An Intemational Joumal computers \& mathematics with applications

# Normal Forms for Nonautonomous Difference Equations 

S. Siegmund<br>Department of Mathematics<br>University of California<br>Berkeley, CA 94720, U.S.A.<br>siegmund@math.berkeley.edu


#### Abstract

We extend Henry Poincare's normal form theory for autonomous difference equations $x_{k+1}=f\left(x_{k}\right)$ to nonautonomous difference equations $x_{k+1}=f_{k}\left(x_{k}\right)$. Poincaré's nonresonance condition $\lambda_{j}-\prod_{i=1}^{n} \lambda_{i}^{q_{i}} \neq 0$ for eigenvalues is generalized to the new nonresonance condition $\lambda_{j} \cap$ $\prod_{i=1}^{n} \lambda_{i}^{q_{i}}=\emptyset$ for spectral intervals. © 2003 Elsevier Science Ltd. All rights reserved.


Keywords-Poincaré normal form theory, Nonautonomous normal forms, Resonance, Nonautonomous difference equations.

## 1. INTRODUCTION

The famous French mathematician Henry Poincare founded the normal form theory for autonomous differential equations $\dot{x}=f(x)$ near a rest point in his thesis in 1879. Soon a parallel theory for autonomous difference equations $x_{k+1}=f\left(x_{k}\right)$ was developed. If the eigenvalues $\lambda_{1}, \ldots, \lambda_{n}$ of the linearization $x_{k+1}=D f\left(x^{0}\right) x_{k}$ at the rest point $x^{0}$ satisfy the nonresonance condition

$$
\begin{equation*}
\lambda_{j} \neq \lambda_{1}^{q_{1}} \cdots \lambda_{n}^{q_{n}}, \tag{1}
\end{equation*}
$$

$j \in\{1, \ldots, n\}, q_{i} \in \mathbb{N}_{0}=\{0,1, \ldots\}, \sum_{i=1}^{n} q_{i} \geq 2$, then the difference equation can be formally linearized.

As an example, we consider the following planar autonomous system:

$$
\begin{aligned}
x_{k+1} & =2 x_{k}, \\
y_{k+1} & =\lambda y_{k}+x_{k}^{2},
\end{aligned}
$$

with $\lambda \in(0, \infty)$. We are looking for a near-identity transformation

$$
H(x, y)=\binom{x}{y}+h_{2}(x, y)
$$

which eliminates the second order nonlinearity $\binom{0}{x^{2}}$ and we choose $h_{2} \in \operatorname{span}\left\{\binom{x^{2}}{0},\binom{0}{x^{2}},\binom{x y}{0},\binom{0}{x y}\right.$, $\left.\binom{y^{2}}{0},\binom{0}{y^{2}}\right\}$. It is not difficult to show that the transformed equation has no second-order nonlinearity if and only if the so-called homological equation

$$
A h_{2}(x, y)-h_{2}\left(A\binom{x}{y}\right)=f_{2}(x, y)
$$

[^0]is satisfied with
\[

A=\left($$
\begin{array}{ll}
2 & 0 \\
0 & \lambda
\end{array}
$$\right) \quad and \quad f_{2}(x, y)=\binom{0}{x^{2}} .
\]

It is solvable if and only if $\lambda \neq 4$ and with its unique solution we get

$$
H(x, y)=\binom{x}{y+\frac{1}{\lambda-4} x^{2}} .
$$

In this simple example, the transformed equation $x_{k+1}=2 x_{k}, y_{k+1}=\lambda y_{k}$ is linear. In general, the elimination of second-order nonlinearities produces higher-order nonlinearities, and the process has to be iterated. The resulting transformation is the composition of the transformations of each elimination step, and it is nonlinear but is constructed by solving a sequence of linear equations.
In this article, we consider nonautonomous invertible difference equations

$$
\begin{equation*}
x_{k+1}=f_{k}\left(x_{k}\right) \tag{2}
\end{equation*}
$$

not in the vicinity of a rest point as Poincare did it in the autonomous case, but in the vicinity of an arbitrary reference solution $v^{0}: \mathbb{Z} \rightarrow \mathbb{R}^{N}$. For some $p \geq 2$ we assume $f_{k}: D_{f_{k}} \subset$ $\mathbb{R}^{N} \rightarrow f_{k}\left(D_{f_{k}}\right) \subset \mathbb{R}^{N}$ to be a $C^{p}$ diffeomorphism for every $k \in \mathbb{Z}=\{0, \pm 1, \ldots\}$. We will extend Poincaré's normal form theory by showing that if the linearization $x_{k+1}=D f_{k}\left(v_{k}^{0}\right) x_{k}$ of (2) along the reference solution $v^{0}$ has invertible coefficient matrices $D f_{k}\left(v_{k}^{0}\right) \in \mathbb{R}^{N \times N}, k \in \mathbb{Z}$, and satisfies a nonresonance condition, then system (2) is locally $C^{p}$ equivalent to a system $x_{k+1}=g_{k}\left(x_{k}\right)$ in normal form; i.e., with zero reference solution, block diagonal linear part $x_{k+1}=D g_{k}(0) x_{k}$ and all nonresonant Taylor coefficients of $g$ up to order $p$ are zero.

We therefore have to use a proper replacement of the "linear algebra" for autonomous systems (i.e., eigenvalues and eigenspaces) in our nonautonomous situation. A spectral theory for nonautonomous difference equations is developed in [1]. The dichotomy spectrum of the linearized difference equation $x_{k+1}=D f_{k}\left(v_{k}^{0}\right) x_{k}$ consists of at most $N$ closed intervals of the positive real line $\mathbb{R}^{+}=(0, \infty)$; in general, the spectrum may be empty or unbounded. It is nonempty and compact, i.e., consists of $n$ compact intervals with $1 \leq n \leq N$, if the system has bounded growth. A linear system $x_{k+1}=A_{k} x_{k}$ has bounded growth if its evolution operator $\Phi$ satisfies the estimate $\|\Phi(k, \ell)\| \leq K a^{|k-\ell|}$ for $k, \ell \in \mathbb{Z}$ with constants $K, a \geq 1$. Bounded growth is equivalent to the boundedness of the coefficients and their inverses [1, Lemma 2.3], and hence, the linearized difference equation $x_{k+1}=D f_{k}\left(v_{k}^{0}\right) x_{k}$ has bounded growth if $\left\|A_{k}\right\| \leq M$ and $\left\|A_{k}^{-1}\right\| \leq M$ for $k \in \mathbb{Z}$ with some constant $M \geq 0$ and $A_{k}=D f_{k}\left(v_{k}^{0}\right)$.

For simplicity, we assume in the following that the linearized equation has bounded growth, although the theory could also be developed in the general case.

## 2. PRELIMINARIES

Let $\phi(\cdot ; \ell, \xi): I_{\ell, \xi} \rightarrow \mathbb{R}^{N}$ denote the unique maximal solution of the initial value problem (2), $x_{\ell}=\xi$ for $\xi \in D_{f_{\ell}}$ where $I_{\ell, \xi}$ is a $\mathbb{Z}$-interval (i.e., an intersection of an interval with $\mathbb{Z}$ ) containing $\ell$ such that the solution identity

$$
\phi(k+1 ; \ell, \xi)=f_{k}(\phi(k ; \ell, \xi)), \quad \text { for } k, k+1 \in I_{\ell, \xi},
$$

holds. We have $\phi(\ell ; \ell, \xi)=\xi$ and $\phi(k ; \phi(m ; \ell, \xi))=\phi(k+m ; \ell, \xi)$ for $m, k+m \in I_{\ell, \xi}$.
There is no straightforward way to define a notion of conjugacy for nonautonomous difference equations. What do we mean by this? Two autonomous difference equations $x_{k+1}=f^{1}\left(x_{k}\right)$ and $x_{k+1}=f^{2}\left(x_{k}\right)$ in $\mathbb{R}^{N}$ are said to be conjugate if there exists a homcomorphism $H: \mathbb{R}^{N} \rightarrow \mathbb{R}^{N}$ such that the flows $\phi_{1}(\cdot ; \xi)$, respectively, $\phi_{2}(\cdot ; \eta)$ satisfy the conjugacy relation $H\left(\phi_{1}(k ; \xi)\right)=$ $\phi_{2}(k ; H(\xi))$ for all $\xi \in \mathbb{R}^{N}, k \in I_{\xi}$; i.e., $H$ maps solutions of the first equation onto solutions of
the second equation and vice versa for $H^{-1}$. Now if we would define a conjugacy between two nonautonomous difference equations $x_{k+1}=f_{k}^{1}\left(x_{k}\right)$ and $x_{k+1}=f_{k}^{2}\left(x_{k}\right)$ by the same property, but now with a $k$-dependent $H$, then for every $\ell \in I_{k, x}$

$$
H_{k}(x):=\phi_{2}\left(k ; \ell, \phi_{1}(\ell ; k, x)\right)
$$

would establish a conjugacy; i.e., $H$ maps solutions of the first equation onto solutions of the second equation and vice versa with

$$
H_{k}^{-1}(x):=\phi_{1}\left(k ; \ell, \phi_{2}(\ell ; k, x)\right) .
$$

In the nonautonomous situation, we need, therefore, some additional conditions which ensure that qualitative behaviour-at least for a single reference solution-is preserved under the transformation.

It is easy to see in the autonomous situation that for a conjugacy, periodic solutions, limit sets, and invariant sets of the first equation are bijectively mapped onto periodic solutions, limit sets, and invariant sets, respectively, of the second equation and that (asymptotic) stability, attractivity, and instability of bounded solutions are preserved under the conjugacy. In most cases this is enough, but note that the assumption of boundedness of solutions is essential for the preservation of stability. For example: the two linear systems $x_{k+1}=x_{k}+1, y_{k+1}=(1 / 2) y_{k}$ and $x_{k+1}=x_{k}+1, y_{k+1}=2 y_{k}$ are conjugate via $H(x, y)=\left(x, y 2^{2 x}\right)$, but the first system is stable and the second is unstable. To preserve the stability of an unbounded solution $v_{k}$, it would be necessary to pose some uniformity condition on $H$, e.g., $\lim _{x \rightarrow 0} H\left(v_{k}+x\right)=H\left(v_{k}\right)$ uniformly in $k \in \mathbb{Z}$. Such a uniformity condition is exactly what we need in the nonautonomous situation to define a meaningful notion of $C^{p}$ equivalence.

Consider difference equations together with reference solutions

$$
\begin{array}{ll}
x_{k+1}=f_{k}\left(x_{k}\right), & v^{0}: \mathbb{Z} \rightarrow \mathbb{R}^{N}, \\
x_{k+1}=g_{k}\left(x_{k}\right), & w^{0}: \mathbb{Z} \rightarrow \mathbb{R}^{N}, \tag{4}
\end{array}
$$

where $f_{k}$ and $g_{k}$ are $C^{p}$ diffeomorphisms, i.e., $f_{k} \in \operatorname{Diff}^{p}\left(D_{f_{k}}, f_{k}\left(D_{f_{k}}\right)\right), g_{k} \in \operatorname{Diff}^{p}\left(D_{g_{k}}, g_{k}\left(D_{g_{k}}\right)\right)$, $p \geq 0$. We assume that uniform neighbourhoods of the reference solutions are contained in the corresponding sets of definition, i.e., there exist $\varepsilon>0$ and $\delta>0$ such that

$$
B_{\varepsilon}\left(v_{k}^{0}\right) \subset D_{f_{k}} \quad \text { and } \quad B_{\delta}\left(w_{k}^{0}\right) \subset D_{g_{k}}, \quad \text { for } k \in \mathbb{Z}
$$

where $B_{\varepsilon}\left(x^{0}\right):=\left\{x \in \mathbb{R}^{N}:\left\|x-x^{0}\right\|<\varepsilon\right\}$. Define $U_{\varepsilon}\left(v^{0}\right):=\left\{(k, x) \in \mathbb{Z} \times \mathbb{R}^{N}: x \in B_{\varepsilon}\left(v_{k}^{0}\right)\right\}$.
Definition 1. Consider systems (3) and (4). If there exist $\varepsilon^{\prime}, \delta^{\prime}$ with $0<\varepsilon^{\prime} \leq \varepsilon$ and $0<\delta^{\prime} \leq \delta$ together with functions

$$
H: U_{\varepsilon^{\prime}}\left(v^{0}\right) \rightarrow \mathbb{R}^{N}, \quad H^{-1}: U_{\delta^{\prime}}\left(w^{0}\right) \rightarrow \mathbb{R}^{N},
$$

then $H$ is called a local $C^{p}$ equivalence between system (3) with solution $v^{0}$ and system (4) with solution $w^{0}$, if the following statements are valid.
(A) For each $k \in \mathbb{Z}$, the mappings

$$
\begin{gathered}
H_{k}: B_{\varepsilon^{\prime}}\left(v_{k}^{0}\right) \rightarrow H_{k}\left(B_{\varepsilon^{\prime}}\left(v_{k}^{0}\right)\right) \subset B_{\delta}\left(w_{k}^{0}\right), \\
H_{k}^{-1}: B_{\delta^{\prime}}\left(w_{k}^{0}\right) \rightarrow H_{k}^{-1}\left(B_{\delta^{\prime}}\left(w_{k}^{0}\right)\right) \subset B_{\varepsilon}\left(v_{k}^{0}\right)
\end{gathered}
$$

are $C^{p}$ diffeomorphisms (or homeomorphisms if $p=0$ ) with

$$
H_{k}\left(H_{k}^{-1}(x)\right)=x \quad \text { and } \quad H_{k}^{-1}\left(H_{k}(x)\right)=x
$$

for all $x$ for which the compositions are defined.
(B) If $v_{k}$ is a solution of (3) in $U_{\varepsilon^{\prime}}\left(v^{0}\right)$, then $H_{k}\left(v_{k}\right)$ is a solution of (4). If $w_{k}$ is a solution of (4) in $U_{\delta^{\prime}}\left(w^{0}\right)$, then $H_{k}^{-1}\left(w_{k}\right)$ is a solution of (3).
(C) The reference solutions are mapped uniformly onto each other,

$$
\begin{aligned}
\lim _{x \rightarrow 0} H_{k}\left(v_{k}^{0}+x\right) & =w_{k}^{0}, & & \text { uniformly in } k \in \mathbb{Z} \\
\lim _{x \rightarrow 0} H_{k}^{-1}\left(w_{k}^{0}+x\right) & =v_{k}^{0}, & & \text { uniformly in } k \in \mathbb{Z} .
\end{aligned}
$$

Lemma 2. Consider systems (3) and (4) together with a solution $v^{0}: I \rightarrow \mathbb{R}^{N}$ of (3) which is defined on some $\mathbb{Z}$-interval $I$. Then a mapping $w: J \rightarrow \mathbb{R}^{N}$ defined on a $\mathbb{Z}$-interval $J \subset I$ is a solution of the difference equation

$$
\begin{equation*}
x_{k+1}=g_{k}\left(x_{k}+v_{k}^{0}\right)-f_{k}\left(v_{k}^{0}\right) \tag{5}
\end{equation*}
$$

if and only if $w+v^{0}: J \rightarrow \mathbb{R}^{N}$ is a solution of the difference equation (4).
Proof. Since $v^{0}$ is a solution of (3) one has for $k, k+1 \in J$

$$
w_{k+1}=g_{k}\left(w_{k}+v_{k}^{0}\right)-f_{k}\left(v_{k}^{0}\right) \Leftrightarrow w_{k+1}+v_{k+1}^{0}=g_{k}\left(w_{\dot{k}}+v_{k}^{0}\right),
$$

and the claim is proved.

## 3. NORMAL FORMS

We consider a difference equation together with a reference solution

$$
\begin{equation*}
x_{k+1}=f_{k}\left(x_{k}\right), \quad v^{0}: \mathbb{Z} \rightarrow \mathbb{R}^{N} \tag{6}
\end{equation*}
$$

which satisfy the following conditions.

- Smoothness: $f_{k} \in \operatorname{Diff}^{p}\left(D_{f_{k}}, f_{k}\left(D_{f_{k}}\right)\right)$ for a $p \geq 2$.
- Set of definition: $\varepsilon:=\inf \left\{\varepsilon^{\prime} \geq 0: B_{\varepsilon^{\prime}}\left(v_{k}^{0}\right) \subset D_{f_{k}}\right.$ for all $\left.k \in \mathbb{Z}\right\}>0$.
- Linearity: $x_{k+1}=D f_{k}\left(v_{k}^{0}\right) x_{k}$ has bounded growth.
- Nonlinearity: $\left\|D^{j} f_{k}\left(v_{k}^{0}\right)\right\| \leq M$ for $2 \leq j \leq p$ and all $k \in \mathbb{Z}$.

We will simplify system (6) in three steps.

## Step 1: Trivialization of the Reference Solution.

Recall Lemma 2. If $f \equiv g$, then system (5) reduces to

$$
\begin{equation*}
x_{k+1}=f_{k}\left(x_{k}+v_{k}^{0}\right)-f_{k}\left(v_{k}^{0}\right), \tag{7}
\end{equation*}
$$

which is usually called the difference equation of perturbed motion of (6) w.r.t. the solution $v^{0}$. Obviously, (7) has the zero solution, and because of Lemma 2, w:J‘I $\rightarrow \mathbb{R}^{N}$ is a solution of (7) if and only if $w+v^{0}$ is a solution of (6), and hence, the mappings

$$
\begin{aligned}
R_{k}: B_{\varepsilon}\left(v_{k}^{0}\right) \rightarrow \mathbb{R}^{N}, & & x \mapsto x-v_{k}^{0}, \\
R_{k}^{-1}: B_{\varepsilon}(0) \rightarrow \mathbb{R}^{N}, & & x \mapsto x+v_{k}^{0},
\end{aligned}
$$

define a $C^{\infty}$ equivalence between (6) with reference solution $v^{0}$ and system (7) with zero reference solution. We rewrite (7) as

$$
\begin{equation*}
x_{k+1}=A_{k}^{*} x_{k}+F_{k}^{*}\left(x_{k}\right) \tag{8}
\end{equation*}
$$

where $A_{k}^{*}=D f_{k}\left(v_{k}^{0}\right)$ is the linear part and $F_{k}^{*}\left(x_{k}\right)=f_{k}\left(x_{k}+v_{k}^{0}\right)-f_{k}\left(v_{k}^{0}\right)-D f_{k}\left(v_{k}^{0}\right) x_{k}$ is the nonlinearity. Obviously, $U_{\varepsilon}(0)=\mathbb{R} \times B_{\varepsilon}(0)$ is contained in the set of definition of the right-hand
side of (8). Note that this simple transformation is a powerful nonautonomous tool. It is of no use in a purely autonomous framework, since (8) in general is nonautonomous.
Step 2: Block Diagonalization of the Linear Part.
In [2, Reduction Theorem] it is shown that there exists a kinematic similarity $S: \mathbb{Z} \rightarrow \mathbb{R}^{N \times N}$ between the linearization $x_{k+1}=A_{k}^{*} x_{k}$ of (8) and a linear system

$$
\begin{equation*}
x_{k+1}=A_{k} x_{k} \tag{9}
\end{equation*}
$$

such that $A: \mathbb{Z} \rightarrow \mathbb{R}^{N \times N}$ is in block diagonal form

$$
A_{k}=\left(\begin{array}{ccc}
A_{k}^{1} & & \\
& \ddots & \\
& & A_{k}^{n}
\end{array}\right)
$$

and each block $A^{i}: \mathbb{Z} \rightarrow \mathbb{R}^{N_{i} \times N_{i}}, i=1, \ldots, n$, corresponds to a spectral interval $\lambda_{i}$. System (9) also has bounded growth, the dichotomy spectra $\Sigma\left(A^{*}\right)$ and $\Sigma(A)$ are the same, and they equal $\lambda_{1} \cup \cdots \cup \lambda_{n}$.
Lemma 3. There exist $\varepsilon^{\prime}, \sigma, \sigma^{\prime}$ with $0<\varepsilon^{\prime} \leq \varepsilon$ and $0<\sigma^{\prime} \leq \sigma$ such that for $k \in \mathbb{Z}$, the mappings

$$
\begin{aligned}
B_{\varepsilon^{\prime}}(0) \rightarrow B_{\sigma}(0), & x \mapsto S_{k}^{-1} x, \\
B_{\sigma^{\prime}}(0) \rightarrow B_{\varepsilon}(0), & x \mapsto S_{k} x,
\end{aligned}
$$

define a $C^{\infty}$ equivalence between (8) and the difference equation

$$
\begin{equation*}
x_{k+1}=A_{k} x_{k}+F_{k}\left(x_{k}\right) \tag{10}
\end{equation*}
$$

with $F_{k}(x)=S_{k+1}^{-1} F_{k}^{*}\left(S_{k} x\right)$ and $\left[A_{k}+F_{k}\right] \in \operatorname{Diff}^{p}\left(\dot{B_{\sigma}}(0),\left[\dot{A}_{k}+F_{k}\right]\left(B_{\sigma}(0)\right)\right)$. Moreover, $\left\|D^{j} F_{k}(0)\right\| \leq M^{\prime}$ for all $k \in \mathbb{Z}$ and all $j \in\{2, \ldots, p\}$ with some $M^{\prime} \geq 0$.
Proof. Due to [2, Corollary 2.1], the kinematic similarity satisfies $S_{k+1}=A_{k}^{*} S_{k} A_{k}^{-1}$. Let $v_{k}$ be a solution of (8). Then $w_{k}:=S_{k}^{-1} v_{k}$ satisfies for all $k \in \mathbb{Z}$

$$
w_{k+1}=S_{k+1}^{-1}\left[A_{k}^{*} v_{k}+F_{k}^{*}\left(v_{k}\right)\right]=A_{k} w_{k}+S_{k+1}^{-1} F_{k}^{*}\left(S_{k} w_{k}\right)
$$

i.e., $w_{k}$ is a solution of (10) with $F_{k}\left(x_{k}\right)=S_{k+1}^{-1} F_{k}^{*}\left(S_{k} x_{k}\right)$. The remaining claims of the lemma follow with the definitions $\sigma:=\sigma^{\prime}:=|S|^{-1} \varepsilon, \varepsilon^{\prime}:=\min \left\{\varepsilon,|S|^{-1} \sigma\right\}$, and $M^{\prime}:=|S|^{p+1} M$ where

$$
|S|:=\max \left\{\sup _{k \in \mathbb{Z}}\left\|S_{k}\right\|, \sup _{k \in \mathbb{Z}}\left\|S_{k}^{-1}\right\|\right\}<\infty,
$$

due to the boundedness of a kinematic similarity.

## Step 3: Elimination of Nonresonant Taylor Components.

This is the crucial step. We will eliminate Taylor components of the nonlinearity which correspond to the blocks $A_{k}^{i} \in \mathbb{R}^{N_{i} \times N_{i}}, i=1, \ldots, n$, of the linear part $A_{k}$. Therefore, define $E^{i}:=\mathbb{R}^{N_{i}}$, $i=1, \ldots, n$, and write $F=\left(F^{1}, \ldots, F^{n}\right)$ with the component functions $F^{i}: D_{F} \rightarrow E^{i}$. Let $\Phi^{i}$ denote the evolution operator of the linear block system $x_{k+1}^{i}=A_{k}^{i} x_{k}^{i}$.

In order to present the ideas, we first motivate the construction of the transformation and the nonresonance condition. For simplicity assume, therefore, that system (10) is globally defined, i.e., $D_{F}=\mathbb{Z} \times \mathbb{R}^{N}$ and that each solution exists on $\mathbb{Z}$; this can be achieved by cutting of $F$ outside the neighbourhood $U_{\delta}(0)$ of the zero solution. Now for all $k \in \mathbb{Z}$ we can expand $F_{k}$ into a Taylor series at $x=0$,

$$
F_{k}(x)=\sum_{q \in \mathbb{N}_{0}^{n}: 2 \leq|q| \leq p} \frac{1}{q!} D^{q} F_{k}(0) \cdot x^{q}+o\left(\|x\|^{p}\right),
$$

$q=\left(q_{1}, \ldots, q_{n}\right) \in \mathbb{N}_{0}^{n}$ a multi-index, $q!=q_{1}!\cdots q_{n}!, x^{q}=\left(x^{1}, \ldots, x^{n}\right)^{q}=\left(x^{1}\right)^{q_{1}} \cdots\left(x^{n}\right)^{q_{n}}$, $|q|=q_{1}+\cdots+q_{n}$. Now we are looking for a condition under which a $C^{p}$ transformation exists which eliminates the $j^{\text {th }}$ component $(1 / q!) D^{q} F_{k}^{j}(0) \cdot x^{q}$ of a summand in the Taylor expansion. Therefore, choose and fix a $j \in\{1, \ldots, n\}$ and a multi-index $q \in \mathbb{N}_{0}^{n}$ with $2 \leq|q| \leq p$. For simplicity we assume that the Taylor coefficients of $F$ at $x=0$ up to order $|q|-1$ are already eliminated; i.e.,

$$
\begin{equation*}
D^{\hat{q}} F_{k}(0)=0, \quad \text { for all } k \in \mathbb{Z}, \quad \text { and all } \hat{q} \in \mathbb{N}_{0}^{n}, \quad \text { with }|\hat{q}| \leq|q|-1 \tag{11}
\end{equation*}
$$

We define a new function $G: \mathbb{Z} \times \mathbb{R}^{N} \rightarrow \mathbb{R}^{N}$ by

$$
G_{k}(x):=F_{k}(x)-\left(0, \ldots, 0, \frac{1}{q!} D^{q} F_{k}^{j}(0) \cdot x^{q}, 0, \ldots, 0\right)
$$

To derive some necessary conditions for the existence of an equivalence we assume now that a near identity $C^{p}$ equivalence $H_{k}(x)=x+h_{k}(x)$ between (10) with zero reference solution and the difference equation

$$
\begin{equation*}
x_{k+1}=A_{k} x_{k}+G_{k}\left(x_{k}\right) \tag{12}
\end{equation*}
$$

with zero reference solution exists, where $h: \mathbb{Z} \times \mathbb{R}^{N} \rightarrow \mathbb{R}^{N}$ is a mapping with $h_{k}(0)=0$ and $D h_{k}(0)=0$ for $k \in \mathbb{Z}$. We will make some observations which will help us to construct an explicit candidate for a $C^{p}$ equivalence.

First, we will assign a difference equation to the values of the transformation $H$ along a fixed solution $\phi(\cdot ; m, \xi)$ of (10).

ObSERVATION 1. For each initial condition $(m, \xi) \in \mathbb{Z} \times \mathbb{R}^{N}$ the mapping $h_{k}(\phi(k ; m, \xi))$ is a solution of

$$
\begin{equation*}
x_{k+1}=A_{k} x_{k}+G_{k}\left(x_{k}+\phi(k ; m, \xi)\right)-F_{k}(\phi(k ; m, \xi)) \tag{13}
\end{equation*}
$$

Observation 1 is a simple but powerful consequence of Lemma 2. Next, we expose a connection between $D^{q} h_{k}(0)$ and $\left.D_{\xi}^{q}\left[h_{k}(\phi(k ; m, \xi))\right]\right|_{\xi=0}$.
ObSERVATION 2. For all $k, m \in \mathbb{Z}$ and $\eta=\left(\eta^{1}, \ldots, \eta^{n}\right) \in E^{1} \times \cdots \times E^{n}=\mathbb{R}^{N}$ we have

$$
\left.D_{\xi}^{q}\left[h_{k}(\phi(k ; m, \xi))\right]\right|_{\xi=0} \cdot \eta^{q}=D^{q} h_{k}(0) \cdot\left[\Phi^{1}(k, m) \eta^{1}\right]^{q_{1}} \cdots\left[\Phi^{n}(k, m) \eta^{n}\right]^{q_{n}}
$$

This can be seen by calculating the partial derivatives which is easily possible since, by (11), the Taylor coefficients of $F$ and $G$ are zero up to order $|q|-1$. The evolution operators $\Phi^{i}$ of the linear block systems $x_{k+1}^{i}=A_{k}^{i} x_{k}^{i}$ come into play because of

$$
\left.D_{\xi^{i}} \phi(k ; m, \xi)\right|_{\xi=0}=\left(0, \ldots, 0, \Phi^{i}(k, m), 0, \ldots, 0\right) \in L\left(E^{i} ; \mathbb{R}^{N}\right)
$$

Now, we replace the $\eta^{i}$ in Observation 2 by $\Phi^{i}(m, k) \zeta^{i}$, and with the identity $\left[\Phi^{i}(k, m)\right]^{-1}=$ $\Phi^{i}(m, k)$, we get the following proposition.
ObSERVATion 3. For all $k, m \in \mathbb{Z}$ and $\zeta \in \mathbb{R}^{N}$ we have

$$
D^{q} h_{k}(0) \cdot \zeta^{q}=D_{\xi}^{q}\left[h_{k}(\phi(k ; m, \xi))\right] l_{\xi=0} \cdot\left[\Phi^{1}(m, k) \zeta^{1}\right]^{q_{1}} \cdots\left[\Phi^{n}(m, k) \zeta^{n}\right]^{q_{n}}
$$

Now, we have a relationship between the Taylor coefficient $D^{q} h_{k}(0)$ of $h$ and the partial derivative $\left.D_{\xi}^{q}\left[h_{k}(\phi(k ; m, \xi))\right]\right|_{\xi=0}$. Observation 1 implies that $h_{k}(\phi(k ; m, \xi))$ is a solution of (13). Then, by differentiation, one can show the following proposition.

Observation 4. The function $D_{\xi}^{q}\left[h_{k}(\phi(k ; m, \xi))\right]_{\xi=0}$ is a solution of

$$
\begin{equation*}
x_{k+1}=A_{k} x_{k}+c_{k}, \tag{14}
\end{equation*}
$$

the variational equation of (13) in $L^{q}\left(E^{1}, \ldots, E^{n} ; \mathbb{R}^{N}\right)$, where

$$
c_{k}=-\left(0, \ldots, 0, D^{q} F_{k}^{j}(0), 0, \ldots, 0\right) \cdot\left[\Phi^{1}(k, m)\right]^{q_{1}} \cdots\left[\Phi^{n}(k, m)\right]^{q_{n}} .
$$

So far, we have assumed that a $C^{p}$ equivalence $H_{k}(x)=x+h_{k}(x)$ between (10) and (12) exists and by Observation 3, the Taylor coefficient $(1 / q!) D^{q} h_{k}(0) \cdot x^{q}$ is a function of a special solution $\left.D_{\xi}^{q}\left[h_{k}(\phi(k ; m, \xi))\right]\right|_{\xi=0}$ of the difference equation (14) and the known evolution operators $\Phi^{i}$.

From now on, we want to use this information to construct a candidate for a $C^{p}$ equivalence between (10) and (12). We make the ansatz

$$
H_{k}(x)=x+\frac{1}{q!} D^{q} h_{k}(0) \cdot x^{q} ;
$$

i.e., $h_{k}(x)=(1 / q!) D^{q} h_{k}(0) \cdot x^{q}$ has only one nontrivial Taylor coefficient.

We make use of Observations 3 and 4 in the way that we choose a special solution $z$ of (14) and interpret $z_{k}$ as $\left.D_{\xi}^{q}\left[h_{k}(\phi(k ; m, \xi))\right]\right|_{\xi=0}$ for an arbitrary but fixed $m \in \mathbb{Z}$. With Observation 3 and our ansatz, this yields

$$
\begin{equation*}
H_{k}(x)=x+\frac{1}{q!} z_{k} \cdot\left[\Phi^{1}(m, k) x^{1}\right]^{q_{1}} \cdots\left[\Phi^{n}(m, k) x^{n}\right]^{q_{n}} . \tag{15}
\end{equation*}
$$

Which solution $z$ of (14) should we choose? To satisfy the condition (C) of Definition 1, it is necessary that $\lim _{x \rightarrow 0} H_{k}(x)=0$ uniformly in $k \in \mathbb{Z}$, and this is satisfied if $z_{k} \cdot\left[\Phi^{1}(m, k) \cdot\right]^{q_{1}} \ldots$ [ $\left.\Phi^{n}(m, k) \cdot\right]^{q_{n}}$ is bounded for $k \in \mathbb{Z}$. Using [2], one knows that the borders of the spectral intervals $\lambda_{i}=\left[a_{i}, b_{i}\right]$ yield the exponential growth rates of the evolution operators $\Phi^{i}$ of $x_{k+1}^{i}=A_{k}^{i} x_{k}^{i}$. Now it is the exponential growth rate of $z$ we have to take care of. Here a key lemma comes in play.
Lemma 8. Consider the $j^{\text {th }}$ component of (14)

$$
\begin{equation*}
x_{k+1}^{j}=A_{k}^{j} x_{k}^{j}+c_{k}^{j} . \tag{16}
\end{equation*}
$$

(A) Assume that the spectral intervals $\lambda_{i}=\left[a_{i}, b_{i}\right]$ satisfy the condition

$$
\begin{equation*}
a_{j}>b_{1}^{q_{1}} \cdots b_{n}^{q_{n}} \tag{17}
\end{equation*}
$$

Choose a $\gamma \in\left(b_{1}^{q_{1}} \cdots b_{n}^{q_{n}}, a_{j}\right)$. Then $z_{k}^{j}:=-\sum_{\ell=k}^{\infty} \Phi^{j}(k, \ell+1) c_{\ell}^{j}$ is the unique solution of (16) with the exponential growth rate $\gamma^{k}$ for $k \rightarrow \infty$, i.e., $\left\|z_{k}^{j}\right\| \leq C^{\prime} \gamma^{k}$ for all $k \geq 0$ with some $C^{\prime} \geq 0$.
(B) Assume the condition

$$
\begin{equation*}
b_{j}<a_{1}^{q_{1}} \cdots a_{n}^{q_{n}} . \tag{18}
\end{equation*}
$$

Choose a $\gamma \in\left(b_{j}, a_{1}^{q_{1}} \cdots a_{n}^{q_{n}}\right)$. Then $z_{k}^{j}:=\sum_{\ell=-\infty}^{k-1} \Phi^{j}(k, \ell+1) c_{\ell}^{j}$ is the unique solution of (16) with the exponential growth rate $\gamma^{k}$ for $k \rightarrow-\infty$.
Proof. We prove only (A). For every $\varepsilon>1$ we get with [2, Corollary 3.1] a constant $K \geq 1$ with

$$
\begin{aligned}
\left\|c_{\ell}^{j}\right\| & \leq\left\|D^{q} F_{\ell}^{j}(0)\right\| \cdot\left\|\Phi^{1}(\ell, m)\right\|^{q_{1}} \cdots\left\|\Phi^{n}(\ell, m)\right\|^{q_{n}} \\
& \leq M K^{|q|} \begin{cases}\left(b_{1}^{q_{1}} \cdots b_{n}^{q_{n}} \varepsilon^{|q|}\right)^{\ell-m}, & \text { for } \ell \geq m \\
\left(a_{1}^{q_{1}} \cdots a_{n}^{q_{n}} \varepsilon^{-|q|}\right)^{\ell-m}, & \text { for } \ell \leq m .\end{cases}
\end{aligned}
$$

The rest follows from [3, Lemma 3.4].

Now, we assume that (17) or (18) holds (both together cannot hold) and choose the following special solution $z=\left(z^{1}, \ldots, z^{n}\right)$ of (14):

$$
z_{k}^{i}:= \begin{cases}0, & \text { if } i \neq j, \\ -\sum_{\ell=k}^{\infty} \Phi^{j}(k, \ell+1) c_{\ell}^{j}, & \text { if } i=j \text { and (17) holds, } \\ \sum_{\ell=-\infty}^{k-1} \Phi^{j}(k, \ell+1) c_{\ell}^{j}, & \text { if } i=j \text { and (18) holds. }\end{cases}
$$

Using (15) and the identity $\Phi^{j}(\ell, m) \cdot \Phi^{j}(m, k)=\Phi^{j}(\ell, k)$, our explicit candidate $H_{k}(x)=x+h_{k}(x)$ for a $C^{p}$ equivalence is defined by

$$
h_{k}^{i}(x)=\left\{\begin{align*}
& 0, \text { if } i \neq j,  \tag{19}\\
& \sum_{\ell=k}^{\infty} \Phi^{j}(k, \ell+1) \frac{1}{q!} D^{q} F_{\ell}^{j}(0) \cdot\left[\Phi^{1}(\ell, k) x^{1}\right]^{q_{1}} \cdots\left[\Phi^{n}(\ell, k) x^{n}\right]^{q_{n}} \\
&-\sum_{\ell=-\infty}^{k-1} \Phi^{j}(k, \ell+1) \frac{1}{q!} D^{q} F_{\ell}^{j}(0) \cdot\left[\Phi^{1}(\ell, k) x^{1}\right]^{q_{1}} \cdots\left[\Phi^{n}(\ell, k) x^{n}\right]^{q_{n}} \\
& \text { if } i=j \text { and (18) holds, holds. }
\end{align*}\right.
$$

Let us have a closer look at conditions (17) and (18). For two compact intervals $[a, b]$ and $[c, d]$ in $\mathbb{R}^{+}$, we introduce a multiplication $[a, b] \cdot[c, d]:=[a c, b d]$ and $[a, b]^{r}:=\left[a^{r}, b^{r}\right]$ for $r \in \mathbb{N}$; furthermore, we will use the relations $[a, b]<[c, d]: \Leftrightarrow b<c$ and $[a, b]>[c, d]: \Leftrightarrow a>d$. With this notation, conditions (17) and (18) for the spectral intervals $\lambda_{i}=\left[a_{i}, b_{i}\right]$ are equivalent to

$$
\lambda_{j}>\lambda_{1}^{q_{1}} \cdots \lambda_{n}^{q_{n}}, \quad \text { respectively, } \quad \lambda_{j}<\lambda_{1}^{q_{1}} \cdots \lambda_{n}^{q_{n}} .
$$

So for our explicit candidate of $H$ to be well defined, we have to assume that one of these two conditions is satisfied and this is equivalent to the so-called nonresonance condition

$$
\begin{equation*}
\lambda_{j} \cap \prod_{i=1}^{n} \lambda_{i}^{q_{i}}=\emptyset . \tag{20}
\end{equation*}
$$

If condition (20) does not hold, then we have a resonance of order $|q|$ and the term $(0, \ldots, 0$, $\left.(1 / q!) D^{q} F_{k}^{j}(0) \cdot x^{q}, 0, \ldots, 0\right)$ is called resonant.
If the linear part $A$ of system (10) does not depend on $k$, then the dichotomy spectrum $\Sigma(A)$ consists of the absolute values $\lambda_{1}, \ldots, \lambda_{n}$ of the eigenvalues of $A$, and the nonresonance condition (20) for the spectral intervals reduces to Poincare's discrete nonresonance condition $\lambda_{j} \neq$ $\lambda_{1}^{q_{1}} \cdots \lambda_{n}^{q_{n}}$ for real eigenvalues.
Now, we prove that $H$ is indeed a $C^{p}$ equivalence, which eliminates a nonresonant term. In contrast to condition (11) in the motivation, we now allow the right-hand side of the difference equation to have nontrivial Taylor coefficients of arbitrary order.
Theorem 9. Consider the difference equation (10). Let $j \in\{1, \ldots, n\}$ be an index and $q \in \mathbb{N}_{0}^{n}$, $2 \leq|q| \leq p$, a multi-index. Assume that the nonresonance condition (20) holds for the spectral intervals $\lambda_{1}, \ldots, \lambda_{n}$. Then a local $C^{p}$ equivalence $H$ exists, which eliminates the $j^{\text {th }}$ Taylor component ( $1 / q!) D^{q} F_{k}^{j}(0) \cdot x^{q}$ belonging to the multi-index $q$ and leaves fixed all other Taylor coefficients up to order $|q|$.

That is, equation (10) is locally $C^{p}$ equivalent to an invertible difference equation

$$
\begin{equation*}
x_{k+1}=A_{k} x_{k}+G_{k}\left(x_{k}\right), \tag{21}
\end{equation*}
$$

with zero reference solution, $\left[A_{k}+G_{k}\right] \in \operatorname{Diff}^{p}\left(B_{\delta}(0),\left[A_{k}+G_{k}\right]\left(B_{\delta}(0)\right)\right)$ with a $\delta=\delta(j, q, A)>0$ and for all $\hat{q} \in \mathbb{N}_{0}^{n}$ with $1 \leq|\hat{q}| \leq|q|$ and all $i \in\{1, \ldots, n\}, k \in \mathbb{Z}$, the identity

$$
D^{\hat{q}} G_{k}^{i}(0) \equiv \begin{cases}D^{\hat{q}} F_{k}^{i}(0), & \text { for } \hat{q} \neq q \text { or } i \neq j \\ 0, & \text { for } \hat{q}=q \text { and } i=j\end{cases}
$$

holds. There exist $\sigma^{\prime}$, $\delta^{\prime}$ with $0<\sigma^{\prime} \leq \sigma, 0<\delta^{\prime} \leq \delta$ and the local near-identity $C^{p}$ equivalence $H: \mathbb{Z} \times B_{\sigma^{\prime}}(0) \rightarrow B_{\delta}(0),(k, x) \mapsto x+h_{k}(x)$ between (10) and (21) with respect to the zero solutions is defined through (19). The inverse transformation $H^{-1}: \mathbb{Z} \times B_{\delta^{\prime}}(0) \rightarrow B_{\sigma}(0)$ has the form

$$
H_{k}^{-1}(x)=x-h_{k}(x)+\psi_{k}(x)
$$

with a continuous mapping $\psi: \mathbb{Z} \times B_{\delta^{\prime}}(0) \rightarrow \mathbb{R}^{N}$, which satisfies the limiting relation $\lim _{x \rightarrow 0} \times$ $\left(\psi_{k}(x) /\|x\|^{|q|^{2}-1}\right)=0$ uniformly in $k \in \mathbb{Z}$. Moreover, for every $k \in \mathbb{Z}$ one has the estimates

$$
\begin{aligned}
\left\|H_{k}(x)-H_{k}(\bar{x})\right\| \leq 2\|x-\bar{x}\|, & \text { for all } x, \bar{x} \in B_{\sigma^{\prime}}(0) \\
\left\|H_{k}^{-1}(x)-H_{k}^{-1}(\bar{x})\right\| & \leq 2\|x-\bar{x}\|,
\end{aligned} \quad \text { for all } x, \bar{x} \in B_{\delta^{\prime}}(0)
$$

Proof. The proof is divided into seven steps. In the first step, we show that $h$ is well defined. The smoothness of $H$ is examined in the second step. In the third step, we construct the inverse transformation $H^{-1}$, and in the following step the Lipschitz estimates for $H$ and $H^{-1}$ are shown. The explicit form of $H^{-1}$ is elaborated in Step 5. The difference equation $x_{k+1}=A_{k} x_{k}+G_{k}\left(x_{k}\right)$ is constructed in the sixth step and it is shown that $G_{k}$ coincides up to order $|q|$ with $F_{k}$ except for the $j^{\text {th }}$ component of the Taylor component belonging to the multi-index $q$; this component is eliminated in $G_{k}$. In the final step, it is proved that $H$ is a local $C^{p}$ equivalence between (10) and (21) with respect to the zero reference solutions.
STEP 1. The mapping $h: \mathbb{Z} \times \mathbb{R}^{N} \rightarrow \mathbb{R}^{N}$ is well defined and the estimate

$$
\begin{equation*}
\left\|h_{k}(x)\right\| \leq C\|x\|^{|q|} \tag{22}
\end{equation*}
$$

holds with a constant $C=C(j, q, A) \geq 0$.
Proof of Step 1. The nonresonance condition (20) implies one of the two estimates (17) or (18) for the spectral intervals $\lambda_{i}=\left[a_{i}, b_{i}\right], i=1, \ldots, n$. For every spectral interval $\lambda_{i}=\left[a_{i}, b_{i}\right]$ choose two numbers $\alpha_{i}$ and $\beta_{i}$ with $0<\alpha_{i}<a_{i}$ and $b_{i}<\beta_{i}$ such that $\alpha_{j}>\beta_{1}^{q_{1}} \cdots \beta_{n}^{q_{n}}$ if (17) holds, and $\beta_{j}<\alpha_{1}^{q_{1}} \ldots \alpha_{n}^{q_{n}}$ if (18) holds. Then as a consequence of [2, Corollary 3.1], there exists a $K=K(j, q, A) \geq 1$ with

$$
\begin{array}{ll}
\left\|\Phi^{i}(k, \ell)\right\| \leq K \beta_{i}^{k-\ell}, & \text { for } k \geq \ell \\
\left\|\Phi^{i}(k, \ell)\right\| \leq K \alpha_{i}^{k-\ell}, & \text { for } k \leq \ell
\end{array}
$$

for $i=1, \ldots, n$. For all $k, \ell \in \mathbb{Z}, x \in \mathbb{R}^{N}$ one has

$$
\begin{aligned}
& \left\|\Phi^{j}(k, \ell+1) \frac{1}{q!} D^{q} F_{k}^{j}(0) \cdot\left[\Phi^{1}(\ell, k) x^{1}\right]^{q_{1}} \cdots\left[\Phi^{n}(\ell, k) x^{n}\right]^{q_{n}}\right\| \\
& \leq\left\|\Phi^{j}(k, \ell+1)\right\| \cdot \frac{1}{q!} M \cdot\left\|\Phi^{1}(\ell, k)\right\|^{q_{1}} \cdots\left\|\Phi^{n}(\ell, k)\right\|^{q_{n}} \cdot\left\|x^{1}\right\|^{q_{1}} \cdots\left\|x^{n}\right\|^{q_{n}} \\
& \leq \frac{1}{q!} M K^{|q|+1}\left\|x^{1}\right\|^{q_{1}} \cdots\left\|x^{n}\right\|^{q_{n}} \begin{cases}\beta_{j}^{-1}\left(\beta_{j} \alpha_{1}^{-q_{1}} \cdots \alpha_{n}^{-q_{n}}\right)^{k-\ell}, & \text { if } k \geq \ell, \\
\alpha_{j}^{-1}\left(\alpha_{j} \beta_{1}^{-q_{1}} \cdots \beta_{n}^{-q_{n}}\right)^{k-\ell}, & \text { if } k \leq \ell,\end{cases}
\end{aligned}
$$

and the claim follows.
STEP 2. For every $k \in \mathbb{Z}$ the mapping $H_{k}: \mathbb{R}^{N} \rightarrow \mathbb{R}^{N}, x \mapsto x+h_{k}(x)$ is $C^{\infty}$.

Proof of Step 2. Arguing for each component separately, the proof of this claim reduces to the verification that $h_{k}^{j}: \mathbb{R}^{N} \rightarrow E^{j}$ is $C^{\infty}$ for every $k \in \mathbb{Z}$. Assume that (17) holds. Similarly to Step 1 , one can show that $h_{k}^{j}$ is differentiable and one gets for all $k \in \mathbb{Z}$ and $x, \xi \in \mathbb{R}^{N}$ the derivative

$$
\begin{aligned}
D h_{k}^{j}(x) \cdot \xi= & \sum_{\ell=k}^{\infty} D_{x}\left[\Phi^{j}(k, \ell+1) \frac{1}{q!} D^{q} F_{\ell}^{j}(0) \cdot\left[\Phi^{1}(\ell, k) x^{1}\right]^{q_{1}} \cdots\left[\Phi^{n}(\ell, k) x^{n}\right]^{q_{n}}\right] \cdot \xi \\
= & \sum_{i=1, \ldots, n: q_{i} \geq 1} q_{i} \sum_{\ell=k}^{\infty} \Phi^{j}(k, \ell+1) \frac{1}{q!} D^{q} F_{\ell}^{j}(0) \\
& \times\left[\Phi^{1}(\ell, k) x^{1}\right]^{q_{1}} \cdots\left[\Phi^{i}(\ell, k) \xi^{i}\right] \cdot\left[\Phi^{i}(\ell, k) x^{i}\right]^{q_{i}-1} \cdots\left[\Phi^{n}(\ell, k) x^{n}\right]^{q_{n}} d s
\end{aligned}
$$

and the estimate

$$
\begin{equation*}
\left\|D h_{k}^{j}(x)\right\| \leq|q| C\|x\|^{|q|-1} . \tag{23}
\end{equation*}
$$

Now, $D h_{k}^{j}(x)$ is again differentiable and the second derivative operates on $\xi, \eta \in \mathbb{R}^{N}$ through

$$
\begin{gathered}
D^{2} h_{k}^{j}(x) \cdot \xi \cdot \eta=\sum_{i, m=1, \ldots, n: q_{i}, q_{m} \geq 1, i \neq m} q_{i} q_{m} \sum_{\ell=k}^{\infty} \Phi^{j}(k, \ell+1) \frac{1}{q!} D^{q} F_{\ell}^{j}(0) \\
\times\left[\Phi^{1}(\ell, k) x^{1}\right]^{q_{1}} \cdots\left[\Phi^{i}(\ell, k) \xi^{i}\right] \cdot\left[\Phi^{i}(\ell, k) x^{i}\right]_{i}^{q_{i}-1} \\
\cdots\left[\Phi^{m}(\ell, k) \eta^{m}\right] \cdot\left[\Phi^{m}(\ell, k) x^{m}\right]^{q_{m}-1} \cdots\left[\Phi^{n}(\ell, k) x^{n}\right]^{q_{n}} \\
+\sum_{i=1, \ldots, n: q_{i} \geq 2} q_{i}\left(q_{i}-1\right) \sum_{\ell=k}^{\infty} \Phi^{j}(k, \ell+1) \frac{1}{q!} D^{q} F_{\ell}^{j}(0) \cdot\left[\Phi^{1}(\ell, k) x^{1}\right]^{q_{1}} \\
\cdots\left[\Phi^{i}(\ell, k) \xi^{i}\right] \cdot\left[\Phi^{i}(\ell, k) \eta^{i}\right] \cdot\left[\Phi^{i}(\ell, k) x^{i}\right]^{q_{i}-2} \cdots\left[\Phi^{n}(\ell, k) x^{n}\right]^{q_{n}}
\end{gathered}
$$

and implies the estimate

$$
\begin{equation*}
\left\|D^{2} h_{k}^{j}(x)\right\| \leq|q|^{2} C\|x\|^{|q|-2} . \tag{24}
\end{equation*}
$$

Mathematical induction yields the existence of the derivatives $D^{m} h_{k}^{j}: \mathbb{R}^{N} \rightarrow L^{m}\left(\mathbb{R}^{N} ; E^{j}\right)$ for $k \in \mathbb{Z}$ and $m=1, \ldots, p$. For $m>p$, the mapping $D^{m} h_{k}^{j}$ is zero, and for this reason $h_{k}^{j}$ and, therefore also, $H_{k}$ is $C^{\infty}$.
STEP 3. There exist $\sigma_{0}^{\prime}, \delta_{0}^{\prime}, \delta_{0}$ with $0<\sigma_{0}^{\prime} \leq \sigma_{0}:=\min \left\{\sigma,(2|q| C)^{-1 /(|q|-1)}\right\}, 0<\delta_{0}^{\prime} \leq \delta_{0}$, such that for arbitrary $k \in \mathbb{Z}$,

$$
H_{k}: B_{\sigma_{0}^{\prime}}(0) \rightarrow H_{k}\left(B_{\sigma_{0}^{\prime}}(0)\right) \subset B_{\delta_{0}}(0)
$$

is a $C^{p}$ diffeomorphism and such that a $C^{p}$ diffeomorphism

$$
H_{k}^{-1}: B_{\delta_{0}^{\prime}}(0) \rightarrow H_{k}^{-1}\left(B_{\delta_{0}^{\prime}}(0)\right) \subset B_{\sigma_{0}}(0)
$$

exists and the identities

$$
H_{k}^{-1}\left(H_{k}(x)\right)=x \quad \text { and } \quad H_{k}\left(H_{k}^{-1}(x)\right)=x
$$

are valid for all $x \in B_{\sigma_{0}^{\prime}}(0)$, respectively, $x \in B_{\delta_{0}^{\prime}}(0)$.
Proof of Step 3. With (24), [4, Proposition 2.5.6, pp. 119-121] implies the claim.
STEP 4. For every $k \in \mathbb{Z}$ we have

$$
\begin{align*}
\left\|H_{k}(x)-H_{k}(\bar{x})\right\| \leq 2\|x-\bar{x}\|, & \text { for } x, \bar{x} \in B_{\sigma_{0}^{\prime}}(0),  \tag{25}\\
\left\|H_{k}^{-1}(x)-H_{k}^{-1}(\bar{x})\right\| \leq 2\|x-\bar{x}\|, & \text { for } x, \bar{x} \in B_{\delta_{0}^{\prime}}(0) . \tag{26}
\end{align*}
$$

Proof of Step 4. First we prove the Lipschitz continuity of $h_{k}: B_{\sigma_{0}^{\prime}}(0) \rightarrow \mathbb{R}^{N}$. Estimate (23) for the derivative of $h_{k}^{j}$ implies for all $k \in \mathbb{Z}$ and $x \in B_{\sigma_{0}^{\prime}}(0)$ the inequality $\left\|D h_{k}(x)\right\| \leq 1 / 2$, and hence,

$$
\begin{equation*}
\left\|h_{k}(x)-h_{k}(\bar{x})\right\| \leq \frac{1}{2}\|x-\bar{x}\| \tag{27}
\end{equation*}
$$

which implies (25). To prove the Lipschitz estimate for $H^{-1}$ we use (27) to show for $k \in \mathbb{Z}$ and $y, \bar{y} \in B_{\sigma_{0}^{\prime}}(0)$ the estimate

$$
\|y-\bar{y}\|-\frac{1}{2}\|y-\bar{y}\| \leq\|y-\bar{y}\|-\left\|h_{k}(y)-h_{k}(\bar{y})\right\|
$$

and it follows that

$$
\frac{1}{2}\|y-\bar{y}\| \leq\left\|H_{k}(y)-H_{k}(\bar{y})\right\| .
$$

Step 3 implies for $x, \bar{x} \in B_{\delta_{0}^{\prime}}(0)$ the identities $H_{k}\left(H_{k}^{-1}(x)\right)=x$ and $H_{k}\left(H_{k}^{-1}(\bar{x})\right)=\bar{x}$ and with $y:=H_{k}^{-1}(x), \bar{y}:=H_{k}^{-1}(\bar{x})$ and one gets estimate (26).
STEP 5. For $k \in \mathbb{Z}$, the mapping $H_{k}^{-1}: B_{\delta_{0}^{\prime}}(0) \rightarrow \mathbb{R}^{N}$ is of the form

$$
H_{k}^{-1}(x)=x-h_{k}(x)+\psi_{k}(x)
$$

with a continuous mapping $\cdot \psi_{k}: B_{\delta_{0}^{\prime}}(0) \rightarrow \mathbb{R}^{N}$, which satisfies

$$
\begin{equation*}
\lim _{x \rightarrow 0} \frac{\left\|\psi_{k}(x)\right\|}{\|x\|^{|q|^{2}-1}}=0, \quad \text { uniformly in } k \in \mathbb{Z} \tag{28}
\end{equation*}
$$

Proof of Step 5. For $k \in \mathbb{Z}$ the inverse of $H_{k}$ can be given explicitely with the Neumann-series (see, e.g., [4, p. 117])

$$
H_{k}^{-1}(x):=\sum_{i=0}^{\infty}\left(-h_{k}\right)^{i}(x), \quad \text { for } x \in B_{\delta_{0}^{\prime}}(0)
$$

With the mapping $\psi_{k}: B_{\delta_{0}^{\prime}}(0) \rightarrow \mathbb{R}^{N}, x \mapsto \sum_{i=2}^{\infty}\left(-h_{k}\right)^{i}(x)$ one has for arbitrary $x \in B_{\delta_{0}^{\prime}}(0)$ the identity

$$
H_{k}^{-1}(x)=x-h_{k}(x)+\psi_{k}(x) .
$$

To show the limiting relation (28) one has to apply twice estimate (22) together with (27) to get for all $x \in B_{\delta_{0}^{\prime}}(0)$ and $i \geq 2$ the following inequalities:

$$
\begin{aligned}
\left\|\left(-h_{k}\right)^{i}(x)\right\| & \leq C\left\|\left(-h_{k}\right)^{i-1}(x)\right\|^{|q|} \leq C^{2}\left\|\left(-h_{k}\right)^{i-2}(x)\right\|^{|q|^{2}} \\
& \leq C^{2}\left(\frac{1}{2}\right)^{(i-2)|q|^{2}}\|x\|^{|q|^{2}}
\end{aligned}
$$

and this implies

$$
\left\|\psi_{k}(x)\right\| \leq \frac{C^{2}}{1-(1 / 2)^{|q|^{2}}}\|x\|^{|q|^{2}}
$$

and therefore, the limiting relation (28).
STEP 6. Define $\delta:=\min \left\{\delta_{0}^{\prime}, \sigma / 2, \sigma_{0}^{\prime} / 2 \tilde{M}\right\}, \delta^{\prime}:=\delta$, and $\sigma^{\prime}:=\min \left\{\sigma_{0}^{\prime}, \delta / 2\right\}$ where $\tilde{M}>0$ is a constant with

$$
\left\|A_{k}+D F_{k}(x)\right\| \leq \tilde{M}, \quad \text { for } x \in B_{\sigma}(0)
$$

If $v$ is a solution of (10) in $B_{\sigma^{\prime}}(0)$, then $H_{k}\left(v_{k}\right)$ is a solution of $x_{k+1}=\tilde{G}_{k}\left(x_{k}\right)$ with $\tilde{G}_{k} \in$ $\operatorname{Diff}^{p}\left(B_{\delta}(0), \tilde{G}_{k}\left(B_{\delta}(0)\right)\right)$ and

$$
\tilde{G}_{k}\left(x_{k}\right)=H_{k+1}\left(A_{k} H_{k}^{-1}\left(x_{k}\right)+F_{k}\left(H_{k}^{-1}\left(x_{k}\right)\right)\right)
$$

If $w$ is a solution of $x_{k+1}=\tilde{G}_{k}\left(x_{k}\right)$ in $B_{\delta^{\prime}}(0)$, then $H_{k}^{-1}\left(w_{k}\right)$ is a solution of (10). Moreover, $\tilde{G}_{k}$ has the form $\tilde{G}_{k}\left(x_{k}\right)=A_{k} x_{k}+G_{k}\left(x_{k}\right)$ and for the components of the Taylor coefficients of $G_{k}: B_{\delta}(0) \rightarrow E^{1} \times \cdots \times E^{n}=\mathbb{R}^{N}$ the following identities hold:

$$
D^{\hat{q}} G_{k}^{i}(0) \equiv \begin{cases}D^{\hat{q}} F_{k}^{i}(0), & \text { for } \hat{q} \neq q \text { or } i \neq j \\ 0, & \text { for } \hat{q}=q \text { and } i=j\end{cases}
$$

for all $\hat{q} \in \mathbb{N}_{0}^{n}$ with $|\hat{q}| \leq|q|$ and all $i \in\{1, \ldots, n\}$.
Proof of Step 6. The Lipschitz estimates (25),(26) and $\left\|\left[A_{k} x+F_{k}(x)\right]-\left[A_{k} \bar{x}+F_{k}(\bar{x})\right]\right\| \leq$ $\tilde{M}\|x-\bar{x}\|$ for $x, \bar{x} \in B_{\sigma}(0)$ imply that $\tilde{G}_{k}$ is a composition of $C^{p}$ diffeomorphisms on $B_{\delta}(0)$.

Let $v_{k} \in B_{\sigma^{\prime}}(0)$ be a solution of (10). Then $H_{k}\left(v_{k}\right) \in B_{\delta}(0)$ is a solution of $x_{k+1}=\tilde{G}_{k}(x)$, and similarly if $w_{k} \in B_{\delta^{\prime}}(0)$ is a solution of $x_{k+1}=\tilde{G}_{k}(x)$, then $H_{k}^{-1}\left(w_{k}\right) \in B_{\sigma}(0)$ is a solution of (10).

Now, we write the components $\tilde{G}^{i}: \mathbb{Z} \times B_{\delta}(0) \rightarrow E^{i}, i=1, \ldots, n$, of the right-hand side of the transformed difference equation as a sum of terms up to order $|q|$ and terms of higher order. The most important relation to do this is the following connection between $h_{k+1}$ and $h_{k}$. For all $k \in \mathbb{Z}$ and $x \in B_{\sigma}(0)$, one has the identity

$$
h_{k+1}\left(A_{k} x\right)=A_{k} h_{k}(x)-\left(0, \ldots, 0, \frac{1}{q!} D^{q} F_{k}^{j}(0) \cdot x^{q}, 0, \ldots, 0\right)
$$

Taylor-expanding $\tilde{G}_{k}$ near $x=0$, one gets the identity

$$
\begin{aligned}
\tilde{G}_{k}(x) & =A_{k} x-A_{k} h_{k}(x)+F_{k}(x)+h_{k+1}\left(A_{k} x\right)+o\left(\|x\|^{|q|}\right) \\
& =A_{k} x+F_{k}(x)-\left(0, \ldots, 0, \frac{1}{q!} D^{q} F_{k}^{j}(0) \cdot x^{q}, 0, \ldots, 0\right)+o\left(\|x\|^{|q|}\right)
\end{aligned}
$$

and the claim follows.
STEP 7. The mapping $H: \mathbb{Z} \times B_{\sigma^{\prime}}(0) \rightarrow B_{\delta}(0)$ is a $C^{p}$ equivalence between systems (10) and (21) with respect to the zero solutions with the inverse transformation $H^{-1}: \mathbb{Z} \times B_{\delta^{\prime}}(0) \rightarrow B_{\sigma}(0)$.
Proof of Step 7. We only have to verify the properties of the definition of a $C^{p}$ equivalence. Use Steps 3, 4, and 6.
Corollary 10. Let $A_{k}$ and $F_{k}$ be periodic in $k$ with a period $\kappa \geq 1$; i.e., for all $k \in \mathbb{Z}$, the identities

$$
A_{k+\kappa}=A_{k} \quad \text { and } \quad F_{k+\kappa}=F_{k}
$$

hold. Then $H$ from Theorem 9 is also periodic in $k$ with period $\kappa$. Especially, if (10) is autonomous, then $H$ is independent of $k$.
Proof. For $\ell \in \mathbb{Z}$ and $\xi \in \mathbb{R}^{N}$, the mapping $\Phi(k+\kappa, \ell+\kappa) \xi$ is the unique solution of the initial value problem $x_{k+1}=A_{k+\kappa} x_{k}, x(\ell)=\xi$ and also $\Phi(k, \ell) \xi$ is the unique solution of the same initial value problem $x_{k+1}=A_{k} x_{k}, x(\ell)=\xi$, and therefore, the identity $\Phi(k+\kappa, \ell+\kappa)=\Phi(k, \ell)$ holds for all $k, \ell \in \mathbb{Z}$. Moreover, the $\kappa$-periodicity of $F$ in $k$ implies the relation $D^{q} F_{k+\kappa}^{j}(0)=D^{q} F_{k}^{j}(0)$, and in case of (17) one gets the equality

$$
h_{k+\kappa}^{j}(x)=\sum_{\ell=k+\kappa}^{\infty} \Phi^{j}(k+\kappa, \ell+1) D^{q} F_{\ell}^{j}(0) \cdot\left[\Phi^{1}(\ell, k+\kappa) x^{1}\right]^{q_{1}} \cdots\left[\Phi^{n}(\ell, k+\kappa) x^{n}\right]^{q_{n}}=h_{k}^{j}(x),
$$

and the claim follows.
Now, it is easy to get our main result on normal forms. Combining the three steps, we immediately get the following theorem.

Theorem 11. Normal Form. Consider a difference equation

$$
\begin{equation*}
x_{k+1}=f_{k}(x) \tag{29}
\end{equation*}
$$

together with a reference solution $v^{0}: \mathbb{Z} \rightarrow \mathbb{R}^{N}$. Assume that
(A) a neighbourhood $U_{\varepsilon}\left(v^{0}\right)$ is contained in $D_{f}$ for some $\varepsilon>0$,
(B) $f_{k} \in \operatorname{Diff}^{p}\left(B_{\varepsilon}\left(v_{k}\right), f_{k}\left(B_{\varepsilon}\left(v_{k}\right)\right)\right)$ for a $p \geq 2$,
(C) the linearization $x_{k+1}=D f_{k}\left(v_{k}^{0}\right) x_{k}$ of (29) along $v^{0}$ has bounded growth, and therefore, [1] the dichotomy spectrum consists of $n, 1 \leq n \leq N$, compact intervals $\lambda_{i}=\left[a_{i}, b_{i}\right], i=$ $1, \ldots, n$, and
(D) higher-order terms of $f$ in $x$ along $v^{0}$ are uniformly bounded in $k$; i.e., there is an $M>0$ such that

$$
\left\|D^{j} f_{k}\left(v_{k}^{0}\right)\right\| \leq M, \quad \text { for all } k \in \mathbb{Z} \text { and all } j \in\{2, \ldots, p\}
$$

Then (29) is locally $C^{p}$ equivalent to a difference equation

$$
\begin{equation*}
x_{k+1}=g_{k}\left(x_{k}\right) \tag{30}
\end{equation*}
$$

with zero reference solution and (30) is in normal form; i.e., it holds that
(A) $g_{k} \in \operatorname{Diff}^{p}\left(B_{\delta}(0), g_{k}\left(B_{\delta}(0)\right)\right)$ for some $\delta>0$,
( $\mathrm{B}^{\prime}$ ) the linearization $x_{k+1}=D g_{k}(0) x_{k}$ of (30) along the zero solution has the same dichotomy spectrum as the linearization of (29) along $v^{0}$ and additionally is block-diagonalized, each block corresponds to a spectral interval $\lambda_{i}$, and
$\left(\mathrm{C}^{\prime}\right)$ all nontrivial Taylor components of $g$ of order 2 to $p$ are resonant; i.e., for every $j \in$ $\{1, \ldots, n\}$ and $q \in \mathbb{N}_{0}^{n}, 2 \leq|q| \leq p$ with

$$
\lambda_{j} \cap \prod_{i=1}^{n} \lambda_{i}^{q_{i}}=\emptyset,
$$

we have $D^{q} g_{k}^{j}(0)=0$ for $k \in \mathbb{Z}$.
We apply the normal form theorem to an example. It is the same example which we used at the beginning to explain Poincarés normal form theory. Therefore, consider again

$$
\begin{aligned}
x_{k+1} & =2 x_{k}, \\
y_{k+1} & =\lambda y_{k}+x_{k}^{2},
\end{aligned}
$$

with $\lambda \in(0, \infty)$. The spectral intervals of the first and second equations are the one-point sets $\lambda_{1}=\{2\}$ and $\lambda_{2}=\{\lambda\}$, respectively, consisting of the eigenvalues of the linear part. We want to eliminate the quadratic term $x_{k}^{2}$ in the second component of the difference equation, i.e., $j=2$ and $q=(2,0)$. For $\lambda<4$ the condition $\lambda_{2}<\left(2 \lambda_{1}+0 \lambda_{2}\right)$ holds, so we have no resonance and get

$$
\begin{aligned}
h_{k}^{2}(x, y) & =-\sum_{\ell=-\infty}^{k-1} \Phi^{2}(k, \ell+1) \cdot\left[\Phi^{1}(\ell, k) x\right]^{2} \\
& =-\sum_{\ell=-\infty}^{k-1} \lambda^{k-\ell-1} \cdot 4^{\ell-k} \cdot x^{2}=\frac{1}{\lambda-4} x^{2}
\end{aligned}
$$

and therefore, $H$ is (we get the same $h_{2}$ for $\lambda<4$ )

$$
H_{k}(x, y)=\binom{x}{y}+\binom{0}{h_{k}^{2}(x, y)}=\binom{x_{1}}{y+\frac{1}{\lambda-4} x^{2}} .
$$

This is the same result as we calculated above with Poincaré's method.

Table 1.

| Autonomous Poincaré Theory | Nonautonomous Theory |
| :---: | :---: |
| $x_{k+1}=A x_{k}+f\left(x_{k}\right)$ <br> Linear part $A \in \mathbb{R}^{N \times N}$ in block diagonal form $A=\operatorname{diag}\left(A^{1}, \ldots, A^{n}\right)$. Eigenvalues $\mu_{1}, \ldots, \mu_{N}$ of $A$ or eigenvalue real parts $\lambda_{1}, \ldots, \lambda_{n}$ of blocks $A^{1}, \ldots, A^{n}$. | $x_{k+1}=A_{k} x_{k}+f_{k}\left(x_{k}\right)$ <br> Linear part $A_{k} \in \mathbb{R}^{N \times N}$ in block diagonal form $A_{k}=\operatorname{diag}\left(A_{k}^{1}, \ldots, A_{k}^{n}\right)$. Compact dichotomy spectrum $\lambda_{1}=\left[a_{1}, b_{1}\right], \ldots, \lambda_{n}=\left[a_{n}, b_{n}\right]$ of blocks $A_{k}^{1}, \ldots, A_{k}^{n}$. |
| Elimination of a Taylor coefficient: algebraically <br> Solve a linear homological equation $L^{j, q} h^{j, q}=D^{q} F^{j}(0)$ <br> with a linear operator $L^{j, q}$ on a finite-dimensional space of monomials. <br> $\leadsto$ Solve a linear equation. | Elimination of a Taylor coefficient: analytically <br> Solve a linear difference equation $x_{k+1}^{j}=A_{k}^{j} x_{k}^{j}-D^{q} F_{k}^{j}(0) \cdot \Phi(k, m)^{q}$ <br> where the solution is unique with a prescribed growth rate. <br> $\leadsto$ Solve a linear difference equation. |
| Nonresonance condition $\begin{gathered} \lambda_{j}-\lambda_{1}^{q_{1}} \cdots \lambda_{n}^{q_{n}} \neq 0 \\ j=1, \ldots, n, q_{i} \in \mathbb{N}_{0}, 2 \leq\|q\| \leq p . \end{gathered}$ | Nonresonance condition (new) $\begin{gathered} \lambda_{j} \cap \lambda_{1}^{q_{1}} \cdots \lambda_{n}^{q_{n}}=\emptyset \\ j=1, \ldots, n, q_{i} \in \mathbb{N}_{0}, 2 \leq\|q\| \leq p . \end{gathered}$ |

## 4. CONCLUSION

We extended Poincare's normal form theory for autonomous difference equations to the class of nonautonomous differential equations in the vicinity of an arbitrary reference solution. Poincare's nonresonance condition for the eigenvalues of the linearization is generalized to a new nonresonance condition for the spectral intervals. A comparison of the new normal form theory with Poincare's method is contained in Table 1.

Normal forms traditionally are an important tool in bifurcation theory (see, e.g., [5-8]). We hope to stimulate the development of a nonautonomous bifurcation theory for difference equations.

## REFERENCES

1. B. Aulbach and S. Siegmund, A spectral theory for nonautonomous difference equations, In Proceedings of the Fifth International Conference on Difference Equations and Applications, Gordon \& Breach, (2000).
2. S. Siegmund, Block diagonalization of linear difference equations, Journal of Difference Equations and Applications 8, 177-189, (2002).
3. B. Aulbach, The fundamental existence theorem on invariant fiber bundles, Journal of Difference Equations and Applications 3, 501-537, (1998).
4. R. Abraham, J.E. Marsden and T. Ratiu, Manifolds, Tensor Analysis, and Applications, Applied Mathematical Sciences, Volume 75, Springer, New York, (1983).
5. V.I. Arnold, Geometrical Methods in the Theory of Ordinary Differential Equations, A Series of Comprehensive Studies in Mathematics, Volume 250, Springer, New York, (1983).
6. V.I. Arnold, Dynamical Systems V, Encyclopaedia of Mathematical Sciences, Volume 5, Springer, Berlin, (1994).
7. S.-N. Chow and J.K. Hale, Methods of Bifurcation Theory, Grundlehren der mathematischen Wissenschaften, Volume 251, Springer, Berlin, (1996).
8. S. Wiggins, Introduction to Applied Nonlinear Dynamical Systems and Chaos, Texts in Applied Mathematics, Volume 2, Springer, New York, (1990).
9. A. Katok and B. Hasselblatt, Introduction to the Modern Theory of Dynamical Systems, pp. 56-104, Cambridge University Press, (1995).
10. S. Siegmund, Spektral-Theorie, glatte Faserungen und Normalformen für Differentialgleichungen vom Cara-théodory-Typ, Dissertation, University of Augsburg, Germany, (1999).
11. S. Siegmund, Dichotomy spectrum for nonautonomous differential equations, Journal of Dynamics and Differential Equations 14 (1), 243-258, (2002).
12. S. Sternberg, Local contractions and a theorem of Poincare, American Journal of Mathematics 79, 809-824, (1957).
13. S. Sternberg, On the structure of local homeomorphisms of Euclidian $n$-space. II., American Journal of Mathematics 80, 623-631, (1958).
14. S. Sternberg, The structure of local homeomorphisms. III, American Journal of Mathematics 81, 578-604, (1959).
15. F. Takens, Partially hyperbolic fixed points, Topology 10, 133-147, (1971).

[^0]:    $0898-1221 / 03 / \$$ - see front matter (C) 2003 Elsevier Science Ltd. All rights reserved. Typeset by $\mathcal{A}_{\mathcal{M}} \mathcal{S}-\mathrm{TE}_{\mathrm{E}}$ PII: S0898-1221(03)00085-3

