$. This sequence can converge only to a fixed point of the continuous map $r \mapsto \max _{|\zeta|=r}\left|G_{a}(\zeta)\right|$, which is, inside $B_{\zeta_{0}}^{\circ}$, solely $r=0$.

Estimate the right-hand side of map (3). Using $|\lambda(a)|=\sqrt{a-1},|d(a)|=1 /(\sqrt{4 a-5})$ and $\left|g_{a}(z)\right| \leq a(2|\lambda(a)|+1)|z|^{2}=a(2 \sqrt{a-1}+1)|z|^{2}$, we obtain

$$
\left|\lambda(a) z+d(a) g_{a}(z)\right| \leq|z|\left(\sqrt{a-1}+\frac{a(2 \sqrt{a-1}+1)}{\sqrt{4 a-5}}|z|\right)<|z|
$$

for every $z \neq 0$, provided $|z|<\zeta_{0}:=(\sqrt{4 a-5}(1-\sqrt{a-1})) /(a(2 \sqrt{a-1}+1))$.
To obtain an estimation of the real variables $u$ and $v$, we use the expression $z=$ $\langle p(a), U\rangle=d((\lambda-1) u+v)$. Supposing $|u|,|v| \leq \delta$, we obtain

$$
|z| \leq|d||(\lambda-1) u+v|=|d| \sqrt{\frac{4 a-5}{4} u^{2}+\left(v-\frac{u}{2}\right)^{2}} \leq \delta \frac{\sqrt{a+1}}{\sqrt{4 a-5}}
$$

therefore, if $\delta \leq \zeta_{0}(\sqrt{4 a-5}) /(\sqrt{a+1})$, then $|z| \leq \zeta_{0}$. Set $\varepsilon(a)=\zeta_{0}(\sqrt{4 a-5}) /(\sqrt{a+1})$. Then points, whose coordinates satisfy $|u|,|v| \leq \varepsilon(a)$, are in the region of attraction of the fixed point.

It is easy to see the set $\mathcal{N}$ shrinks to the fixed point $(A, A)$ as $a$ tends to 2 , since $\lim _{a \rightarrow 2} \varepsilon(a)=0$. Consequently, close to the critical parameter value, the neighbourhood, obtained by linearization, is not suitable for reliable numerical methods. In fact, the smaller the neighbourhood, the less efficient, and more time-consuming the numerical part of the proof. Furthermore, the linearization does not provide an attractive neighbourhood
at the critical parameter value $a=2$; therefore, we need another approach to construct a neighbourhood $\mathcal{N}$ for parameter values close to 2 .

In the subsequent section, we use the normal form of the Neimark-Sacker bifurcation and create a neighbourhood whose size is independent of $a$. Actually, the first method with the linearization become rather compute-intensive at about the parameter range (1.99, 2), but we will use the second technique with the normal form in a bigger parameter range, namely for $a \in(1.95,2]$. The normal form technique provides a significantly larger neighbourhood than the first method can do for parameters close to the critical value, so the second method is more efficient even for $a \in(1.95,1.99]$, too.

## 4. Transforming to normal form

In this section, first, we give a general method to construct an attractive neighbourhood around a fixed point, which undergoes a supercritical Neimark-Sacker bifurcation at $a_{0}$. This neighbourhood is suitable for parameters close to the critical value $a_{0}$, ie for $a \in\left[a_{0}-\beta_{0}, a_{0}\right]$ with some fixed $\beta_{0}>0$. We follow the steps of finding the normal form of the Neimark-Sacker bifurcation, according to Kuznetsov [9].

Suppose, we have a map

$$
\begin{equation*}
x \mapsto F_{a}(x), \tag{4}
\end{equation*}
$$

where $x \in \mathbb{R}^{2}, F_{a}$ is smooth and $a \in \mathbb{R}$ is the parameter. Furthermore, we have a fixed point $\tilde{x}=\tilde{x}(a)$, which undergoes a supercritical Neimark-Sacker bifurcation at $a_{0}$. Fix some $\beta_{0}>0$. According to Kuznetsov, if $|\lambda(a)|<1$ for all $a \in\left[a_{0}-\beta_{0}, a_{0}\right.$ ), then map (4) can be transformed into the following form:

$$
\begin{equation*}
z \mapsto G(z)=\lambda(a) z+G_{2}(z, a) \tag{5}
\end{equation*}
$$

where $z \in \mathbb{C}$, and $G_{2}$ is smooth (cf. Section 3).
We can write the smooth $G(z)$ as a formal Taylor series in two complex variables ( $z$ and $\bar{z}):$

$$
G(z)=\lambda(a) z+\sum_{k+l=2}^{4} \frac{g_{k l}}{k!l!} z^{k} \bar{z}^{l}+R_{1}
$$

where $g_{k l}=g_{k l}(a)$ and $R_{1}=R_{1}(z, \bar{z}, a)=O\left(|z|^{5}\right)$. Then, with smooth and invertible functions, we transform map (5) into the normal form of the bifurcation:

$$
\begin{equation*}
w \mapsto \lambda(a) w+c_{1}(a) w^{2} \bar{w}+R_{2} \tag{6}
\end{equation*}
$$

where $R_{2}=R_{2}(w, \bar{w}, a)=O\left(|w|^{4}\right)$. If we show that there exists $\rho_{0}>0$, such that for every $0<|w|<\rho_{0}$ and $a \in\left[a_{0}-\beta_{0}, a_{0}\right]$ the following holds:

$$
\begin{equation*}
\left|\lambda(a) w+c_{1}(a) w^{2} \bar{w}+R_{2}\right|<|w| \tag{7}
\end{equation*}
$$

then we obtain that $B_{\rho_{0}}=\left\{w \in \mathbb{C}:|w|<\rho_{0}\right\}$ is in the region of attraction of the fixed point 0 of map (6). Since $|\lambda| \leq 1$ for $a \leq a_{0}$, and the bifurcation is supercritical, ie $\Re c_{1}\left(a_{0}\right) / \lambda\left(a_{0}\right)<0$, it is easy to see that inequality (7) holds for all sufficiently small $\rho_{0}$ and $\beta_{0}$.

Our aim is to obtain an explicit value for $\rho_{0}$ assuming $\beta_{0}$ is given. Furthermore, $\rho_{0}$ needs to be as big as possible, because of the computer-assisted part of the proof. Consequently, the estimation of the higher order terms $\left(R_{2}\right)$ is the most essential part of the method, just like in the linearized case. Note that, in the end, we need to derive a $\{z:|z|<\varepsilon\}$-type neighbourhood, related to original map (5).

To obtain the normal form, we look for a smooth invertible function $h: \mathbb{C} \rightarrow \mathbb{C}$ in a neighbourhood of $0 \in \mathbb{C}$ which transforms map (5) with the new coordinate $w=h^{-1}(z)$ into the following form:

$$
\begin{equation*}
w \mapsto h^{-1}(G(h(w)))=\lambda(a) w+c_{1}(a) w^{2} \bar{w}+R_{2} . \tag{8}
\end{equation*}
$$

According to Kuznetsov [9], such a function can be found in the form:

$$
\begin{equation*}
h(w)=w+\frac{h_{20}}{2} w^{2}+h_{11} w \bar{w}+\frac{h_{02}}{2} \bar{w}^{2}+\frac{h_{30}}{6} w^{3}+\frac{h_{12}}{2} w \bar{w}^{2}+\frac{h_{03}}{6} \bar{w}^{3}, \tag{9}
\end{equation*}
$$

where $h_{i j}=h_{i j}(a)$. To this transformation, we need the non-resonance condition

$$
\left(\frac{\lambda\left(a_{0}\right)}{\left|\lambda\left(a_{0}\right)\right|}\right)^{k} \neq 1, \quad \text { where } k \in 1,2,3,4
$$

Clearly, $h$ has an inverse in a small neighbourhood of $0 \in \mathbb{C}$, and $h^{-1}$ can be written in the following form:

$$
\begin{equation*}
h^{-1}(z)=h_{0}^{-1}(z)+R_{3} \tag{10}
\end{equation*}
$$

where

$$
h_{0}^{-1}(z)=z+\sum_{2 \leq k+l \leq 4} \tilde{h}_{k l} z^{k} z^{l}
$$

$R_{3}=R_{3}(z, \bar{z}, a)=O\left(|z|^{5}\right)$ and $\tilde{h}_{k l}=\tilde{h}_{k l}(a)$. The coefficient $\tilde{h}_{k l}$ can be obtained by substituting $w=h^{-1}(z)$ into $z=h(w)$ and equating the coefficients of the same type up to fourth order. The $h_{i j}$ was obtained in a similar manner: we need to choose the coefficients so that the second- and third-order terms (apart from $\left.w^{2} \bar{w}\right)$ of $h^{-1}(G(h(w)))$ are eliminated. The formulas can be found in the Appendix. Notice that $h_{k l}$ and consequently $\tilde{h}_{k l}$ depend only on the at most third-order terms of $G$.

First, we will give a finite-order polynomial estimation on the functions $G, h$ and $h^{-1}$ :

$$
\begin{aligned}
|h(w)| & \leq|w|+h_{2}|w|^{2}+h_{3}|w|^{3} \\
|G(z)| & \leq|z|+g_{2}|z|^{2}+g_{3}|z|^{3}+g_{4}|z|^{4}+R_{10}|z|^{5} \\
\left|h^{-1}(z)\right| & \leq|z|+\tilde{h}_{2}|z|^{2}+\tilde{h}_{3}|z|^{3}+\tilde{h}_{4}|z|^{4}+R_{30}|z|^{5}
\end{aligned}
$$

where the coefficients are independent of $a$. With them we can give an estimation on $R_{2}$, ie the higher order terms of the composition $h^{-1}(G(h(w)))$. Clearly the Taylor expansion of $h$ is finite, but generally the other two Taylor expansions have infinitely many terms. So the at least fifth-order terms are estimated in $R_{10}|z|^{5}$ and $R_{30}|z|^{5}$. For the lower order terms, we have explicit formulae, and they could be estimated by interval arithmetic. As for the higher order terms, it is essential to be able to say how large can be the moduli of


Figure 3. The size of the domains of $h, G$ and $h^{-1}$ provided that $|w|<\rho_{0}$.
$h(w), G(h(w))$ and $h^{-1}(G(h(w)))$ in (8) if $|w|<\rho_{0}$ is assumed, since the estimation of the remaining terms of a Taylor expansion highly depends on the size of the neighbourhood on which it needs to be valid. The radii $\rho_{1}, \rho_{2}$ and $\rho_{3}$ must be chosen so that $h\left(B_{\rho_{0}}\right) \subseteq B_{\rho_{1}}$, $G\left(B_{\rho_{1}}\right) \subseteq B_{\rho_{2}}$ and $h^{-1}\left(B_{\rho_{2}}\right) \subseteq B_{\rho_{3}}$ (see Figure 3); consequently, during the study of $G$ and $h^{-1}$ we can assume that the domains are in $B_{\rho_{1}}$ and $B_{\rho_{2}}$, respectively.

After gaining an estimation on $R_{2}$, we show that inequality (7) holds for $|w| \in\left(0, \rho_{0}\right]$. From this result a neighbourhood in the $z$-plane can easily be obtained: the set $B_{\varepsilon}=$ $\{z \in \mathbf{C}:|z| \leq \varepsilon\}$ is inside the attractive neighbourhood of the fixed point of map (5) if $h^{-1}\left(B_{\varepsilon}\right) \subseteq B_{\rho_{0}}$, ie $B_{\varepsilon}$ is mapped inside the region of attraction of map (6).

Here, we emphasize that for our calculations the only thing we need to know is the at most fourth-order terms of the function $G(z)$ and an $R_{10}|z|^{5}$-type estimation of the at least fifth-order terms of $G(z)$.

Until this point in the section, we described the method for a general $F_{a}(x)$. Now, we turn our attention to the specific $F_{a}(x, y)=(y, a y(1-x))$ from (1).

The main results of this section are the following two propositions. We prove only Proposition 4.2 as the whole argument can be repeated to get an attracting neighbourhood when only $a \in[1.95,2]$ is assumed. The differences appear only in concrete values in the given estimations. Details of Proposition 4.1 can be found on our website [12].

Proposition 4.1: For all fixed $a \in[1.95,2]$, the $\operatorname{set}\{z \in \mathbb{C}:|z| \leq 0.013\}$ belongs to the basin of attraction of the fixed point 0 of $G(z)$.

Proposition 4.2: For all fixed $a \in[1.995,2]$, the set $\{z \in \mathbb{C}:|z| \leq 0.014\}$ belongs to the basin of attraction of the fixed point 0 of $G(z)$.

Proof: Throughout the proof we suppose $a \in\left[a_{0}-\beta_{0}, a_{0}\right]$, where $\beta_{0}=0.005$ and $a_{0}=2$. In our calculations, we use symbolic computation and built-in interval arithmetic tools of Wolfram Mathematica v. 11.

### 4.1. Estimation of the lower order terms in $G$, $h$ and $h^{-1}$

Throughout this section, we need an estimation of the coefficients of the lower order terms in $G, h$ and $h^{-1}$, such that these estimations are independent of $a \in[1.995,2]$. We use interval arithmetic tools to compute them for $a \in[1.995,2]$.

In our particular case, the function $G(z)$ can be written in the following form:

$$
G(z)=\lambda(a) z+\sum_{k+l=2} \frac{g_{k l}}{k!l!} z^{k} \bar{z}^{l},
$$

since $G$ has only at most second-order terms. Furthermore, we use (9) and (10). We look for constants satisfying the following inequalities:

$$
\begin{aligned}
& g_{2} \geq \max _{a \in\left[a_{0}-\beta_{0}, a_{0}\right]}\left(\frac{\left|g_{20}\right|}{2}+\left|g_{11}\right|+\frac{\left|g_{02}\right|}{2}\right), \\
& h_{2} \geq \max _{a \in\left[a_{0}-\beta_{0}, a_{0}\right]}\left(\frac{\left|h_{20}\right|}{2}+\left|h_{11}\right|+\frac{\left|h_{02}\right|}{2}\right), \\
& h_{3} \geq \max _{a \in\left[a_{0}-\beta_{0}, a_{0}\right]}\left(\frac{\left|h_{30}\right|}{6}+\frac{\left|h_{12}\right|}{2}+\frac{\left|h_{03}\right|}{6}\right), \\
& \tilde{h}_{n} \geq \max _{a \in\left[a_{0}-\beta_{0}, a_{0}\right]}\left(\sum_{i+j=n}\left|\tilde{h}_{i j}\right|\right),
\end{aligned}
$$

where $n=2,3,4$. With interval arithmetic it can be shown that $g_{2}=3.47, h_{2}=2.9, h_{3}=$ 4.7, $\tilde{h}_{2}=h_{2}=2.9, \tilde{h}_{3}=8.2$ and $\tilde{h}_{4}=30$ fulfil the requirements. From the definition of these constants, we obtain the following estimations:

$$
\begin{align*}
&|h(w)| \leq h^{\text {max }}(|w|) \\
&|G(z)| \leq|w|+h_{2}|w|^{2}+h_{3}|w|^{3}, \\
&\left|h_{0}^{-1}(z)\right| \leq \tilde{h}_{0}^{\text {max }}(|z|):=|z|+g_{2}|z|^{2}, \\
&\left|h^{-1}(z)\right| \leq|z|+\tilde{h}_{2}^{\text {max }}(|z|):=|z|+\tilde{h}_{2}|z|^{2}+\tilde{h}_{3}|z|^{3}+\tilde{h}_{4}|z|^{3}+\tilde{h}_{4}|z|^{4}+R_{30}|z|^{5}, \tag{11}
\end{align*}
$$

if in the last equation $R_{30}$ satisfies $\left|R_{3}\right| \leq R_{30}|z|^{5}$. We will determine $R_{30}$ later.
From the definition of $h_{2}$ and $h_{3}$, we also obtain

$$
\begin{equation*}
|w|-h_{2}|w|^{2}-h_{3}|w|^{3} \leq|h(w)| . \tag{12}
\end{equation*}
$$

Consequently, assuming $|w| \leq \rho$, we can make the following estimation:

$$
\begin{equation*}
|w| \leq \eta(\rho)|h(w)|, \tag{13}
\end{equation*}
$$

with

$$
\eta(\rho):=\frac{1}{1-h_{2} \rho-h_{3} \rho^{2}} .
$$

We choose $\rho_{0}=0.015, \rho_{1}=h^{\max }\left(\rho_{0}\right), \rho_{2}=G^{\max }\left(\rho_{1}\right)$ and from (12) it is clear that $B_{\rho_{2}}$ cannot be mapped outside of the circle with radius 0.018 ; consequently, this value is a suitable choice for $\rho_{3}$.

### 4.2. The domain of $h$ and $h^{-1}$

Now, we show that $h$ is injective in $\overline{B_{1 / 9}} \subseteq \mathbb{C}$, and $h^{-1}$ is defined on $\overline{B_{1 / 16}}$. Let $z \in \mathbb{C}$, $a \in[1.995,2]$ be fixed and denote $H_{a, z}: \mathbb{C} \ni w \mapsto w+z-h(w) \in \mathbb{C}$. With this notation $H_{a, z}=w$ if and only if $h(w)=z$.

$$
\begin{aligned}
\left|H_{a, z}\left(w_{1}\right)-H_{a, z}\left(w_{2}\right)\right| & =\left|w_{1}-h\left(w_{1}\right)-w_{2}+h\left(w_{2}\right)\right| \\
& \leq\left|w_{1}-w_{2}\right|\left(h_{2}\left(\left|w_{1}\right|+\left|w_{2}\right|\right)+h_{3}\left(\left|w_{1}\right|^{2}+\left|w_{1}\right| \cdot\left|w_{2}\right|+\left|w_{2}\right|^{2}\right)\right)
\end{aligned}
$$

If $|w| \leq \delta_{1}$ and $|z| \leq \delta_{2}$, then

$$
\left|H_{a, z}\left(w_{1}\right)-H_{a, z}\left(w_{2}\right)\right| \leq\left|w_{1}-w_{2}\right|\left(2 \delta_{1} h_{2}+3 \delta_{1}^{2} h_{3}\right)
$$

and

$$
\left|H_{a, z}(w)\right| \leq \delta_{2}+\delta_{1}^{2} h_{2}+\delta_{1}^{3} h_{3}
$$

Choosing $\delta_{1}=\frac{1}{9}$ and $\delta_{2}=\frac{1}{16}$ the map $H_{a, z}$ is a contraction mapping $\overline{B_{1 / 9}}$ into itself. Consequently, for every $z \in \overline{B_{1 / 16}}$ there exists only one $w=w(z) \in \overline{B_{1 / 9}}$ such that $h(w(z))=z$, ie $h^{-1}$ is defined on $\overline{B_{1 / 16}}$.

It is clear that $\rho_{0}, \rho_{3}<\delta_{1}$ and $\rho_{1}, \rho_{2}<\delta_{2}$, where $\rho_{0}, \rho_{1}, \rho_{2}, \rho_{3}$ were chosen at the end of the previous subsection.

### 4.3. The estimation of the higher order terms in $h^{-1}$

Now, we turn our attention to the estimation of $R_{3}$ in (10), which consists of the fifth and higher order terms of $h^{-1}$. We need an estimation $\left|R_{3}(z)\right|<R_{30}|z|^{5}$, assuming $|z| \leq$ $\rho_{2}$. But first, we give an estimation of type $\left|R_{3}(h(w))\right|<C|w|^{5}$, assuming $|w| \leq \rho_{3}$ (see Figure 3). Using the definition of $h_{0}^{-1}$ and $h$, it follows from (10) that

$$
R_{3}(h(w))=w-h_{0}^{-1}(h(w))=\sum_{5 \leq k+l \leq 12} r_{3}^{k l}(a) w^{k} \bar{w}^{l},
$$

since it is a 12th-order polynomial of $w$ and $\bar{w}$, which has only 5 th and higher order terms.
Consider the composition

$$
\tilde{R}_{3}(|w|)=\tilde{h}_{0}^{\max }\left(h^{\max }(|w|)\right)=\sum_{j=5}^{12} r_{3}^{j}|w|^{j}
$$

of the real functions $h^{\max }, \tilde{h}_{0}^{\max }$. It is clear that $\sum_{k+l=j}\left|r_{3}^{k l}\right| \leq r_{3}^{j}$ holds for $5 \leq j \leq 12$. Consequently,

$$
\left|R_{3}(h(w))\right| \leq \sum_{5 \leq k+l \leq 12}\left|r_{3}^{k l}(a)\right||w|^{k+l} \leq \sum_{5 \leq j \leq 12} r_{3}^{j}|w|^{j} \leq \sum_{5 \leq j \leq 12} r_{3}^{j} \rho_{3}^{j-5}|w|^{5},
$$

assuming $|w| \leq \rho_{3}$. Using (13), we gain $|w|<\eta\left(\rho_{3}\right)|z|$, and

$$
\left|R_{3}(z)\right| \leq \sum_{5 \leq j \leq 12} r_{3}^{j} \rho_{3}^{j-5}\left(\eta\left(\rho_{3}\right)\right)^{5}|z|^{5} \leq 1070|z|^{5},
$$

therefore, $R_{30}=1070$ is a suitable choice.

### 4.4. The estimation of the higher order terms in the normal form

Now we turn our attention to $R_{2}$, which estimates the at least fourth-order terms of $h^{-1}(G(h(w)))$. To obtain a better estimation, we handle the fourth-order terms $\left(R_{24}\right)$ and the higher order ones $\left(R_{25}\right)$ separately. Set $R_{2}(w)=R_{24}(w)+R_{25}(w)$.

The fourth-order coefficients $r_{2}^{k l}(a)$ (where $\left.k+l=4\right)$ of $h^{-1}(G(h(w)))$ can be calculated explicitly; the formulae of $\left|r_{2}^{k l}(a)\right|$ can be found in Appendix. With interval arithmetic, it can be shown that $\sum_{k+l=4}\left|r_{2}^{k l}(a)\right| \leq 40$; consequently, $\left|R_{24}(w)\right| \leq 40|w|^{4}$.

As for the higher order terms, we use $h^{\max }, G^{\max }$ and $\tilde{h}^{\max }$ (from (11)), similarly to the estimation of $R_{3}$. Consider the composition

$$
\tilde{R}_{2}(|w|)=\tilde{h}^{\max }\left(G^{\max }\left(h^{\max }(|w|)\right)\right)=\sum_{j=1}^{30} r_{2}^{j}|w|^{j} .
$$

It is clear that for $|w| \leq \rho_{0}$ the following holds:

$$
\left|R_{25}(w)\right| \leq \sum_{k=5}^{30} r_{2}^{k}(|w|) \rho_{0}^{k-4}|w|^{4} \leq 90|w|^{4}
$$

Combining these two results, for $|w| \leq \rho_{0}$, we obtain

$$
\left|R_{2}(w)\right| \leq\left|R_{24}(w)\right|+\left|R_{25}(w)\right| \leq 130|w|^{4},
$$

and consequently $R_{20}=130$.

### 4.5. The attracting neighbourhood

Now, with our previous estimation on $R_{2}$ we can finish our proof. Since

$$
\left|\lambda(a) w+c_{1}(a) w^{2} \bar{w}+R_{2}\right| \leq|w|\left(\left.\left.| | \lambda\left|+\tilde{c}_{1}\right| w\right|^{2}\left|+R_{20}\right| w\right|^{3}\right),
$$

where $\tilde{c}_{1}=|\lambda(a)| / \lambda(a) c_{1}(a)$ and $\lambda=\lambda(a)$, we only need to prove

$$
\begin{equation*}
\left.\left.\left||\lambda|+\tilde{c}_{1}\right| w\right|^{2}\left|+R_{20}\right| w\right|^{3}<1 \tag{14}
\end{equation*}
$$

for every $|w| \leq \rho_{0}$ and $a \in\left[a_{0}-\beta_{0}, a_{0}\right]$.
To this end, we show that the following inequality holds with a suitable $R_{4}>0$ :

$$
\left.\left.\left||\lambda|+\tilde{c}_{1}\right| w\right|^{2}\left|\leq|\lambda|+\left(\mathfrak{R} \tilde{c}_{1}\right)\right| w\right|^{2}+R_{4}|w|^{3},
$$

or equivalently

$$
\begin{equation*}
0 \leq 2 R_{4}|\lambda|-\left(\Im_{\Im} \tilde{c}_{1}\right)^{2}|w|+2 R_{4}\left(\Re \tilde{c}_{1}\right)|w|^{2}+R_{4}^{2}|w|^{3}, \tag{15}
\end{equation*}
$$

for every $|w| \leq \rho_{0}$ and $a \in\left[a_{0}-\beta_{0}, a_{0}\right]$. For $a \in\left[a_{0}-\beta_{0}, a_{0}\right]$ we can make the following estimations with interval arithmetic on the coefficients in (15) depending on $a$ : $\mathfrak{R} \tilde{c}_{1}$ and $\Im \tilde{c}_{1}$ are negative, $\left|\mathfrak{i} \tilde{c}_{1}\right| \leq 2.1,\left|\Im \tilde{c}_{1}\right| \leq 3.5$ and $|\lambda| \geq 0.99$. From this it is clear that for $|w| \leq \rho_{0}$
the choice $R_{4}=0.1$ will be suitable. Therefore, the left-hand side of inequality (14) can be written in the following form:

$$
\begin{aligned}
\left.\left.\left||\lambda|+\tilde{c}_{1}\right| w\right|^{2}\left|+R_{20}\right| w\right|^{3} & \leq\left(|\lambda|+\Re \tilde{c}_{1}|w|^{2}\right)+\left(R_{4}+R_{20}\right)|w|^{3} \\
& \leq 1+\left(\Re \tilde{c}_{1}+\left(R_{4}+R_{20}\right)|w|\right)|w|^{2}
\end{aligned}
$$

which is less than 1 , provided

$$
|w|<\frac{-\mathfrak{R} \tilde{c}_{1}}{R_{4}+R_{20}}
$$

Using the fact that $\left|\mathfrak{i} \tilde{c}_{1}\right| \geq 2$, we obtain $\left(-\mathfrak{R} \tilde{c}_{1}\right) /\left(R_{4}+R_{20}\right)>\rho_{0}$; therefore, inequality (14) holds for every $|w|<\rho_{0}$. From $\left|h_{a}^{-1}(z)\right|<h_{i n v}^{\max }(|z|)$, inequality $|z|<\varepsilon_{G}:=0.014$ implies $|w|=\left|h^{-1}(z)\right|<\rho_{0}$, so the proof of Proposition 4.2 is complete.

To obtain a neighbourhood in the real coordinate system $(u, v)$ we use $z=\langle p(a), U\rangle$, just like in the linearized case. Note that the set $\left\{z \in \mathbb{C}:|z| \leq \varepsilon_{G}\right\}$ will be transformed into an ellipse-shaped neighbourhood in the $u v$-plane.

## 5. Graph representation

In the computer-assisted part, we follow the method in [4]. In this section (for $\frac{3}{2}<a \leq 2$ ) we associate map (1) with a directed graph, which reflects the behaviour of the map up to a given resolution. Therefore, we can derive the properties of our dynamical system through the study of this graph. More precisely our aim is to show with the help of this graph that every point of $S \backslash \mathcal{N}$ enters the attracting neighbourhood $\mathcal{N}$ of the nontrivial fixed point constructed in the previous sections.

Let $D$ be a subset of $\mathbb{R}^{n}$. A set $\mathfrak{S}$ is called a cover of $D$, if the elements of $\mathfrak{S}$ are subsets of $\mathbb{R}^{n}$ and $\cup_{s \in \mathfrak{S}^{\mathcal{S}}} \supset D$. Let a map $f: D_{f} \subseteq \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$, a subset $D \subseteq D_{f}$ and a cover $\mathfrak{S}$ of $D$ be given. The directed graph $G(V, E)$ is called a graph representation of $f$ on $D$ with respect to $\mathfrak{S}$, if there exists a bijection $i: V \rightarrow \mathfrak{S}$, such that the following implication is true for all $u, v \in V$ :

$$
f(i(u) \cap D) \cap i(v) \cap D \neq \emptyset \Rightarrow(u, v) \in E .
$$

The meaning of the implication in the previous definition is the following. If we can get with $\operatorname{map} f$ from an element $s_{1}$ of the cover to another (possibly the same) element $s_{2}$ of it, ie there exists $x \in s_{1}$ and $y \in s_{2}$ such that $f(x)=y$, then there is an edge between the vertices corresponding to the two sets, more precisely $(u, v) \in E$ for $s_{1}=i(u)$ and $s_{2}=i(v)$. The reverse implication is not necessarily true, namely if there is a directed edge between the vertices $u$ and $v$, it is not sure there exists $x \in s_{1}$ such that $f(x) \in s_{2}$, where $s_{1}$ and $s_{2}$ are the corresponding sets to $u$ and $v$.

It is easy to see the implication above can be reformulated as follows. For every $u \in V$

$$
\begin{equation*}
f(i(u) \cap D) \subseteq \bigcup_{v \in K_{u}} i(v) \cap D \tag{16}
\end{equation*}
$$

where $K_{u}$ denotes the set of vertices, into which there is an edge from $u$ in graph $G$. So the sets corresponding to vertices in $K_{u}$ need to form a cover of the image of $i(u)$. From this it can be seen the graph representation can be regarded as some kind of upper estimation of
the original map $f$. The finer the cover is, the better the graph representation approximates the map. Therefore, if we would like to determine the possible location of the image of a point $P \in \mathbb{R}^{n}$ under $f$, we can do it with the help of the graph, since $f(P) \in \cup_{v \in K_{u}} i(v) \cap D$ for $P \in i(u)$. This means iterating $f$ the point $P$ can move forward only along the edges, ie it can move from an element of the cover to another one only if there is an edge between the two vertices corresponding to them. Consequently, we can draw conclusions regarding the possible location of the iterates of a point studying only the graph. In the following we take the liberty to handle the elements of the cover as vertices and vice versa, omitting the use of $i$.

The construction of the graph representation in our case is the following. For a fixed $k \in \mathbb{N}$ we divide the unit square $[0,1] \times[0,1]$ parallel to the sides into $2^{k} \times 2^{k}$ pieces of small closed squares with side length $r=2^{-k}$. According to Proposition 2.5 we only need to consider the squares lying in $\left[\frac{1}{16}, \frac{7}{8}\right]^{2}$. The cover $\mathfrak{S}$ of $S$ consists of these sets. The small squares correspond to the vertices of the graph. As for the edges, for every small square $s$ we construct a rectangle with reliable numerical methods which contains $f(s)$. If the rectangle intersects the small square $s_{2}$, then there is an edge from $s$ to $s_{2}$. It is clear that this construction satisfies relation (16). Note that we considered only that part of the rectangle obtained by the numerical method, which lies inside the square $\left[\frac{1}{16}, \frac{7}{8}\right]^{2}$, but this is not a restriction, since the studied set $\tilde{S}$ is invariant under map (1), so getting out of the unit square is only the consequence of the numerical method and the 'upper estimation' nature of the graph representation. Note also that instead of map (1) we use the second iterate of it, since the formula is still compact enough not to cause big overestimation in interval arithmetic, and it considerably speeds up the calculations.

In this paper, we suppose a graph is always finite. A graph is strongly connected if there are $u v$ and $v u$ (directed) paths for every $u \neq v$ vertices of the graph. We use the following decomposition of a directed graph [see 2].

Proposition 5.1: The vertices of a directed graph can be classified, and the classes can be ordered such that

- the subgraphs spanned by the classes are strongly connected, and
- for every directed edge between these classes, the class of the tail of this edge precedes the class of the head of it;
moreover, the partition above is unique.
The aforementioned classes are called the strongly connected components (SCC) of the graph. A strongly connected component is called non-essential, if it consists of one vertex without a loop. Otherwise, we call it essential.

From the graph representation and from Proposition 5.1, it is clear what happens to an arbitrary point of $S$ during the iterations. Starting from a small square containing this point it moves to an other (possibly the same) small square along a directed edge. If we are not in an essential SCC, we step out of this small square not returning to here afterwards because of the ordering of the SCCs. If we are in an essential SCC, it can happen that the point stays here forever, or the point steps out of this SCC, but in this case it cannot return to this SCC anymore.

Since during the partition we obtain finitely many small squares, and consequently, the graph is finite; it is straightforward that for every point of $S$ there exists an essential strongly connected component, which the point enters during the iteration and never leaves it. So it is true for every $x \in S$ that it enters an essential SCC with finitely many steps and stays here afterwards ; therefore, we only need to study the essential SCCs.

Our aim is to show that those essential SCCs, in which the points of $S$ can get stuck, are in the attracting neighbourhood $\mathcal{N}$ of the fixed point $(A, A)$, which neighbourhood was constructed analytically in the previous sections. It is important to note that, it is possible for some essential SCC that none of the points of $S$ can get stuck here. Actually, this would be the case close to the trivial fixed point $(0,0)$, since it is a saddle; that shows the necessity of Proposition 2.5 and $\tilde{S}$.

As a next step, we refine the partition as follows. We divide the small squares into four smaller squares that have a side length half as long as before, determine their images with reliable numerical methods and construct the SCCs again. Because of the properties of interval arithmetic (inclusion isotonicity: $I_{1} \subseteq I_{2} \Rightarrow F\left(I_{1}\right) \subseteq F\left(I_{2}\right)$, where $F$ is the interval-extension of $f$, [see 14]), if there is an edge between two new small squares, then there must be an edge between their predecessors with the same orientation. We come to the conclusion that during the refinement, an essential SCC can arise only from a former essential SCC ; therefore, it is really enough to trace merely the essential SCCs. Note that, with the refinements the graph representation becomes a more and more accurate approximation of the represented map, so an essential SCC can fall apart into smaller pieces, and it even can happen that none of the small squares born from a former essential SCC compose a new essential SCC, ie this cycle in the graph is only the consequence of the 'upper estimation' nature of the graph representation. We continue these refinement steps, until all the remaining SCCs are inside the region of attraction of the fixed point obtained in the previous section. If it occurs in finitely many steps, our main theorem is proven.

Finally, instead of checking after every refinement, whether the remaining SCCs are in the analytically constructed attracting neighbourhood $\mathcal{N}$, we can remove all the small squares lying entirely in $\mathcal{N}$ before the first refinement. In that case for a fixed $a$, the main theorem will be proved, if the set of the new SCCs will be empty after a refinement. We show the correctness of this method.

- If we erase a vertex which is a non-essential SCC, it has no effect at all compared to our former method (when checking after every refinement).
- If we remove a whole essential SCC, it also has no substantial effect, because during the checking it always would be in the attracting neighbourhood.
- The only significant change happens when we erase only a proper subgraph of an essential SCC. Consider such an SCC and colour blue the vertices we want to remove (and do not remove them yet). Delete the directed edges stemming from a blue vertex, then form the SCCs (referred to as new SCCs later on) of the new graph and order them such that the blue vertices are at the end of the ordering. (It can be done, since there are no edges from coloured to uncoloured vertices.)
(i) An uncoloured vertex can be in a new essential SCC; in that case they remain under study after the removal of the blue vertices, just as they would be in the original method.

```
Algorithm Proving the global stability of \((A, A)\) for the logistic map
    procedure Log2d
        \(V \leftarrow\) the initial partition
                                    \(r=2^{-10}\)
        \(E \leftarrow\) the edges construct them with reliable num. method
        \(C \leftarrow \operatorname{SCC}\) of directed \(\operatorname{graph}(V, E)\)
        remove the nonessential SCCs from \(V\)
        remove the SCC at the origin from \(V\) if possible
        remove the initial attracting neighbourhood from \(V\)
        repeat
            \(V \leftarrow \operatorname{refine}(V) \quad r \leftarrow r / 2\)
            \(E \leftarrow\) the edges
            \(C \leftarrow\) SCC of directed \(\operatorname{graph}(V, E)\)
            remove the nonessential SCCs from \(V\)
        until \(\left|B_{1}\right|=\emptyset\)
    end procedure
```

(ii) However, if an uncoloured vertex is a non-essential SCC it will be erased (as we keep only the essential SSCs), unlike in the method without deleting the vertices of the attracting neighbourhood, but this is not a problem because every point of this vertex enters a new SCC or a blue vertex (because of the ordering) in finitely many steps, so this vertex really can be deleted.

Note that the aforementioned method can be regarded as a proof, since the graph problems are finite, so the computer can work on them punctually, moreover the method used during the construction of edges was executed with reliable numerical methods, therefore if we have sufficiently much time, then we could reconstruct by hands the parts which were executed by the computer, and we would come to the same conclusion, if our estimation is as good as the computer's.

The program code and the outputs can be found on the link [12].

## 6. Completion of the proof

In the previous sections, we obtained an attracting neighbourhood and then a method to prove the global stability of the nontrivial fixed point for a fixed $a \in[1.5,2]$. In this section, we show how to modify our method to handle not only a single value of $[1.5,2]$ but a small subinterval of that, instead.

Let $[a]=\left[a^{-}, a^{+}\right] \subseteq[1.5,2]$ be a fixed small interval. First, we need a new attracting neighbourhood $\mathcal{N}([a])$, such that for every $a \in[a]$, the attracting neighbourhood $\mathcal{N}(a)$ contains this set, ie $\cap_{a \in[a]} \mathcal{N}(a) \supset \mathcal{N}([a])$. To this end, we need to take into consideration the displacement of the fixed point and the change in the size of the neighbourhood. Secondly, during the construction of the edges of the graph representation the number $a$ have to be replaced by the interval [a], since (1) and consequently its second iterate depends on $a$. So while studying the image of a small square $k_{1}$, we need to study it for every $a \in[a]$, ie

Table 1. The partition of the parameter range.

|  | Parameter | Size of slices | Shape of $\mathcal{N}$ | Parameters of $\mathcal{N}$ |
| :--- | :---: | :---: | :---: | :---: |
| $\mathcal{I}_{1}$ | $[1.5,1.95]$ | $2^{-10}$ | Rectangle | $7 \cdot 10^{-3}$ |
| $\mathcal{I}_{2}$ | $[1.95,1.995]$ | $2^{-13}$ | Ellipse | $\varepsilon_{G}=0.0138$ |
| $\mathcal{I}_{3}$ | $[1.995,2]$ | $2^{-16}$ | Ellipse | $\varepsilon_{G}=0.0146$ |

we take a set of small squares during the estimation of the image set such that they cover $f_{a}^{2}\left(k_{1}\right)$ for every $a \in[a]$.

We divide the interval [1.5,2] into subintervals $\mathcal{I}_{1}=[1.5,1.95], \mathcal{I}_{2}=[1.95,1.995]$ and $\mathcal{I}_{3}=[1.995,2]$, then divide further these intervals into smaller subintervals with length $2^{-10}, 2^{-13}$, and $2^{-16}$ respectively (see Table 1).

For small intervals in $\mathcal{I}_{1}$ we use the linearized map and the square-shaped neighbourhoods with side length $2 \varepsilon(a)$ (Proposition 3.1). It is easy to see the size of this set and the location of the fixed point also change as $a$ changes. However, it can be shown that $\varepsilon(a) \geq 0.007$ for every $a \in[1,1.95]$, so considering this value fixed, we need to handle only the displacement of the fixed point.

For small intervals in $\mathcal{I}_{2}$ and $\mathcal{I}_{3}$ we use the bifurcation normal form; therefore, the size of the ellipse-shaped neighbourhood is fixed (Propositions 4.1 and 4.2), so we only need to consider the displacement of the fixed point.

## 7. The algorithm

During the calculation of edges of the graph representation, we use the second iterate of original map (1):

$$
\begin{equation*}
\binom{x_{0}}{y_{0}} \mapsto\binom{x_{2}}{y_{2}}=\binom{a y_{0}\left(1-x_{0}\right)}{a^{2} y_{0}\left(1-x_{0}\right)\left(1-y_{0}\right)} . \tag{17}
\end{equation*}
$$

Regarding the examined parameter domain $\left[a^{-}, a^{+}\right]$and the sides $\left[x_{i}^{-}, x_{i}^{+}\right]$and $\left[y_{i}^{-}, y_{i}^{+}\right]$ of the squares as intervals, simply, we could use interval arithmetic tools, such as IntLab to compute the image of a small square. However, the map is quite simple, so we can accelerate this method as follows. Notice that $x_{2}^{-}=a^{-} y_{0}^{-}\left(1-x_{0}^{+}\right)$, so we only need to force the computer to use a downward rounding in order to guarantee that the obtained below estimation is really not larger than the possible first coordinates of the image of any point from the initial square. Similarly, we can estimate $x_{2}^{+}$, but at this time we use upward rounding. As for the $y_{2}^{-}$and $y_{2}^{+}$, remark that instead of $y_{0}^{-}\left(1-y_{0}^{+}\right)$in $y_{2}^{-}$we can use $\min \left\{y_{0}^{-}\left(1-y_{0}^{-}\right), y_{0}^{+}\left(1-y_{0}^{+}\right)\right\}$because $y_{0}$ denotes the same number in expression (17), and the function $x(1-x)$ is monotone on intervals which do not contain $\frac{1}{2}$ in the interior (and it is fulfilled in the partition). In $y_{2}^{+}$we replace the minimum by maximum, after that we proceed just like in the case of $x_{2}$.

We implemented our program in MATLAB, and used the built-in digraph function to construct the directed graph from the edge list and the conncomp function to divide the graph into strongly connected components.

Now, we can run our algorithm with parameters summarized in Table 1. As an example, for the parameter slice $\left[2-\frac{61}{2^{13}}, 2-\frac{60}{2^{13}}\right]$ we show the evolution of the remaining SCCs during the first 4 iterations on Figure 4.


Figure 4. The remaining vertices for $[a]=\left[2-\frac{61}{2^{13}}, 2-\frac{60}{2^{13}}\right]$. (a) Before the first refinement, (b) after 2 refinements and (c) after 4 refinements.

The program ran successfully ; therefore, we established the nontrivial fixed point that is globally attracting for $a \in[1.5,2]$. Combining this with Proposition 2.4 and the asymptotic stability for $a \in(1,2]$ the proof of Theorem 1.1 is completed.

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## Appendix. The formulae of the coefficients

$$
\begin{aligned}
& \lambda(a)=\frac{1}{2}+\frac{1}{2} i \sqrt{4 a-5} \\
& g_{20}(a)=-a+\frac{i}{\sqrt{4 a-5}}, \quad g_{11}(a)=\frac{i a}{\sqrt{4 a-5}}, \quad g_{02}(a)=a+\frac{i}{\sqrt{4 a-5}} \\
& h_{20}(a)=\frac{4 a}{4 a-5+i \sqrt{4 a-5}}, \quad h_{11}(a)=\frac{4 i a}{\sqrt{4 a-5}(-i+\sqrt{4 a-5})^{2}}, \\
& h_{02}(a)=\frac{a\left(i-\frac{1}{\sqrt{4 a-5}}\right)}{i-i a+\sqrt{4 a-5}}, \\
& h_{30}(a)=-\frac{12 i a^{2}}{-2(5 i+\sqrt{4 a-5})+a(13 i-\sqrt{4 a-5}+2 a(-2 i+\sqrt{4 a-5}))} \\
& h_{12}(a)=\frac{16 a^{2}(2-2 i \sqrt{4 a-5}+a(-5+2 a-\sqrt{4 a-5}))}{\sqrt{4 a-5}(-i+\sqrt{4 a-5})^{4}(-7 i-3 \sqrt{4 a-5}+a(7 i-i a+2 \sqrt{4 a-5}))} \\
& h_{03}(a)=\frac{96 a^{2}(-2+a+i \sqrt{4 a-5})(a-1+i \sqrt{4 a-5})^{-1}}{\sqrt{4 a-5}(1+i \sqrt{4 a-5})\left(12 i(a-1)-7 \sqrt{4 a-5}+\sqrt{(4 a-5)^{3}}\right)} \\
& \tilde{h}_{20}(a)=-\frac{h_{20}}{2}, \quad \tilde{h}_{11}(a)=-h_{11}, \quad \tilde{h}_{02}(a)=-\frac{h_{02}}{2} \\
& \tilde{h}_{30}(a)=\frac{1}{6}\left(3 h_{20}^{2}-h_{30}+3 h_{11} \bar{h}_{02}\right) \\
& \tilde{h}_{21}(a)=\frac{1}{2}\left(3 h_{11} h_{20}+h_{02} \bar{h}_{02}+2 h_{11} \bar{h}_{11}\right) \\
& \tilde{h}_{12}(a)=\frac{1}{2}\left(2 h_{11}^{2}-h_{12}+h_{02} h_{20}+2 h_{02} \bar{h}_{11}+h_{11} \bar{h}_{20}\right) \\
& \tilde{h}_{03}(a)=\frac{1}{6}\left(3 h_{02} h_{11}-h_{03}+3 h_{02} \bar{h}_{20}\right) \\
& \tilde{h}_{40}(a)=\frac{1}{24}\left(-15 h_{20}^{3}+10 h_{20} h_{30}-30 h_{11} h_{20} \bar{h}_{02}-3 h_{02} \bar{h}_{02}^{2}+4 h_{11} \bar{h}_{03}-12 h_{11} \bar{h}_{02} \bar{h}_{11}\right) \\
& \tilde{h}_{31}(a)=\frac{1}{6}\left(-15 h_{11} h_{20}^{2}+4 h_{11} h_{30}-12 h_{11}^{2} \bar{h}_{02}+3 h_{12} \bar{h}_{02}-6 h_{02} h_{20} \bar{h}_{02}+h_{02} \bar{h}_{03}\right. \\
& \left.-12 h_{11} h_{20} \bar{h}_{11}-6 h_{02} \bar{h}_{02} \bar{h}_{11}-6 h_{11} \bar{h}_{11}^{2}+3 h_{11} \bar{h}_{12}-3 h_{11} \bar{h}_{02} \bar{h}_{20}\right) \\
& \tilde{h}_{22}(a)=\frac{1}{4}\left(-12 h_{11}^{2} h_{20}+3 h_{12} h_{20}-3 h_{02} h_{20}^{2}+h_{02} h_{30}+h_{03} \bar{h}_{02}-9 h_{02} h_{11} \bar{h}_{02}\right. \\
& -12 h_{11}^{2} \bar{h}_{11}+4 h_{12} \bar{h}_{11}-6 h_{02} h_{20} \bar{h}_{11}-6 h_{02} \bar{h}_{11}^{2}+2 h_{02} \bar{h}_{12}-3 h_{11} h_{20} \bar{h}_{20} \\
& \left.-3 h_{02} \bar{h}_{02} \bar{h}_{20}-6 h_{11} \bar{h}_{11} \bar{h}_{20}\right) \\
& \tilde{h}_{13}(a)=\frac{1}{6}\left(-6 h_{11}^{3}+6 h_{11} h_{12}+h_{03} h_{20}-9 h_{02} h_{11} h_{20}-3 h_{02}^{2} \bar{h}_{02}+3 h_{03} \bar{h}_{11}\right. \\
& -18 h_{02} h_{11} \bar{h}_{11}-6 h_{11}^{2} \bar{h}_{20}+3 h_{12} \bar{h}_{20}-3 h_{02} h_{20} \bar{h}_{20}-12 h_{02} \bar{h}_{11} \bar{h}_{20} \\
& \left.-3 h_{11} \bar{h}_{20}^{2}+h_{11} \bar{h}_{30}\right) \\
& \tilde{h}_{04}(a)=\frac{1}{24}\left(4 h_{03} h_{11}-12 h_{02} h_{11}^{2}+6 h_{02} h_{12}-3 h_{02}^{2} h_{20}-12 h_{02}^{2} \bar{h}_{11}+6 h_{03} \bar{h}_{20}\right. \\
& \left.-18 h_{02} h_{11} \bar{h}_{20}-15 h_{02} \bar{h}_{20}^{2}+4 h_{02} \bar{h}_{30}\right)
\end{aligned}
$$

$$
\begin{aligned}
&\left|r_{2}^{40}(a)\right|=\frac{a^{3} \sqrt{1+a+\frac{4}{9(-5+4 a)}+\frac{54-10 a(3+a)}{9(-9+a(12+(-5+a) a))}}}{4+a\left(-6+a+a^{2}\right)} \\
&\left|r_{2}^{31}(a)\right| \\
&= \frac{a^{3} \sqrt{4-6 a+21 a^{2}-242 a^{3}+741 a^{4}-1035 a^{5}+824 a^{6}-426 a^{7}+148 a^{8}-32 a^{9}+4 a^{10}}}{(-1+a)^{2}(-5+4 a) \sqrt{(-1+a)(1+a)(-4+a(2+a))(-9+a(12+(-5+a) a))}} \\
&\left|r_{2}^{22}(a)\right|= \frac{a^{3} \sqrt{1+2 a+10 a^{2}+6 a^{3}+220 a^{4}-434 a^{5}+222 a^{6}-20 a^{7}+a^{8}}}{\sqrt{-5+4 a}(-1+a)^{3}(1+a)(-4+a(2+a))} \\
&\left|r_{2}^{13}(a)\right|= a^{3}\left(-256+64 a+1676 a^{2}-2498 a^{3}-95 a^{4}+2796 a^{5}\right. \\
&\left.-2219 a^{6}+187 a^{7}+730 a^{8}-550 a^{9}+200 a^{10}-40 a^{11}+4 a^{12}\right)^{\frac{1}{2}} \\
&\left((-1+a)^{\frac{5}{2}}(-5+4 a)(-4+a(2+a)) \sqrt{(1+a)(-9+a(12+(-5+a) a))}\right)^{-1} \\
&\left|r_{2}^{04}(a)\right|= \frac{a^{3} \sqrt{-11+a(-6+a(27+a(-17+4 a)))}}{(-4+a(2+a)) \sqrt{(-1+a)(-5+4 a)(-9+a(12+(-5+a) a))}}
\end{aligned}
$$


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