# Morse decompositions for delay-difference equations 

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#### Abstract

Scalar difference equation $x_{k+1}=f\left(x_{k}, x_{k-d}\right)$ with delay $d \in \mathbb{N}$ are well-motivated from applications e.g. in the life sciences or discretizations of delay-differential equations. We investigate their global dynamics by providing a (nontrivial) Morse decomposition of the global attractor. Under an appropriate feedback condition on the second variable of $f$, our basic tool is an integer-valued Lyapunov functional. Keywords: delay-difference equation; global attractor; Morse decomposition; discrete Lyapunov functional 2010 Mathematics Subject Classification: 37B25; 37C70; 34K28; 39A11; 39A30; 92A15


## 1 Introduction

This paper studies the global dynamics of scalar difference equations

$$
\begin{equation*}
x_{k+1}=f\left(x_{k}, x_{k-d}\right) \tag{1.1}
\end{equation*}
$$

involving an arbitrary delay $d \in \mathbb{N}$, where $f$ fulfills a suitable positive or negative feedback condition in the second variable. Such problems are of interest not only as time-discretizations of delay-differential equations, but they also intrinsically arise in a multitude of models in the life sciences (see [5, 6] for references). Particularly in the latter applications, much work so far concentrated on the problem to provide (sufficient) conditions for the global asymptotic stability of a positive equilibrium of (1.1) (cf. [6, 8, 9]). Nevertheless, being a discrete time model it is no surprise that much more complicated dynamics can be expected.

[^0]Under quite natural and frequently met assumptions the delay-difference equations (1.1) are dissipative. This means their forward dynamics eventually enters a bounded subset of the state space. Whence, a global attractor $\mathcal{A}$ exists, which contains all bounded entire solutions and therefore all dynamically relevant objects, like for instance equilibria, periodic solutions, as well as homo- or heteroclinics. The dynamics on the attractor itself can be highly nontrivial due to e.g. a cascade of period doubling bifurcations and even chaotic dynamics might arise [10]. Beyond the pure existence of a global attractor, one is rather interested in an as detailed as possible picture of its interior structure. An adequate tool for this endeavor is a Morse decomposition of $\mathcal{A}$ due to the following reasons:

- It allows to disassemble an attractor into finitely many invariant, compact subsets (the Morse sets) and their connecting orbits,
- the recurrent dynamics in $\mathcal{A}$ occurs entirely in the Morse sets,
- outside the Morse sets the dynamics of (1.1) on $\mathcal{A}$ is gradient-like.

On the one hand, identifying a nontrivial Morse decomposition of the attractor (along with the connecting orbits) provides a more detailed picture of the long-term behavior of (1.1), since every solution is attracted by exactly one Morse set. On the other hand, obtaining a Morse decomposition is a difficult task and requires further tools. In our case, this is an integervalued (or discrete) Lyapunov functional, which roughly speaking counts the number of sign changes and decreases along solutions. This allows to quantize solutions to (1.1) in terms of their oscillation rates. Such a concept is not new and actually turned out to be very fruitful to understand the global behavior of other finite and infinite dimensional dynamical systems. For instance, discrete Lyapunov functionals are used to obtain convergence to equilibria in tridiagonal ODEs [21] and scalar parabolic equations [15], but also a Poincaré-Bendixson theory for a class of ODEs [11] in $\mathbb{R}^{n}, n>2$, reaction-diffusion equations [3] and delay-differential equations [13]. Finally, a Morse decomposition of global attractors for delay-differential equations is constructed in $[12,16]$. In conclusion, the existence of such a discrete Lyapunov functional imposes a serious constraint on the possible long-term behavior of various systems.

All the above applications have in common to address problems in continuous time, that is differential equations. A first approach to tackle difference equations via discrete Lyapunov functionals is due to Mallet-Paret and Sell [14]. In showing that such an integer-valued functional $V$ decreases along forward solutions, they lay the foundations of our present work. Yet, [14] is primarily motivated by time-discretizations of delay-differential equations, while we are furthermore interested in applications being time-discrete right from the beginning by means of models originating e.g. in life sciences. Note
that [14] prove the decay of $V$ for a larger class of difference equations than (1.1), which combines (1.1) with (cyclic) tridiagonal systems (see also our Remark 3.4). Nonetheless, up to our knowledge, this paper is the first contribution using a discrete Lyapunov functional to actually understand the dynamics of discrete-time models.

Our detailed setting is as follows: The right-hand side $f: J^{2} \rightarrow J$ of (1.1) is assumed to be of class $C^{1}$ with a closed interval $J \subseteq \mathbb{R}$. Since (1.1) is equivalent to the first order difference equation

$$
y_{k+1}=F\left(y_{k}\right), \quad F(y(0), \ldots, y(d)):=\left(\begin{array}{c}
y(1)  \tag{1.2}\\
y(2) \\
\vdots \\
f(y(d), y(0))
\end{array}\right)
$$

the natural state space for (1.1) is the cartesian product $J^{d+1} \subseteq \mathbb{R}^{d+1}$. Indeed, the mapping $F: J^{d+1} \rightarrow J^{d+1}$ generates a discrete semi-dynamical system via its iterates $F^{k}, k \geq 0$. It is one-to-one, if and only if

$$
\begin{equation*}
f(x, \cdot): J \rightarrow \mathbb{R} \quad \text { is injective for all } x \in J \tag{1.3}
\end{equation*}
$$

and $F: J^{d+1} \rightarrow J^{d+1}$ is onto, if every $f(x, \cdot): J \rightarrow J, x \in J$, has this property. An entire solution to (1.2) is a sequence $\left(y_{k}\right)_{k \in \mathbb{Z}}$ in $J^{d+1}$ satisfying the identity $y_{k+1} \equiv F\left(y_{k}\right)$ on $\mathbb{Z}$, and we speak of an entire solution through $\eta \in J^{d+1}$, provided also $y_{0}=\eta$ holds.

Under appropriate conditions on the function $f$, the equation (1.2) is dissipative, i.e. there exists a bounded subset $A \subset J^{d+1}$ (the absorbing set) such that for every bounded $B \subseteq J^{d+1}$ there exists a $K=K(B) \geq 0$ with

$$
F^{k}(B) \subseteq A \quad \text { for all } k \geq K
$$

The global attractor $\mathcal{A} \subseteq \bar{A}$ for (1.2), is a compact, invariant and nonempty set attracting all bounded $B \subseteq J^{d+1}$. In our case, the existence of the global attractor is equivalent to the existence of a bounded absorbing set [19]. The global attractor is unique and allows the dynamical characterization

$$
\mathcal{A}=\left\{\xi \in J^{d+1}: \text { there is a bounded entire solution of (1.2) through } \xi\right\}
$$

(cf. [19, Lemma 2.18]). It is invariantly connected (cf. [19, Proposition 2.20]). In case $J=\mathbb{R}$ the state space $\mathbb{R}^{d+1}$ of (1.2) is a Banach space, $\mathcal{A}$ is connected (cf. [19, Prop 2.19]) and contains a fixed point (cf. [19, Theorem 2.29]).

The bounded entire solutions constituting the global attractor $\mathcal{A}$ are uniquely determined under the injectivity condition (1.3) and thus the continuous bijection $\left.F\right|_{\mathcal{A}}: \mathcal{A} \rightarrow \mathcal{A}$ is a homeomorphism on the compact set $\mathcal{A}$. Hence, the iterates $\left.F\right|_{\mathcal{A}} ^{k}, k \in \mathbb{Z}$, are a dynamical system on $\mathcal{A}$ and for each $\xi \in \mathcal{A}$ the corresponding $\alpha$ - and $\omega$-limit sets are denoted by $\alpha(\xi), \omega(\xi) \subseteq \mathcal{A}$.


Figure 1.1: The mappings $f$ and $g$ from Example 1.1

In the general, non-invertible case, for a bounded entire solution $\left(y_{k}\right)_{k \in \mathbb{Z}}$ of (1.2), we introduce the $\alpha$-limit set of $y$, defined by

$$
\alpha(y)=\left\{\xi \in \mathcal{A}: \text { there exists } k_{n} \rightarrow-\infty, \text { such that } y_{k_{n}} \rightarrow \xi, \text { as } n \rightarrow \infty\right\} .
$$

Coming to our central notion, following [1], a Morse decomposition of the global attractor $\mathcal{A}$ is a finite ordered collection

$$
\mathcal{M}_{1}<\mathcal{M}_{2}<\cdots<\mathcal{M}_{m}
$$

of pairwise disjoint, compact and invariant Morse sets $\mathcal{M}_{1}, \ldots, \mathcal{M}_{m} \subseteq \mathcal{A}$ such that for all $\xi \in \mathcal{A}$ and all bounded entire solution $\left(y_{k}\right)_{k \in \mathbb{Z}}$ through $\xi$ it holds that there exist $i \geq j$ with

$$
\begin{aligned}
& \left(m_{1}\right) \alpha(y) \subseteq \mathcal{M}_{i} \text { and } \omega(\xi) \subseteq \mathcal{M}_{j} \\
& \left(m_{2}\right) \quad i=j \Rightarrow \xi \in \mathcal{M}_{i}\left(\text { thus, } y_{k} \in \mathcal{M}_{i} \text { for every } k \in \mathbb{Z}\right)
\end{aligned}
$$

The connecting orbits constitute the connecting sets

$$
\mathcal{C}_{j}^{i}:=\left\{\xi \in \mathcal{A}: \alpha(y) \subseteq \mathcal{M}_{i} \text { and } \omega(\xi) \subseteq \mathcal{M}_{j}\right\} \quad \text { for all } j<i
$$

Together with the Morse sets $\mathcal{M}_{1}, \ldots, \mathcal{M}_{m}$ they form a partition of the global attractor $\mathcal{A}$. Moreover, since $\mathcal{A}$ includes every limit point of an equation, due to $\left(m_{1}\right)$ indeed the Morse sets contain all limit sets.

A simple illustration is given in the following example:
Example 1.1. Consider the scalar difference equation (without delay)

$$
x_{k+1}=f\left(x_{k}\right), \quad f(x):=\frac{2 x}{\sqrt{1+3 x^{2}}}
$$

having the equilibria $-1,0,1$ and the global attractor $[-1,1]$. Using graphical iteration (see Fig. 1.1 (left)) or from the explicit representation $f^{k}(x)=$ $\frac{2^{k} x}{\sqrt{1+4^{k}\left(x^{2}-1\right)}}$ one readily obtains the Morse sets and connecting sets

$$
\mathcal{M}_{1}=\{-1,1\}, \quad \mathcal{M}_{2}=\{0\}, \quad \mathcal{C}_{1}^{2}=(-1,0) \cup(0,1)
$$

The components of $\mathcal{M}_{1}$ and $\mathcal{C}_{1}^{2}$ are invariant; $\mathcal{M}_{1}$ is not minimal. The slightly modified difference equation (see Fig. 1.1 (right))

$$
x_{k+1}=g\left(x_{k}\right), \quad g(x):=\frac{-2 x}{\sqrt{1+3 x^{2}}}
$$

has the same global attractor, Morse and connecting sets. Now the Morse set $\mathcal{M}_{1}$ is a 2-periodic orbit, while $(-1,0),(0,1)$ fail to be invariant.

The paper is organized as follows: In many cases the discrete Lyapunov functional is applied to a variational equation. This allows to argue that it decreases along the difference of two solutions. Hence, the next two sections include both preparatory work and crucial related results on linear difference equations; for instance Lemma 2.1 provides a description of the eigenvalues distribution for the linearized equation. In Section 3 we formulate our main assumptions, the feedback conditions and introduce an integer-valued Lyapunov functional, whose properties are given by Proposition 3.2 and Theorem 3.3. Our main result is Theorem 4.1, which yields a Morse decomposition of the global attractor of (1.2) (and in turn (1.1)) under appropriate assumptions on the right-hand side $f$. Most of the arguments are given in the negative feedback case, which seems to have more applications and requires slightly more involved proofs. For the convenience of the reader, the proofs for positive feedback are given in the Appendix. Section 5 contains applications of the theory, mainly from the life sciences, which underline that both the positive and negative feedback case are relevant. Finally some open questions and perspectives are addressed in Section 6.

We conclude with our notation: A discrete interval $\mathbb{I}$ is the intersection of a real interval with the integers and $\mathbb{I}^{\prime}:=\{k \in \mathbb{I}: k+1 \in \mathbb{I}\}$. Special cases are the positive half-axis $\mathbb{Z}_{+}:=\{k \in \mathbb{Z}: k \geq 0\}$ and the negative half-axis $\mathbb{Z}_{-}:=\{k \in \mathbb{Z}: k \leq 0\}$.

## 2 Spectrum of the linearization

Let us consider an autonomous linear delay-difference equation

$$
\begin{equation*}
x_{k+1}=a x_{k}+b x_{k-d} \tag{2.1}
\end{equation*}
$$

with real coefficients $a>0, b \neq 0$. Our further analysis requires a detailed understanding of the associated characteristic equation

$$
\begin{equation*}
\lambda^{d+1}-a \lambda^{d}-b=0 \tag{2.2}
\end{equation*}
$$

whose solutions occur as complex-conjugated pairs $\lambda_{j}, \overline{\lambda_{j}}$, where for real roots, clearly $\lambda_{j}=\overline{\lambda_{j}}$. By convention, suppose that $\operatorname{Im} \lambda_{j} \geq 0$.

The subsequent lemma thoroughly describes the distribution of the roots of (2.2) based on the following sectors of the complex plane:

$$
S_{j}^{+}:=\left\{z \in \mathbb{C}: \frac{(2 j-1) \pi}{d}<\arg z<\frac{2 j \pi}{d+1}\right\} \quad \text { for } j \in \mathbb{N}, j<\frac{d+1}{2}
$$

and

$$
S_{j}^{-}:=\left\{z \in \mathbb{C}: \frac{2 j \pi}{d}<\arg z<\frac{(2 j+1) \pi}{d+1}\right\} \quad \text { for } j \in \mathbb{N}_{0}, j<\frac{d}{2}
$$

Lemma 2.1 (Huszár [4]). Let $m_{0}:=-d^{d}\left(\frac{a}{d+1}\right)^{d+1}$.
(a) If $b>0$, then for all $j \in \mathbb{N}, j<\frac{d+1}{2}$, there exists exactly one pair of characteristic roots $\lambda_{j}, \overline{\lambda_{j}}$ of (2.2) with $\lambda_{j} \in S_{j}^{+}$, and all other roots are real. For even $d$ there exists a unique real root, namely $\lambda_{+}>0$ and for odd $d$ there exist exactly two real roots, namely $\lambda_{-}<0<\lambda_{+}$.
(b) If $b<0$, then for all $j \in \mathbb{N}, j<\frac{d}{2}$, there exists exactly one pair of characteristic roots $\lambda_{j}, \overline{\lambda_{j}}$ of (2.2), such that $\lambda_{j} \in S_{j}^{-}$.

1) If $m_{0} \leq b<0$, then all other roots are real. Exactly two of those are positive and they are denoted by $0<\lambda_{+, 2} \leq \lambda_{+, 1}$.
2) If $b<m_{0}$, then there is also a pair of complex roots $\lambda_{0}, \overline{\lambda_{0}}$, such that $\lambda_{0} \in S_{0}^{-}$, and then there exists no positive characteristic root.
3) If $d$ is odd, then there is no negative real root, otherwise there is exactly one, which is denoted by $\lambda_{-}$.
(c) If $b \neq 0$, and $z_{1}=r_{1} e^{i \varphi_{1}}$ and $z_{2}=r_{2} e^{i \varphi_{2}}$ are two distinct roots of (2.2) such that $0 \leq \varphi_{1}<\varphi_{2} \leq \pi$ holds, then $r_{1}>r_{2}$.

Note that the above lemma also implies that the only possible multiple (at most double) root which possibly occurs, is the pair of two positive roots $\lambda_{+, 1}$ and $\lambda_{+, 2}$ which may coincide. All the other roots are simple. We refer to Figures 2.1 and 2.2, in which the number and sign of the real roots of (2.2) is presented for all possible cases.

Proof. By means of the variable change $\lambda=a x$ one transforms the characteristic equation (2.2) into the trinomial equation

$$
x^{d+1}-x^{d}+p=0, \quad p:=-\frac{b}{a^{d+1}}
$$

Then one applies [4, 1. §] to obtain the parts of statements (a) and (b) concerning the distribution of complex roots. Moreover, according to 8. § of the same paper there is no multiple complex root. This determines also the number of real roots (counting multiplicity) in each case. Then applying Descartes' sign rule yields immediately our conclusion about real roots. From $[4,2 . \S]$ we readily have statement (c).


Figure 2.1: Roots of (2.2) (dots) and sectors (grey) from Huszár's Lemma 2.1 for even delays $d$

$m_{0} \leq b<0$

$\lambda_{-}<0<\lambda_{+}$
$0<\lambda_{2,+}<\lambda_{1,+}$
Figure 2.2: Roots of (2.2) (dots) and sectors (grey) from Huszár's Lemma 2.1 for odd delays $d$

## 3 A Lyapunov functional

Our overall analysis is based on the following assumptions on the right-hand side $f: J \times J \rightarrow J$ of (1.1):
$\left(H_{1}\right) f$ is continuously differentiable and 0 is an inner point of $J ;$
$\left(H_{2}\right) D_{1} f\left(x_{d}, x_{0}\right)>0$ for all $x_{d}, x_{0} \in J ;$
moreover there exists a $\delta^{*} \in\{-1,1\}$, fixed from now on, such that:
$\left(H_{3}\right) \delta^{*} D_{2} f(0,0)>0, \delta^{*} f\left(0, x_{0}\right) x_{0}>0$ for all $x_{0} \in J \backslash\{0\}$.
Note that for $\delta^{*}=1$ one speaks of positive feedback, while $\delta^{*}=-1$ means negative feedback in (1.1).

It is central to have an ambient integer-valued Lyapunov functional for a variational equation along solutions to (1.1) at hand. For sequences $\left(a_{k}\right)_{k \in \mathbb{Z}}$, $\left(b_{k}\right)_{k \in \mathbb{Z}}$ in $\mathbb{R}$ we consider the nonautonomous linear delay-difference equation

$$
\begin{equation*}
x_{k+1}=a_{k} x_{k}+b_{k} x_{k-d} \tag{3.1}
\end{equation*}
$$

as well as the associated first order system

$$
y_{k+1}=A_{k} y_{k}, \quad A_{k}:=\left(\begin{array}{ccccc}
0 & 1 & 0 & \cdots & 0  \tag{3.2}\\
0 & 0 & 1 & \cdots & 0 \\
\vdots & \vdots & \ddots & \ddots & \vdots \\
0 & 0 & \cdots & 0 & 1 \\
b_{k} & 0 & \cdots & 0 & a_{k}
\end{array}\right) \text {, }
$$

or equivalently,

$$
\left(\begin{array}{c}
y_{k+1}(0)  \tag{3.3}\\
y_{k+1}(1) \\
\vdots \\
y_{k+1}(d-1) \\
y_{k+1}(d)
\end{array}\right)=\left(\begin{array}{c}
y_{k}(1) \\
y_{k}(2) \\
\vdots \\
y_{k}(d) \\
a_{k} y_{k}(d)+b_{k} y_{k}(0)
\end{array}\right) .
$$

Remark 3.1. The assumptions $\left(H_{1}\right)-\left(H_{3}\right)$ allow to write every nonlinear difference equation (1.1) in the form (3.1) with

$$
a_{k}=\int_{0}^{1} D_{1} f\left(h x_{k}, x_{k-d}\right) \mathrm{d} h, \quad b_{k}=\int_{0}^{1} D_{2} f\left(0, h x_{k-d}\right) \mathrm{d} h .
$$

Moreover, for an entire solution $\left(x_{k}\right)_{k \in \mathbb{Z}}$ of (1.1), $a_{k}>0$ holds for all $k \in \mathbb{Z}$ (for which $a_{k}$ is defined), and as $b_{k}=f\left(0, x_{k-d}\right) / x_{k-d}$ for any $x_{k-d} \neq 0$ and $b_{k}=D_{2} f(0,0)$ for $x_{k-d}=0$, one has $\delta^{*} b_{k}>0$ for all $k \in \mathbb{Z}$.

Analogously, every equation (1.2) can be rewritten in the form (3.2) with

$$
\begin{equation*}
a_{k}=\int_{0}^{1} D_{1} f\left(h y_{k}(d), y_{k}(0)\right) \mathrm{d} h, \quad b_{k}=\int_{0}^{1} D_{2} f\left(0, h y_{k}(0)\right) \mathrm{d} h \tag{3.4}
\end{equation*}
$$

where $a_{k}>0, \delta^{*} b_{k}>0$ hold for all $k \in \mathbb{Z}$ from the domain of the solution $y$.
Before proceeding, we abbreviate $\mathbb{R}_{*}^{n}:=\mathbb{R}^{n} \backslash\{0\}$ for $n \in \mathbb{N}$. It is crucial to note that $\mathbb{R}_{*}^{d+1}$ is positively invariant under the matrix (3.2) provided $a_{k}>0$ and $\delta^{*} b_{k}>0$ hold on $\mathbb{Z}$. Let us also define the function

$$
\text { sc: } \bigcup_{n \in \mathbb{N}_{0}} \mathbb{R}_{*}^{n+1} \rightarrow \mathbb{N}_{0},
$$

where sc $y$ is the number of sign changes in $y=(y(0), y(1), \ldots, y(n)) \in \mathbb{R}_{*}^{n+1}$, that is

$$
\text { sc } y:=\max \left\{\ell \in \mathbb{N} \mid \exists 0 \leq i_{0}<\cdots<i_{\ell} \leq n: y\left(i_{j-1}\right) y\left(i_{j}\right)<0 \forall 1 \leq j \leq \ell\right\} .
$$

Also, by definition, set sc $y:=0$, if there is no such $\ell$. Based on this, we define the integer-valued functions

$$
V^{+}: \mathbb{R}_{*}^{d+1} \rightarrow 2 \mathbb{N}_{0}, \quad V^{+}(y):= \begin{cases}\text { sc } y, & \text { if sc } y \text { is even }, \\ \text { sc } y+1, & \text { if sc } y \text { is odd },\end{cases}
$$

and

$$
V^{-}: \mathbb{R}_{*}^{d+1} \rightarrow 2 \mathbb{N}_{0}+1, \quad V^{-}(y):= \begin{cases}\text { sc } y, & \text { if sc } y \text { is odd } \\ \text { sc } y+1, & \text { if sc } y \text { is even. }\end{cases}
$$

Finally we introduce the sets $\mathcal{R}^{+}$resp. $\mathcal{R}^{-}$of regular points via

$$
\mathcal{R}^{ \pm}:=\left\{\begin{array}{l|l}
y \in \mathbb{R}^{d+1} \left\lvert\, \begin{array}{l}
y(i)=0 \Rightarrow y(i-1) y(i+1)<0 \text { for } 0 \leq i \leq d \\
\text { with } y(d+1):=\delta^{*} y(0) \text { and } y(-1):=\delta^{*} y(d)
\end{array}\right.
\end{array}\right\}
$$

for the positive (resp. negative) feedback case.
The subsequent Proposition 3.2 and Theorem 3.3 give some nice properties of the functionals $V^{ \pm}$that will be essential in the proof of our main result. We state them for both positive and negative feedback, however, we only present the proofs for the negative feedback case $\delta^{*}=-1$. For positive feedback, the proofs are rather similar to the ones presented below, and require only straightforward modifications in the arguments (cf. Appendix).
Proposition 3.2. If $\left(y^{n}\right)_{n \in \mathbb{N}}$ is a sequence in $\mathbb{R}_{*}^{d+1}$ with limit $y \in \mathbb{R}_{*}^{d+1}$, then the following statements hold:
(a) $V^{ \pm}(y) \leq \liminf _{n \rightarrow \infty} V^{ \pm}\left(y^{n}\right)$,
(b) $V^{ \pm}(y)=\lim _{n \rightarrow \infty} V^{ \pm}\left(y^{n}\right)$, if $y \in \mathcal{R}^{ \pm}$.

Proof. (a) follows from the lower semi-continuity of the function sc on $\mathbb{R}_{*}^{d+1}$.
(b) is shown only in the negative feedback case, since the positive feedback case can be handled by straightforward modifications in the argument.

First note that there exists $n_{0} \in \mathbb{N}$ such that for $n>n_{0}, y^{n}(j)$ has the same sign as $y(j)$, where $j \in\{0, \ldots, d\}$ is an index for which $y(j) \neq 0$. Moreover, from $y \in \mathcal{R}^{-}$it follows that for $j \in\{1, \ldots, d-1\}$ with $y(j)=0$,

$$
\operatorname{sgn} y(j-1)=\operatorname{sgn} y^{n}(j-1)=-\operatorname{sgn} y(j+1)=-\operatorname{sgn} y^{n}(j+1) \neq 0
$$

and

$$
\operatorname{sc}(y(j-1), y(j), y(j+1))=\operatorname{sc}\left(y^{n}(j-1), y^{n}(j), y^{n}(j+1)\right)=1
$$

holds for all $n>n_{0}$. This in turn yields that for $n>n_{0}$ and $0 \leq i<j \leq d$ such that $y(i) \neq 0 \neq y(j)$ one has

$$
\begin{equation*}
\operatorname{sc}(y(i), y(i+1), \ldots, y(j))=\operatorname{sc}\left(y^{n}(i), y^{n}(i+1), \ldots, y^{n}(j)\right), \tag{3.5}
\end{equation*}
$$

which in particular means that sc $y^{n}=\operatorname{sc} y$ and thus also $V^{-}(y)=V^{-}\left(y^{n}\right)$ hold whenever $y(0) \neq 0 \neq y(d)$.

Since $y \in \mathcal{R}^{-}$excludes the possibility of $y(0)$ and $y(d)$ both being 0 at once, there remains the case when exactly one of them is 0 . We may assume $y(d)=0$ (the case $y(0)=0$ is analogous). Again, by regularity of $y$ we have that $y(0) y(d-1)>0$, which together with (3.5) yields that on the one hand

$$
\operatorname{sc}(y(0), y(1), \ldots y(d-1))=\operatorname{sc}\left(y^{n}(0), y^{n}(1), \ldots y^{n}(d-1)\right)
$$

holds for $n>n_{0}$ and on the other hand this number is even, say $2 m$. It is now clear that sc $y=2 m \leq \operatorname{sc} y^{n} \leq 2 m+1$ and therefore $V^{-}(y)=V^{-}\left(y^{n}\right)$ is satisfied for all $n>n_{0}$, which completes the proof.

Theorem 3.3. If $\left(y_{k}\right)_{k \in \mathbb{I}}$ is a solution of (3.2) with $a_{k}>0$ and $\delta^{*} b_{k}>0$ for all $k \in \mathbb{I}$, then
(a) $V^{ \pm}\left(y_{k+1}\right) \leq V^{ \pm}\left(y_{k}\right)$ holds for any $k \in \mathbb{I}^{\prime}$,
(b) if $k \in \mathbb{I}$ is such that $k+4 d+2 \in \mathbb{I}$ and moreover $y_{k} \neq 0$ and $V^{ \pm}\left(y_{k}\right)=$ $V^{ \pm}\left(y_{k+4 d+2}\right)$ hold, then $y_{k+4 d+2} \in \mathcal{R}^{ \pm}$.

According to Remark 3.1 this applies to solutions of (1.2) as well. We additionally remark that Theorem 3.3 remains valid in the nonautonomous situation when $f$ depends on $k$ and our assumptions hold for every $k \in \mathbb{I}$.

The next proof is valid for any $d \in \mathbb{N}$ (using notations $y(d+1):=\delta^{*} y(0)$ and $y(-1):=\delta^{*} y(d)$ where necessary in the $d=1$ case). However it is worth mentioning that Theorem 3.3 almost trivially holds for $d=1$, and some parts of the following proof can be omitted even in the case $d=2$.

Proof of Theorem 3.3. For better readability, we give the proof for negative feedback $\delta^{*}=-1$, the positive feedback case can be found in the Appendix.
(a) Since the last $d$ components of $y_{k} \in \mathbb{R}_{*}^{d+1}$ are the same as the first $d$ coordinates of $y_{k+1}$, clearly sc $y_{k+1} \leq \mathrm{sc} y_{k}+1$ holds. Assuming now to the contrary that $V^{-}\left(y_{k+1}\right)>V^{-}\left(y_{k}\right)$ holds for some $k$ implies that sc $y_{k+1}=\operatorname{sc} y_{k}+1$ and sc $y_{k}$ is odd. The former combined with $a_{k}>0$ yields that $y_{k}(0) \neq 0$ and the oddness of sc $y_{k}$ implies that the last nonzero coordinate of $y_{k}$ (say $\left.y_{k}(j), 1 \leq j \leq d\right)$ has the opposite sign as $y_{k}(0)$. We may assume w.l.o.g. that $y_{k}(j)<0<y_{k}(0)$ (the other case is analogous). In particular this implies $y_{k}(d) \leq 0$. Using $a_{k}>0$ and $b_{k}<0$ one obtains readily from (3.3) that $y_{k+1}(d)=a_{k} y_{k}(d)+b_{k} y_{k}(0)$ is negative. All these together give

$$
\begin{aligned}
\operatorname{sc} y_{k+1} & =\operatorname{sc}\left(y_{k+1}(0), y_{k+1}(1), \ldots, y_{k+1}(d-1), y_{k+1}(d)\right) \\
& =\operatorname{sc}\left(y_{k}(1), y_{k}(2), \ldots, y_{k}(d), y_{k+1}(d)\right) \leq \operatorname{sc} y_{k}
\end{aligned}
$$

where the last inequality is due to $\operatorname{sgn} y_{k}(j)=\operatorname{sgn} y_{k+1}(d)$. This is a contradiction, which proves statement (a).
(b) Before we prove the statement, we need to introduce some auxiliary notions. Let us say that a component $y(j)(0 \leq j \leq d)$ of a vector $y \in \mathbb{R}^{d+1}$ is irregular if $y(j)=0$, and $y(j-1) y(j+1) \geq 0$, where $y(-1):=-y(d)$ and $y(d+1):=-y(0)$. Using this terminology, $y \in \mathcal{R}^{-}$(i.e. $y$ is regular) holds if and only if $y$ has no irregular component. Furthermore, we call a vector $(y(i), \ldots, y(j))$ with $0 \leq i \leq j \leq d$ an irregular block (of zeros) in $y \in \mathbb{R}^{d+1}$ if $y(i)=\cdots=y(j)=0$, and moreover it has maximal length in the sense that either $i=0$ or $i \geq 1$ and $y(i-1) \neq 0$ hold, and similarly either $j=d$ or
$j<d$ and $y(j+1) \neq 0$ hold. The dimension of the block will be regarded as the length of it. Note also that consecutive zero components are irregular by definition. The proof of (b) consists of several, yet elementary, steps. From now on we always assume that $V^{-}\left(y_{k}\right)=V^{-}\left(y_{k+4 d+2}\right)$, which implies, in the light of statement (a), that $V^{-}\left(y_{k}\right)=V^{-}\left(y_{k+\ell}\right)$ holds for all $0 \leq \ell \leq 4 d+2$.
Step 1. If for some index $k, y_{k+1}(j)=0$ is irregular $(0 \leq j<d)$, then so is $y_{k}(j+1)=0$ irregular.

This is trivial in case $1 \leq j \leq d-2$.
If $y_{k+1}(d-1)=0$ is irregular (i.e. $y_{k+1}(d-1)=0$ and $y_{k+1}(d-$ 2) $y_{k+1}(d) \geq 0$ hold), then on the one hand $y_{k}(d)=0$ holds and on the other hand, $b_{k}<0$ implies $\operatorname{sgn} y_{k+1}(d)=-\operatorname{sgn} y_{k}(0)$. Combining this with $y_{k}(d-1)=y_{k+1}(d-2)$ yields $\operatorname{sgn}\left(y_{k}(d-1) y_{k}(d+1)\right)=-\operatorname{sgn}\left(y_{k}(d-\right.$ 1) $\left.y_{k}(0)\right)=\operatorname{sgn}\left(y_{k+1}(d-2) y_{k+1}(d)\right)$, meaning $y_{k}(d)$ is also irregular in $y_{k}$.

For the case when $y_{k+1}(0)=0$ is irregular, assume to the contrary that $y_{k}(1)=0$ is regular, i.e. $y_{k}(0) y_{k}(2)<0$. W.l.o.g. we may assume $y_{k}(0)<$ $0<y_{k}(2)$. Irregularity of $y_{k+1}(0)$ combined with $y_{k+1}(1)=y_{k}(2)>0$ yields $y_{k+1}(d) \leq 0$. Now, from $a_{k}>0$ and $b_{k}<0$ it follows that $y_{k}(d)<0$. These all together imply that sc $y_{k}$ is an even number, moreover sc $y_{k}=\operatorname{sc} y_{k+1}+1$, which gives $V^{-}\left(y_{k}\right)=V^{-}\left(y_{k+1}\right)+2$, a contradiction.

Step 2. If $y_{k+1}(d)=0$, then there are two possibilities. Either $y_{k}(0)=$ $y_{k}(d)=0$, or $y_{k}(0) y_{k}(d)>0$. Assume now that the latter holds and let $0 \leq j<d$ be the smallest integer such that $y_{k}(j+1) \neq 0$. As $y_{k}(0) y_{k}(d)>0$, hence sc $y_{k}$ is even, so $V^{-}\left(y_{k}\right)=V^{-}\left(y_{k+1}\right)$ holds only if $\operatorname{sgn} y_{k}(j+1)=$ $\operatorname{sgn} y_{k}(0)$, which in turn yields also

$$
\begin{equation*}
\operatorname{sgn} y_{k}(j+1)=\operatorname{sgn} y_{k}(d)=\operatorname{sgn} y_{k+1}(d-1) \neq 0 \tag{3.6}
\end{equation*}
$$

If in addition $y_{k+1}(d)=0$ is irregular, then $0 \geq y_{k+1}(0) y_{k+1}(d-1)=$ $y_{k}(1) y_{k}(d)$ holds, which in the light of (3.6) yields that $j \neq 0$, or what is equivalent, $y_{k}(1)=0$.

To sum up, if $y_{k+1}(d)=0$ is irregular, then either $y_{k}(0)=y_{k}(d)=0$ or there exists some $1 \leq j \leq d-1$, such that

$$
\begin{align*}
y_{k}(0) & \neq 0 \\
y_{k}(1) & =\cdots=y_{k}(j)=0 \\
y_{k}(j+1) & \neq 0  \tag{3.7}\\
y_{k}(d) & \neq 0 \\
\operatorname{sgn} y_{k}(0) & =\operatorname{sgn} y_{k}(j+1)=\operatorname{sgn} y_{k}(d)
\end{align*}
$$

hold, which in particular means that $y_{k}(0)$ is followed by an irregular block of zeros of length $j \geq 1$.

Step 3. As a consequence of the first two steps, if $y_{k} \in \mathcal{R}^{-}$and $V^{-}\left(y_{k}\right)=$ $V^{-}\left(y_{k+1}\right)$, then $y_{k+1} \in \mathcal{R}^{-}$. Therefore, in order to prove the statement, it is sufficient to show that there exists $0 \leq \ell \leq 4 d+2$ such that $y_{k+\ell} \in \mathcal{R}^{-}$.

Step 4. Note that if $y_{k}(d)=0$ and $y_{k} \neq 0 \in \mathbb{R}^{d+1}$, then since $a_{j}>0$ for all $j \geq k$, there exists $k+1 \leq k_{1} \leq k+d$ such that $y_{k_{1}}(d) \neq 0$. Thus we may assume that $k_{1}, k \leq k_{1} \leq k+d$ is such that $y_{k_{1}}(d) \neq 0$.
Step 5. This in particular implies that $y_{k_{1}}(d)$ is regular. The rest of the vector $y_{k_{1}}$ may contain several irregular blocks, which are separated by at least one regular coordinate from each other. The coordinate $y_{k_{1}}(0)$ may also be zero, or even irregular.

Next we will study how the irregular blocks of $y_{k_{1}+d+1}$ can be described by the irregular blocks of $y_{k_{1}}$. Let us use notation $k_{2}=k_{1}+d+1$.
Step 5.1. If there exists $1 \leq j \leq d$ so that $y_{k_{1}}(0)=\cdots=y_{k_{1}}(j-1)=$ 0 and $y_{k_{1}}(j) \neq 0$, then, as $a_{j}>0$ for all $j \geq k, \operatorname{sgn} y_{k_{1}+j}(d)=\cdots=$ $\operatorname{sgn} y_{k_{1}+j}(d-j)=\operatorname{sgn} y_{k_{1}}(d)$ holds, and consequently one has $\operatorname{sgn} y_{k_{2}}(0)=$ $\cdots=\operatorname{sgn} y_{k_{2}}(j-1)=\operatorname{sgn} y_{k_{1}}(d) \neq 0$.
Step 5.2. If $y_{k_{1}}(i)=\cdots=y_{k_{1}}(j)=0$ is an irregular block with $1 \leq i<$ $j<d$ such that $y_{k_{1}}(i-1) y_{k_{1}}(j+1)<0$, then for $k^{\prime}=k_{1}+i-1$ one has that $y_{k^{\prime}}(1)=\cdots=y_{k^{\prime}}(j-i+1)=0$ and $y_{k^{\prime}}(0) y_{k^{\prime}}(j-i+2)<0$ hold. Due to the argument seen in Step 2, and that neither $y_{k^{\prime}}(d)=y_{k^{\prime}}(0)=0$, nor (3.7) holds (with $k=k^{\prime}$ ), we obtain that $y_{k^{\prime}+1}(d)$ is regular, which implies in this case that it is also nonzero. Going a bit further, to $k^{\prime \prime}=k^{\prime}+j-i+2=k_{1}+j+1$, one infers that $\operatorname{sgn} y_{k^{\prime \prime}}(d)=\cdots=\operatorname{sgn} y_{k^{\prime \prime}}(d-j+i-1) \neq 0$. Then it is straightforward that $\operatorname{sgn} y_{k_{2}}(i)=\cdots=\operatorname{sgn} y_{k_{2}}(j) \neq 0$ holds. Observe that in case $i=j, y_{k_{1}}(i)=0$ would not be irregular due to assumption $y_{k_{1}}(i-1) y_{k_{1}}(i+1)<0$.

Step 5.3. Thus it remains to consider the case when $y_{k_{1}}(i)=\cdots=y_{k_{1}}(j)=$ 0 is an irregular block with $1 \leq i \leq j<d$ such that $y_{k_{1}}(i-1) y_{k_{1}}(j+1)>0$. Choosing $k^{\prime}=k_{1}+i-1$ just as in the previous step one obtains $y_{k^{\prime}}(1)=$ $\cdots=y_{k^{\prime}}(j-i+1)=0$ and $y_{k^{\prime}}(0) y_{k^{\prime}}(j-i+2)>0$.

Now there are two possibilities. Either $y_{k^{\prime}+1}(d) \neq 0$ and then we have the same situation as in Step 5.2 , i.e. $\operatorname{sgn} y_{k_{2}}(i)=\cdots=\operatorname{sgn} y_{k_{2}}(j) \neq 0$, or $y_{k^{\prime}+1}(d)=0$. In the latter case one can easily see that $\operatorname{sgn} y_{k_{2}}(i-1)=\cdots=$ $\operatorname{sgn} y_{k_{2}}(j)=0$, moreover, if $i \geq 2$, then $y_{k_{2}}(i-2) y_{k_{2}}(j+1)<0$ holds. Note that the index of the last components of the corresponding irregular blocks in $y_{k_{1}}$ and $y_{k_{2}}$ coincide (it is $j$ ).
Step 6. Conversely, Steps 1, 2 and 5 combined show that if $y_{k_{2}}$ has an irregular block, whose last component is at index $j$ (with $0 \leq j<d$ ), i.e. $y_{k_{2}}(j)=0$ is irregular and $y_{k_{2}}(j+1) \neq 0$, then $y_{k_{1}}(j)=0$ is irregular and $y_{k_{1}}(j+1) \neq 0$.

Our aim is to apply the arguments presented in Step 5 now for $y_{k_{3}}$, where $k_{3}$ is to be defined soon. For this reason we need to distinguish two cases with respect to the regularity of $y_{k_{2}}(d)$.

If it is regular, then it may be zero or nonzero. If it is nonzero, then Step 5 shows that all irregular blocks of $y_{k_{2}}$ are of the type considered in

Steps 5.1 and 5.2 (now with $k_{2}$ instead of $k_{1}$ ). If $y_{k_{2}}(d)=0$, then from its regularity it follows that $y_{k_{2}}(0) \neq 0$, so all irregular blocks of zeros are of the type studied in Step 5.2. Note that in that step we did not use that $y_{k_{1}}(d) \neq 0$. For regular $y_{k_{2}}(d)$ let us define $k_{3}=k_{2}$.

If $y_{k_{2}}(d)=0$ is irregular, then $y_{k_{2}-1}(0)=y_{k_{1}}(d) \neq 0$ implies that (3.7) holds with $k=k_{2}-1$ and some $1 \leq j \leq d-1$. Then for $k_{3}=k_{2}+j+1 \leq$ $k_{2}+d$ one infers that $y_{k_{3}}(d) \neq 0$, and all irregular blocks are of the types handled in Steps 5.1 and 5.2.

Applying now the arguments of Step 5.1 and 5.2 for $y_{k_{3}}$ instead of $y_{k_{1}}$ we obtain that, for $k_{4}=k_{3}+d+1, y_{k_{4}}(i)$ is regular for all $0 \leq i<d$. We claim that the last coordinate $y_{k_{4}}(d)$ is also regular. Arguing by way of contradiction, if $y_{k_{4}}(d)$ is irregular, then according to Step 2 there are two possibilities: either $y_{k_{4}-1}(0)=y_{k_{4}-1}(d)=0$ or (3.7) holds for $k=k_{4}-1$. In the former case $y_{k_{4}}(d-1)=y_{k_{4}}(d)=0$ are both irregular, which is a contradiction. If the latter holds, then $y_{k_{4}-1}(1)=0$ is irregular and $y_{k_{4}-1}(0) \neq 0$, thus from Step 1 we obtain that $y_{k_{3}+1}(d)=y_{k_{4}-1-(d-1)}(d)=0$ is also irregular and $y_{k_{3}+1}(d-1)=y_{k_{4}-1}(0) \neq 0$. Thus, (3.7) must hold for $k=k_{3}$, as well, which contradicts the fact that every irregular block of $y_{k_{3}}$ is of the types seen in Step 5.1 and 5.2. This proves that $y_{k_{4}}(d)$ is regular.

Thus, we proved that $y_{k_{4}} \in \mathcal{R}^{-}$is regular for some $0 \leq k_{4} \leq k+4 d+2$. From Step 3 we obtain that $y_{k+4 d+2} \in \mathcal{R}^{-}$.

Assertion (a) of Theorem 3.3 is also true, if we extend $V^{ \pm}$into the origin as $V^{+}(0)=0$ resp. $V^{-}(0)=1$. In this case we may also allow $a_{k} \geq 0$ and $\delta^{*} b_{k} \geq 0$ instead of strict inequalities. The proof of this slightly modified statement is essentially identical to the one presented above.

Remark 3.4. With the modified assumptions, Theorem 3.3 (a) is a special case of a general result by Mallet-Paret and Sell [14, Theorem 2.1]. Yet, there are two reasons for presenting an independent proof: First, our simpler situation allows a much shorter argument. Second, in the more general setting of [14] extra (smallness) conditions on the sequences $a_{k}$ and $b_{k}$ are necessary. As their motivation was to give the Lyapunov property for timediscretization of delay-differential equations and tridiagonal ODEs, those extra conditions were always met for sufficiently small step-sizes. However, our motivation comes not only from time-discretization of continuous-time equations, therefore it is important that we can omit such extra conditions.

## 4 The Morse decomposition

From now on, we suppose that the conditions $\left(H_{1}\right)-\left(H_{3}\right)$ are fulfilled. Furthermore, (1.2) is assumed to possess a global attractor $\mathcal{A}$.

If $a:=D_{1} f(0,0)$ and $b:=D_{2} f(0,0)$, then $\left(H_{1}\right)-\left(H_{3}\right)$ imply $a, \delta^{*} b>0$.

Throughout this section we will treat the positive and negative feedback cases in a parallel way. Both statements and their proofs are formulated in this manner. The only exception in this regard is the proof of Proposition 4.2, whose positive feedback case is postponed to the Appendix.

Before stating our main theorem we introduce some notations concerning the linearization of equation (1.2), which becomes

$$
y_{k+1}=A y_{k}, \quad A:=\left(\begin{array}{ccccc}
0 & 1 & 0 & \cdots & 0  \tag{4.1}\\
0 & 0 & 1 & \cdots & 0 \\
\vdots & \vdots & \ddots & \ddots & \vdots \\
0 & 0 & \cdots & 0 & 1 \\
b & 0 & \cdots & 0 & a
\end{array}\right)
$$

or what is equivalent,

$$
\left(\begin{array}{c}
y_{k+1}(0)  \tag{4.2}\\
y_{k+1}(1) \\
\vdots \\
y_{k+1}(d-1) \\
y_{k+1}(d)
\end{array}\right)=\left(\begin{array}{c}
y_{k}(1) \\
y_{k}(2) \\
\vdots \\
y_{k}(d) \\
a y_{k}(d)+b y_{k}(0)
\end{array}\right)
$$

Denote by $M^{*}$ the number of eigenvalues (counting multiplicities) of $A$ with absolute values strictly greater than 1 . Moreover, let

$$
N_{+}^{*}:= \begin{cases}M^{*}, & \text { if } 0 \text { is hyperbolic or } M^{*}=0 \\ M^{*}+1, & \text { otherwise }\end{cases}
$$

and

$$
N_{-}^{*}:= \begin{cases}M^{*}, & \text { if } 0 \text { is hyperbolic or } M^{*}=1 \\ M^{*}+1, & \text { otherwise }\end{cases}
$$

Our main result is the next theorem:
Theorem 4.1 (Morse decomposition). If (1.2) has a global attractor $\mathcal{A}$ and the assumptions $\left(H_{1}\right)-\left(H_{3}\right)$ hold, then the following collection of sets is a Morse decomposition of $\mathcal{A}$ :

$$
\begin{aligned}
& \mathcal{M}_{n}=\{ \xi \in \mathcal{A} \backslash\{0\}:(1.2) \text { has a bounded entire solution }\left(y_{k}\right)_{k \in \mathbb{Z}} \\
&\text { through } \left.\xi \text { with } V^{ \pm}\left(y_{k}\right) \equiv n \text { on } \mathbb{Z} \text { and } 0 \notin \alpha(y) \cup \omega(\xi)\right\}, \\
& \mathcal{M}_{N_{ \pm}^{*}}=\{ \xi \in \mathcal{A} \backslash\{0\}:(1.2) \text { has a bounded entire solution }\left(y_{k}\right)_{k \in \mathbb{Z}} \\
&\text { through } \left.\xi \text { with } V^{ \pm}\left(y_{k}\right) \equiv N_{ \pm}^{*} \text { on } \mathbb{Z}\right\} \cup\{0\} \\
&\text { for } \left.0 \leq n \leq d+1, n \in 2 \mathbb{N}_{0} \backslash\left\{N_{+}^{*}\right\} \text { (resp. } n \in\left(2 \mathbb{N}_{0}+1\right) \backslash\left\{N_{-}^{*}\right\}\right) .
\end{aligned}
$$

Note that due to definitions of $V^{ \pm}$and $N_{ \pm}^{*}, \mathcal{M}_{N_{ \pm}^{*}} \neq\{0\}$ can only occur if the origin is non-hyperbolic or $N_{ \pm}^{*}=d+1$, and $\bar{d}$ is odd (resp. even) in the positive (resp. negative) feedback case.

The proof of Theorem 4.1 is based on several tools, whose logical structure is borrowed from [16]. However, there are certain differences: For instance, the finite dimensional state space of (1.2) allows to simplify various arguments based on the Arzelà-Ascoli theorem. Another simplification in our finite dimensional setting is the nonexistence of superexponentially decaying solutions established in Lemma 4.5. On the flip side, the nonconnectedness of orbits gives rise to technical difficulties in some arguments.

The next two results play a significant role in the sequel and show how bounded solutions (on either of the half-axes) can be characterized with the aid of $N_{ \pm}^{*}$ and functional $V^{ \pm}$. Although they correspond to continuous time results in $[16,12]$, we give a detailed proof.

## Proposition 4.2.

(a) If $\left(y_{k}\right)_{k \leq 0}$ is a nontrivial, bounded solution of (4.1) in $\mathbb{R}^{d+1}$, then $V^{ \pm}\left(y_{k}\right) \leq N_{ \pm}^{*}$ for all $k \leq 0$.
(b) If $\left(y_{k}\right)_{k \geq 0}$ is a nontrivial, bounded solution of (4.1) in $\mathbb{R}^{d+1}$, then $V^{ \pm}\left(y_{k}\right) \geq N_{ \pm}^{*}$ for all $k \geq 0$.

Proof. We present only the proof (a) in the negative feedback case, i.e. when $\left(H_{1}\right)-\left(H_{3}\right)$ with $\delta^{*}=-1$ hold. The proof of $(\mathrm{b})$ is analogous. The proof for the positive feedback case is also rather similar, but simpler, as there are less cases for the distribution of eigenvalues of (4.1) (cf. Lemma 2.1).

The eigenvalues of (4.1) are precisely the solutions of the characteristic equation (2.2). If $M^{*}$ is equal to 0,1 or $d+1$, then the statement is trivial. Otherwise, according to Lemma 2.1, $M^{*}$ is an even number, say $M^{*}=$ $2 n^{\prime}+2$, where $0 \leq n^{\prime} \leq \frac{d-2}{2}$.

We use the notations and ordering for eigenvalues introduced in Lemma 2.1. For a non-real eigenvalue $\lambda_{j}=r_{j} e^{i \varphi_{j}} \in S_{j}^{-}$(with $\varphi_{j} \in(0, \pi)$, $j \in \mathbb{Z}_{+}, j<\frac{d}{2}$ ) any linear combination of the real eigensolutions corresponding to $\lambda_{j}$ can be written in the form $c_{j} z_{j}$, where $c_{j} \in \mathbb{R}$ and

$$
z_{j, k}=\left(\begin{array}{c}
z_{j, k}(0)  \tag{4.3}\\
z_{j, k}(1) \\
\vdots \\
z_{j, k}(d-1) \\
z_{j, k}(d)
\end{array}\right)=\left(\begin{array}{c}
r_{j}^{k} \sin \left(k \varphi_{j}+\omega_{j}\right) \\
r_{j}^{k+1} \sin \left((k+1) \varphi_{j}+\omega_{j}\right) \\
\vdots \\
r_{j}^{k+d-1} \sin \left((k+d-1) \varphi_{j}+\omega_{j}\right) \\
r_{j}^{k+d} \sin \left((k+d) \varphi_{j}+\omega_{j}\right)
\end{array}\right)
$$

with some $\omega_{j} \in[0,2 \pi)$. From $\varphi_{j} \in(0, \pi)$ it is clear that sc $z_{j, k}$ equals the number of sign changes the sine function has on the interval $\left(k \varphi_{j}+\omega_{j}\right.$, $\left.(k+d) \varphi_{j}+\omega_{j}\right)$. According to the definition of the sector $S_{j}^{-}$, the length of
this interval is $d \varphi_{j} \in\left(2 j \pi, \frac{d}{d+1}(2 j+1) \pi\right)$. Now, from the definition of $V^{-}$ we obtain that $V^{-}\left(z_{j, k}\right)=2 j+1$ for all $k \in \mathbb{Z}_{-}$.

As for the eigensolutions corresponding to real eigenvalues, let us first recall that determined by the value of $b$, either there are exactly two positive eigenvalues of (4.1) $\left(\lambda_{+, 2} \leq \lambda_{+, 1}\right)$, and in that case there is no eigenvalue in the sector $S_{0}^{-}$, or there are no positive eigenvalues and there exists a pair of complex eigenvalues $\lambda_{0}, \overline{\lambda_{0}}$, where $\lambda_{0} \in S_{0}^{-}$. In the former case let $z_{0}^{1}, z_{0}^{2}$ and $z_{0}^{3}$ be defined by $z_{0, k}^{1}(j)=\lambda_{+, 1}^{k+j}, z_{0, k}^{2}(j)=\lambda_{+, 2}^{k+j}$ and $z_{0, k}^{3}(j)=(k+j) \lambda_{+, 2}^{k+j}$ for all $k \in \mathbb{Z}_{-}$and $j \in \mathbb{Z}_{+}, j \leq d$, respectively. Clearly, for any $i=1,2,3$ and for all $k \in \mathbb{Z}_{-}, k<-d$, every component of $z_{0, k}^{i}$ has the same sign, and in particular $V^{-}\left(z_{0, k}^{i}\right)=1$ for all $k<-d$. Concerning negative eigenvalues, if $d$ is even, say $d=2 \ell$ for some $\ell \in \mathbb{N}$, then any eigensolution corresponding to the unique negative eigenvalue $-r_{\ell}:=\lambda_{-}<0$ can be written in the form $c_{\ell} z_{\ell}$, where $c_{\ell} \in \mathbb{R}$ and $z_{\ell, k}(j)=\left(-r_{\ell}\right)^{k+j}$ for all $k \in \mathbb{Z}_{-}$and $j \in \mathbb{Z}_{+}, j \leq d$. Obviously sc $z_{\ell, k}=d$ and therefore also $V^{-}\left(z_{\ell, k}\right)=d+1=2 \ell+1$ hold for all $k \in \mathbb{Z}_{-}$. According to Lemma 2.1 there are no other real eigenvalues.

First, let us consider the hyperbolic case. Assume that $M^{*}=2 n^{\prime}+2$ for some integer $0 \leq n^{\prime} \leq \frac{d-2}{2}$. Then $N_{-}^{*}=M^{*}=2 n^{\prime}+2$. Let $y$ be an arbitrary solution of (4.2) which is bounded in backward time. This means that there exist an integer $0 \leq n \leq n^{\prime}$ and appropriate constants $\omega_{j} \in[0,2 \pi)$, $c_{0,1}, c_{0,2}, c_{0,3}, c_{j} \in \mathbb{R}$ with $j=0, \ldots, n$, such that

$$
\begin{equation*}
y=\sum_{j=1}^{3} c_{0, j} z_{0}^{j}+\sum_{j=0}^{n} c_{j} z_{j} \tag{4.4}
\end{equation*}
$$

where those coefficients are understood to be automatically set to zero, that correspond to eigensolutions which do not exist in the given case (e.g. there exists no positive eigenvalue).

We claim that $V^{-}\left(y_{k}\right) \leq V^{-}\left(z_{n, k}\right)=2 n+1<2 n^{\prime}+2=N_{-}^{*}$ holds for all $k \leq k_{1}$, where $k_{1}$ is to be defined soon. Then by monotonicity of $V^{-}$(proved in Theorem 3.3) we get statement (a). If $c_{0}, \ldots, c_{n}$ are all zeros, then the claim trivially holds. Otherwise we may assume w.l.o.g. that $c_{n} \neq 0$.

If $c_{0} \neq 0$, then there are no positive eigenvalues of (4.1), so the first sum in (4.4) is zero, and we readily obtain that for $n=0, V^{-}\left(y_{k}\right)=V^{-}\left(z_{n, k}\right)=$ $1<2 n^{\prime}+2=N_{-}^{*}$ holds.

Thus we may assume $n \geq 1$. From Lemma 2.1 one can easily obtain that there exists $\vartheta_{n}>0$, sufficiently small, such that the inequalities

$$
\begin{align*}
\vartheta_{n} & <\min \left\{\frac{\pi-\varphi_{n}}{2}, \frac{\varphi_{n}}{2}\right\}  \tag{4.5}\\
d \varphi_{n}+2 \vartheta_{n} & <(2 n+1) \pi \quad \text { and }  \tag{4.6}\\
2 n \pi & <d \varphi_{n}-2 \vartheta_{n} \tag{4.7}
\end{align*}
$$

hold. Moreover, by the Lemma 2.1 (c) we can choose $k_{0}<0$ so small that

$$
\begin{equation*}
\sum_{j=0}^{n-1}\left|c_{j}\right| r_{j}^{k}+\left|c_{0,1}\right| \lambda_{+, 1}^{k}+\left|c_{0,2}\right| \lambda_{+, 2}^{k}+\left|k c_{0,3}\right| \lambda_{+, 2}^{k}<\left|c_{n}\right| r_{n}^{k} \sin \vartheta_{n} \tag{4.8}
\end{equation*}
$$

holds for all $k \leq k_{0}$. Let $k_{1}=k_{0}-d$.
For simplicity, let us say that the $\ell$-th component of $z_{n, k}$, i.e. $z_{n, k}(\ell)=$ $r_{n}^{k} \sin \left((k+\ell) \varphi_{n}+\omega_{n}\right)$ is small if $\left|\sin \left((k+\ell) \varphi_{n}+\omega_{n}\right)\right| \leq \sin \vartheta_{n}$, otherwise the component will be said to be big. Inequality (4.8) yields that if $z_{n, k}(\ell)$ is big for some $\ell=0,1, \ldots, d$ and $k \leq k_{1}$, then $\operatorname{sgn} y_{k}(\ell)=\operatorname{sgn} z_{n, k}(\ell) \neq 0$ holds. Note that inequality (4.5) guarantees that there is no $k \leq k_{1}$ and component $0 \leq \ell<d$ such that $z_{n, k}(\ell)$ and $z_{n, k}(\ell+1)$ are both small. Similarly, (4.6)-(4.7) imply that $z_{n, k}(0)$ and $z_{n, k}(d)$ cannot be small at the same time. Moreover, by (4.5) one obtains that if $z_{n, k}(\ell)$ is small for some $k \leq k_{1}$ and $\ell \in\{1, \ldots, d-1\}$, then $\operatorname{sgn} z_{n, k}(\ell-1)=-\operatorname{sgn} z_{n, k}(\ell+1)$.

Combining them means that, if $k \leq k_{1}$ and $z_{n, k}(0)$ and $z_{n, k}(d)$ are both big, sc $y_{k}=\operatorname{sc} z_{n, k}$ holds, and $V^{-}\left(y_{k}\right)=V^{-}\left(z_{n, k}\right)=2 n+1<N_{-}^{*}$ follows.

There remains the case, when exactly one of $z_{n, k}(0)$ and $z_{n, k}(d)$ is small. W.l.o.g. we may assume that $z_{n, k}(d)$ is small, $z_{n, k}(0)$ is big and positive. Then, as $z_{n, k}(d-1)$ is also big,

$$
\operatorname{sc}\left(y_{k}(0), \ldots, y_{k}(d-1)\right)=\operatorname{sc}\left(z_{n, k}(0), \ldots, z_{n, k}(d-1)\right)
$$

follows. Thus either $V^{-}\left(y_{k}\right) \leq V^{-}\left(z_{n, k}\right)$ holds or sc $y_{k}=\mathrm{sc} z_{n, k}+1$ is even. We show that the latter case leads to contradiction. Should it hold, then $y_{k}(0)$ and $z_{n, k}(0)$ would be both positive, $y_{k}(d-1)$ and $z_{n, k}(d-1)$ would be both negative, while $y_{k}(d)>0$ and $z_{n, k}(d) \leq 0$. Now, since $z_{n, k}(d) \leq 0$ is small, $\left((k+d) \varphi_{n}+\omega_{n}\right) \in\left(-\vartheta_{n}, 0\right] \cup\left[\pi, \pi+\vartheta_{n}\right) \bmod 2 \pi$ holds. On the other hand from $z_{n, k}(d-1)<0$ and (4.5) it follows that $\left((k+d) \varphi_{n}+\omega_{n}\right) \in\left(-\vartheta_{n}, 0\right]$ $\bmod 2 \pi$, which implies by inequalities (4.6) and (4.7) that $z_{n, k}(0)<0$. This is a contradiction, so our proof is complete for the hyperbolic case.

If the origin is non-hyperbolic, and $M^{*} \geq 2$, then $N_{-}^{*}=M^{*}+1$. Keeping our notations from above, there is a pair of eigenvalues $\lambda_{n^{\prime}+1}, \overline{\lambda_{n^{\prime}+1}}$ on the unit circle, or -1 is a simple eigenvalue (if $d$ is even and $M^{*}=d$ ). In the latter case, statement (a) of the lemma is trivial. In the former case all backward bounded solutions of (4.2) can be written as (4.4), this time with $n \leq n^{\prime}+1$, and an argument similar to the one applied in the hyperbolic case shows that $V^{-}\left(y_{k}\right)=2 n+1 \leq 2 n^{\prime}+3=N_{-}^{*}$ for all $k \in \mathbb{Z}_{-}$.

Proposition 4.3. There exists an open neighborhood $U \subseteq J^{d+1}$ of 0 , such that for all nontrivial solutions $\left(y_{k}\right)_{k \in \mathbb{Z}}$ in $\mathbb{R}^{d+1}$ of the nonlinear equation (1.2) the following statements hold:
(a) If $y_{k} \in \bar{U} \cap \mathcal{A}$ for all $k \leq 0$, then $V^{ \pm}\left(y_{k}\right) \leq N_{ \pm}^{*}$ for all $k \in \mathbb{Z}$.
(b) If $y_{k} \in \bar{U} \cap \mathcal{A}$ for all $k \geq 0$, then $V^{ \pm}\left(y_{k}\right) \geq N_{ \pm}^{*}$ for all $k \in \mathbb{Z}$.

Proof. The proof relies mainly on Proposition 4.2. We only prove statement (a), since the proof of part (b) is analogous.

We argue by way of contradiction. If the statement is not true, then there exists a sequence of nontrivial, bounded entire solutions $y^{n}: \mathbb{Z} \rightarrow J^{d+1}$ of equation (1.2) such that $\sup _{k \leq 0}\left\|y_{k}^{n}\right\| \rightarrow 0$, as $n \rightarrow \infty$ and $V^{ \pm}\left(y_{k_{n}}^{n}\right)>N_{ \pm}^{*}$ holds for some $k_{n} \in \mathbb{Z}$.

Since $y^{n}$ is bounded, there exists an integer $j_{n} \leq k_{n}$, such that $\left\|y_{k}^{n}\right\|<$ $2\left\|y_{j_{n}}^{n}\right\|$ holds for all integers $n \in \mathbb{Z}_{+}$and $k \leq k_{n}$.

As equation (1.2) is autonomous and $V^{ \pm}$is monotone non-increasing, we may assume w.l.o.g. that $j_{n}=0$ for all $n \in \mathbb{Z}_{+}$. Recall that equation (1.2) can be written in the form (3.2), (3.4) and set

$$
z_{k}^{n}:=\frac{1}{\left\|y_{0}^{n}\right\|} y_{k}^{n} .
$$

Keeping the notations of Remark 3.1, $y^{n}$ is a bounded entire solution of (3.2),(3.4) for all $n \in \mathbb{Z}_{+}$, where $y, A, a$ and $b$ should be replaced throughout by $y^{n}, A^{n}, a^{n}$ and $b^{n}$, respectively. From the linearity of equation (3.2) one obtains that $z^{n}$ is an entire solution of

$$
z_{k+1}^{n}=A_{k}^{n} z_{k}^{n} \quad \text { for all } n \in \mathbb{Z}_{+} .
$$

By definition, $\left\|z_{k}^{n}\right\| \leq 2$ holds for all $n \in \mathbb{Z}_{+}$, so by the Cantor diagonalization argument, there exists a subsequence $\left(n_{\ell}\right)_{\ell=0}^{\infty}$ and a bounded sequence $\left(z_{k}\right)_{k=-\infty}^{0}$ such that $z_{k}^{n_{\ell}} \rightarrow z_{k}$, as $\ell \rightarrow \infty$, for all integers $k \leq 0$.

On the other hand, by our assumptions, for any fixed integer $k \leq 0$, $y_{k}^{n_{\ell}} \rightarrow 0$, as $\ell \rightarrow \infty$. Thus $a_{k}^{n_{\ell}} \rightarrow a=D_{1} f(0,0)$ and $b_{k}^{n_{\ell}} \rightarrow b=D_{2} f(0,0)$ hold as $\ell \rightarrow \infty$.

It follows then that $\left(z_{k}\right)_{k \leq 0}$ is a bounded solution of the linear equation

$$
z_{k+1}=A z_{k},
$$

where $A$ is defined by (4.1), moreover it is nontrivial, since $\left\|z_{0}\right\|=\left\|z_{0}^{n}\right\|=1$ for all $n \in \mathbb{Z}_{+}$. Thus Proposition 4.2 can be applied to obtain that $V^{ \pm}\left(z_{k}\right) \leq$ $N_{ \pm}^{*}$ holds for all $k \leq 0$. Furthermore, Theorem 3.3 (b) yields that there exists $k_{1} \leq 0$ such that $z_{k} \in \mathcal{R}^{ \pm}$for all $k \leq k_{1}$. Finally, since $z_{k_{1}}^{n_{\ell}} \rightarrow z_{k_{1}}$, as $\ell \rightarrow \infty$, it follows from Proposition 3.2 (b) that

$$
\lim _{\ell \rightarrow \infty} V^{ \pm}\left(z_{k_{1}}^{n_{\ell}}\right)=V^{ \pm}\left(z_{k_{1}}\right) \leq N_{ \pm}^{*},
$$

contradicting our assumption $N_{ \pm}^{*}<V^{ \pm}\left(y_{k}^{n_{\ell}}\right)=V^{ \pm}\left(z_{k}^{n_{\ell}}\right)$ for all $k \leq 0$.
From the Proposition 4.3 we immediately infer the next result.
Corollary 4.4. The following holds for solutions of (1.2):
(a) If $N_{ \pm}^{*}>1$, then there exists no heteroclinic solution towards the trivial solution.
(b) If there exists a homoclinic solution $\left(\tilde{y}_{k}\right)_{k \in \mathbb{Z}}$, then $V^{ \pm}\left(\tilde{y}_{k}\right) \equiv N_{ \pm}^{*}$ on $\mathbb{Z}$.
(c) In particular, if the trivial solution is hyperbolic, then there exists no heteroclinic solution in the positive feedback case, while in case of negative feedback it may only exist if $M^{*}=N_{-}^{*}=1$.

Proof. (a) Fixed-points $y^{*} \neq 0$ of (1.2) have equal, nonzero components, so $V^{ \pm}\left(y^{*}\right) \leq 1$ and $y^{*} \in \mathcal{R}^{ \pm}$. Assume now to the contrary that $\left(\tilde{y}_{k}\right)_{k \in \mathbb{Z}}$ is an entire solution with $\lim _{k \rightarrow \infty} \tilde{y}_{k}=0$ and $\lim _{k \rightarrow-\infty} \tilde{y}_{k}=y^{*}$. Then from Proposition 3.2 (b) we get $\lim _{k \rightarrow-\infty} V^{ \pm}\left(\tilde{y}_{k}\right) \leq 1$, while Proposition 4.3 (b) yields $\lim _{k \rightarrow \infty} V^{ \pm}\left(\tilde{y}_{k}\right) \geq N_{ \pm}^{*}>1$, a contradiction to Proposition 3.3 (a).
(b) This follows immediately from Propositions 4.3 and 3.3 (a).
(c) Assume the trivial solution is hyperbolic. If $M^{*}$ equals 0 or $d+1$, then the statement holds trivially since either the local unstable or the local stable manifold is trivial. Except the case of negative feedback with $M^{*}=N_{-}^{*}=1$, $N_{ \pm}^{*}=M^{*}$ is an odd (resp. even) number in the positive (resp. negative) feedback case. Now, recall that $V^{ \pm}$takes on only even (resp. odd) values, then apply (b) to conclude the proof.

The next technical lemma shows that solutions of (1.2) can neither grow nor decrease faster than exponentially in some neighborhood of the origin.

Lemma 4.5. Let $\left(y_{k}\right)_{k \in \mathbb{I}}$ be a solution of (3.2). If $\tilde{a}_{1}, \tilde{b}_{0}$ and $\tilde{b}_{1}$ are positive reals with $0<a_{k}<\tilde{a}_{1}$ and $\tilde{b}_{0}<\left|b_{k}\right|<\tilde{b}_{1}$ for all $k \in \mathbb{I}$, then there exist positive constants $c$ and $C$ (that may depend on $\tilde{a}_{1}, \tilde{b}_{0}$ and $\tilde{b}_{1}$ ) such that

$$
c\left\|y_{k}\right\| \leq\left\|y_{k+1}\right\| \leq C\left\|y_{k}\right\| \quad \text { if } k \in \mathbb{I}^{\prime} .
$$

Proof. Without loss of generality we may use the $\|\cdot\|_{1}$-norm. Since $A_{k}$ is invertible for all $a_{k} \neq 0 \neq b_{k}$, and

$$
\left\|y_{k+1}\right\|_{1} \leq\left\|A_{k}\right\|_{1}\left\|y_{k}\right\|_{1}=\max \left\{1+\left|a_{k}\right|,\left|b_{k}\right|\right\}\left\|y_{k}\right\|_{1}
$$

and

$$
\left\|y_{k+1}\right\|_{1} \geq \frac{\left\|y_{k}\right\|_{1}}{\left\|A_{k}^{-1}\right\|_{1}}=\min \left\{\left|b_{k}\right|, \frac{\left|b_{k}\right|}{\left|b_{k}\right|+a_{k}}\right\}\left\|y_{k}\right\|_{1}
$$

thus

$$
c:=\min \left\{\tilde{b}_{0}, \frac{\tilde{b}_{0}}{\tilde{b}_{0}+\tilde{a}_{1}}\right\} \quad \text { and } \quad C:=\max \left\{1+\tilde{a}_{1}, \tilde{b}_{1}\right\}
$$

are appropriate choices.
Remark 4.6. From the above proof and the $C^{1}$-smoothness of $f$ it is clear that for any solution $\left(y_{k}\right)_{k \in \mathbb{I}}$ of (3.2), (3.4) on some discrete interval $\mathbb{I}$ and any compact set $K \subset \mathbb{R}^{d+1}$, there exists $C>0$, such that $y_{k} \in K$ for some $k \in \mathbb{I}$ implies $\left\|y_{k+1}\right\| \leq C\left\|y_{k}\right\|$. Moreover, if $K$ is a sufficiently small neighborhood of 0 , then there exists also $c>0$ such that $c\left\|y_{k}\right\| \leq\left\|y_{k+1}\right\|$.

We combine the next three lemmas with Proposition 4.3 in order to prove that the family of sets $\mathcal{M}_{n}$ fulfills the "Morse properties" $\left(m_{1}\right)$ and $\left(m_{2}\right)$.

Lemma 4.7. Suppose that $\left(y_{k}\right)_{k \in \mathbb{Z}}$ is a bounded entire solution of (1.2) through $\xi \in \mathcal{A}$.
(a) If $\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right)=N$, then $V^{ \pm}(\eta)=N$ for any $\eta \in \omega(\xi) \backslash\{0\}$.
(b) If $\lim _{k \rightarrow-\infty} V^{ \pm}\left(y_{k}\right)=K$, then $V^{ \pm}(\eta)=K$ for any $\eta \in \alpha(y) \backslash\{0\}$.

Proof. We only prove statement (a), since (b) can be shown analogously.
For an arbitrary $\eta \in \omega(\xi) \backslash\{0\}$, there exist a monotone sequence $k_{n} \rightarrow \infty$, as $n \rightarrow \infty$, such that $y_{k_{n}} \rightarrow \eta$, and - by the invariance of $\omega(\xi) \backslash\{0\}-$ there exists a bounded entire solution $z$ of (1.2) such that $z_{0}=\eta$ and $z_{k} \in \omega(\xi) \backslash\{0\}$ holds for all $k \in \mathbb{Z}$. By Theorem $3.3(\mathrm{~b})$, there exist integers $\ell_{1}<0<\ell_{2}$ such that $z_{\ell_{1}}$ and $z_{\ell_{2}}$ both belong to $\mathcal{R}^{ \pm}$. Since $z_{\ell_{1}} \in \omega(\xi) \backslash\{0\}$, there exists a monotone sequence $k_{n}^{\prime} \rightarrow \infty$, as $n \rightarrow \infty$, such that $y_{k_{n}^{\prime}} \rightarrow z_{\ell_{1}}$, and from the continuity of the right-hand side of equation (1.2) one infers that $y_{k_{n}+\ell_{2}} \rightarrow z_{\ell_{2}}$, as $n \rightarrow \infty$.

Combining these with Proposition $3.2(\mathrm{~b})$ and with the assumption $\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right)=N$ we get that

$$
V^{ \pm}\left(z_{\ell_{1}}\right)=\lim _{n \rightarrow \infty} V^{ \pm}\left(y_{k_{n}^{\prime}}\right)=N=\lim _{n \rightarrow \infty} V^{ \pm}\left(y_{k_{n}+\ell_{2}}\right)=V^{ \pm}\left(z_{\ell_{2}}\right)
$$

Using that $V^{ \pm}$is non-increasing and that $\ell_{1}<0<\ell_{2}$, and $z_{0}=\eta$ hold, $V^{ \pm}(\eta)=N$ is established.

Lemma 4.8. Suppose that $\left(y_{k}\right)_{k \in \mathbb{Z}}$ is a bounded entire solution of (1.2) through $\xi \in \mathcal{A}$.
(a) If $\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right) \neq N_{ \pm}^{*}$, then either $\omega(\xi)=\{0\}$ or $0 \notin \omega(\xi)$.
(b) If $\lim _{k \rightarrow-\infty} V^{ \pm}\left(y_{k}\right) \neq N_{ \pm}^{*}$, then either $\alpha(y)=\{0\}$ or $0 \notin \alpha(y)$.

Proof. Again, we only prove part (a). Let us assume that $0 \in \omega(\xi)$, but $\omega(\xi) \neq\{0\}$, furthermore, let $U$ be defined as in Proposition 4.3. Then there exists $U_{1}$, an open neighborhood of 0 in $\mathcal{A}$ such that $U_{1} \subset U$ and $\omega(\xi) \nsubseteq \overline{U_{1}}$. This, together with $0 \in \omega(\xi)$ imply that $y_{k}$ enters and leaves $\overline{U_{1}}$ infinitely many times as $k \rightarrow \infty$. Thus there exist positive integer sequences $k_{n}, \sigma_{n}, \tau_{n}$ ( $n \geq 0$ ), such that

$$
\begin{gathered}
y_{k_{n}} \rightarrow 0 \quad \text { and } \quad k_{n} \rightarrow \infty, \quad \text { as } n \rightarrow \infty \\
k_{n}+\tau_{n}<k_{n+1}-\sigma_{n+1}-1 \\
y_{k} \in \overline{U_{1}}, \quad \text { if } \quad k_{n}-\sigma_{n} \leq k \leq k_{n}+\tau_{n}
\end{gathered}
$$

and

$$
y_{k_{n}-\sigma_{n}-1}, y_{k_{n}+\tau_{n}+1} \notin \overline{U_{1}}
$$

hold for all $n \geq 0$. By Remark 4.6 we get $\tau_{n} \rightarrow \infty$ and $\sigma_{n} \rightarrow \infty$, as $n \rightarrow \infty$.
Now, let

$$
z_{k}^{n}=y_{k_{n}-\sigma_{n}+k} \quad \text { and } \quad w_{k}^{n}=y_{k_{n}+\tau_{n}+k} \quad \text { for } k \in \mathbb{Z}
$$

Using the compactness of $\mathcal{A}$ and applying a Cantor diagonalization argument one obtains that there exist a subsequence $\left(n_{\ell}\right)_{\ell \in \mathbb{N}_{0}}$ and bounded entire solutions $\left(z_{k}\right)$ and $\left(w_{k}\right)$ of (1.2), such that $z_{k}^{n_{\ell}} \rightarrow z_{k}$ and $w_{k}^{n_{\ell}} \rightarrow w_{k}$ hold for all $k \in \mathbb{Z}$, as $\ell \rightarrow \infty$. Solutions $z$ and $w$ are nontrivial, since both $z_{-1}$ and $w_{1}$ are from $\mathcal{A} \backslash U_{1}$. On the other hand, $z_{k} \in \overline{U_{1}} \subset \bar{U}$ and $w_{-k} \in \overline{U_{1}} \subset \bar{U}$ hold for $k \in \mathbb{Z}_{+}$. Thus by Proposition 4.3 it follows that

$$
\begin{equation*}
V^{ \pm}\left(w_{k}\right) \leq N_{ \pm}^{*} \leq V^{ \pm}\left(z_{k}\right) \quad \text { for all } k \in \mathbb{Z} \tag{4.9}
\end{equation*}
$$

Observe that $z_{0}=\lim _{\ell \rightarrow \infty} y_{k_{n_{\ell}}-\sigma_{n_{\ell}}}$, hence $z_{0} \in \omega(\xi) \backslash\{0\}$, and similarly $w_{0} \in \omega(\xi) \backslash\{0\}$. Then, from Lemma 4.7 it follows that

$$
V^{ \pm}\left(z_{0}\right)=V^{ \pm}\left(w_{0}\right)=\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right)
$$

This combined with (4.9) (for $k=0$ ) leads to $\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right)=N_{ \pm}^{*}$.
Lemma 4.9. Suppose that $\left(y_{k}\right)_{k \in \mathbb{Z}}$ is a bounded entire solution of (1.2) through $\xi \in \mathcal{A} \backslash\{0\}$.
(a) If $N:=\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right) \neq N_{ \pm}^{*}$, then either $\omega(\xi)=\{0\}$ or $\omega(\xi) \subseteq \mathcal{M}_{N}$.
(b) If $K:=\lim _{k \rightarrow-\infty} V^{ \pm}\left(y_{k}\right) \neq N_{ \pm}^{*}$, then either $\alpha(y)=\{0\}$ or $\alpha(y) \subseteq \mathcal{M}_{K}$.

Proof. In order to prove assertion (a), let us assume that $\omega(\xi) \neq\{0\}$ and that $N:=\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right) \neq N_{ \pm}^{*}$. According to Lemma 4.8, $0 \notin \omega(\xi)$. Let $\eta \in \omega(\xi)$ be arbitrarily chosen, then using the invariance of $\omega(\xi)$, there exists a bounded entire solution $\left(z_{k}\right)$, such that $z_{0}=\eta$ and $z_{k} \in \omega(\xi)$ holds for all $k \in \mathbb{Z}$. From compactness of $\omega(\xi), \alpha(z) \cup \omega(\eta) \subseteq \omega(\xi)$ follows. This implies that $0 \notin \alpha(z) \cup \omega(\eta)$. On the other hand, $z_{k} \in \omega(\xi)$ and Lemma 4.7 together ensure that $V^{ \pm}\left(z_{k}\right)=N$ for all $k \in \mathbb{Z}$, meaning that $\eta \in \mathcal{M}_{N}$ holds.

The proof of statement (b) is analogous.
The next lemma is a step towards the compactness proof of Morse sets.
Lemma 4.10. For every $N \in 2 \mathbb{N}_{0} \backslash N_{+}^{*}$ (resp. $\left.N \in\left(2 \mathbb{N}_{0}+1\right) \backslash N_{-}^{*}\right)$ there exists an open neighborhood $\tilde{U}$ in $\mathcal{A}$ of the origin, such that $\mathcal{M}_{N} \cap \tilde{U}=\emptyset$.

Proof. Assume to the contrary that there is a non-negative even (resp. odd) $N \neq N_{ \pm}^{*}$ and a sequence $\left(\xi^{n}\right)_{n \geq 0}$ in $\mathcal{M}_{N}$ such that $\xi^{n} \rightarrow 0$ as $n \rightarrow \infty$. For each $n \geq 0$ let $y^{n}$ be a bounded entire solution such that $y_{0}^{n}=\xi^{n}$, $0 \notin \alpha\left(y^{n}\right) \cup \omega\left(\xi^{n}\right)$ and $V^{ \pm}\left(y_{k}^{n}\right)=N$ hold for all $k \in \mathbb{Z}$.

We claim that $\omega\left(\xi^{n}\right)$ cannot be a subset of $\bar{U}$, where $U$ is from Proposition 4.3. To see this, suppose that $\omega\left(\xi^{n}\right) \subseteq \bar{U}$. Let $\eta \in \omega\left(\xi^{n}\right)$. By invariance
of $\omega\left(\xi^{n}\right)$ there exist a bounded entire solution $\left(z_{k}\right)$ through $\xi^{n}$, such that $z_{k} \in \omega\left(\xi^{n}\right) \subseteq \bar{U}$ holds for all $k \in \mathbb{Z}$. Thus Porposition 4.3 implies that $V^{ \pm}\left(z_{k}\right)=N_{ \pm}^{*}$ holds for all $k \in \mathbb{Z}$. However, $V^{ \pm}(\eta)=N$ should also hold by Lemma 4.8, which is a contradiction, so $\omega\left(\xi^{n}\right) \nsubseteq \bar{U}$ for all $n \geq 0$.

Assume for definiteness that $N>N_{ \pm}^{*}$. Since $\lim _{n \rightarrow \infty} \xi^{n}=0$, there exists a series of non-negative integers $\left(\sigma_{n}\right)_{n \in \mathbb{N}_{0}}$, such that

$$
y_{k}^{n} \in \bar{U} \quad \text { for } 0 \leq k \leq \sigma_{n} \quad \text { and } \quad y_{\sigma_{n}+1}^{n} \notin \bar{U}
$$

holds for all $n \geq 0$. By Remark 4.6, it follows that $\sigma_{n} \rightarrow \infty$, as $n \rightarrow \infty$. Letting $z_{k}^{n}:=y_{k+\sigma_{n}}^{n}$ for all $k \in \mathbb{Z}$ and applying a Cantor diagonalization argument one obtains that there is subsequence $\left(z^{n_{\ell}}\right)_{\ell \in \mathbb{N}_{0}}$ and an entire solution $\left(z_{k}\right)$ of (1.2) such that $z_{k}^{n_{\ell}} \rightarrow z_{k}$ holds for all $k \in \mathbb{Z}$, as $\ell \rightarrow \infty$. Solution $\left(z_{k}\right)$ is nontrivial because $z_{1} \notin U$. On the other hand, $z_{k} \in \bar{U}$ holds for all $k \leq 0$, so Proposition 4.3 yields that $V^{ \pm}\left(z_{k}\right) \leq N_{ \pm}^{*}$ holds for all $k \in \mathbb{Z}$.

Finally, Theorem 3.3 (b) ensures that there exists some $k^{\prime}>0$, so that $z_{k^{\prime}} \in \mathcal{R}^{ \pm}$. As $y_{\sigma_{n_{\ell}+k^{\prime}}}^{n_{\ell}}=z_{k^{\prime}}^{n_{\ell}} \rightarrow z_{k^{\prime}}(\ell \rightarrow \infty)$, we obtain

$$
N=\lim _{\ell \rightarrow \infty} V^{ \pm}\left(y_{\sigma_{n_{\ell}}+k^{\prime}}^{n_{\ell}}\right)=V^{ \pm}\left(z_{k^{\prime}}\right) \leq N_{ \pm}^{*}
$$

from Proposition $3.2(\mathrm{~b})$, which is a contradiction that proves our statement.
The proof for the case $N<N_{ \pm}^{*}$ is analogous.
Lemma 4.11. The set $\mathcal{M}_{N}$ is closed for all $N \in 2 \mathbb{N}_{0} \cup\left\{N_{+}^{*}\right\}$ (resp. all $\left.N \in\left(2 \mathbb{N}_{0}+1\right) \cup\left\{N_{-}^{*}\right\}\right)$.

Proof. Let $\left(\xi^{n}\right)_{n \geq 0}$ be a sequence in $\mathcal{M}_{N}$ for some $N \in 2 \mathbb{N}_{0} \cup\left\{N_{+}^{*}\right\}$ (resp. $\left.N \in\left(2 \mathbb{N}_{0}+1\right) \cup\left\{N_{-}^{*}\right\}\right)$, such that $\lim _{n \rightarrow \infty} \xi^{n}=\xi$ for some $\xi \in \mathcal{A}$. We claim that $\xi \in \mathcal{M}_{N}$.

If $\xi=0$, then Lemma 4.10 yields that $N=N_{ \pm}^{*}$. By definition of $\mathcal{M}_{N_{ \pm}^{*}}$, $0 \in \mathcal{M}_{N_{ \pm}^{*}}$ also holds, so the statement is proved in this particular case.

Assume now that $\xi \neq 0$. By definition of $\mathcal{M}_{N}$ and using $\xi^{n} \in \mathcal{M}_{N}$, there exist bounded entire solutions $\left(y_{k}^{n}\right)$ of (1.2), such that $y_{0}^{n}=\xi^{n}, y_{k}^{n} \in \mathcal{M}_{N}$ and $V^{ \pm}\left(y_{k}^{n}\right)=N$ hold for all $n \geq 0$ and $k \in \mathbb{Z}$. The compactness of $\mathcal{A}$ and a Cantor diagonalization argument leads to a subsequence $\left(y^{n_{\ell}}\right)_{\ell \in \mathbb{N}}$, and a solution $\left(y_{k}\right)_{k \in \mathbb{Z}}$ of (1.2), so that $y_{k}^{n_{\ell}} \rightarrow y_{k}$ holds for every $k \in \mathbb{Z}$, as $\ell \rightarrow \infty$.

Necessarily, $y_{0}=\xi$ holds. This combined with the compactness of $\mathcal{A}$ yields that $y_{k} \in \mathcal{A} \backslash\{0\}$ for all $k \in \mathbb{Z}$. Applying Theorem 3.3 (b) we obtain integers $\ell_{1}<0<\ell_{2}$ such that $y_{\ell_{1}}$ and $y_{\ell_{2}}$ are from $\mathcal{R}^{ \pm}$, and that

$$
\begin{equation*}
V^{ \pm}\left(y_{\ell_{1}}\right)=\lim _{k \rightarrow-\infty} V^{ \pm}\left(y_{k}\right), \quad \text { and } \quad V^{ \pm}\left(y_{\ell_{2}}\right)=\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right) \tag{4.10}
\end{equation*}
$$

hold. On the other hand, $y_{\ell_{1}}^{n_{\ell}} \rightarrow y_{\ell_{1}}$ and $y_{\ell_{2}}^{n_{\ell}} \rightarrow y_{\ell_{2}}$ hold as $\ell \rightarrow \infty$, so Proposition $3.2(\mathrm{~b})$ assures that $V^{ \pm}\left(y_{\ell_{1}}\right)=V^{ \pm}\left(y_{\ell_{2}}\right)=N$. This combined with (4.10) yields that $V^{ \pm}\left(y_{k}\right)=N$ for all $k \in \mathbb{Z}$. In case $N=N_{ \pm}^{*}, \xi \in \mathcal{M}_{N}$
is readily established. Otherwise, it remains to prove that $0 \notin \alpha(y) \cup \omega(\xi)$. In this case Lemma 4.10 gives an open neighborhood $\tilde{U}$ of 0 in $\mathcal{A}$, such that $y_{k}^{n} \in \mathcal{A} \backslash \tilde{U}$ holds for all $k \in \mathbb{Z}$ and $n \in \mathbb{N}_{0}$. Hence,

$$
y_{k} \in \overline{\mathcal{A} \backslash \tilde{U}}=\mathcal{A} \backslash \tilde{U}
$$

holds for all $k \in \mathbb{Z}$, which leads to $\alpha(y) \cup \omega(\xi) \subseteq \mathcal{A} \backslash \tilde{U}$. This shows that $0 \notin \alpha(y) \cup \omega(\xi)$ and completes the proof.

Now, we are finally in the position to prove our main theorem.
Proof of Theorem 4.1. By definition, the sets $\mathcal{M}_{n}$ are pairwise disjoint and invariant. Since $\mathcal{A}$ is compact, Lemma 4.11 yields that $\mathcal{M}_{n}$ is compact for all possible $n$.

It remains to prove that these sets fulfill the "Morse properties" $\left(m_{1}\right)$ and $\left(m_{2}\right)$, i.e. for all $\xi \in \mathcal{A}$ and any bounded entire solution $\left(y_{k}\right)_{k \in \mathbb{Z}}$ for which $y_{0}=\xi$ holds, there exist $i \geq j$ with $\alpha(y) \subseteq \mathcal{M}_{i}$ and $\omega(\xi) \subseteq \mathcal{M}_{j}$, and in case $i=j$, then $\xi \in \mathcal{M}_{i}$ (thus, $y_{k} \in \mathcal{M}_{i}$ for every $k \in \mathbb{Z}$ ).

In order to show this, let $\xi \in \mathcal{A} \backslash\{0\}$ be arbitrary (the statement is trivial for $\xi=0$ ), and $\left(y_{k}\right)_{k \in \mathbb{Z}}$ be a bounded entire solution of (1.2) for which $y_{0}=\xi$ holds. Furthermore, define

$$
N:=\lim _{k \rightarrow \infty} V^{ \pm}\left(y_{k}\right) \quad \text { and } \quad K:=\lim _{k \rightarrow-\infty} V^{ \pm}\left(y_{k}\right)
$$

Note that from monotonicity of $V^{ \pm}$we get that $N \leq K$.
First, observe that if $N=N_{ \pm}^{*}$, then $\omega(\xi) \subseteq \mathcal{M}_{N_{ \pm}^{*}}$. In order to prove this, choose $\eta \in \omega(\xi)$ arbitrarily. If $\eta=0$, then $\eta \in \mathcal{M}_{N_{ \pm}^{*}}$ holds by definition, thus we may assume now that $\eta \neq 0$. By Lemma 4.7 we obtain that $V^{ \pm}(\eta)=N_{ \pm}^{*}$. Moreover, by the invariance of $\omega(\xi) \backslash\{0\}$ there exists an entire solution $\left(z_{k}\right)$ in $\omega(\xi) \backslash\{0\}$, such that $z_{0}=\eta$. Thus Lemma 4.7 yields that $V^{ \pm}\left(z_{k}\right)=N_{ \pm}^{*}$ for all $k \in \mathbb{Z}$, meaning that $\eta \in \mathcal{M}_{N_{ \pm}^{*}}$ holds.

A similar argument can be applied to prove that $K=N_{ \pm}^{*}$ implies that $\alpha(y) \subseteq \mathcal{M}_{N_{ \pm}^{*}}$ holds.

We will distinguish four cases in terms of the values of $N$ and $K$.
Case 1. If $N=K=N_{ \pm}^{*}$, then $\alpha(y) \cup \omega(\xi) \subseteq \mathcal{M}_{N_{ \pm}^{*}}$ holds by our previous observation. Moreover, from the monotonicity of $V^{ \pm}$it follows that $V^{ \pm}\left(y_{k}\right) \equiv N_{ \pm}^{*}$ on $\mathbb{Z}$. This implies $\xi \in \mathcal{M}_{N_{ \pm}^{*}}$, thus both $\left(m_{1}\right)$ and $\left(m_{2}\right)$ hold.
Case 2. $N=N_{ \pm}^{*}<K$. As already shown, $\omega(\xi) \subseteq \mathcal{M}_{N_{ \pm}^{*}}$ holds in this case.
Moreover, observe that $\alpha(y) \neq\{0\}$. Otherwise Proposition 4.3 would imply $V^{ \pm}\left(y_{k}\right) \leq N_{ \pm}^{*}$ for $k \in \mathbb{Z}$, and thus $K \leq N_{ \pm}^{*}=N$, which is impossible.

Therefore, Lemma 4.9 can be applied to obtain $\alpha(y) \subseteq \mathcal{M}_{K}$, so property $\left(m_{1}\right)$ is fulfilled. Note that $\left(m_{2}\right)$ holds automatically, since the two Morse sets in question, i.e. $\mathcal{M}_{N_{ \pm}^{*}}$ and $\mathcal{M}_{K}$, are different.

Case 3. A similar argument applies in the case when $N<N_{ \pm}^{*}=K$.
Case 4. If $N \neq N_{ \pm}^{*} \neq K$, then Lemma 4.9 yields that either $\omega(\xi)=\{0\}$ or $\omega(\xi) \subseteq \mathcal{M}_{N}$. Similarly, either $\alpha(y)=\{0\}$ or $\alpha(y) \subseteq \mathcal{M}_{K}$ holds. Note that $\omega(\xi)$ and $\alpha(y)$ cannot be both $\{0\}$ in this case, because then Proposition 4.3 would imply $V^{ \pm}\left(y_{k}\right) \equiv N_{ \pm}^{*}$ on $\mathbb{Z}$, contradicting assumption $N \neq N_{ \pm}^{*} \neq K$.

If $\omega(\xi) \neq\{0\} \neq \alpha(y)$, then Lemma 4.9 yields that $\omega(\xi) \subseteq \mathcal{M}_{N}$ and $\alpha(y) \subseteq \mathcal{M}_{K}$ hold, so $\left(m_{1}\right)$ is fulfilled. If $K=N$, then their definition imply that $V^{ \pm}\left(y_{k}\right)=N=K$ for all $k \in \mathbb{Z}$. Since the limit sets are now also assumed to be nontrivial, thus Lemma 4.8 ensures that $0 \notin \alpha(y) \cup \omega(\xi)$, thus $y_{k} \in \mathcal{M}_{N}$ also holds for all $k \in \mathbb{Z}$. This establishes property $\left(m_{2}\right)$.

If $\omega(\xi)=\{0\} \neq \alpha(y)$, then $\omega(\xi) \subseteq \mathcal{M}_{N_{ \pm}^{*}}$ holds by definition. Furthermore, Proposition 4.3 implies that $V^{ \pm}\left(y_{k}\right) \geq N_{ \pm}^{*}$ holds for all $k \in \mathbb{Z}$, and consequently $N_{ \pm}^{*}<N \leq K$. On the other hand Lemma 4.9 yields that $\alpha(y) \subseteq \mathcal{M}_{K}$, so $\left(m_{1}\right)$ holds. Property $\left(m_{2}\right)$ is fulfilled automatically.

An analogous argument applies for the case when $\omega(\xi) \neq\{0\}=\alpha(y)$. We have taken all possible cases into consideration, so our proof is complete.

## 5 Applications

Since the existence of a global attractor is assumed in Theorem 4.1, we begin with a condition for dissipativity:
Lemma 5.1. Suppose $A \in \mathbb{R}^{d \times d}$ and $H: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ satisfy:
(i) There exist reals $a \in(0,1), K \geq 1$ with $\left\|A^{k}\right\| \leq K a^{k}$ for all $k \in \mathbb{Z}_{+}$,
(ii) there exist reals $\beta_{0}, \beta_{1} \geq 0$ so that $\|H(y)\| \leq \beta_{0}+\beta_{1}\|y\|$ for all $y \in \mathbb{R}^{d}$.

If $\beta_{1} \in\left[0, \frac{1-a}{K}\right)$ and $R>\frac{K \beta_{0}}{1-a-K \beta_{1}}$, then $B_{R}(0)$ is an absorbing set for

$$
\begin{equation*}
y_{k+1}=A y_{k}+H\left(y_{k}\right) \tag{5.1}
\end{equation*}
$$

Proof. Due to the variation of constants formula (cf. [18, p. 100, Theorem 3.1.16(a)]) the semigroup induced by (5.1) satisfies

$$
\phi(k ; \xi):=A^{k} \xi+\sum_{l=0}^{k-1} A^{k-l-1} H(\phi(l ; \xi)) \quad \text { for all } k \in \mathbb{Z}_{+}
$$

and our assumptions yield

$$
\begin{aligned}
\|\phi(k ; \xi)\| & \stackrel{(i)}{\leq} K a^{k}\|\xi\|+K \sum_{l=0}^{k-1} a^{k-l-1}\|H(\phi(l ; \xi))\| \\
& \stackrel{(i i)}{\leq} K a^{k}\|\xi\|+K \beta_{0} \sum_{l=0}^{k-1} a^{k-l-1}+K \beta_{1} \sum_{l=0}^{k-1} a^{k-l-1}\|\phi(l ; \xi)\|
\end{aligned}
$$

for $k \in \mathbb{Z}_{+}$. With the formula for the geometric series in mind, this implies

$$
\frac{\|\phi(k ; \xi)\|}{a^{k}} \leq K\|\xi\|+\frac{K \beta_{0}}{1-a}\left(a^{-k}-1\right)+\frac{K \beta_{1}}{a} \sum_{l=0}^{k-1} \frac{\|\phi(l ; \xi)\|}{a^{l}} \quad \text { for all } k \in \mathbb{Z}_{+}
$$

using the Grönwall inequality [18, p. 348, Proposition A.2.1(a)] we obtain $\|\phi(k ; \xi)\| \leq K\left(a+K \beta_{1}\right)^{k}\|\xi\|+\frac{K \beta_{0}}{1-a-K \beta_{1}}\left(1-\left(a+K \beta_{1}\right)^{k}\right) \quad$ for all $k \in \mathbb{Z}_{+}$ and $a+K \beta_{1} \in[0,1)$ leads to $\lim _{\sup }^{k \rightarrow \infty} \boldsymbol{\|}\|(t ; \xi)\| \leq \frac{K \beta_{0}}{1-a-K \beta_{1}}$.

### 5.1 Life sciences

Applications from the life sciences typically require solutions with values in $\mathbb{R}_{+}:=[0, \infty)$. We thus present some dissipative delay-difference equations

$$
\begin{equation*}
x_{k+1}=f_{0}\left(x_{k}, x_{k-d}\right) \tag{5.2}
\end{equation*}
$$

with $C^{1}$-right-hand side $f_{0}: \mathbb{R}_{+} \times \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$fitting into our setting. They commonly have a unique equilibrium $x_{*}>0$. Hence, the shifted equation

$$
x_{k+1}=f\left(x_{k}, x_{k-d}\right), \quad f(x, y):=f_{0}\left(x+x_{*}, y+x_{*}\right)-x_{*}
$$

with $f: J \times J \rightarrow J$ and the closed interval $J:=\left[-x_{*}, \infty\right)$ possesses the trivial equilibrium. Under dissipativity conditions on (5.2) we obtain absorbing sets of the form $[0, R]^{d+1}$ with $R>R_{0}$ for some $R_{0}>0$ and consequently global attractors $\mathcal{A} \subseteq\left[-x_{*}, R-x_{*}\right]^{d+1}$ for (1.2). For brevity, let us introduce $\underline{x}_{*}:=\left(x_{*}, \ldots, x_{*}\right) \in \mathbb{R}^{d+1}$.

In the literature often sufficient conditions for global asymptotic stability of $x_{*}$ are given (cf., for example $[9,8]$ ). In this case, only one Morse set $\mathcal{M}_{N_{ \pm}^{*}}$ is obtained from Theorem 4.1. Let us consequently list several dissipative difference equations that fulfill the assumptions of our main theorem, i.e. $\left(H_{1}\right)-\left(H_{3}\right)$, and indicate bifurcations which lead to more complex dynamics and possibly Morse sets.

Throughout, suppose that $\alpha \in(0,1), \beta>0$.

### 5.1.1 Clarke-type models

Given a $C^{1}$-function $h: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$we denote difference equations

$$
\begin{equation*}
x_{k+1}=\alpha x_{k}+h\left(x_{k-d}\right) \tag{5.3}
\end{equation*}
$$

as Clarke's delayed recruitment models. For such problems with $h \in C^{3}$ it was recently shown in [5] that even if the Schwarzian

$$
S h(x):=\frac{h^{\prime \prime \prime}(x)}{h^{\prime}(x)}-\frac{3}{2}\left(\frac{h^{\prime \prime}(x)}{h^{\prime}(x)}\right)^{2}
$$

| $h$ | $x_{*}$ | $R_{0}$ | GAS |
| ---: | :---: | :---: | :---: |
| $(5.4)$ | - | $\frac{\beta}{1-\alpha}$ | $p \leq 1$ |
| $(5.6)$ for $p<1$ | $\left(\frac{\alpha+\beta-1}{1-\alpha}\right)^{1 / p}$ | $\frac{\beta \eta}{1+\eta^{p}}-\beta_{1} \eta$ | $1-\alpha<\beta<\frac{1-\alpha}{\alpha^{d}(1-\alpha)-\beta_{1}}$ |
| $(5.6)$ for $p=1$ | $\frac{\alpha+\beta-1}{1-\alpha}$ | $\frac{\beta}{1-\alpha}$ | $1-\alpha<\beta<\frac{1-\alpha}{1-\alpha^{1+\alpha}}$ |
| $(5.7)$ | $W\left(\frac{\beta}{1-\alpha}\right)$ | $\frac{\beta}{1-\alpha}$ | $\frac{\beta}{1-\alpha} \leq e$ |

Table 5.1: Clarke-type models (5.3): Sufficient conditions for global asymptotic stability of the equilibrium $x_{*}>0$
is negative, supercritical Neimark-Sacker bifurcations can occur. Thus, an invariant circle around the nontrivial equilibrium contributes to $\mathcal{A}$.

For bounded functions $h$ one obtains the following result on dissipativity.
Proposition 5.2 (dissipativity for (5.3), cf. [17]). If there exists a $K^{+} \geq 0$ such that $h(x) \leq K^{+}$for all $x \in \mathbb{R}_{+}$, then (5.3) is dissipative and $A=[0, R]$ is absorbing for every $R>\frac{K^{+}}{1-\alpha}$.

Example 5.3 (Mackey-Glass equation I). Let $p>0$. For the Mackey-Glass equation (5.3) one considers

$$
\begin{equation*}
h(x)=\frac{\beta}{1+x^{p}} \tag{5.4}
\end{equation*}
$$

and the injectivity condition (1.3) holds. Moreover, $S h(x)<0$ is true on the interval $(0, \infty)$. There exists a unique equilibrium $x_{*}>0$ and Proposition 5.2 guarantees the absorbing set specified in Tab. 5.1. Shifting $x_{*}$ into the origin yields a right-hand side

$$
f(x, y)=\alpha\left(x+x_{*}\right)+\frac{\beta}{1+\left(y+x_{*}\right)^{p}}-x_{*}=\alpha x+\frac{\beta}{1+\left(y+x_{*}\right)^{p}}-\frac{\beta}{1+x_{*}^{p}}
$$

which satisfies $\left(H_{1}\right)-\left(H_{3}\right)$ (negative feedback) with $a=\alpha, b=-\frac{\beta p x_{*}^{p-1}}{\left(1+x_{*}^{p}\right)^{2}}<0$ and a global attractor $\mathcal{A}$ consisting of entire solutions being uniquely defined in backward time.

Our next example requires some dissipativity criterion which applies even though $h$ is unbounded.

Proposition 5.4 (dissipativity for (5.3)). Let $\beta_{0}, \beta_{1} \geq 0$ be reals satisfying $0 \leq \beta_{1} \leq \alpha^{d}(1-\alpha)$. If $h: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$fulfills

$$
\begin{equation*}
h(x) \leq \beta_{0}+\beta_{1} x \quad \text { for all } x \in \mathbb{R}_{+}, \tag{5.5}
\end{equation*}
$$

then (5.3) is dissipative and $A=[0, R]$ is absorbing for every $R>\frac{\beta_{0}}{\alpha^{d}(1-\alpha)-\beta_{1}}$.

Proof. We equip $\mathbb{R}^{d+1}$ with the max-norm and formulate (5.3) as first order system (5.1) in $\mathbb{R}^{d+1}$. Then it is not difficult to see the relations

$$
\left\|A^{k}\right\| \leq \alpha^{-d} a^{k} \quad \text { for all } k \in \mathbb{Z}_{+}, \quad\|H(y)\| \leq \beta+\beta_{1}\|y\| \quad \text { for all } y \in \mathbb{R}_{+}^{d+1}
$$

with $a:=\alpha$ and Lemma 5.1 implies the assertion.
Example 5.5 (Mackey-Glass equation II). Let $p \in(0,1]$. A further difference equation attributed to Mackey and Glass is of the form (5.3) with

$$
\begin{equation*}
h(x)=\frac{\beta x}{1+x^{p}} \tag{5.6}
\end{equation*}
$$

and the injectivity condition (1.3) holds. On the one hand, due to [9, Theorem 2] the trivial solution is globally asymptotically stable for $\alpha+\beta<1$. On the other hand, for $\alpha+\beta>1$ there bifurcates a unique positive equilibrium $x_{*}$ given in Tab. 5.1.

- For $p<1$ choose $\beta_{1} \in\left(0, \alpha^{d}(1-\alpha)\right)$ and from $h^{\prime}(y)=\frac{\beta\left(1+(1-p) y^{p}\right)}{\left(1+y^{p}\right)^{2}} \rightarrow 0$ as $y \rightarrow \infty$ there is an $\eta>0$ such that $h^{\prime}(\eta)=\beta_{1}$. Hence, (5.5) holds with $\beta_{0}=h(\eta)-\beta_{1} \eta>0$ and Proposition 5.4 yields $R>\frac{h(\eta)-\beta_{1} \eta}{\alpha^{d}(1-\alpha)-\beta_{1}}$.
- For $p=1$ the bounded function $y \mapsto \frac{\beta y}{1+y}$ is strictly increasing to $\beta$, we derive from Proposition 5.2 that $R>\frac{\beta}{1-\alpha}$.

The resulting right-hand side $f(x, y)=\alpha x+\frac{\beta\left(y+x_{*}\right)}{1+\left(y+x_{*}\right)^{p}}-\frac{\beta x_{*}}{1+x_{*}^{p}}$ fulfills $\left(H_{1}\right)-$ $\left(H_{3}\right)$ (positive feedback) and $a=\alpha, b=\beta \frac{1+(1-p) x_{*}^{p}}{\left(1+x_{*}^{p}\right)^{2}}>0$. Whence, a global attractor $\mathcal{A}$ containing the equilibria $-\underline{x}_{*}, 0$ exists and consists of bounded entire solutions being unique in backward time.

Example 5.6 (Wazewska-Lasota equation). Let $W$ denote the Lambert-$W$-function. We investigate (5.3) with

$$
\begin{equation*}
h(x)=\beta e^{-x} \tag{5.7}
\end{equation*}
$$

Note that (1.3) holds, as well as a negative Schwarzian $S h$. Moreover, one observes a subcritical flip bifurcation in the equation $x_{k+1}=\alpha x_{k}+h\left(x_{k}\right)$ for critical parameters $\alpha$ such that $\frac{1+\alpha}{1-\alpha}=W\left(\frac{\beta}{1-\alpha}\right)$. Using Proposition 5.2 we obtain a concrete absorbing set from Tab. 5.1, which also lists the unique positive equilibrium $x_{*}>0$. The right-hand side

$$
f(x, y):=\alpha\left(x+x_{*}\right)+\beta e^{-\left(x_{k-d}+x_{*}\right)}-x_{*}=\alpha x+(1-\alpha) W\left(\frac{\beta}{1-\alpha}\right)\left(e^{-y}-1\right)
$$

fulfills $\left(H_{1}\right)-\left(H_{3}\right)$ (negative feedback) with $a=\alpha, b=-\beta e^{-x_{*}}<0$ and a global attractor $\mathcal{A}$ containing unique backward solutions.

| $F$ | $x_{*}$ | GAS |
| ---: | :---: | :---: |
| $(5.9)$ | $\frac{1-\alpha}{\beta}$ | $1-\alpha<\frac{3 d+4}{2(d+1)^{2}}$ |
| $(5.10)$ | $\beta$ | $\beta<\frac{3 d+4}{2(d+1)^{2}}$ |

Table 5.2: Hutchinson-type models (5.8): Sufficient conditions for global asymptotic stability of the equilibrium $x_{*}>0$ due to [8, Proposition 4]

### 5.1.2 Models of Hutchinson-type

Let $F: \mathbb{R}_{+} \rightarrow(0, \infty)$ be of class $C^{1}$. Delay-difference equations

$$
\begin{equation*}
x_{k+1}=x_{k} F\left(x_{k-d}\right) \tag{5.8}
\end{equation*}
$$

are said to be of Hutchinson type.
Example 5.7 (Pielou equation). The equation (5.8) with

$$
\begin{equation*}
F(x)=\frac{1}{\alpha+\beta x} \tag{5.9}
\end{equation*}
$$

is dissipative by [7, Theorem 3.1]. Its unique equilibrium $x_{*}>0$ is asymptotically stable for $\cos \frac{d \pi}{2 d+1}>\frac{1-\alpha}{2}$. A sufficient condition for global asymptotic stability is given in Tab. 5.2. Both above sufficient conditions for asymptotic stability fail for large delays, but when $\frac{d^{d}}{(d+1)^{d+1}}<1-\alpha$, nontrivial solutions oscillate about $x_{*}$ (see [6, pp. 68-69, Theorem 3.4.2(c)]). Furthermore, the shifted equation with

$$
f(x, y):=\frac{x+(\alpha-1) y}{1+\beta y}
$$

fulfills $\left(H_{1}\right)-\left(H_{3}\right)$ (negative feedback) with $a=1, b=\alpha-1<0$ and a global attractor $\mathcal{A}$ containing the equilibria $-\underline{x}_{*}, 0$.

Example 5.8 (Ricker equation). For delayed Ricker equations (5.8) with

$$
\begin{equation*}
F(x)=e^{\beta-x} \tag{5.10}
\end{equation*}
$$

dissipativity was shown in [22, Theorem 3.1] for a more general equation. A sufficient condition for global asymptotic stability of the nontrivial equilibrium is given in Tab. 5.2. Moreover, the shifted equation with

$$
\begin{equation*}
f(x, y):=(x+\beta) e^{-y}-\beta \tag{5.11}
\end{equation*}
$$

fulfills $\left(H_{1}\right)-\left(H_{3}\right)$ (negative feedback) with $a=1, b=-\beta<0$ and a global attractor $\mathcal{A} \subseteq\left[-x_{*}, R-x_{*}\right]^{d+1}$ containing the fixed points $-\underline{x}_{*}, 0$.

Example 5.9 (nonempty Morse sets). Let $d=m p$ for some $m, p \in \mathbb{N}$, $p \geq 2$. We utilize the simple observation that $\tilde{x}^{p}$ is a $p$-periodic solution of equation (1.1) if and only if it is a $p$-periodic solution of the corresponding undelayed equation $x_{k+1}=f\left(x_{k}, x_{k}\right)$.

It is well known that the undelayed Ricker equation $x_{k+1}=x_{k} e^{\beta-x_{k}}$ admits a 2 -periodic solution if and only if $\beta>2$, and that this solution oscillates about $x_{*}=\beta$ (see [20, Proposition 3]). Now, let us fix $\beta>2$. Then there exists $\left(\tilde{x}_{k}^{2}\right)_{k \in \mathbb{Z}}$, a 2-periodic solution of (5.10) for any even $d$. Thus the corresponding first order, $d+1$ dimensional equation (1.2) with (5.11) has a 2 -periodic solution $\left(\tilde{y}_{k}^{2}\right)_{k \in \mathbb{Z}}$, moreover, from the oscillation of $\left(\tilde{x}_{k}^{2}\right)_{k \in \mathbb{Z}}$ follows immediately that $V^{-}\left(\tilde{y}_{k}^{2}\right) \equiv d+1$ holds on $\mathbb{Z}$. In particular, the corresponding Morse set $\mathcal{M}_{d+1}$ contains a 2 -periodic orbit. Note that in this case $d+1=N_{-}^{*}$ holds, meaning that $\mathcal{M}_{N^{*}}$ is nontrivial.

Now assume that $\beta>1+\ln 9$ and let $d=3 m$ for some $m \in \mathbb{N}$. Then according to [20, Proposition 5] and the comment subsequent to it, there exists a 3-periodic solution $\left(\tilde{x}_{k}^{3}\right)_{k \in \mathbb{Z}}$ of $x_{k+1}=x_{k} e^{\beta-x_{k}}$ (by Šarkovs'kiŭ's Theorem [23], also $n$-periodic solutions exist for any $n \in \mathbb{N}$ ). It is easy to see that $\left(\tilde{x}_{k}^{3}\right)_{k \in \mathbb{Z}}$ oscillates about $x_{*}=\beta$. Then one readily obtains that the corresponding $d+1$ dimensional equation (1.2) with (5.11) has a 3 -periodic solution $\left(\tilde{y}_{k}^{3}\right)_{k \in \mathbb{Z}}$, for which $V^{-}\left(\tilde{y}_{k}^{3}\right) \equiv 2 m+1$ holds on $\mathbb{Z}$, and thus the corresponding Morse set $\mathcal{M}_{2 m+1}$ contains a 3-periodic orbit.

As a consequence, if $\beta>1+\ln 9$ and $d=6 n$ for some $n \in \mathbb{N}$, then the Morse set $\mathcal{M}_{6 n+1}=\mathcal{M}_{N_{-}^{*}}$ is nontrivial and $\mathcal{M}_{4 n+1}$ is nonempty for the $d+1$ dimensional Ricker map (1.2) with (5.11).

Note that although one can similarly find oscillatory $n$-periodic solutions for $n m+1$ dimensional Ricker maps with $\beta>1+\ln 9$, it is nontrivial to determine the number of sign changes they have on one period, and therefore it is not clear in which Morse set they are contained.

The above example is closely related to the topic of the first open problem presented in Section 6.

### 5.2 Discretizations

We consider the scalar retarded delay-differential equation

$$
\begin{equation*}
\dot{x}(t)=g(x(t), x(t-r)) \tag{5.12}
\end{equation*}
$$

with a continuously differentiable right-hand side $g: \mathbb{R} \times \mathbb{R} \rightarrow \mathbb{R}$ and delay $r>0$. Its forward Euler discretization of stepsize $\frac{r}{d}, d \in \mathbb{N}$, becomes

$$
\begin{equation*}
x_{k+1}=f\left(x_{k}, x_{k-d}\right), \quad f(\xi, \eta):=\xi+\frac{r}{d} g(\xi, \eta) \tag{5.13}
\end{equation*}
$$

and is of the form (1.1) with $J=\mathbb{R}$.

The Krisztin-Walther equation $\dot{x}(t)=-a x(t)+h(x(t-r))$ is a special case with a $C^{1}$-function $h: \mathbb{R} \rightarrow \mathbb{R}$ satisfying

$$
a \neq h^{\prime}(0), \quad\left\{\begin{array}{l}
h((0, \infty)) \subseteq(0, \infty) \text { and } h((-\infty, 0)) \subseteq(-\infty, 0), \text { or } \\
h((0, \infty)) \subseteq(-\infty, 0) \text { and } h((-\infty, 0)) \subseteq(0, \infty) .
\end{array}\right.
$$

In this situation, the discretization is of Clarke-type

$$
\begin{equation*}
x_{k+1}=\left(1-a \frac{r}{d}\right) x_{k}+\frac{r}{d} h\left(x_{k-d}\right), \tag{5.14}
\end{equation*}
$$

whenever $a \frac{r}{d}<1$ holds.
Proposition 5.10 (dissipativity, cf. [17]). If there exist $K^{-}, K^{+} \geq 0$ such that $-K^{-} \leq h(x) \leq K^{+}$for $x \in \mathbb{R}$, then (5.14) is dissipative and every $A=\left[R^{-}, R^{+}\right]$with $R^{+}>K^{+}$and $R^{-}<-K^{-}$is an absorbing set.

As a prototype for equations (5.14) we study
Example 5.11. Let $J=\mathbb{R}$ and $\alpha \in(0,1), \beta \in \mathbb{R} \backslash\{0\}$. Consider the delay-difference equation

$$
\begin{equation*}
x_{k+1}=\alpha x_{k}+\beta \tanh \left(x_{k-d}\right), \tag{5.15}
\end{equation*}
$$

which according to [9, Theorem 2] has a globally asymptotically stable trivial solution for $\alpha+|\beta|<1$. The injectivity condition (1.3) holds and using Proposition 5.10 we see that $[-R, R]$ is an absorbing set when $R>\frac{|\beta|}{1-\alpha}$ and the backward solutions on the global attractor of (5.15) are unique. One has positive feedback for $\beta>0$, negative feedback for $\beta<0$ and it is

$$
a=\alpha, \quad b=\beta .
$$

In case $\frac{\beta}{1-\alpha} \in(0,1]$ there exists the unique fixed point 0 , while for $\frac{\beta}{1-\alpha}>1$ a symmetric pair of fixed points $x_{-}<0<x_{+}$bifurcates. In the positive feedback case, note that shifting (5.15) into one of the nontrivial equilibria $x_{ \pm}$might yield another Morse decomposition.

## 6 Perspectives

We conclude this paper with raising some open questions and indicating possible future research directions:

- Which Morse sets are nonempty? First note that the Morse set $\mathcal{M}_{N_{ \pm}^{*}}$ is obviously nonempty, as the trivial solution is contained in it. In Example 5.9 we showed one possible method which can provide nonempty Morse sets different from $\mathcal{M}_{N_{ \pm}^{*}}$. However it seems a challenging problem whether it can be shown (maybe under further assumptions on (1.1)) that $\mathcal{M}_{n} \neq \emptyset$ holds for all $n<N_{ \pm}^{*}$, as it was proved by MalletParet in the continuous time case with negative feedback [12].
- Using a generalization of Huszár's Lemma 2.1 by Egerváry [2], there is hope that similar results can be shown for delay-difference equations

$$
x_{k+1}=f\left(x_{k-m}, x_{k-d}\right)
$$

with coprime $m+1, d+1$.

- It would be interesting to obtain a Morse decomposition in the general setting of [14] and in particular for tridiagonal equations.
- Finally, what can be said in the nonautonomous case?


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## Appendix A Proofs for positive feedback

Proof of Theorem 3.2. First note that there exists $n_{0} \in \mathbb{N}$ such that for $n>n_{0}, y^{n}(j)$ has the same sign as $y(j)$, where $j \in\{0, \ldots, d\}$ is an index for which $y(j) \neq 0$. Moreover, from $y \in \mathcal{R}^{+}$it follows that for $j \in\{1, \ldots, d-1\}$ with $y(j)=0$,

$$
\operatorname{sgn} y(j-1)=\operatorname{sgn} y^{n}(j-1)=-\operatorname{sgn} y(j+1)=-\operatorname{sgn} y^{n}(j+1) \neq 0
$$

and

$$
\operatorname{sc}(y(j-1), y(j), y(j+1))=\operatorname{sc}\left(y^{n}(j-1), y^{n}(j), y^{n}(j+1)\right)=1
$$

holds for all $n>n_{0}$. This in turn yields that for $n>n_{0}$ and $0 \leq i<j \leq d$ such that $y(i) \neq 0 \neq y(j)$ one has

$$
\begin{equation*}
\operatorname{sc}(y(i), y(i+1), \ldots, y(j))=\operatorname{sc}\left(y^{n}(i), y^{n}(i+1), \ldots, y^{n}(j)\right) \tag{A.1}
\end{equation*}
$$

which in particular means that sc $y^{n}=\operatorname{sc} y$ and thus also $V^{+}(y)=V^{+}\left(y^{n}\right)$ hold whenever $y(0) \neq 0 \neq y(d)$.

Since $y \in \mathcal{R}^{+}$excludes the possibility of $y(0)$ and $y(d)$ both being 0 at once, there remains the case when exactly one of them is 0 . We may assume $y(d)=0$ (the case $y(0)=0$ is analogous). Again, by regularity of $y$ we have that $y(0) y(d-1)<0$, which together with (A.1) yields that on the one hand

$$
\mathrm{sc}(y(0), y(1), \ldots, y(d-1))=\operatorname{sc}\left(y^{n}(0), y^{n}(1), \ldots, y^{n}(d-1)\right)
$$

holds for $n>n_{0}$ and on the other hand this number is odd, say $2 m+1$. It is now clear that sc $y=2 m+1 \leq \operatorname{sc} y^{n} \leq 2 m+2$ and therefore $V^{+}(y)=V^{+}\left(y^{n}\right)$ holds for all $n>n_{0}$, which completes the proof.

Proof of Theorem 3.3. (a) Since the last $d$ components of $y_{k} \in \mathbb{R}_{*}^{d+1}$ are the same as the first $d$ coordinates of $y_{k+1}$, clearly sc $y_{k+1} \leq \operatorname{sc} y_{k}+1$ holds. Assuming now to the contrary that $V^{+}\left(y_{k+1}\right)>V^{+}\left(y_{k}\right)$ holds for some $k$ implies that sc $y_{k+1}=\operatorname{sc} y_{k}+1$ and sc $y_{k}$ is even. The former combined with $a_{k}>0$ yields $y_{k}(0) \neq 0$ and the evenness of sc $y_{k}$ implies that the last nonzero coordinate of $y_{k}\left(\right.$ say $\left.y_{k}(j), 0 \leq j \leq d\right)$ has the same sign as $y_{k}(0)$. We may assume w.l.o.g. that they are both positive (the negative case is analogous). In particular this implies $y_{k}(d) \geq 0$. Using now the positivity of both $a_{k}$ and $b_{k}$ one obtains readily from (3.3) that $y_{k+1}(d)=a_{k} y_{k}(d)+b_{k} y_{k}(0)$ is also positive. All these together give

$$
\begin{aligned}
\operatorname{sc} y_{k+1} & =\operatorname{sc}\left(y_{k+1}(0), y_{k+1}(1), \ldots, y_{k+1}(d-1), y_{k+1}(d)\right) \\
& =\operatorname{sc}\left(y_{k}(1), y_{k}(2), \ldots, y_{k}(d), y_{k+1}(d)\right) \leq \operatorname{sc} y_{k},
\end{aligned}
$$

where the last inequality is due to $\operatorname{sgn} y_{k}(j)=\operatorname{sgn} y_{k+1}(d)$. This is a contradiction, which proves statement (a).
(b) Before proving the statement, we need to introduce some auxiliary notions. Let us say that a component $y(j)(0 \leq j \leq d)$ of a vector $y \in \mathbb{R}^{d+1}$ is irregular if $y(j)=0$, and $y(j-1) y(j+1) \geq 0$, where $y(-1):=y(d)$ and $y(d+1):=y(0)$. Using this terminology, $y \in \mathcal{R}^{+}$(i.e. $y$ is regular) holds if and only if $y$ has no irregular component. Furthermore, we call a vector $(y(i), \ldots, y(j))$ with $0 \leq i \leq j \leq d$ an irregular block (of zeros) in $y \in \mathbb{R}^{d+1}$, if $y(i)=\cdots=y(j)=0$, and moreover it has maximal length in the sense that either $i=0$ or $i \geq 1$ and $y(i-1) \neq 0$ hold, and similarly either $j=d$ or $j<d$ and $y(j+1) \neq 0$ hold. The dimension of the block will be regarded as the length of it. Note also that consecutive zero components are irregular by definition.

The proof of statement (b) consists of several, yet elementary, steps.
From now on we shall always assume that $V^{+}\left(y_{k}\right)=V^{+}\left(y_{k+4 d+2}\right)$, which implies, in the light of statement (a), that $V^{+}\left(y_{k}\right)=V^{+}\left(y_{k+\ell}\right)$ holds for all $0 \leq \ell \leq 4 d+2$.

Step 1. If for some $k, y_{k+1}(j)=0$ is irregular $(0 \leq j<d)$, then so is $y_{k}(j+1)=0$ irregular.

This is trivial in case $1 \leq j \leq d-2$.
If $y_{k+1}(d-1)=0$ is irregular (i.e. $y_{k+1}(d-1)=0$ and $y_{k+1}(d-$ 2) $\left.y_{k+1}(d) \geq 0\right)$, then on the one hand $y_{k}(d)=0$ holds and on the other hand, $b_{k}>0$ yields that $\operatorname{sgn} y_{k+1}(d)=\operatorname{sgn} y_{k}(0)$. Combining this with $y_{k}(d-1)=$ $y_{k+1}(d-2)$ yields that $\operatorname{sgn}\left(y_{k}(d-1) y_{k}(0)\right)=\operatorname{sgn}\left(y_{k+1}(d-2) y_{k+1}(d)\right)$, meaning that $y_{k}(d)$ is also irregular in $y_{k}$.

For the case when $y_{k+1}(0)=0$ is irregular, assume to the contrary that $y_{k}(1)=0$ is regular, i.e. $y_{k}(0) y_{k}(2)<0$. W.l.o.g. we may assume $y_{k}(0)<$ $0<y_{k}(2)$. Irregularity of $y_{k+1}(0)$ combined with $y_{k+1}(1)=y_{k}(2)>0$ yields $y_{k+1}(d) \geq 0$. Now, from $a_{k}>0$ and $b_{k}>0$ it follows that $y_{k}(d)>0$. These
all together imply that sc $y_{k}$ is an odd number, moreover sc $y_{k}=\operatorname{sc} y_{k+1}+1$, which gives $V^{+}\left(y_{k}\right)=V^{+}\left(y_{k+1}\right)+2$, a contradiction.
Step 2. If $y_{k+1}(d)=0$, then there are two possibilities. Either $y_{k}(0)=$ $y_{k}(d)=0$, or $y_{k}(0) y_{k}(d)<0$. Assume now that the latter holds and let $0 \leq j<d$ be the smallest integer such that $y_{k}(j+1) \neq 0$. As $y_{k}(0) y_{k}(d)<0$, hence sc $y_{k}$ is odd, so $V^{+}\left(y_{k}\right)=V^{+}\left(y_{k+1}\right)$ holds only if $\operatorname{sgn} y_{k}(j+1)=$ $\operatorname{sgn} y_{k}(0)$, which in turn yields also

$$
\begin{equation*}
\operatorname{sgn} y_{k}(j+1)=-\operatorname{sgn} y_{k}(d)=-\operatorname{sgn} y_{k+1}(d-1) \neq 0 \tag{A.2}
\end{equation*}
$$

If in addition $y_{k+1}(d)=0$ is irregular, then $0 \leq y_{k+1}(0) y_{k+1}(d-1)=$ $y_{k}(1) y_{k}(d)$ holds, which in the light of (A.2) yields that $j \neq 0$, or what is equivalent, $y_{k}(1)=0$.

To sum up, if $y_{k+1}(d)=0$ is irregular, then either $y_{k}(0)=y_{k}(d)=0$ or there exists some $1 \leq j \leq d-2$, such that

$$
\begin{align*}
y_{k}(0) & \neq 0 \\
y_{k}(1) & =\cdots=y_{k}(j)=0, \\
y_{k}(j+1) & \neq 0  \tag{A.3}\\
y_{k}(d) & \neq 0 \\
\operatorname{sgn} y_{k}(0) & =\operatorname{sgn} y_{k}(j+1)=-\operatorname{sgn} y_{k}(d)
\end{align*}
$$

hold, which in particular means that $y_{k}(0)$ is followed by an irregular block of zeros of length $j \geq 1$.
Step 3. As a result of the first two steps, if $y_{k} \in \mathcal{R}^{+}$and $V^{+}\left(y_{k}\right)=$ $V^{+}\left(y_{k+1}\right)$, then $y_{k+1} \in \mathcal{R}^{+}$. Therefore, in order to prove the statement, it is sufficient to show that there exists $0 \leq \ell \leq 4 d+2$ such that $y_{k+\ell} \in \mathcal{R}^{+}$.

Steps 4-6 are identical to the corresponding ones presented for negative feedback and the same argument shows that $y_{k+4 d+2} \in \mathcal{R}^{+}$.

Proof of Theorem 4.2. The eigenvalues of (4.1) are precisely the solutions of the characteristic equation (2.2). If $M^{*}$ is equal to 0 or $d+1$, then the statement is trivial. Otherwise, according to Lemma 2.1, $M^{*}$ is an odd number, say $M^{*}=2 n^{\prime}+1$, where $0 \leq n^{\prime} \leq \frac{d-1}{2}$.

We use the notations and ordering for eigenvalues introduced in Lemma 2.1. For a non-real eigenvalue $\lambda_{j}=r_{j} e^{i \varphi_{j}} \in S_{j}^{+}$(with $\varphi_{j} \in(0, \pi)$, $j \in \mathbb{N}, j<\frac{d+1}{2}$ ) any linear combination of the real eigensolutions corresponding to $\lambda_{j}$ can be written in the form $c_{j} z_{j}$, where $c_{j} \in \mathbb{R}$ and

$$
z_{j, k}=\left(\begin{array}{c}
z_{j, k}(0)  \tag{A.4}\\
z_{j, k}(1) \\
\vdots \\
z_{j, k}(d-1) \\
z_{j, k}(d)
\end{array}\right)=\left(\begin{array}{c}
r_{j}^{k} \sin \left(k \varphi_{j}+\omega_{j}\right) \\
r_{j}^{k+1} \sin \left((k+1) \varphi_{j}+\omega_{j}\right) \\
\vdots \\
r_{j}^{k+d-1} \sin \left((k+d-1) \varphi_{j}+\omega_{j}\right) \\
r_{j}^{k+d} \sin \left((k+d) \varphi_{j}+\omega_{j}\right)
\end{array}\right)
$$

with some $\omega_{j} \in[0,2 \pi)$. From $\varphi_{j} \in(0, \pi)$ it is clear that sc $z_{j, k}$ equals the number of sign changes the sine function has on the interval $\left(k \varphi_{j}+\omega_{j},(k+\right.$ d) $\left.\varphi_{j}+\omega_{j}\right)$. According to the definition of the sector $S_{j}^{+}$, the length of this interval is $d \varphi_{j} \in\left((2 j-1) \pi, \frac{d}{d+1} 2 j \pi\right)$. Now, from the definition of $V^{+}$we obtain that $V^{+}\left(z_{j, k}\right)=2 j$ for all $k \in \mathbb{Z}_{-}$.

As for the eigensolutions corresponding to real eigenvalues, any eigensolution corresponding to the leading positive eigenvalue $r_{0}:=\lambda_{+}>0$ can be written in the form $c_{0} z_{0}$, where $c_{0} \in \mathbb{R}$ and $z_{0, k}(j)=r_{0}^{k+j}$ for all $k \in \mathbb{Z}_{-}$ and $j \in \mathbb{Z}_{+}, j \leq d$. Clearly, every component of $z_{0, k}$ has the same sign for any $k \in \mathbb{Z}_{-}$and in particular $V^{+}\left(z_{0, k}\right)=0$ for all $k \in \mathbb{Z}_{-}$. Similarly, if $d$ is odd, say $d=2 \ell-1$ for some $\ell \in \mathbb{N}$, then any eigensolution corresponding to the unique negative eigenvalue $-r_{\ell}:=\lambda_{-}<0$ can be written in the form $c_{\ell} z_{\ell}$, where $c_{\ell} \in \mathbb{R}$ and $z_{\ell, k}(j)=\left(-r_{\ell}\right)^{k+j}$ for all $k \in \mathbb{Z}_{-}$and $j \in \mathbb{Z}_{+}, j \leq d$. Obviously sc $z_{\ell, k}=d$ and therefore also $V^{+}\left(z_{\ell, k}\right)=d+1=2 \ell$ hold for all $k \in \mathbb{Z}_{-}$. According to Lemma 2.1 there are no other real eigenvalues.

First, let us consider the hyperbolic case. Assume that $M^{*}=2 n^{\prime}+1$ for some integer $0 \leq n^{\prime} \leq \frac{d-1}{2}$. Then $N_{+}^{*}=M^{*}=2 n^{\prime}+1$. Let $y$ be an arbitrary solution of (4.2) which is bounded in backward time. This means that there exist an integer $0 \leq n \leq n^{\prime}$ and appropriate constants $c_{0}, c_{j} \in \mathbb{R}$ and $\omega_{j} \in[0,2 \pi)$ for $j=1, \ldots, n$ with $c_{n} \neq 0$, such that

$$
y=\sum_{j=0}^{n} c_{j} z_{j} .
$$

We will show that $V^{+}\left(y_{k}\right) \leq V^{+}\left(z_{n, k}\right)=2 n<2 n^{\prime}+1=N_{+}^{*}$ holds for all $k \leq k_{1}$, where $k_{1}$ is to be defined soon. Then by monotonicity of $V^{+}$ (proved in Theorem 3.3) we get statement (a).

Case $n=0$ is trivial, thus we may assume $n \geq 1$. From Lemma 2.1 one can easily obtain that there exists $\vartheta_{n}>0$, sufficiently small, such that

$$
\begin{align*}
\vartheta_{n} & <\min \left\{\frac{\pi-\varphi_{n}}{2}, \frac{\varphi_{n}}{2}\right\}  \tag{A.5}\\
d \varphi_{n}+2 \vartheta_{n} & <2 n \pi \text { and }  \tag{A.6}\\
(2 n-1) \pi & <d \varphi_{n}-2 \vartheta_{n} \tag{A.7}
\end{align*}
$$

hold. Moreover, by the same lemma we have $\left|r_{n}\right|<\left|r_{j}\right|$ for any $0 \leq j<n$, so we may choose $k_{0}<0$ so small that

$$
\begin{equation*}
n r_{n-1}^{k} \max _{0 \leq j \leq n-1}\left|c_{j}\right|<\left|c_{n}\right| r_{n}^{k} \sin \vartheta_{n} \tag{A.8}
\end{equation*}
$$

holds for all $k \leq k_{0}$. Let $k_{1}=k_{0}-d$.
For simplicity, let us say that the $\ell$-th component of $z_{n, k}$, i.e. $z_{n, k}(\ell)=$ $r_{n}^{k} \sin \left((k+\ell) \varphi_{n}+\omega_{n}\right)$ is small if $\left|\sin \left((k+\ell) \varphi_{n}+\omega_{n}\right)\right| \leq \sin \vartheta_{n}$, otherwise the component will be said to be $b i g$. Inequality (A.8) yields that if $z_{n, k}(\ell)$
is big for some $\ell=0,1, \ldots, d$, then $\operatorname{sgn} y_{k}(\ell)=\operatorname{sgn} z_{n, k}(\ell) \neq 0$ holds for all $k \leq k_{1}$. Note that inequality (A.5) guarantees that there is no $k \leq k_{1}$ and component $0 \leq \ell<d$ such that $z_{n, k}(\ell)$ and $z_{n, k}(\ell+1)$ are both small. Similarly, (A.6)-(A.7) imply that $z_{n, k}(0)$ and $z_{n, k}(d)$ cannot be small at the same time. Moreover, by (A.5) one obtains that if $z_{n, k}(\ell)$ is small for some $k \leq k_{1}$ and $\ell \in\{1, \ldots, d-1\}$, then $\operatorname{sgn} z_{n, k}(\ell-1)=-\operatorname{sgn} z_{n, k}(\ell+1)$.

All these together mean that, if $k \leq k_{1}$ and $z_{n, k}(0)$ and $z_{n, k}(d)$ are both big, sc $y_{k}=\operatorname{sc} z_{n, k}$ holds, and $V^{+}\left(y_{k}\right)=V^{+}\left(z_{n, k}\right)=2 n<N_{+}^{*}$ follows.

There remains the case, when exactly one of $z_{n, k}(0)$ and $z_{n, k}(d)$ is small. W.l.o.g. we assume that $z_{n, k}(d)$ is small, $z_{n, k}(0)$ is big and positive. Then, as $z_{n, k}(d-1)$ is also big, $\operatorname{sc}\left(y_{k}(0), \ldots, y_{k}(d-1)\right)=\operatorname{sc}\left(z_{n, k}(0), \ldots, z_{n, k}(d-1)\right)$ follows. Thus either $V^{+}\left(y_{k}\right) \leq V^{+}\left(z_{n, k}\right)$ holds or $\operatorname{sc} y_{k}=\operatorname{sc} z_{n, k}+1$ is odd. We show that the latter case leads to contradiction. Should it hold, then $y_{k}(0), y_{k}(d-1), z_{n, k}(0)$ and $z_{n, k}(d-1)$ would be all positive, while $y_{k}(d)<0$ and $z_{n, k}(d) \geq 0$. Now, since $z_{n, k}(d) \geq 0$ is small, $\left((k+d) \varphi_{n}+\omega_{n}\right) \in$ $\left[0, \vartheta_{n}\right) \cup\left(\pi-\vartheta_{n}, \pi\right] \bmod 2 \pi$ holds. On the other hand, from $z_{n, k}(d-1)>0$ and (A.5) it follows that $\left((k+d) \varphi_{n}+\omega_{n}\right) \in\left(\pi-\vartheta_{n}, \pi\right] \bmod 2 \pi$, which implies by inequalities (A.6) and (A.7) that $z_{n, k}(0)<0$. This is a contradiction, so our proof is complete for the hyperbolic case.

If the origin is non-hyperbolic, and $M^{*} \geq 1$, then $N_{+}^{*}=M^{*}+1$. Keeping our notations from above, there is a pair of eigenvalues $\lambda_{n^{\prime}+1}, \overline{\lambda_{n^{\prime}+1}}$ on the unit circle, or -1 is a simple eigenvalue (if $d$ is odd and $M^{*}=d$ ). In the latter case, statement (a) of the lemma is trivial. In the former case all solutions of (4.2) that are bounded in backward time can be written in the form $y=\sum_{j=0}^{n} c_{j} z_{j}$, where $n \leq n^{\prime}+1, c_{j} \in \mathbb{R}, c_{n} \neq 0$, and $z_{j}$ are of the form (A.4). Now, the same argument presented in the hyperbolic case can be applied to prove that $V^{+}\left(y_{k}\right)=2 n \leq 2 n^{\prime}+2=N_{+}^{*}$ for all $k \in \mathbb{Z}_{-}$.

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