## Eindhoven University of Technology

## MASTER

## Numerical and experimental analysis of multiple Chua circuits

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# Numerical and experimental analysis of multiple Chua circuits 

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## Master's thesis

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## Abstract

Nonlinear systems can have multiple equilibria or limit cycles and bifurcations can change the qualitative features of a system, which can result in very complex behavior. Some examples are the global weather, turbulence of fluids or chemical reactions.

The objective of this thesis is to achieve a better understanding of the complicated behavior of complex chaotic systems including synchronization. To provide this insight one of many possible complex systems is chosen and analyzed. The system that is considered in this thesis is the electrical Chua circuit. This seemingly simple circuit, the only nonlinearity is given by a piecewiselinear characteristic, is however capable of generating bifurcation and chaos phenomena.

The rich dynamic behavior of a Chua circuit has been analyzed by determining the stability properties of the equilibria and a Poincaré map has been derived to explain the possible bifurcations. With this powerful tool the qualitative behavior of the Chua circuit can be easily visualized. Further it is shown that the conditions, which are required for chaos, can be fulfilled.

To compare numerical results with experimental results an experimental setup has been designed and built. The parameters of the setup are estimated using two identification methods. The first method is based on measurements of the separate components of the circuit, for the second method a switching nonlinear Kalman filter is designed and implemented. The obtained numerical results are qualitatively comparable with experiments, however it is impossible to follow a chaotic trajectory if no information of the experimental trajectory is used.

Furthermore the possibility of experimental synchronization has been studied. Two synchronization methods, master-slave synchronization and mutual synchronization, have been applied. To be able to specify asymptotic synchronization of non-identical systems a form of practical synchronization is introduced. This is necessary because it is impossible to built identical circuits, due to tolerances in electrical components.
In experiments master-slave synchronization has been achieved by applying a negative feedback to the slave system. In the case of mutual synchronization the circuits are diffusively coupled. Besides synchronization of two systems, synchronization in a network of four symmetrically coupled circuits has been investigated. It is shown that, under certain conditions, it is possible to achieve a partial synchronized situation, which is also experimentally confirmed.

## Samenvatting

Niet-lineaire systemen kunnen meerdere evenwichtspunten of periodieke oplossingen bevatten en bifurcaties kunnen het kwalitatieve gedrag van een systeem veranderen. Dit heeft tot gevolg dat er zeer complex gedrag in een niet-lineair systeem kan optreden. Enkele voorbeelden hiervan zijn het weer, turbulentie in vloeistoffen en chemische reacties.

De doelstelling van dit onderzoek is het verkrijgen van een beter praktisch inzicht in het complexe gedrag van niet-lineaire chaotische systemen inclusief synchronisatie met andere systemen. Om dit inzicht te verkrijgen is een representatief systeem gekozen en onderzocht. Het systeem dat in dit onderzoek wordt beschouwd, is het elektrische Chua circuit. Dit ogenschijnlijk simpele systeem, de enige niet-lineariteit wordt beschreven door een stuksgewijze lineaire karakteristiek, vertoont bifurcaties en chaotisch gedrag.

Het variërende gedrag van het Chua circuit is geanalyseerd door het bepalen van de stabiliteits eigenschappen van de evenwichtspunten en door middel van de afleiding van een Poincaré afbeelding zijn de mogelijke bifurcaties te verklaren. Met de Poincaré afbeelding kan het kwalitatieve gedrag van het Chua circuit eenvoudig zichtbaar gemaakt worden. Eveneens is aangetoond dat aan de condities, die vereist zijn voor chaos, voldaan kan worden.

Om numerieke en experimentele resultaten met elkaar te kunnen vergelijken is er een experimentele opstelling ontworpen en gebouwd. De parameters van de opstelling zijn geschat door middel van het toepassen van twee identificatie methoden. De eerste methode is gebaseerd op metingen van de afzonderlijke componenten van het circuit, voor de twee methode is een schakelend niet-lineair Kalman filter ontworpen en geïmplementeerd. De verkregen numerieke oplossingen zijn kwalitatief vergelijkbaar met de experimenten, maar het is onmogelijk een chaotische oplossing te volgen indien er geen gebruik gemaakt wordt van meetdata.

Tevens is de mogelijkheid van synchronisatie experimenteel onderzocht. Twee synchronisatie methoden zijn toegepast, master-slave synchronisatie en synchronisatie op basis van wederzijdse koppeling. In het geval van master-slave synchronisatie wordt het gedrag van de slave opgelegd door de master. Om asymptotische synchronisatie van niet-identieke systemen te kunnen specificeren is er een vorm van praktische synchronisatie geïntroduceerd. Dit is noodzakelijk, omdat het onmogelijk is identieke circuits te fabriceren vanwege toleranties van de elektrische componenten. Master-slave synchronisatie in experimenten is verkregen door toepassing van een negatieve terugkoppeling op het slave systeem. In het geval van wederzijdse synchronisatie zijn alle circuits met elkaar gekoppeld. Behalve synchronisatie van twee systemen is ook de mogelijkheid van synchronisatie in een netwerk van vier symmetrisch gekoppelde circuits onderzocht. Er is aangetoond dat, onder bepaalde voorwaarden, er partiële synchronisatie optreedt, dit is tevens experimenteel bevestigd.

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Numerical and experimental analysis of multiple Chua circuits

## Chapter 1

## Introduction

### 1.1 Complex systems

The term complex system refers to a system of parts coupled in a nonlinear fashion. Such a system may be discrete, e.g. difference equations, or it may be continuous as in a system of differential equations. Because they are nonlinear, complex systems are more than the sum of their parts because a linear system is subject to the principle of superposition, and hence is literally the sum of its parts. A linear system can have only one or an infinite number of equilibrium points. Nonlinear systems can have multiple equilibria and limit cycles. Bifurcations can change the qualitative features of a complex system such as the number of equilibrium points or limit cycles or their stability properties. This is caused by parametric variation in the model.
Complex systems are found in many fields. Some examples include the global weather, turbulence of fluids, chemical reactions and the well-known nonlinear van der Pol oscillator (van der Pol, 1927).

### 1.2 Chaotic systems

Complex systems are often referred to as chaotic systems. However this is not correct, a chaotic system, in contrast to complex systems, has to satisfy specific properties, which are usually very hard to verify.
The word chaos was introduced by Yorke (Li and Yorke, 1975) several years after the publication of Lorenz (Lorenz, 1963). Although Lorenz' model is considered as one of the first chaotic models, chaos was already noticed by Poincaré. In the year 1887 the king of Sweden sponsored a mathematical competition for the resolution of the question of how stable is the solar system, a variation of the three-body problem. Poincaré found that the evolution of such a system is often 'chaotic' in the sense that a small perturbation in the initial state such as a slight change in one body's initial position might lead to a radically different later state (Poincaré, 1890).
However after the work of Lorenz chaos has developed into a new research area. Smale developed the so-called horseshoe map (Smale, 1967). A new view on the turbulence problem was given by (Ruelle and Takens, 1971), which led to the so-called 'strange' attractor in the phase space. A system inspired by the work of Lorenz is the Rössler system (Rössler, 1976), which has only one spiral instead of two. Another well-known example from biology is the population map or logistic equation (May, 1976), which was investigated further by Feigenbaum. Feigenbaum developed the base of what is now known as the Feigenbaum number. He described that in systems where period doubling occur, the period doubling follows a constant relation (Feigenbaum, 1978).
To illustrate the very complex behavior of chaotic systems we look at the, at first sight simple, logistic equation $x_{t+1}=f\left(x_{t}\right)$ with $f(x)=\alpha x(1-x), x \in[0,1]$ for some values of $\alpha$, see figure 1.1.


Figure 1.1: 140 iterations of $f(x)=\alpha x(1-x)$ for four different values of $\alpha$. If $1<\alpha<3$ we have one stable fixed point, i.e., after a number of iterations $x_{t+1}=x_{t}$. If the value of $\alpha>3$ this fixed point loses its stability and a period doubling bifurcation occurs, if $\alpha$ is increased further several period doubling bifurcations occur and for $3.570<\alpha<4$ the dynamics are very complicated. The trajectories appear to wander around randomly.

### 1.3 Synchronization

The Dutch scientist Christiaan Huygens was probably the first who described synchronization. He discovered that two pendulum clocks hanging from a beam exhibit anti-phase synchronization, i.e. the two pendulums move in opposite directions (Huygens, 1986). Huygens used the pendulum clock to solve the problem of finding longitude at sea, but ironically his discovery that the pendulums influenced each other led the Royal Society to lose faith in pendulum clocks as a solution to the longitude problem.
Synchronization can be understood as the adjustment of rhythms of oscillating objects due to weak interaction (Pikovsky et al., 2001). The requirement of (self-sustaining) oscillating objects gives an indication that complex systems are needed to make synchronization possible. Synchronization is encountered in several fields e.g. (van der Pol and van der Mark, 1928; Winfree, 1980; Gray, 1994; Pikovsky et al., 2001; Nijmeijer and Rodrigues-Angeles, 2003).
However since the publication of (Pecora and Carroll, 1990), chaotic synchronization has been receiving much interest in literature. Some examples of applications of chaotic synchronization can be found in (Dedieu et al., 1993; Terry et al., 1999; Ticos et al., 2000). Synchronization of
electrical circuits is considered in (Heagy et al., 1994; Wu and Chua, 1995; Chua et al., 1996; Matías et al., 1997). A different view on chaotic synchronization, where synchronization is viewed as an observer problem, is given in (Nijmeijer and Mareels, 1997; Santoboni et al., 2001).

### 1.4 Objective

The main goal of this thesis is to get more insight in the complicated behavior of complex (chaotic) systems including synchronization, not only at a theoretical level, but also based on experimental results. To achieve this we divide the objective into the following aspects

- Analysis of a single complex system.
- Development and comparison of numerical and experimental results.
- Synchronization of multiple complex systems; when are two coupled complex systems synchronized?

First of all we need to choose a complex system to investigate. As already mentioned there are many possible choices. However we want to include experimental results and that gives some practical restrictions. Although the drill-string setup in the DCT-lab (Mihajlovic, 2005) is a complex mechanical system, it is not a practical setup to be used for synchronization due to the large dimensions. Biological cells for instance are much more appealing to be used for synchronization. But the main disadvantage of biological systems is to keep the systems and environmental conditions at an adequate level. Also the modeling of biological cells seems still an issue. Therefore we choose an electrical system. This has some great advantages, they are relatively cheap to build and the dimensions can be kept small. A well-known complex electrical system is the so-called Chua circuit (Matsumoto, 1984). This electrical circuit is a hybrid system, due to a piecewiselinear characteristic, and is capable of generating bifurcation and chaos. Although a mechanical equivalent of the Chua circuit is proposed in (Awrejcewicz and Calvisi, 2002), it is questionable of the proposed device will be able to reproduce the diverse dynamical behavior of the electrical circuit, cf. (Verhees, 2004), and therefore the electrical circuit will be used for experiments.

### 1.5 Outline

The outline of this thesis is as follows. In chapter 2 we start with a numerical analysis of a single Chua circuit. By varying a bifurcation parameter several trajectories are generated. We analyze these trajectories by using the equilibrium points of the system and a Poincaré map. This complex Poincaré map is approximated by a simple 1D Poincaré map. With this map the bifurcations are visualized and explained. At the end of this chapter a definition of chaos is presented. With this definition it is possible to show that a Chua system is chaotic in a mathematical context.
In the third chapter the different aspects of building a Chua circuit are explained. Two identification methods are used to obtain parameters of the fabricated realization. The first identification method is based on measurements of the used components and measurements of separate segments of the total circuit. Secondly a Kalman filter is used to estimate states and parameters for the complete circuit. Different strategies are compared and the results are discussed. In chapter 4 the experiments, which validate the numerical results, are described.
Synchronization of Chua circuits is treated in chapter 5. Two synchronization methods are presented here. First master-slave synchronization is explained and applied both numerically and experimentally. Secondly mutual synchronization is applied on a network consisting of maximal four circuits, again both numerically and experimentally.
Finally we end with conclusions and recommendations in chapter 6.

## Chapter 2

## Dynamics of the Chua circuit

The Chua circuit, first discovered by computer simulations (Matsumoto, 1984) and experimentally validated by (Zhong and Ayrom, 1985), is a simple autonomous electrical circuit capable of generating bifurcation and chaos phenomena. The circuit consists of all but one linear elements, a nonlinear resistor with a piecewise-linear characteristic.
The three equations of motion for the circuit are given by

$$
\begin{align*}
C_{1} \dot{v}_{1} & =G\left(v_{2}-v_{1}\right)-f\left(v_{1}\right) \\
C_{2} \dot{v}_{2} & =G\left(v_{1}-v_{2}\right)+i_{L}  \tag{2.1}\\
L \dot{i}_{L} & =-v_{2}-R_{0} i_{L}
\end{align*}
$$

with $G=\frac{1}{R}$ and the function $f\left(v_{1}\right)$ is defined as

$$
\begin{equation*}
f\left(v_{1}\right)=G_{b} v_{1}+\frac{1}{2}\left(G_{a}-G_{b}\right)\left(\left|v_{1}+B_{p}\right|-\left|v_{1}-B_{p}\right|\right) . \tag{2.2}
\end{equation*}
$$

In these equations the state variables $v_{1}$ and $v_{2}$ are the voltages across the capacitors, $C_{1}$ and $C_{2}$, $i_{L}$ is the current flowing through the inductor $L$, which has an internal resistance $R_{0} . G_{a}$ and $G_{b}$ are the conductances of the piecewise characteristic for $\left|v_{1}\right|<B_{p}$ and $\left|v_{1}\right| \geq B_{p}$ respectively. $B_{p}$ is the voltage of the breakpoint. A schematic representation of a Chua oscillator can be seen in figure 2.1. After adding the resistor, $R_{0}$, the circuit has been referred to as a Chua oscillator. Another common name is the 'double scroll attractor', which will be explained later. However all names will be used throughout the text for the circuit as given in (2.1).
An example of a piecewise linear resistor is given in figure 2.2.


Figure 2.1: Schematic layout of a Chua circuit.


Figure 2.2: $v-i$ characteristic of the piecewise linear resistor.

### 2.1 Parameters

The variable resistor $R$, see figure 2.1, will be used as a bifurcation parameter, while all other parameters are considered fixed. In the remainder of this chapter a typical parameter set is chosen as in table 2.1.
Before we present a rigorous analysis of (2.1) some possible trajectories, with the chosen parameters, are given in the next figures to illustrate the variety of the circuit. In figure 2.3 (a), (b), (c) and (d) the periodic solutions are given for both negative and positive initial conditions. It is also possible to choose initial conditions such that the trajectory is unbounded. This is due to the fact that system (2.1) is not a passive system, see (Willems, 1972; van der Steen and Nijmeijer, 2006).

Table 2.1: Parameter values for the circuit.

| Component | Value |
| :--- | ---: |
| $C_{1}$ | $10[\mathrm{nF}]$ |
| $C_{2}$ | $100[\mathrm{nF}]$ |
| $L$ | $22[\mathrm{mH}]$ |
| $R_{0}$ | $22[\Omega]$ |
| $G_{a}$ | $-0.75[\mathrm{mS}]$ |
| $G_{b}$ | $-0.40[\mathrm{mS}]$ |
| $B_{p}$ | $1.75[\mathrm{~V}]$ |



(h) $R=1680[\Omega]$.

Figure 2.3: Trajectories for different values of $R$.

### 2.2 Equilibrium points

There are three equilibrium points of (2.1), which are given by

$$
\begin{align*}
& e q_{1}=\left[\begin{array}{lll}
\frac{B_{p}\left(-G_{a}+G_{b}\right)\left(G R_{0}+1\right)}{G+G_{b} G R_{0}+G_{b}} & \frac{B_{p} G\left(-G_{a}+G_{b}\right) R_{0}}{G+G_{b} G R_{0}+G_{b}} & -\frac{B_{p} G\left(-G_{a}+G_{b}\right)}{G+G_{b} G R_{0}+G_{b}}
\end{array}\right]  \tag{2.3}\\
& e q_{2}=\left[\begin{array}{lll}
0 & 0 & 0
\end{array}\right]  \tag{2.4}\\
& e q_{3}
\end{align*}=\left[\begin{array}{lll}
-\frac{B_{p}\left(-G_{a}+G_{b}\right)\left(G R_{0}+1\right)}{G+G_{b} G R_{0}+G_{b}} & -\frac{B_{p} G\left(-G_{a}+G_{b}\right) R_{0}}{G+G_{b} G R_{0}+G_{b}} & \frac{B_{p} G\left(-G_{a}+G_{b}\right)}{G+G_{b} G R_{0}+G_{b}} \tag{2.5}
\end{array}\right] .
$$

The feasibility of these equilibrium points depends on the value of state $v_{1}$, i.e., eq1 and eq3 are only feasible if $v_{1} \geq B_{p}$ or $v_{1} \leq-B_{p}$ respectively, leading to the following condition for the feasibility of these equilibrium points

$$
\begin{equation*}
\frac{B_{p}\left(-G_{a}+G_{b}\right)\left(G R_{0}+1\right)}{G+G_{b} G R_{0}+G_{b}} \geq B_{p} \tag{2.6}
\end{equation*}
$$

Rewriting condition (2.6) and substituting the parameters values gives a lower and upper bound on the variable resistor $R$

$$
\begin{align*}
\left\lvert\, \frac{G_{b}}{1+G_{b} R_{0}}\right. & \leq G \leq\left|G_{a}\right|  \tag{2.7}\\
1333[\Omega] & \leq R \leq 2478[\Omega] \tag{2.8}
\end{align*}
$$

Using (2.3) and (2.7) the lower bound of $v_{1}$, which occurs at $G=\left|G_{a}\right|=-G_{a}$, is given by

$$
\begin{equation*}
\frac{B_{p}\left(-G_{a}+G_{b}\right)\left(-G_{a} R_{0}+1\right)}{-G_{a}-G_{b} G_{a} R_{0}+G_{b}} \Rightarrow \frac{-B_{p} G_{a}^{2} R_{0}}{G_{a}+G_{a} G_{b} R_{0}-G_{b}}+B_{p} \tag{2.9}
\end{equation*}
$$

The nominator and denominator of (2.9) are both negative, since $B_{p}, R_{0}, G_{a}^{2}>0, G_{a}<G_{b}<0$ and $\left|G_{a}\right| \gg\left|G_{a} G_{b} R_{0}\right|$.
The theoretical upper bound of $v_{1}$ is given by

$$
\begin{equation*}
\lim _{G \rightarrow\left|\frac{G_{b}}{1+G_{b} R_{0}}\right|} \frac{B_{p}\left(-G_{a}+G_{b}\right)\left(G R_{0}+1\right)}{G+G_{b} G R_{0}+G_{b}}=+\infty . \tag{2.10}
\end{equation*}
$$

Therefore the equilibrium point (2.3) is located at a value $v_{1}>B_{p}$.

### 2.2.1 Stability of the equilibrium points

The stability of the equilibrium points is determined by the linearization of (2.1) around the equilibrium points $(2.3),(2.4)$ and (2.5). The stability properties of the third equilibrium point are equal to the first equilibrium point because the linearization around (2.3) and (2.5) are equal. The linearization around the first equilibrium point and the origin for varying $R$ are given in figure 2.4 (a) and (b) respectively.

It is immediately clear that the origin is always an unstable equilibrium point, while there is a region where the real part of all three eigenvalues are negative for the first (and third) equilibrium point. The crossing of the imaginary axis is numerically determined to find the stable region for the parameter $R$.
By means of simulation it is shown that a trajectory of system (2.1) converges to one of those equilibrium points, if $R$ is chosen between $2007-2478$, for a large set of initial conditions. A prove might be obtained using the work of (Leonov et al., 1996).
Further it follows from figure 2.4 (a) that there are two pure imaginary eigenvalues and one real eigenvalue at $R$ equal to 2007 [ $\Omega$ ]. At this point a Hopf bifurcation, see (Guckenheimer and Holmes, 1983), occurs, i.e., the equilibrium point loses its stability and a stable limit cycle appears.


Figure 2.4: Stability diagram for $1333[\Omega]<R<2478[\Omega]$. The arrows indicate the direction of the eigenvalues for an increasing value of $R$. The equilibrium point $e q_{1}$ is stable for $2007[\Omega]<$ $R<2478[\Omega]$.

### 2.3 Poincaré map

To explain figure 2.3 in more detail we use the strategy described in (Chua et al., 1986). In this work a one dimensional approximation of a two dimensional Poincaré map is described. With this tool it is possible to visualize the periodic solutions as well as the 'birth' and 'death' of the double scroll in a rather straightforward manner.

### 2.3.1 Dimensionless model

The dimensionless model of (2.1) is given by

$$
\begin{align*}
\dot{x} & =\alpha(y-h(x)) \\
\dot{y} & =x-y+z  \tag{2.11}\\
\dot{z} & =-\beta y-\gamma z
\end{align*}
$$

where

$$
\begin{array}{rlrl}
h(x) & =x+f(x)=m_{1} x+\frac{1}{2}\left(m_{0}-m_{1}\right)(|x+1|-|x-1|), \\
x & =\frac{v_{1}}{B_{p}}, & \alpha=\frac{C_{1}}{C_{2}}, & a=R G_{a}, \\
y & =\frac{v_{2}}{B_{p}}, & \beta=\frac{R^{2} C_{2}}{L}, & b=R G_{b}, \\
z & =i \frac{R}{B_{p}}, & \gamma=\frac{R R_{0} C_{2}}{L}, & \tau=\frac{t}{R C_{2}} .
\end{array}
$$

The breakpoints are located at $x=-1$ and $x=1$, while the slope is equal to $m_{0}=a+1<0$ at the inner segment and $m_{1}=b+1>0$ at the outer segment. Therefore $h(x)$ can be rewritten

$$
\begin{align*}
h(x) & =m_{1} x+\left(m_{0}-m_{1}\right), & & x \geq 1 \\
& =m_{0} x, & & |x|<1  \tag{2.12}\\
& =m_{1} x-\left(m_{0}-m_{1}\right), & & x \leq 1 .
\end{align*}
$$

Now system (2.11) can be described by five parameters $\alpha, \beta, \gamma, m_{0}$ and $m_{1}$.
The vector field corresponding to (2.11), $\xi: \mathbb{R}^{3} \rightarrow \mathbb{R}^{3}$, satisfies the following properties, cf. (Chua et al., 1986):

1. $\xi$ is symmetric with respect to the origin, i.e., $\xi(-x,-y,-z)=-\xi(x, y, z)$.
2. There a two planes $U_{1}$ and $U_{-1}$ which are symmetric with respect to the origin and they partition $\mathbb{R}^{3}$ into three closed regions $D_{1}, D_{0}$ and $D_{-1}$, see figure 2.5.
3. In each region $D_{i},(i=-1,0,1)$, the vector field $\xi$ is affine, i.e.,

$$
D \xi(x, y, x)=M_{i}, \text { for }(x, y, z) \in D_{i}
$$

where $D \xi$ denotes the Jacobian matrix of $\xi(\mathbf{x})$ and $M_{i}$ a $3 \times 3$ real constant matrix.
4. $\xi$ has three equilibrium points, one at the origin $O$, and one in the interior of $D_{i}, i=-1,1$ denoted by $P^{ \pm}$.
5. Each matrix $M_{i}$ has a pair of complex conjugate eigenvalues $\tilde{\sigma}_{0} \pm j \tilde{\omega}_{0}$ for $M_{0}$ and $\tilde{\sigma}_{1} \pm j \tilde{\omega}_{1}$ for $M_{-1}$ and $M_{1}$, where $\tilde{\omega}_{0}>0$ and $\tilde{\omega}_{1}>0$ and a real eigenvalue $\tilde{\gamma}_{0} \neq 0$ for $M_{0}$ and $\tilde{\gamma}_{1} \neq 0$ for $M_{-1}$ and $M_{1}$.
6. The eigenspace associated with either the real or complex eigenvalues at each equilibrium point is not parallel to $U_{1}$ or $U_{-1}$.

$$
\begin{aligned}
& A=L_{0} \cap L_{1} \\
& B=L_{1} \cap L_{2} \\
& C=E^{r}(O) \cap U_{1} \\
& D=E^{r}\left(P^{+}\right) \cap U_{1} \\
& E=L_{0} \cap L_{2} \\
& F=\mathbf{x} \in L_{2}: \xi(\mathbf{x}) \| L_{2} \\
& L_{0}=E^{c}(O) \cap U_{1} \\
& L_{1}=E^{c}\left(P^{+}\right) \cap U_{1} \\
& L_{2}=\mathbf{x} \in U_{1}: \xi(\mathbf{x}) \| U_{1}
\end{aligned}
$$



Figure 2.5: Eigenspaces of the equilibria. $E^{c}(0)=2 \mathrm{D}$ eigenspace spanned by the real and imaginary part of the complex eigenvalue $\tilde{\sigma}_{0} \pm j \tilde{\omega}_{0}, E^{r}(0)=1 \mathrm{D}$ eigenspace corresponding to the real eigenvalue $\tilde{\gamma}_{0} . E^{c}\left(P^{+}\right)=2 \mathrm{D}$ eigenspace spanned by the real and imaginary part of the complex eigenvalue $\tilde{\sigma}_{1} \pm j \tilde{\omega}_{1}, E^{r}\left(P^{+}\right)=1 \mathrm{D}$ eigenspace corresponding to the real eigenvalue $\tilde{\gamma}_{1}$.

### 2.3.2 Reference frame

A reference frame is introduced by two coordinate transformations for each of the linear regions $D_{0}$ and $D_{1}$ :

$$
\begin{align*}
& \Psi_{0}: D_{0} \rightarrow \mathbb{R}^{3},  \tag{2.13}\\
& \Psi_{1}: D_{1} \rightarrow \mathbb{R}^{3} \tag{2.14}
\end{align*}
$$

These affine transformations reduce $M_{0}$ and $M_{1}$ to the real Jordan form (Chua et al., 1986) and simultaneously transform the planes $U_{ \pm 1}$ in a simple form. In the new coordinate system we have

$$
\begin{align*}
\Psi_{0}(0) & =0  \tag{2.15}\\
\Psi_{0}\left(U_{1}\right) & =V_{0}=\{(x, y, z): x+z=1\}  \tag{2.16}\\
\Psi_{0}\left(U_{-1}\right) & =V_{0}^{-}=\{(x, y, z): x+z=-1\}  \tag{2.17}\\
\frac{1}{\tilde{\omega}_{0}} D \Psi_{0}\left(\xi\left(\Psi_{0}^{-1}\right)\right) & =\xi_{0}(\mathbf{x})=\left[\begin{array}{ccc}
\sigma_{0} & -1 & 0 \\
1 & \sigma_{0} & 0 \\
0 & 0 & \gamma_{0}
\end{array}\right] \mathbf{x}  \tag{2.18}\\
\Psi_{1}(P) & =0  \tag{2.19}\\
\Psi_{1}\left(U_{1}\right) & =V_{1}=\{(x, y, z): x+z=1\}  \tag{2.20}\\
\frac{1}{\tilde{\omega}_{1}} D \Psi_{1}\left(\xi\left(\Psi_{1}^{-1}\right)\right) & =\xi_{1}(\mathbf{x})=\left[\begin{array}{ccc}
\sigma_{1} & -1 & 0 \\
1 & \sigma_{1} & 0 \\
0 & 0 & \gamma_{1}
\end{array}\right] \mathbf{x} \tag{2.21}
\end{align*}
$$

where $\sigma_{0}=\frac{\tilde{\sigma}_{0}}{\tilde{\omega}_{0}}, \gamma_{0}=\frac{\tilde{\gamma}_{0}}{\tilde{\omega}_{0}}, \sigma_{1}=\frac{\tilde{\sigma}_{1}}{\tilde{\omega}_{1}}, \gamma_{1}=\frac{\tilde{\gamma}_{1}}{\tilde{\omega}_{1}}, \mathbf{x}=[x, y, z]^{T}$ and $P=P^{+}$. Although we still use the notation $(x, y, z)$ it are different coordinate systems for both reference frames. However due to the subscript 0 and 1 no confusion should arise. The important points $A, B, C, D, E$ and $F$ from figure 2.5 are also transformed in the new reference frame and denoted by

$$
\begin{array}{lll}
A_{0}=\Psi_{0}(A), & B_{0}=\Psi_{0}(B), & C_{0}=\Psi_{0}(C) \\
D_{0}=\Psi_{0}(D), & E_{0}=\Psi_{0}(E), & F_{0}=\Psi_{0}(F) \\
A_{1}=\Psi_{1}(A), & B_{1}=\Psi_{1}(B), & C_{1}=\Psi_{1}(C) \\
D_{1}=\Psi_{1}(D), & E_{1}=\Psi_{1}(E), & F_{1}=\Psi_{1}(F) \tag{2.23}
\end{array}
$$

The explicit coordinates of (2.22) and (2.23) can be found in appendix A.

### 2.3.3 Poincaré map

In figure 2.6 we consider again figure 2.5, but now with two typical arbitrary trajectories and their corresponding trajectories in the two reference frames. The upper trajectory $\Gamma_{1}$ originates from a point on $U_{1}$, moves downward and turns around, before hitting $U_{-1}$ and returns to $U_{1}$ after a certain amount of time. The trajectory continues upward, turns around and returns again on $U_{1}$. This typical trajectory defines a Poincaré map from some subset $S \subset U_{1}$ into $S$.
The Poincaré map can be decomposed into two components. A half return map, which maps the initial point on $U_{1}$ to the first return point on $U_{1}$, and a second half return map, which maps the first return point to the second return point on $U_{1}$.
The other trajectory, $\Gamma_{2}$ also starts at $U_{1}$, but it crosses the plane $U_{-1}$ before it turns around and hits $U_{-1}$ a second time. Due to the odd-symmetry of the vector field each return point $\mathbf{x}$ on $U_{-1}$ can be identified by its reflected image $-\mathbf{x}$ in $U_{1}$. Also the portion below $U_{-1}$ can be identified with a corresponding version of $\Gamma_{1}$ above $U_{1}$. This means that both typical trajectories define the same Poincaré map, which is the composition of the two half return maps. However the half returns maps cannot be calculated explicitly. Therefore the simplified but equivalent reference frames are used to compute the half return maps.


Figure 2.6: Original system with two typical trajectories and their images in the $D_{0}$ and $D_{1}$ units of the transformed system.

### 2.3.4 First half return map

The line $L_{2}$ in figure 2.6 is mapped into a straight line $L_{20}$ in the $D_{0}$ unit through the points $B_{0}$ and $E_{0}$. The vector field $\xi_{0}(\mathbf{x})$ penetrates $V_{0}$ from above for all $x$ to the right of $L_{2}$, while the vector field $\xi_{0}(\mathbf{x})$ crosses $V_{0}$ from below for $x$ to the left of $L_{2}$. This means that any trajectory starting inside the triangular region

$$
\begin{equation*}
\triangle A_{0} B_{0} E_{0}=\left\{\mathbf{x} \in V_{0} \mid \mathbf{x} \text { is bounded in the triangle } A_{0} B_{0} E_{0}\right\} \tag{2.24}
\end{equation*}
$$

penetrates $V_{0}$ from above and moves down. But the z-axis $=\Psi_{0}\left(E^{r}(0)\right)$ corresponds to the unstable eigenvector $E^{r}(0)$ and therefore the trajectory moves upwards to $V_{0}$ again. This can be defined as one part of the half return map

$$
\begin{equation*}
\pi_{0}^{+}(\mathbf{x}): \triangle A_{0} B_{0} E_{0} \rightarrow V_{0} \tag{2.25}
\end{equation*}
$$

via the image

$$
\begin{equation*}
\pi_{0}^{+}(\mathbf{x})=\varphi_{0}^{T}(\mathbf{x}) \tag{2.26}
\end{equation*}
$$

where

$$
\varphi_{0}^{T}\left(\mathbf{x}_{0}\right)=\left[\begin{array}{ccc}
e^{\sigma_{0} T} \cos (T) & -e^{\sigma_{0} T} \sin (T) & 0  \tag{2.27}\\
e^{\sigma_{0} T} \sin (T) & e^{\sigma_{0} T} \cos (T) & 0 \\
0 & 0 & e^{\gamma_{0} T}
\end{array}\right]\left[\begin{array}{l}
x_{0} \\
y_{0} \\
z_{0}
\end{array}\right]
$$

the flow from $\mathbf{x}$ to the first return point on $V_{0}$ at a first return time $T>0$, where

$$
\begin{equation*}
T=T(\mathbf{x})=\inf \left\{t>0 \mid \varphi_{0}^{t}(\mathbf{x}) \in V_{0}\right\} \tag{2.28}
\end{equation*}
$$

It is assumed that $\varphi_{0}^{t}(\mathbf{x})$ does not cross $V_{0}^{-}$before time $T$.
If a trajectory starts from a point in the wedge, see figure 2.6,

$$
\begin{equation*}
\angle A_{0} B_{0} E_{0}=\left\{\mathbf{x} \in V_{0} \mid \mathbf{x} \text { lies within the wedge extension of } \triangle A_{0} B_{0} E_{0}\right\} \tag{2.29}
\end{equation*}
$$

it must go downward again and the trajectory intersects with $V_{0}^{-}$. This corresponds to the first part of the trajectory $\Gamma_{2}$ and the following map can be defined

$$
\begin{equation*}
\pi_{0}^{-}(\mathbf{x}): \angle A_{0} B_{0} E_{0} \backslash \triangle A_{0} B_{0} E_{0} \rightarrow V_{0}^{-} \tag{2.30}
\end{equation*}
$$

via the image

$$
\begin{equation*}
\pi_{0}^{-}(\mathbf{x})=\varphi_{0}^{T}(\mathbf{x}) \tag{2.31}
\end{equation*}
$$

where

$$
\begin{equation*}
T=T(\mathbf{x})=\inf \left\{t>0 \mid \varphi_{0}^{t}(\mathbf{x}) \in V_{0}^{-}\right\} \tag{2.32}
\end{equation*}
$$

is now the time when the trajectory first hits $V_{0}^{-}$. By making use of the reflected image in $V_{0}$ the first half return map is defined as

$$
\pi_{0}(\mathbf{x})= \begin{cases}\pi_{0}^{+}(\mathbf{x}) & \mathbf{x} \in \triangle A_{0} B_{0} E_{0}  \tag{2.33}\\ -\pi_{0}^{-}(\mathbf{x}) & \mathbf{x} \in \angle A_{0} B_{0} E_{0} \backslash \triangle A_{0} B_{0} E_{0}\end{cases}
$$

In order to calculate (2.33) a local coordinate system is introduced for points inside $\angle A_{0} B_{0} E_{0}$. Each point on $V_{0}$ can be defined by its $x$ and $y$ coordinate since the $z$-coordinate on $V_{0}$ is given by $z=1-x$. A coordinate system $(u, v)$ is chosen such that all points $\mathbf{x} \in \angle A_{0} B_{0} E_{0}$ are specified in terms of $(u, v)$. These coordinates are chosen as a weighted sum of the four corner points $A_{0}, B_{0}, E_{0}$ and $F_{0}$

$$
\begin{equation*}
\mathbf{x}_{0}(u, v)=u\left[v A_{0}+(1-v) E_{0}\right]+(1-u)\left[v B_{0}+(1-v) F_{0}\right] \tag{2.34}
\end{equation*}
$$

where $0 \leq u<\infty$ and $0 \leq v \leq 1$. This is visualized in figure 2.7. Note that all points inside


Figure 2.7: Local $(u, v)$ coordinate system.
$\triangle A_{0} B_{0} E_{0}$ have $0 \leq u<1$, while points outside this triangle have $1 \leq u<\infty$. We can rewrite (2.24) and (2.29) in an equivalent form as follows

$$
\begin{align*}
\triangle A_{0} B_{0} E_{0} & =\left\{\mathbf{x}_{0}(u, v) \mid(u, v) \in[0,1] \times[0,1]\right\}  \tag{2.35}\\
\angle A_{0} B_{0} E_{0} & =\left\{\mathbf{x}_{0}(u, v) \mid(u, v) \in[0, \infty) \times[0,1]\right\} \tag{2.36}
\end{align*}
$$

For a given $\mathbf{x}_{0}(u, v)$ the return map $\pi_{0}^{+}\left(\mathbf{x}_{0}\right)$ is now given by

$$
\pi_{0}^{+}\left(\mathbf{x}_{0}\right)=e^{\sigma_{0} t}\left[\begin{array}{cc}
\cos (t) & -\sin (t)  \tag{2.37}\\
\sin (t) & \cos (t)
\end{array}\right] \mathbf{x}_{0}(u, v)
$$

where $t$ is the first return time and can be calculated by an inverse return time function using the $v$ coordinate

$$
\begin{equation*}
u^{+}(v, t)=\frac{\left\langle\varphi_{0}^{t}\left(B_{0 v}\right), h\right\rangle-1}{\left\langle\varphi_{0}^{t}\left(B_{0 v}-A_{0 v}\right), h\right\rangle} \tag{2.38}
\end{equation*}
$$

where $h=\left[\begin{array}{lll}1 & 0 & 1\end{array}\right]^{T}$ is the normal vector from to origin to $V_{0},\langle$,$\rangle is the inner product, A_{0 v}$ is the location of a point along the line $\overline{E_{0} A_{0}}, v$ units from $E_{0}$ and $B_{0 v}$ is the location of a point along the line $\overline{F_{0} B_{0}}, v$ units from $F_{0}$.
Using the $u$ coordinate the first return time can be calculated by

$$
\begin{equation*}
t=\inf \left\{t \geq 0 \mid u^{+}(v, t)=u\right\} \tag{2.39}
\end{equation*}
$$

The return map $\pi_{0}^{-}\left(\mathbf{x}_{0}\right)$ can be calculated in a similar way using the following equations

$$
\begin{gather*}
u^{-}(v, t)=\frac{\left\langle\varphi_{0}^{t}\left(B_{0 v}\right), h\right\rangle+1}{\left\langle\varphi_{0}^{t}\left(B_{0 v}-A_{0 v}\right), h\right\rangle}  \tag{2.40}\\
t=\inf \left\{t \geq 0 \mid u^{-}(v, t)=u\right\} \tag{2.41}
\end{gather*}
$$

### 2.3.5 Second half return map

The second return map is defined in the same way as the first half return map. The only difference is that the stability properties of the equilibrium point $P$ are opposite of the equilibrium point at the origin. Therefore we look at the inverse image, such that the stable eigenvector $E^{r}(P)$ and the unstable eigenvectors of $E^{c}(P)$ become unstable and stable respectively. This means that a trajectory crossing the $V_{1}$ plane from above inside the angular region $\angle A_{1} B_{1} E_{1}$ moves upwards to hit $V_{1}$. Hence the second half return map can be defined as

$$
\begin{equation*}
\pi_{1}(\mathbf{x}): \angle A_{1} B_{1} E_{1} \rightarrow V_{1} \tag{2.42}
\end{equation*}
$$

via the inverse image

$$
\begin{equation*}
\pi_{1}(\mathbf{x})=\varphi_{1}^{-T}(\mathbf{x}) \tag{2.43}
\end{equation*}
$$

where $\varphi_{1}^{-T}(\mathbf{x})$ is the flow from $\mathbf{x}$ to the first return point where the trajectory hits $V_{1}$ at a reverse time $-T<0$ where

$$
\begin{equation*}
T=T(x)=\inf \left\{t>0 \mid \varphi_{1}^{-t}(x) \in V_{1}\right\} \tag{2.44}
\end{equation*}
$$

To link the first and second half return maps together we need a relation between the $D_{0}$ unit and the $D_{1}$ unit. Therefore a connection map is defined as follows

$$
\begin{equation*}
\Phi=\left(\left.\Psi_{1}\right|_{U_{1}}\right) \circ\left(\left.\Psi_{0}\right|_{U_{1}}\right)^{-1} \tag{2.45}
\end{equation*}
$$

where $\left.\Psi_{1}\right|_{U_{1}}$ and $\left.\Psi_{0}\right|_{U_{1}}$ denote the restrictions on $U_{1}$. Since the coordinates of points on the $V_{0}$ and the $V_{1}$ planes are explicitly known, it is possible to transform the coordinates of one unit into the other.
The total 2D Poincaré map can now be defined as

$$
\begin{equation*}
\pi: V_{1}^{\prime} \rightarrow V_{1}^{\prime} \tag{2.46}
\end{equation*}
$$

where

$$
\begin{equation*}
V_{1}^{\prime}=\left\{(x, y) \in V_{1} \mid x \leq 1\right\} \tag{2.47}
\end{equation*}
$$

and is given by

$$
\pi(\mathbf{x})= \begin{cases}\pi_{1}^{-1} \Phi \pi_{0} \Phi^{-1}(\mathbf{x}) & \mathbf{x} \in \angle A_{1} B_{1} E_{1}  \tag{2.48}\\ \Phi \pi_{0} \Phi^{-1} \pi_{1}^{-1}(\mathbf{x}) & \mathbf{x} \in V_{1}^{\prime} \backslash \angle A_{1} B_{1} E_{1}\end{cases}
$$

### 2.3.6 Approximation of the Poincaré map

Instead of using the complex Poincaré map (2.48) a simpler approximation can be made. For the parameter values given in table 2.1 it holds that the real eigenvalue $\left|\gamma_{1}\right|$ is relatively large compared to $\sigma_{1}$. This means that all trajectories originating on either side of the eigenspace $E^{c}(P)$ will converge rapidly to this eigenspace and eventually cross the plane $U_{1}$ at the line $L_{1}$ in figure 2.5. Therefore we define a 1D approximation of the Poincaré map $\pi$ by restricting its domain to a line $L_{1}^{\prime}$, which starts at point $B$, passes through $A$ and extends to infinity. The domain of the 1D approximation map is chosen as another semi-infinite line segment. This line is constructed through the points $M=\Psi_{1}^{-1}[1,0,0]^{T}$ and $P^{+}$, its extension through $N$ and then deleting the line segment between $M$ and $P^{+}$, which defines the following Poincaré map

$$
\begin{equation*}
\pi^{*}: \overline{P^{+} N_{\infty}} \rightarrow \overline{P^{+} N_{\infty}} \tag{2.49}
\end{equation*}
$$



Figure 2.8: Geometrical interpretation of the one dimensional approximated Poincaré map.

In figure 2.8 the $D_{1}$ unit is shown where the plane $W$, which corresponds to the line segments $\overline{P^{+} D}$ and $\overline{N D}$ from figure 2.5, is transformed into the rectangular region

$$
\begin{equation*}
\left.W_{1}=\{x, y, z) \in \mathbb{R}^{3} \mid x \leq 0, y=0\right\} \tag{2.50}
\end{equation*}
$$

In terms of local coordinates $(u, v)$, points on the line $\overline{B_{1} A_{1 \infty}}$ are identified by a single coordinate $u$ since $v=1$ on this line

$$
\mathbf{x}(u)=\mathbf{x}_{1}(u, 1) \begin{cases}0 \leq u \leq 1 & \mathbf{x}(u) \in \overline{B_{1} A_{1}}  \tag{2.51}\\ 1<u<\infty & \mathbf{x}(u) \in \overline{A_{1} A_{\infty}}\end{cases}
$$

All trajectories originating from the line $\overline{B_{1} A_{1 \infty}}$ remain in the $x-y$ plane, since this is the eigenspace $\Psi_{1}\left(E^{c}(P)\right.$ ), spiraling inwards (in backward time) and hit the line $\overline{0 N_{1}}$ a distance $X(u)$ from 0 after a specific time interval $\pi+\arg (\mathbf{x}(u))$. We can rewrite this as

$$
\begin{equation*}
X(u)=|\mathbf{x}(u)| e^{-\sigma_{1}(\pi+\arg (\mathbf{x}(u)))} \geq 0 \tag{2.52}
\end{equation*}
$$

where

$$
\begin{equation*}
\arg (\mathbf{x})=\arctan \left(\frac{\mathbf{x}_{y}(u)}{\mathbf{x}_{x}(u)}\right) \tag{2.53}
\end{equation*}
$$

It is shown in (Chua et al., 1986) that the inverse return time functions $u^{+}(1, t)(2.38)$ and $u^{-}(1, t)$ (2.40) are both monotone functions and hence have an unique inverse. This means that every point
$X(u) \geq 0$ on the line $\overline{N_{1} X(0)}$ maps uniquely into a point on $\overline{B_{1} A_{1 \infty}}$ via the flow $\varphi_{1}^{t}$, where $X(0)$ is the limiting point which maps into $B_{1}$. Any point $d_{1}$ between $X(0)$ and 0 must map into a point $d_{2}$ where

$$
\begin{equation*}
d_{2}=e^{2 \pi \sigma_{1}} d_{1} \tag{2.54}
\end{equation*}
$$

since it does not intersect the line $\overline{B_{1} A_{1 \infty}}$.
Each point on the line $\overline{B_{1} A_{1 \infty}}$ maps into one of the two spirals, defined by (2.33) and transformed into the $D_{1}$ unit as shown in figure 2.8 with coordinates $\mathbf{y}(u)$. It is also shown in (Chua et al., 1986) that under the flow $\varphi_{1}^{t}$ the point $\mathbf{y}(u)$ intersects with the plane $W_{1}$ with two spirales at coordinates

$$
\begin{equation*}
Y^{\prime}(u)=\left(-|\mathbf{y}(u)| e^{\sigma_{1}(\pi-\arg (\mathbf{y}(u)))}, 0, \mathbf{y}_{z}(u) e^{\gamma_{1}(\pi-\arg (\mathbf{y}(u)))}\right) \tag{2.55}
\end{equation*}
$$

where

$$
\begin{equation*}
\arg (\mathbf{y})=\arctan \left(\frac{\mathbf{y}_{y}(u)}{\mathbf{y}_{x}(u)}\right) \tag{2.56}
\end{equation*}
$$

If $\gamma_{1}$ is relatively large, the two spirales are squeezed into a thin line close to the line $\overline{N_{\infty} 0}$ and the points $Y^{\prime}(u)$ can be approximated with points on the line $\overline{N_{\infty} 0}$ by

$$
\begin{equation*}
Y(u)=|\mathbf{y}(u)| e^{\sigma_{1}(\pi-\arg (\mathbf{y}(u)))}, \quad 0 \leq u \leq \infty \tag{2.57}
\end{equation*}
$$

To calculate the approximated Poincaré map (2.49) for $X(u)>X(0)$ we use

$$
(X(u), Y(u))=\left\{\begin{array}{ll}
\left(X\left(u^{+}(1, t)\right), Y\left(u^{+}(1, t)\right)\right), & 0 \leq t<\infty \text { for } 0 \leq u<1  \tag{2.58}\\
\left(X\left(u^{-}(1, t)\right), Y\left(u^{-}(1, t)\right)\right), & 0<t<\infty \text { for } 1 \leq u<\infty
\end{array}\right\}
$$

For points between $X(0)$ and 0 we make use of (2.54)

$$
\begin{equation*}
Y(u)=e^{2 \pi \sigma_{1}} X(u), \quad u<0 \tag{2.59}
\end{equation*}
$$

The approximated Poincaré (2.49) map can also be identified in figure 2.5. The point $\frac{B^{\prime} \text { corre- }}{\overline{P+B^{\prime}}}$ sponds with $X(0)$. For each point $\mathbf{x} \in \overline{P^{+} B^{\prime}}$ the poincaré map is a linear map from $\overline{P^{+} B^{\prime}}$ onto $\overline{P^{+} \pi^{*}\left(B^{\prime}\right)}$. For points $\mathbf{x} \in \overline{B^{\prime} N_{\infty}}, \pi^{*}$ is a nonlinear continuous map from $\overline{B^{\prime} N_{\infty}}$ into $\overline{P^{+} N_{\infty}}$.

### 2.4 1D Poincaré map

With the Poincaré map and a value for the variable resistor $R$ the qualitative behavior of system (2.1) can be predicted, without explicitly integrating the equations. The approximated Poincaré map is given by (2.58) and (2.59) for a given parameter set. After the transient of (2.1) the behavior can be visualized by iterating this map.
An important point of the map is located at $X(1)$ which corresponds to the point $A$. If a trajectory starts at the line $\overline{B A}$, which is part of the triangular region (2.24), it does not intersect with the plane $U_{-1}$, see figure 2.6, and we call this trajectory a single scroll trajectory. If a trajectory starts at the line $\overline{A A_{\infty}}$, which is part of the wedge extension (2.29), the trajectory crosses the $U_{-1}$ plane before it returns and is naturally called a double scroll trajectory. The 'birth' and 'death' of double scroll trajectories can therefore be explained using the Poincaré map.
To compare the approximated Poincaré map with the trajectories of (2.1) as in figure 2.3 the same resistor values are used to compute the maps. In figure 2.9 the map for $R=1985[\Omega]$ is shown. The boxed area in the upper left corner is a magnification of the rectangular box in the middle of the figure. Further the point $X(1)$ is located at $X=1.1223$, indicated by a vertical dashed line. It can be seen that the map has a stable fixed point $X(u)=\pi^{*}(X(u))=Y(u)$ to the left of the point $X(1)$. This means that a single scroll period one trajectory should occur, which is indeed the case in figure 2.3 (a). A single scroll period two solution, indicated by $X_{2}=\pi^{*}\left(X_{1}\right), X_{1}=\pi^{*}\left(X_{2}\right)$ with $X_{1}$ and $X_{2}$ both to the left of $X(1)$ can be seen in figure 2.10 . In figures 2.11 and 2.12 the existence of single scroll periodic solutions with period four and period three are visualized.
The transition from a single scroll to the double scroll, i.e. the 'birth' of the double scroll, is
shown in figure 2.13. With $R=1942[\Omega]$ the maximum value of $Y$ on the interval $[0 X(1)]$ is equal to $\mathrm{X}(1)$. This means that $Y$ maps precisely through the boundary point $A_{1}$, if $R$ is decreased further the point $Y_{\max }$ on the interval $[0 X(1)]$ maps to the right of $X(1)$ and a double scroll scroll trajectory occurs.
Another special case is observed at $R=1778[\Omega]$, see figure 2.14. A double scroll trajectory occurs since $Y_{\max }$ on the interval $[0 X(1)]$ is greater than $X(1)$, but the minimum value $Y_{\min }$ on the interval $[X(1) \infty]$ is equal to zero. This implies that the spiral $A_{1 \infty}^{\prime} C$ in figure 2.8 goes through the point $D_{1}$, because $Y=0$. Recall from figure 2.6 that $D_{1}$ is the point on $V_{1}$ corresponding to the stable eigenvector $\Psi_{1}\left(E^{r}(P)\right)$ of $P^{+}$. Therefore the double scroll trajectory goes through $P$ and spirales outwards until it hits $U_{1}$, this is referred to as a hole-filling orbit.
A different periodic two limit cycle is shown in figure 2.15. Again $X_{2}=\pi^{*}\left(X_{1}\right), X_{1}=\pi^{*}\left(X_{2}\right)$, but now $X_{1}<X(1)$ and $X_{2}>X(1)$. So we have one period of the limit cycle as a single scroll and the other period as the double scroll. This corresponds to the trajectory of figure $2.3(\mathrm{~g})$.
Finally in figure 2.16 the 'death' of the double scroll is visualized. $Y_{\max }$ on the interval $[0 X(1)]$ maps to a second fixed point which corresponds to an unstable saddle-type closed orbit cycle, cf. (Chua et al., 1986). When the attractor collides with the unstable closed orbit a boundary crisis (Grebogi et al., 1982) occurs, resulting in the loss of the attractor (and its basin) and trajectories diverge to $\pm \infty$ for every initial condition.
Since this is physically impossible the model for the nonlinear resistor (2.2) could be extended with a second breakpoint $B_{p 2}>B_{p}$ and a positive slope $G_{c}$ to introduce passivity if $\left|v_{1}\right|>B_{p 2}$,

$$
\begin{equation*}
f\left(v_{1}\right)=G_{c} v_{1}+\frac{1}{2}\left(G_{a}-G_{b}\right)\left(\left|v_{1}+B_{p}\right|-\left|v_{1}-B_{p}\right|+\frac{1}{2}\left(G_{b}-G_{c}\right)\left(\left|v_{1}+B_{p 2}\right|-\left|v_{1}-B_{p 2}\right|\right)\right. \tag{2.60}
\end{equation*}
$$

A new characteristic, with $B_{p 2}=13[\mathrm{~V}]$ and $G_{c}=4.5[\mathrm{mS}]$, is shown in figure 2.17. This has no effect on the actual attractor and the unstable closed orbit, because normally $\left|v_{1}\right|<B_{p 2}$. However if the boundary crisis occurs the resulting trajectory is a large stable limit cycle, see (Matsumoto, 1984).


Figure 2.9: Period one limit cycle for $R=1985[\Omega]$.


Figure 2.10: Period two limit cycle for $R=1975[\Omega]$.


Figure 2.11: Period four limit cycle for $R=1967[\Omega]$.


Figure 2.12: Period three limit cycle for $R=1960[\Omega]$.


Figure 2.13: Birth of the double scroll for $R=1942[\Omega]$.


Figure 2.14: Hole-filling orbit for $R=1778[\Omega]$.


Figure 2.15: Double scroll period two limit cycle for $R=1700[\Omega]$.


Figure 2.16: Death of the double scroll for $R=1609[\Omega]$.


Figure 2.17: $v-i$ characteristic of the extended piecewise linear resistor.

### 2.5 Chaos

It is already mentioned that the Chua circuit is also capable of generating chaos. In this section the approximated Poincaré map is used to explain why chaotic behavior is possible in the circuit. In figure 2.18 another special case is shown. The point $X(1)$ is now a fixed point, i.e., it holds that $Y(1)=\pi^{*}(X(1))=X(1)$. This implies that a trajectory which starts at the point $A$ converges to the origin since it lies in the stable eigenspace of $E^{c}(0)$ and continues along the unstable eigenspace $E^{r}(0)$ until it hits the plane $U_{1}$ in point $C$, see figure 2.5 . Point $C$ corresponds to $C_{1}$ in figure 2.8 and since $Y(1)=X(1)$, the trajectory continuing from $C_{1}$ intersects the plane $W_{1}$ at $Y^{\prime}(1)$, whose projection $Y(1)$ is equal to $X(1)$. This trajectory is a homoclinic orbit of the origin. A homoclinic orbit is an orbit which tends to a fixed point under backward iteration and lands on the same fixed point under forward iteration. This orbit plays a crucial role to show that the double scroll attractor is a chaotic system.


Figure 2.18: Homoclinic orbit for $R=1635.7[\Omega]$ and $G_{b}=-0.5[\mathrm{mS}]$.

First we need to characterize chaos in such a way that it is easy to verify. Therefore we use the definition of chaos for discrete time dynamics as is given in (Devaney, 1986):
Let $V$ be a set. $F: V \rightarrow V$ is said to be chaotic on $V$ if

1. $f$ has sensitive dependence on initial conditions.
2. $f$ is topologically transitive.
3. periodic points are dense in $V$.

If a system has sensitive dependence on initial conditions than initially nearby points can evolve quickly into very different states. This implies that a chaotic system is unpredictable.
A function $f: J \rightarrow J$ is said to be topologically transitive if for any pair of open sets $U, V \subset J$ there exists $k>0$ such that $f^{k}(U) \cap V \neq \emptyset$, i.e., since the intersection is a non-empty set some points from $U$, under iteration, intersect with the set $V$. So a dynamical system, which is topologically transitive, cannot be decomposed into two disjoint open sets which are invariant under the map. Therefore it is not possible to decompose a chaotic system into two subsystems.
However there is still regular behavior possible in a chaotic system since the periodic points are dense in $V$.

To verify that system (2.1) is chaotic we make use of a theorem posed by Šilnikov, cf. (Guckenheimer and Holmes, 1983; Silva, 1993).
Šilnikov considered a three dimensional flow, containing a homoclinic trajectory to a saddle point,
see figure 2.19.
Under the condition that the real eigenvalue $(\gamma)$ has larger magnitude than the real part of the complex eigenvalues $(\sigma \pm j \omega)$, there are horseshoes present in the return maps defined near the homoclinic orbit. The proof for this can be found in (Guckenheimer and Holmes, 1983).


Figure 2.19: Abstraction of a homoclinic orbit to a saddle point.

The horseshoe is an example of a hyperbolic limit set (Guckenheimer and Holmes, 1983). The Smale horseshoe map is defined by the map $f: S \rightarrow \mathbb{R}^{2}$, where $S=[0,1] \times[0,1]$ is the unit square in the plane, shown in figure 2.20 .


Figure 2.20: Abstraction of the Smale horseshoe map.

If $f$ is iterated most points leave $S$, however the points that remain in $S$ for all time form a set $\Lambda=\left\{x \mid f^{i}(x) \in S,-\infty<i<\infty\right\}$. This set is a Cantor set, see (Devaney, 1986). The horseshoe map has the following properties (Guckenheimer and Holmes, 1983):
The map $f$ has an invariant Cantor set $\Lambda$ such that:

1. $\Lambda$ contains a countable set of periodic orbits of arbitrarily long periods.
2. $\Lambda$ contains an uncountable set of bounded nonperiodic motions.
3. $\Lambda$ contains a dense orbit.

Now we are able to verify that the double scroll is indeed chaotic. There exists a parameter set such that there is a homoclinic orbit through the origin, (see figure 2.18), further it holds that $\left|\sigma_{0}\right|<\gamma_{0}$ and therefore Šilnikovs conditions hold. The points two and three in the definition of a chaotic system follow directly from the Cantor set $\Lambda$, (since $\Lambda$ contains a dense orbit it is also topologically transitive (Devaney, 1986)). It remains to verify that system (2.1) is sensitive to
initial conditions.
We try to demonstrate this through a simulation. System (2.1) is integrated twice with slightly different initial conditions. In the first simulation the initial conditions are chosen as $\left[\begin{array}{ll}1 & 1\end{array} 0\right]^{T}$, for the second simulation the initial condition is chosen as $\left[1+1 \cdot 10^{-11} 10\right]^{T}$. The results for the first state, $v_{1}$, are shown in figure 2.21. After $t=0.012$ it is visible that the solutions diverge. This shows that system (2.1) is sensitive to initial conditions.
Another property indicating chaos in system (2.1) is the existence of the period three limit cycle as shown in figure 2.12, see (Šarkovskii, 1964; Li and Yorke, 1975).


Figure 2.21: Time signals for $v_{1}$ of both simulations.

## Chapter 3

## Building and identification of a Chua circuit

In this chapter the electrical realization of model (2.1) is presented. The different aspects of this realization are explained. After the design the identification of the circuit is described. Different identification methods are compared and the results are discussed.

### 3.1 Realization

Consider again the schematic circuit in figure 3.1. The circuit basically consists of two capacitors, an inductor, a variable resistor and the so-called negative resistor. In the next paragraphes the actual electronic counterparts of the schematic components are introduced. The main part of the layout of the circuit is based on the methodology presented in (Kennedy, 1992). The nominal values of the electronic components are chosen to match with the parameter values in table 2.1.


Figure 3.1: Schematic layout of a Chua circuit.

### 3.1.1 Variable resistor and capacitors

The variable resistor, $R$, is composed of a precision metal film resistor with a nominal value of $1500[\Omega]$ and a precision potentiometer with a range from $0-500[\Omega]$. In this way the variable resistor is more accurate in the region of interest. The dielectric capacitors $C_{1}$ and $C_{2}$ are chosen with nominal values of $0.01[\mu F]$ and $0.1[\mu F]$ respectively.

### 3.1.2 Nonlinear resistor

There are several ways to obtain the piecewise linear function, $f\left(v_{1}\right)(2.2)$, using electronic components, e.g. using transistors (Matsumoto et al., 1986), operational amplifiers (Op Amps) (Kennedy, 1992) or operational transconductance amplifiers (Cruz and Chua, 1992). In this design the nonlinear resistor is constructed with operational amplifiers and resistors.
A negative resistance can be synthesized using three normal positive resistances and a voltagecontrolled voltage source (VCVS). A VCVS is an ideal circuit element with two important properties: The voltage at the output is a linear function of the potential difference at the input and no current flows in, or out, the input terminals. Since a VCVS is an ideal circuit element an Op Amp is used as an approximation of a VCVS. The output of an Op amp depends linearly on the potential difference $V_{d}$ at the input, however the output $V_{o}$ of an Op Amp saturates when the absolute difference of the input crosses a threshold voltage. This can be seen in figure 3.2, it is assumed that there is no voltage offset $\left(V_{d}=0[V] \mapsto V_{o}=0[V]\right)$. When the resistors are connected as shown in figure 3.3 the following relations (Kennedy, 1992) hold


Figure 3.2: Operational Amplifier and voltage transfer characteristic.


Figure 3.3: Negative resistor converter and $v-i$ characteristic.

$$
\begin{align*}
s_{0} & =\frac{1}{R_{1}}  \tag{3.1}\\
B_{p}^{+} & =\frac{R_{2}+(1+A) R_{3}}{A\left(R_{2}+R_{3}\right)} V_{s a t}^{+}  \tag{3.2}\\
s_{1} & =\frac{(1-A) R_{2}+R_{3}}{R_{1}\left[R_{2}+(1+A) R_{3}\right]}  \tag{3.3}\\
B_{p}^{-} & =\frac{R_{2}+(1+A) R_{3}}{A\left(R_{2}+R_{3}\right)} V_{\text {sat }}^{-} \tag{3.4}
\end{align*}
$$

where $A$ is the gain of the Op Amp and $V_{\text {sat }}$ the saturation voltage.
When the gain $A$ is large and $R_{2}$ is equal to $R_{1}$ these equations can be approximated by

$$
\begin{align*}
s_{0} & =\frac{1}{R_{1}}  \tag{3.5}\\
B_{p}^{+} & \approx \frac{R_{3}}{R_{2}+R_{3}} V_{\text {sat }}^{+}  \tag{3.6}\\
s_{1} & \approx-\frac{1}{R_{3}}  \tag{3.7}\\
B_{p}^{-} & \approx \frac{R_{3}}{R_{2}+R_{3}} V_{s a t}^{-} \tag{3.8}
\end{align*}
$$

It follows from figure 3.3 that as soon as the voltage exceeds $B_{p}$ the resistance becomes positive. This prevents that after the 'death' of the double scroll the voltage becomes unbounded.
When two negative resistors are connected in parallel the piecewise-linear resistor (2.2) is obtained.

### 3.1.3 Inductor

The inductor in the circuit is constructed using resistors, an Op Amp and a capacitor. With these components it is possible to synthesize an inductor. The inductor used in a previous realization (Hees, 2004) did not give satisfactory results, therefore a different layout is used. This new layout is shown in figure 3.4 and based on the work of (Pegna et al., 2000). Under the assumption that the gain $A \rightarrow \infty$, the value of $L$ can be approximated by

$$
\begin{equation*}
L \approx C_{3} R_{0}\left(R-R_{0}\right) \tag{3.9}
\end{equation*}
$$

where $R_{0}$ acts as the internal resistance of the inductor.


Figure 3.4: Inductor layout.

### 3.1.4 Total circuit

With the parts described in the previous sections the complete circuit can be realized. However some additional components are used to be able to measure the voltage across the two capacitors $C_{1}$ and $C_{2}$. In order to use data acquisition devices, such as Siglab, or other electronic devices (loudspeakers or other Chua circuits) without influencing the circuit itself so-called voltage followers are added. A practical addition are two 'jumpers' that divide the circuit into three separate subsystems. In this way it is possible to measure only the nonlinear resistor or a resonance circuit consisting of the inductor $L$ and capacitor $C_{2}$. Another practical addition are bypass capacitors connected across the power supplies of the Op Amps to maintain a steady DC voltage. The complete realization can be seen in figure 3.5. The circuit layout with component list can be found in appendix B.


Figure 3.5: Realization of Chua circuit.

### 3.2 Identification using measurements

### 3.2.1 Variable resistor and capacitors

The variable resistor is measured using a standard multi-meter. A linear characteristic is assumed and approximated by a least square estimate

$$
\begin{equation*}
R=49.3(\text { turn }+ \text { number } / 100)+0.8 \tag{3.10}
\end{equation*}
$$

The turn is the number of rotations of the potentiometer $(0-10)$ and number is the scale division


Figure 3.6: Characteristic of the variable resistor.
(0-100) in one turn.
The capacitors $C_{1}, C_{2}$ and $C_{3}$ are measured with a Fluke RCL meter before assembly and con-
sidered fixed hereafter.

$$
\begin{align*}
C_{1} & =10.90[n F]  \tag{3.11}\\
C_{2} & =97.93[n F]  \tag{3.12}\\
C_{3} & =10.92[n F] \tag{3.13}
\end{align*}
$$

### 3.2.2 Nonlinear resistor

By removing one of the jumpers the nonlinear resistor part can be isolated from the total circuit. To measure the current through the nonlinear resistor a differential amplifier ( $A D 620$ ) in combination with a resistor $R_{m}$ is used. With the amplified voltage over the resistor $R_{m}$ and using Ohm's law it is possible to obtain the $v-i$ characteristic. A sinusoidal signal is used as input. Several experiments are performed and averaged to determine the positive and negative breakpoints as well as the slopes $G_{a}$ and $G_{b}$. The main observation is that the slopes $G_{b}$ in the positive and negative parts are equal. Furthermore there is a little difference between the breakpoints. This is due to different saturation points for positive and negative voltages of the used operational amplifiers.
One of the disadvantages of this measuring method is that the slopes are very sensitive to the total amplification factor, which consists of the differential amplifier itself times the value of $R_{m}$. The values in table 3.1 are determined using a (ideal) total amplification factor of 100 . The theoretical values for $G_{a}$ and $G_{b}$ are $-0.7591[\mathrm{mS}]$ and $-0.4096[\mathrm{mS}]$ respectively.
Another interesting phenomenon is illustrated in figure 3.8. It seems that there is some sort of hysteresis present in the nonlinear resistor, although this is investigated further in (Criens, 2005; Wouters, 2005) no specific explanation is found.


Figure 3.7: Measured characteristic of the nonlinear resistor with a $30[\mathrm{~Hz}]$ sinusoidal input.

Table 3.1: Parameters for the nonlinear resistor.

| Parameter | Value | Parameter | Value |
| :--- | ---: | :--- | ---: |
| $G_{a}$ | $-0.7665[\mathrm{mS}]$ | $B p^{+}$ | $1.81[\mathrm{~V}]$ |
| $G_{b}^{+}$ | $-0.4119[\mathrm{mS}]$ | $B p^{-}$ | $-1.75[\mathrm{~V}]$ |
| $G_{b}^{-}$ | $-0.4119[\mathrm{mS}]$ |  |  |



Figure 3.8: Measured characteristic of the nonlinear resistor with a $3500[\mathrm{~Hz}]$ sinusoidal input.

### 3.2.3 Inductor

The theoretical values, based on the used components and (3.9), are $0.0245[H]$ for the inductor and $22.40[\Omega]$ for the internal resistance $R_{0}$.
By removing the other jumper a resonance circuit is created, which makes it possible to perform a frequency response measurement.
The resonance circuit is shown in figure 3.9, where $V_{\text {in }}$ and $V_{\text {out }}$ are measured. $V_{\text {out }}$ is used as measurement for the current through the circuit. We assume again that the total amplification factor is equal to 100 , which gives the following relation for the total impedance of the resonance circuit

$$
\begin{equation*}
Z=\frac{V_{\text {in }}}{I_{\text {in }}} \approx 100 \frac{V_{\text {in }}}{V_{\text {out }}} . \tag{3.14}
\end{equation*}
$$

The theoretical impedance is given by

$$
\begin{align*}
Z & =\frac{\left(R_{0}+j \omega L\right) \frac{1}{j \omega C_{2}}}{R_{0}+j \omega L+\frac{1}{j \omega C_{2}}}+R_{m}  \tag{3.15}\\
Z & =\frac{C_{2} L R_{m} s^{2}+\left(L+R_{m} C_{2} R_{0}\right) s+R_{0}+R_{m}}{C_{2} L s^{2}+C_{2} R_{0} s+1} \tag{3.16}
\end{align*}
$$

Under the assumption that $C_{2}$ is known, we can fit a frequency response function and use the denominator of (3.16) to calculate $L$ and $R_{0}$. In this way the effect of the amplification factor is canceled since a gain only affects the nominator. The frequency response function is shown in figure 3.10 and the values for $L$ and $R_{0}$ are given in table 3.2.

Table 3.2: Parameters for the inductor.

| Parameter | Value |
| :--- | ---: |
| $L$ | $0.0243[\mathrm{H}]$ |
| $R_{0}$ | $37.29[\Omega]$ |



Figure 3.9: Resonance circuit.


Figure 3.10: Frequency response function.

### 3.3 Identification using filtering

A popular technique to estimate states (and parameters) in the presence of measurement (and system) noise is Kalman filtering. Kalman filters for linear systems produce optimal - i.e., the filter minimizes the variance of the difference between actual and estimated state - unbiased and consistent estimations of the systems states. Kalman filters can also be applied to nonlinear systems and are known as Extended Kalman filters. The main drawback of an extended filter is however that the filter is not anymore optimal in the above sense.

### 3.3.1 Augmented Continuous-Discrete filter

To identify the state of the double scroll circuit some remarks must be made. The differential equations are continuous in time, while measurements can only be taken at discrete times. This means that a continuous model is needed to estimate the states between measurements and a discrete update law should be used whenever a new measurement is available. This can be done by using a so-called continuous-discrete Kalman filter.

### 3.3.2 Parameters

Besides the states also some parameters are not accurately known, e.g. the slopes of the nonlinear resistor and the internal resistance $R_{0}$ of the inductor. The value of the variable resistor $R$ can be measured, however internal resistances of the print board itself are unknown. Since $R$ is used as bifurcation parameter we want to know its value as accurate as possible. Therefore this parameter should be estimated as well.
The nonlinear resistor is characterized by the slopes, $G_{a}, G_{b}^{+}$and $G_{b}^{-}$and the breakpoints $B_{p}^{+}$and $B_{p}^{-}$. We assume that the breakpoints are fixed, while the slopes cannot be measured accurately enough due to uncertainties in the measurement device. Therefore the state $\mathbf{x}$ is augmented to $\left[v_{1} v_{2} i G R_{0}-G_{a}^{-}-G_{b}^{+}-G_{b}^{-}\right]^{T}$ where $G=\frac{1}{R}$.

### 3.3.3 Algorithm

In this section the filter algorithm, based on (Gelb, 2001), is summarized. We start with a system and measurement model given by

$$
\begin{align*}
\dot{\mathbf{x}}(t) & =f(\mathbf{x}(t), t)+\mathbf{w}(t) & & \mathbf{w}(t) \sim N(0, Q(t))  \tag{3.17}\\
\mathbf{z}_{k} & =\mathbf{h}_{k}\left(\mathbf{x}\left(t_{k}\right)\right)+\mathbf{v}_{k} \quad k=1,2, \ldots & & v_{k} \sim N\left(0, R_{k}\right)  \tag{3.18}\\
t_{k} & =k \Delta t & & \tag{3.19}
\end{align*}
$$

where $\Delta t$ is the time between sample $k$ and $k+1 . \mathbf{w}(t)$ and $\mathbf{v}_{k}$ represents the process and measurement noise respectively. In absence of information about this noise, it is assumed to be white noise with zero mean and a normal distribution. The covariance matrices of the process and measurement noise are $Q(t)$ and $R_{k}$. In the ideal case, having a perfect model and no measurement noise, $\mathbf{w}(t)$ and $\mathbf{v}_{k}$ are absent.
The state estimation and error covariance propagation are as follows

$$
\begin{align*}
\dot{\hat{\mathbf{x}}}(t) & =f(\hat{\mathbf{x}}(t), t)  \tag{3.20}\\
\dot{P}(t) & =F(\hat{\mathbf{x}}(t), t) P(t)+P(t) F^{T}(\hat{\mathbf{x}}(t), t)+Q(t) \tag{3.21}
\end{align*}
$$

where $\hat{\mathbf{x}}$ is the filtered estimate of $\mathbf{x}$ and $F(\hat{\mathbf{x}}(t), t)$ is the Jacobian matrix defined as

$$
\begin{equation*}
F(\hat{\mathbf{x}}(t), t)=\left.\frac{\partial f(\mathbf{x}(t), t)}{\partial \mathbf{x}(t)}\right|_{\mathbf{x}(t)=\hat{\mathbf{x}}(t)} \tag{3.22}
\end{equation*}
$$

The covariance matrix $P(t)$ is a symmetric matrix representing an approximation of the true covariance.
When a new measurement is available the state estimation and error covariance matrix are updated using the following update laws

$$
\begin{align*}
\hat{\mathbf{x}}_{k}(+) & =\hat{\mathbf{x}}_{k}(-)+K_{k}\left[\mathbf{z}_{k}-h_{k}\left(\hat{\mathbf{x}}_{k}(-)\right)\right]  \tag{3.23}\\
P_{k}(+) & =\left[I-K_{k} H_{k}\left(\hat{\mathbf{x}}_{k}(-)\right)\right] P_{k}(-)  \tag{3.24}\\
K_{k} & =P_{k}(-) H_{k}^{T}\left(\hat{\mathbf{x}}_{k}(-)\right)\left[H_{k}\left(\hat{\mathbf{x}}_{k}(-)\right) P_{k}(-) H_{k}^{T}\left(\hat{\mathbf{x}}_{k}(-)\right)+R_{k}\right]^{-1} \tag{3.25}
\end{align*}
$$

where $\hat{\mathbf{x}}_{k}(-)$ is the estimate of $\mathbf{x}(t)$ at time $t_{k}$ before measurement $\mathbf{z}_{k}$ is taken into account and $\hat{\mathbf{x}}_{k}(+)$ is the estimation at time $t_{k}$ after this measurement. $K_{k}$ is the gain matrix and $I$ is the identity matrix. $P_{k}(-)$ and $P_{k}(+)$ are the discrete error covariance before and after the update. $H_{k}\left(\hat{\mathbf{x}}_{k}(-)\right)$ is the Jacobian matrix defined as

$$
\begin{equation*}
H_{k}\left(\hat{\mathbf{x}}_{k}(-)\right)=\left.\frac{\partial h_{k}\left(\mathbf{x}\left(t_{k}\right)\right)}{\partial \mathbf{x}\left(t_{k}\right)}\right|_{\mathbf{x}\left(t_{k}\right)=\hat{\mathbf{x}}_{k}(-)} \tag{3.26}
\end{equation*}
$$

The following model for the electric circuit is used (see chapter 2)

$$
\begin{align*}
C_{1} \dot{v}_{1} & =G\left(v_{2}-v_{1}\right)-f\left(v_{1}\right) \\
C_{2} \dot{v}_{2} & =G\left(v_{1}-v_{2}\right)+i_{L}  \tag{3.27}\\
L \dot{i}_{L} & =-v_{2}-R_{0} i_{L},
\end{align*}
$$

which can be divided into three linear parts by writing out $f\left(v_{1}\right)$ :

$$
\begin{align*}
v_{1} & \geq B_{p}^{+}  \tag{3.28}\\
C_{1} \dot{v}_{1} & =G\left(v_{2}-v_{1}\right)-\left(G_{b}^{+} v_{1}+\left(G_{a}-G_{b}^{+}\right) B_{p}^{+}\right) \\
B_{p}^{-} & <v_{1}<B_{p}^{+}  \tag{3.29}\\
C_{1} \dot{v}_{1} & =G\left(v_{2}-v_{1}\right)-G_{a} v_{1} \\
v_{1} & \leq B_{p}^{-}  \tag{3.30}\\
C_{1} \dot{v}_{1} & =G\left(v_{2}-v_{1}\right)-\left(G_{b}^{-} v_{1}+\left(G_{a}-G_{b}^{-}\right) B_{p}^{-}\right)
\end{align*}
$$

The advantage of this model, compared to the dimensionless model (2.11), is that a measurement sample can be used directly, whereas the dimensionless time $\tau$ depends on the parameter $R$, which is to be estimated.
Since the piecewise linear model consists of three different linear systems the Jacobian of $f(\mathbf{x}(t), t)$ (3.22) is not continuous and a modification is required. There are multiple ways to solve the Ricatti equation (3.21).
In the next section we introduce an extension of a 'switching' Kalman filter as presented in (Cruz and Nijmeijer, 2000; Cruz et al., 2001). Although this filter switches between linear parts it is not an optimal filter. Three different methods will be used and compared.

### 3.3.4 Methods

The first method consists of one Ricatti equation. Although the linearization of $f(\hat{\mathbf{x}}(t), t)$, see (3.27), is evaluated at the current state of the system, it gives three different Jacobian matrices depending on the state $v_{1}$. This gives a 'jump' in the Ricatti equation.
Instead of one Ricatti equation, three Ricatti equations can be used to prevent the jump. There are two ways to solve these equations. One way, method 2 , is to solve all three equations at the same time and chose the right update gain vector, which depends on the state $v_{1}$, to update the state. The Ricatti equations are updated using the three different gain vectors. The problem here is however that the linearization of $f(\hat{\mathbf{x}}(t), t)$ is not a good approximation for two of the three regions. Therefore the update for two Ricatti equations is not correct.
The other way, method 3, is to integrate the system and just one Ricatti equation between two samples. Then update this Ricatti equation and holding the other two at their previous values. Now the linearization of $f(\hat{\mathbf{x}}(t), t)$ is correct, but the Ricatti equations are not continuous in time anymore.

### 3.3.5 Results

To compare the three methods one data set is used and the settings for $Q$ and $R_{k}$ are kept constant. The results can be seen in figure 3.11. In this figure a part of a measured time signal of $v_{1}$ together with the errors between measurement and estimation for the three methods is shown. In the unstable area, the area between the two horizontal lines in the measured time signal, the error remains large compared to the stable regions. In table 3.3 the RMS value of the error is given. It can be seen that the three Ricatti equations perform a little bit better. For the filter, where one Ricatti equation is solved between two samples, the obtained parameters are presented in table 3.4.

### 3.4 Comparison filters for measurement and process noise

### 3.4.1 Measurement noise

To compare the two filters with the three Ricatti equations we investigate the influence of different levels of measurement and process noise. For the measurement noise a diagonal matrix $R_{k}=$ $\operatorname{diag}\left(\left[R_{1} R_{2}\right]\right)$ is used. To compare the influence different values for $R_{k}$ are evaluated ranging from $1^{-8}$ till $1^{-4}$. The comparison is made by using the 2 and infinity norm of the error between measurement and prediction of $v_{1}$. It can be seen in the figure 3.12 that if $R_{k}$ is too small the measurements are not considered reliable, also the filter is not stable anymore.
When $R_{k}$ is too big, the filter tries to follow all measurement samples and the model information is ignored. There seems to be an optimum around $R_{k}=1^{-6}$. For all simulations the process noise is chosen as $Q=\operatorname{diag}\left(\left[1^{-5} 1^{-5} 1^{-5} 000000\right]\right)$.

### 3.4.2 Process noise

To see the influence of process noise on the differential equations $Q$ is varied, while the process noise for the parameter equations is set to zero. $R_{k}$ is chosen as $R_{k}=\operatorname{diag}\left(\left[1^{-6} 1^{-6}\right]\right)$ in this case. It follows that if there is too much process noise the model is not considered reliable and again the filter tries to follow every measurement sample.

### 3.5 Discussion

The described identification methods can be used to obtain values for the parameters. Several simple experiments are used to determine parameters of different parts of the circuit. Further three Kalman filters are used to estimate parameters based on measurements of the whole circuit. However if the obtained parameters from measurements or the Kalman filter are used in a numerical simulation, the solution always differs from experimental data. There are several reasons why this happens. We know that system (2.1) is a chaotic system and is therefore sensitive to parameter values. The initial condition of the current is unknown, which also plays an important role in the final trajectory and the tolerances of the used solver cause unwanted rounding errors. An illustrative example is visualized in figure 3.14. To give an explanation of the radical difference in figure 3.14 we take a closer look at the experiment and the Kalman filter. If we look at one period of the periodic solution, see figure 3.15, it can be seen that the error in the two outer regions is relatively small compared to the error in middle region, especially when the trajectory goes from left to right or visa versa.
However there is always a nonzero error and this means that trajectories will always diverge from a measured trajectory if we do not use measurement data to correct the computed numerical trajectory.
The closed orbit shown in figure 3.14 is not robust to parameter variations and therefore it is practically impossible to obtain this periodic solution by numerical integration of system (2.1). The result is the trajectory as shown in figure 3.14 (c).
So we need measurement information to follow experimental trajectories and therefore we end this chapter with some remarks about the Kalman filter. It is already mentioned that the error in the two outer regions is smaller than the error in middle region. This indicates that results obtained by the Kalman filter might improve if the model for the region $B_{p}^{-}<v_{1}<B_{p}^{+}$is changed. There is a very small current offset at zero volt in the measured nonlinear resistor, see figure 3.7, which is not modeled. Further we approximated the location of the breakpoints by averaging several experiments and used this as the switching between the three regions. A filter that is capable of estimating the breakpoints (see (Santoboni and Nijmeijer, 2001)) could improve the results as well.
Finally, and probably the most important, the dynamic behavior of the electric components is neglected in the used model (3.27). The influence of the operational amplifiers is already pointed out in (Kennedy, 1992). This can be investigated further by performing circuit simulations, such
as SPICE (Quarles et al., 1993).
In the next chapters we will focus more on the qualitative behavior instead of the quantitative behavior.

Table 3.3: Root mean square values of the error signals.

| Method 1 | Method 2 | Method 3 |
| :---: | :---: | :---: |
| $0.0464[\mathrm{~V}]$ | $0.0320[\mathrm{~V}]$ | $0.0262[\mathrm{~V}]$ |



Figure 3.11: Experimental data and error between measurement and estimation of $v_{1}$ for the three methods. The first error is for one Ricatti equation. The second error is obtained by solving all three Ricatti equations at once. The last error plot is obtained by solving one Ricatti equation between two samples.

Table 3.4: Parameters obtained by Kalman filter, std stands for the standard deviation.

| Parameter | Mean value | std |
| :--- | ---: | ---: |
| $R_{0}$ | $25.70[\Omega]$ | $0.20[\Omega]$ |
| $G_{a}$ | $-0.7385[\mathrm{mS}]$ | $0.0049[\mathrm{mS}]$ |
| $G_{b}^{+}$ | $-0.4018[\mathrm{mS}]$ | $0.0013[\mathrm{mS}]$ |
| $G_{b}^{-}$ | $-0.4020[\mathrm{mS}]$ | $0.0024[\mathrm{mS}]$ |



Figure 3.12: Influence of measurement noise levels.


Figure 3.13: Influence of process noise levels.


Figure 3.14: Difference between experiment, Kalman filter and simulation.


Figure 3.15: $\diamond$ indicate measured data points and $*$ indicate filtered data.

## Chapter 4

## Experiments

With the electronic realization as described in the previous chapter several experiments are conducted. In this chapter the qualitative numerical results obtained in chapter 2 are validated experimentally. For the experiments the same approach as in chapter 2 is followed. This means we start at $R \approx 2000[\Omega]$ and then decrease the resistor value. The value of $R$ is obtained by the number of rotations of the potentiometer and using (3.10).

### 4.1 Single scroll trajectories

One of the disadvantages of the setup is that is not possible to choose initial conditions. The setup is initialized when the voltage supply of the operational amplifiers is turned on. Although the initial conditions are random, often the voltage across capacitor $C_{1}$ is negative. Therefore we present the single scroll solutions only for negative voltages $v_{1}$. However all solutions are also observed for positive $v_{1}$.
The bifurcation point at which the stable equilibrium point loses its stability occurs at $R \approx 1902[\Omega]$. In figure 4.1 projections of $v_{2}$ versus $v_{1}$ can be seen. In these figures the trajectories intersect with each other due to the projection of the three dimensional state space on the $v_{1}-v_{2}$ plane. The projection of the period three limit cycle is probably the most interesting, since this is a strong indication that the experimental setup is capable of generating chaos.
The 'birth' of the double scroll occurs at $R \approx 1823[\Omega]$. This is visualized in figure 4.1 (e). In this experiment only one of over two hundred oscillations takes place in the positive part of the $v_{1}-v_{2}$ plane.

### 4.2 Double scroll trajectories

When we further decrease the resistor value, the double scroll is observed. In figure 4.2 several periodic solutions are shown. The so-called hole filling trajectory is also experimentally observed for $R \approx 1772[\Omega]$.
Another special trajectory occurs at $R \approx 1542[\Omega]$. To make that more clear the measured data for this experiment is used in the described Kalman filter to obtain parameter values. With the identified parameters the 1D Poincaré map can be calculated.
This is done by using four different time series in the Kalman filter. For the second iteration the final parameter estimates of the first time series are used as initial conditions. After the fourth iteration the parameter estimates, see table 4.1, are used to compute the 1D Poincaré map. It follows from this map that the trajectory is related to the homoclinic orbit. This homoclinic orbit is another confirmation that the experimental setup is a chaotic system.
It can also be seen in figure 4.4 that this happens just before the double scroll collides with the unstable closed orbit. If the resistor value is decreased further, the trajectory is not bounded by the attractor anymore and $v_{1}$ exceeds the saturation level of the operational amplifiers. The
negative resistance becomes positive (passive) and the resulting trajectory is a large stable limit cycle, see chapter 2.


Figure 4.1: Projections of $v_{2}$ versus $v_{1}$ for different values of $R$.


Figure 4.2: Projections of $v_{2}$ versus $v_{1}$ for different values of $R$.

Table 4.1: Parameters used to compute the 1D Poincaré map.

| Parameter | Value | Parameter | Value |
| :--- | ---: | :--- | ---: |
| $C_{1}$ | $10.90[\mathrm{nF}]$ | $C_{2}$ | $97.93[\mathrm{nF}]$ |
| $G_{a}$ | $-0.7333[\mathrm{mS}]$ | $G_{b}^{+}$ | $-0.4005[\mathrm{mS}]$ |
| $R$ | $1593[\Omega]$ | $R_{0}$ | $25.64[\Omega]$ |
| $L$ | $24.39[\mathrm{mH}]$ | $B_{p}^{+}$ | $1.8191[\mathrm{~V}]$ |



Figure 4.3: Hole filling orbit for $R \approx 1772[\Omega]$.


Figure 4.4: Homoclinic orbit and estimation of the Poincaré map.

## Chapter 5

## Synchronization

As last part of this thesis synchronization of multiple Chua circuits is investigated.
With synchronization we mean that two or more circuits follow the same trajectory. This might seem strange for chaotic systems, since trajectories of two identical systems with nearly identical initial conditions diverge exponentially. However it is possible for chaotic systems to synchronize with each other under certain conditions.
In this chapter two different methods are used, master-slave synchronization and mutual synchronization. In the case of mutual synchronization an article, which has been written during this graduation project, is included. This article is accompanied with additional numerical simulations.

### 5.1 Master-slave synchronization

Synchronization of chaotic systems can be found in (Pecora and Carroll, 1990; Carroll and Pecora, 1991). The general scheme for synchronization here is to take a driving system, create a subsystem and drive this and a duplicate of this subsystem, called a response system, with signals from the drive system. A related method is master-slave synchronization. This method consists of a master and a slave system, which follows the master. This can be obtained if we take two identical systems

$$
\begin{align*}
& \dot{\mathbf{x}}_{1}=\mathbf{F}\left(\mathbf{x}_{1}\right)  \tag{5.1}\\
& \dot{\mathbf{x}}_{2}=\mathbf{F}\left(\mathbf{x}_{2}\right)
\end{align*}
$$

and apply an unidirectional coupling to the second system

$$
\begin{align*}
& \dot{\mathbf{x}}_{1}=\mathbf{F}\left(\mathbf{x}_{1}\right)  \tag{5.2}\\
& \dot{\mathbf{x}}_{2}=\mathbf{F}\left(\mathbf{x}_{2}\right)+f\left(\mathbf{x}_{1}, \mathbf{x}_{2}\right) .
\end{align*}
$$

Asymptotic synchronization occurs when $\lim _{t \rightarrow \infty}\left|\mathbf{x}_{1}(t)-\mathbf{x}_{2}(t)\right|=0$. This can be achieved by using continuous feedback (Ott et al., 1990; Pyragas, 1992). A schematic overview can be seen in figure 5.1. Consider (5.2) where some of the state variables can be measured. The chaotic systems are unidirectionally coupled in such a way that the difference between the signals $\mathbf{x}_{1 i}$ and $\mathbf{x}_{2 i}$ is used as a linear control signal, i.e.

$$
\begin{equation*}
f\left(\mathbf{x}_{1}, \mathbf{x}_{2}\right)=K\left(\mathbf{x}_{1 i}-\mathbf{x}_{2 i}\right) \tag{5.3}
\end{equation*}
$$

which is applied to the second system as a negative feedback. The parameter $K>0$ is an adjustable gain. An important feature of (5.3) is that it does not change the solution of (5.1).

### 5.1.1 Numerical results

We use the same parameters as in table 2.1 and choose $R$ equal to $1775[\Omega]$ such that the circuit operates at the double scroll attractor. The following master-slave configuration, as described in


Figure 5.1: Schematic overview of master-slave synchronization using feedback.
(Kapitaniak, 2000), is used:

$$
\begin{align*}
C_{1} \dot{v}_{11} & =G\left(v_{12}-v_{11}\right)-f\left(v_{11}\right) \\
C_{2} \dot{v}_{12} & =G\left(v_{11}-v_{12}\right)+i_{L 1} \\
\dot{i}_{L 1} & =-v_{12}-R_{0} i_{L 1} \\
C_{1} \dot{v}_{21} & =G\left(v_{22}-v_{21}\right)-f\left(v_{21}\right)+K_{1}\left(v_{11}-v_{21}\right)  \tag{5.4}\\
C_{2} \dot{v}_{22} & =G\left(v_{21}-v_{22}\right)+i_{L 2}+K_{2}\left(v_{12}-v_{22}\right) \\
L i_{L 2} & =-v_{22}-R_{0} i_{L 2}
\end{align*}
$$

with $K_{i}=\frac{1}{R_{c i}}$ the coupling strength and $R_{c i}$ are variable resistors, see figure 5.2. The asymptotic synchronized situation is defined as

$$
\begin{equation*}
\lim _{t \rightarrow \infty}\left|\mathbf{x}_{1}(t)-\mathbf{x}_{2}(t)\right|=0 \tag{5.5}
\end{equation*}
$$

with $\mathbf{x}_{i}=\left[\begin{array}{lll}v_{1 i} & v_{2 i} & i_{L i}\end{array}\right]$.
First synchronization between identical systems is considered. We consider coupling through $K_{1}$ or $K_{2}$ respectively. It can be seen in figure 5.3 that synchronization occurs if $R_{c 1}$ does not exceed $4100[\Omega]$. The upper bound for synchronization using $K_{2}$ is approximately $R_{c 2}=2500[\Omega]$. If $K_{i}$ is increased the circuits synchronize in a shorter time period.


Figure 5.2: Master-slave synchronization for Chua circuits with unidirectional coupling.


Figure 5.3: $R_{c 1}=4100 \Omega$. Synchronization error for identical systems.


Figure 5.4: $R_{c 2}=2500 \Omega$. Synchronization error for identical systems.

### 5.1.2 Non-identical systems

In practice the circuits are build using standard components and due to tolerances of these components every circuit is slightly different. Therefore perfect synchronization is not possible and we introduce a form of practical asymptotic synchronization as

$$
\begin{equation*}
\lim _{t \rightarrow \infty}\left|\mathbf{x}_{1}(t)-\mathbf{x}_{2}(t)\right| \leq \delta, \tag{5.6}
\end{equation*}
$$

for some fixed, preferably small, $\delta>0$.
To simulate the effect of component tolerances some of the parameters are changed, see table 5.1. The results are shown in the figures 5.5 and 5.6.

Table 5.1: Changed parameters for system 1 and 2.

| System 1 | Value | System 2 | Value |
| :--- | ---: | :--- | ---: |
| $C_{1}$ | $10.90[\mathrm{nF}]$ | $C_{1}$ | $10.80[\mathrm{nF}]$ |
| $C_{2}$ | $97.93[\mathrm{nF}]$ | $C_{2}$ | $99.60[\mathrm{nF}]$ |
| $R$ | $1775[\Omega]$ | $R$ | $1778[\Omega]$ |



Figure 5.5: $R_{c 1}=10 \Omega$. Synchronization error for different systems.


Figure 5.6: $R_{c 2}=10 \Omega$. Synchronization error for different systems.

### 5.1.3 Experimental results

For the experiments the value of the variable resistor is fixed for both systems at $1775[\Omega]$. In the first experiment both circuits operate as free systems. Although the projections seems similar, see figure 5.7, there is no synchronization. This is shown in figure 5.8. If there is synchronization, it can be seen as a diagonal line in the phase portrait.
In the second experiment the upper bound for $R_{c 1}$ is determined, while for the third experiment $R_{c 1}$ is chosen as low as possible to create the strongest coupling factor. The same procedure is followed for $R_{c 2}$.
It can be seen that synchronization occurs for a larger range of $R_{c 1}$ than $R_{c 2}$ and the synchronization error is also smaller. Despite the different values for the resistors again the results are qualitative comparable with the obtained simulation results.
Another interesting phenomenon can be seen in figure 5.12. It seems that there occurs intermittency. Although most of the time the slave is synchronized, it bursts out of synchronization for very short time intervals. This occurs also in numerical simulations if the coupling factor $R_{c 2}$ exceeds $2500[\Omega]$. Finally we remark that if the coupling strength is decreased the systems start to desynchronize and above a certain threshold value the slave is no longer bounded by the attractor.


Figure 5.7: Projections of both systems without coupling.


Figure 5.8: Synchronization phase portrait and error without coupling.


Figure 5.9: Synchronization phase portrait and error for $R_{c 1} \approx 3288[\Omega]$.


Figure 5.10: Synchronization phase portrait and error for $R_{c 1} \approx 2[\Omega]$.


Figure 5.11: Synchronization phase portrait and error for $R_{c 2} \approx 200[\Omega]$.


Figure 5.12: Synchronization phase portrait and error for $R_{c 2} \approx 2[\Omega]$.

### 5.2 Mutual synchronization

In this section a different synchronization scheme is applied, instead of a unidirectional coupling (5.2) we use a bidirectional coupling

$$
\begin{align*}
& \dot{\mathbf{x}}_{1}=\mathbf{F}\left(\mathbf{x}_{1}\right)+f\left(\mathbf{x}_{1}, \mathbf{x}_{2}\right) \\
& \dot{\mathbf{x}}_{2}=\mathbf{F}\left(\mathbf{x}_{2}\right)+f\left(\mathbf{x}_{1}, \mathbf{x}_{2}\right) . \tag{5.7}
\end{align*}
$$

A schematic overview can be seen in figure 5.13. Besides synchronization of two circuits we also investigate synchronization of multiple circuits. When a network of systems is considered it is possible to achieve so-called partial synchronization in the network. Under partial synchronization some systems in the network synchronize, while other do not. Recently partial synchronization in networks of identical systems is receiving particular interest in literature (Hasler et al., 1998; Zhang et al., 2001; Pogromsky et al., 2002). In the following article we present obtained experimental synchronization results for different topologies, with a maximum of four Chua circuits.


Figure 5.13: Schematic overview of mutual synchronization using diffusive coupling.

# PARTIAL SYNCHRONIZATION OF DIFFUSIVELY COUPLED CHUA SYSTEMS: AN EXPERIMENTAL CASE STUDY 

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#### Abstract

In this paper partial synchronization of diffusively coupled Chua systems is presented. Partial synchronization is defined as the situation where some circuits synchronize with each other, while others do not. An experimental setup, consisting of maximal four Chua circuits operating in the double scroll regime, is used to show the existence of linear invariant manifolds corresponding to the partial synchronized state. Copyright (c) 2006 IFAC


Keywords: Synchronization, Electrical networks, Chaotic behavior, Network topology, Circuits

## 1. INTRODUCTION

Synchronization of coupled dynamical systems receives much attention in literature. One of the reasons for this is that synchronization can be found in several fields such as nature (Strogatz and Stewart, 1993), brain dynamics (Gray, 1994) and robotics (Nijmeijer and Rodrigues-Angeles, 2003). Also, the potential use of synchronization in communication and coordination forms a major reason for this interest. Recently partial synchronization in networks of identical systems is receiving particular interest. Some examples of partial synchronization can be found in (Hasler et al., 1998; Zhang et al., 2001; Pogromsky et al., 2002).
Although there are many papers describing global synchronization of a network of coupled Chua circuits (Wu and Chua, 1995; Matías et al., 1997; Sánchez et al., 2000), less attention sofar has been devoted to experimental results for bidirectional coupled systems. In this paper attention will be drawn to partial synchronization of Chua circuits. Partial synchronization is defined as the situation
where some circuits synchronize with each other, while others do not. It is shown that under certain conditions it is possible to obtain partial synchronization of diffusively coupled Chua circuits. An experimental setup consisting of four coupled Chua circuits is build to show the possibility of partial synchronization. The experimental results obtained qualitatively confirm simulation results. The remainder of this paper is organized as follows. In section 2 some preliminaries about the used notation are given. Further passive and convergent systems are described and the conditions for partial synchronization are stated. Section 3 deals with the experimental setup that is used. In section 4 and 5 synchronization of two and three diffusively coupled systems is shown, while in section 6 global and partial synchronization of four diffusively coupled systems is presented and discussed. Finally conclusions are drawn in section 7.

## 2. PRELEMINARIES

First a mathematical description for a network of coupled systems is introduced by adopting the notation used in (Pogromsky et al., 2002). A general system description for $k$ identical systems is given by

$$
\begin{equation*}
\dot{\mathbf{x}}_{i}=\mathbf{f}\left(\mathbf{x}_{i}\right)+\mathbf{B u} u_{i}, \quad \mathbf{y}_{i}=\mathbf{C} \mathbf{x}_{i} \tag{1}
\end{equation*}
$$

where $\mathbf{f}$ is a vector field, $i=1, \ldots, k, \mathbf{x}_{i}(t) \in \mathbb{R}^{n}$ is the state of the $i$ th system, $\mathbf{u}_{i}(t) \in \mathbb{R}^{m}$ and $\mathbf{y}_{i}(t) \in \mathbb{R}^{m}$ are the input and output of the $i$ th system, while $\mathbf{B}, \mathbf{C}$ are matrices of appropriate dimension.
The $k$ systems are coupled through linear outputs

$$
\begin{equation*}
\mathbf{u}_{i}=-\gamma_{i 1}\left(\mathbf{y}_{i}-\mathbf{y}_{1}\right)-\gamma_{i 2}\left(\mathbf{y}_{i}-\mathbf{y}_{2}\right)-\cdots-\gamma_{i k}\left(\mathbf{y}_{i}-\mathbf{y}_{k}\right) \tag{2}
\end{equation*}
$$

By defining the symmetric $k \times k$ matrix $\boldsymbol{\Gamma}$ as

$$
\boldsymbol{\Gamma}=\left(\begin{array}{cccc}
\sum_{i=2}^{k} \gamma_{1 i} & -\gamma_{12} & \cdots & -\gamma_{1 k}  \tag{3}\\
-\gamma_{21} & \sum_{i=1, i \neq 2}^{k} \gamma_{2 i} & \cdots & -\gamma_{2 k} \\
\vdots & \vdots & \ddots & \vdots \\
-\gamma_{k 1} & -\gamma_{2 k} & \cdots & \sum_{i=1}^{k-1} \gamma_{1 i}
\end{array}\right)
$$

the collection of $k$ systems, with the matrix $\boldsymbol{\Gamma}$ and feedback $\mathbf{u}_{i}$, can be rewritten as

$$
\begin{equation*}
\dot{\mathbf{x}}=\mathbf{F}(\mathbf{x})+\left(\mathbf{I}_{k} \otimes \mathbf{B}\right) \mathbf{u}, \quad \mathbf{y}=\left(\mathbf{I}_{k} \otimes \mathbf{C}\right) \mathbf{x} \tag{4}
\end{equation*}
$$

with the feedback

$$
\begin{equation*}
\mathbf{u}=-\left(\boldsymbol{\Gamma} \otimes \mathbf{I}_{m}\right) \mathbf{y} \tag{5}
\end{equation*}
$$

where $\mathbf{x}=\operatorname{col}\left(\mathbf{x}_{1}, \ldots, \mathbf{x}_{k}\right), \mathbf{F}(\mathbf{x})=\operatorname{col}\left(\mathbf{f}\left(\mathbf{x}_{1}\right)\right.$, $\left.\ldots, \mathbf{f}\left(\mathbf{x}_{k}\right)\right) \in \mathbb{R}^{k n}, \mathbf{y}=\operatorname{col}\left(\mathbf{y}_{1}, \ldots, \mathbf{y}_{k}\right)$ and $\mathbf{u}=\operatorname{col}\left(\mathbf{u}_{1}, \ldots, \mathbf{u}_{k}\right) \in \mathbb{R}^{k m}$. The notation $\operatorname{col}\left(\mathbf{x}_{1}, \ldots, \mathbf{x}_{k}\right)$ stands for the column vector composed of the vectors $\mathbf{x}_{1}, \ldots, \mathbf{x}_{k}$. The notation $\otimes$ stands for the Kronecker product.

A system

$$
\begin{equation*}
\dot{\mathbf{x}}=\mathbf{f}(\mathbf{x}, \mathbf{u}), \quad \mathbf{y}=\mathbf{h}(\mathbf{x}) \tag{6}
\end{equation*}
$$

is called passive, see (Willems, 1972), if the following inequality holds

$$
\begin{equation*}
\frac{d}{d t} V(\mathbf{x})=\frac{\partial V(\mathbf{x})}{\partial \mathbf{x}} \mathbf{f}(\mathbf{x}, \mathbf{u}) \leq \mathbf{y}^{T} \mathbf{u} \tag{7}
\end{equation*}
$$

where $V(\mathbf{x})$ is a nonnegative function (storage function) defined on $\mathbb{R}^{n}$, for which $V(0)=0$. If the dissipation inequality (7) is satisfied only for x lying outside some ball

$$
\begin{equation*}
\dot{V}(\mathbf{x}, \mathbf{u}) \leq \mathbf{y}^{T} \mathbf{u}-H(\mathbf{x}) \tag{8}
\end{equation*}
$$

where the function $H: \mathbb{R}^{n} \rightarrow \mathbb{R}$ is nonnegative outside some ball

$$
\begin{equation*}
\exists \rho>0, \quad|\mathbf{x}| \geq \rho \Rightarrow H(\mathbf{x}) \geq 0 \tag{9}
\end{equation*}
$$

then the system is semipassive, see (Pogromsky et al., 2002).

Consider a system

$$
\begin{equation*}
\dot{\mathbf{z}}=\mathbf{q}(\mathbf{z}, w(t)) \tag{10}
\end{equation*}
$$

with $\mathbf{z} \in \mathbb{R}^{l}$, driven by an external signal $w(t)$ taking values from a compact set. The system (10) is called convergent if for any bounded input $w(t)$ the solution of (10) converges to a solution $\mathbf{z}_{w}(t)$, in other words, the solution of (10) will forget their specific initial condition. If there exists a positive definite symmetric $l \times l$ matrix $\mathbf{P}$ such that all eigenvalues $\lambda_{i}(\mathbf{Q})$ of the symmetric matrix
$\mathbf{Q}(\mathbf{z}, w)=\frac{1}{2}\left[\mathbf{P}\left(\frac{\partial \mathbf{q}}{\partial \mathbf{z}}(\mathbf{z}, w)\right)+\left(\frac{\partial \mathbf{q}}{\partial \mathbf{z}}(\mathbf{z}, w)\right)^{T} \mathbf{P}\right]$
are negative and separated from zero, such that

$$
\lambda_{i}(\mathbf{Q}(\mathbf{z}, w)) \leq \epsilon<0
$$

with $\epsilon>0$ and $i=1 \ldots l$ for all $\mathbf{z}, w \in \mathbb{R}^{l}$, then system (10) is convergent, cf (Pavlov et al., 2004).
If the network contains repeating patterns, the permutation of some elements of $\boldsymbol{\Gamma}$ leave the network invariant. Such a permutation matrix $\boldsymbol{\Pi}$ is a symmetry for the network if $\Pi$ commutes with $\boldsymbol{\Gamma}$, i.e. $\boldsymbol{\Pi} \boldsymbol{\Gamma}-\boldsymbol{\Gamma} \boldsymbol{\Pi}=0$. A permutation matrix $\Pi$ commuting with $\boldsymbol{\Gamma}$ defines a linear invariant manifold of the closed loop system (4) and (5) as

$$
\begin{equation*}
\operatorname{ker}\left(\mathbf{I}_{k n}-\boldsymbol{\Pi} \otimes \mathbf{I}_{n}\right) \tag{13}
\end{equation*}
$$

The stability of such manifolds depends on the asymptotic stability of sets. Due to converse Lyapunov theorem, e.g. (Lin et al., 1996), the asymptotic stability of a set is equivalent to the existence of a scalar storage function $V$, which is zero only on the set and decays along the trajectories otherwise. In the context of the coupled systems (1,2) a Lyapunov function should be found as a sum of two functions. The first function depends on the input-output relations of the systems (1), while the second function depends on the interacting due to the coupling of systems.
Under the assumption that the matrix $\mathbf{C B}$ is nonsingular (and positive definite) a linear coordinate transformation $\mathbf{x}_{i} \rightarrow\left(\mathbf{y}_{i}, \mathbf{z}_{i}\right)$ exists such that

$$
\begin{equation*}
\dot{\mathbf{z}}_{i}=\mathbf{q}\left(\mathbf{z}_{i}, \mathbf{y}_{i}\right), \quad \dot{\mathbf{y}}_{i}=\mathbf{g}\left(\mathbf{z}_{i}, \mathbf{y}_{i}\right)+\mathbf{C B u} u_{i} \tag{14}
\end{equation*}
$$

where $\mathbf{z}_{i} \in \mathbb{R}^{n-m}$ and $\mathbf{q}$ and $\mathbf{g}$ are vector functions. Then the stability of the manifolds given by (13) is determined by the following theorem.
Theorem 1 (Pogromsky et al., 2002). Let $\lambda^{\prime}$ be the minimal eigenvalue of $\boldsymbol{\Gamma}$ under restriction that the eigenvectors of $\boldsymbol{\Gamma}$ are taken from the set range $\left(\mathbf{I}_{k}-\boldsymbol{\Pi}\right)$. Suppose that:

1. Each individual system (1) is strictly semipassive with respect to the input $\mathbf{u}_{i}$ and output
$\mathbf{y}_{i}$ with a radially unbounded storage function $V\left(\mathbf{x}_{\mathbf{i}}, \mathbf{u}_{\mathbf{i}}\right)$.
2. There exists a positive definite matrix $\mathbf{P}$ such that inequality (12) holds for some $\epsilon>0$ for the matrix $\mathbf{Q}$ defined as in (11) for $\mathbf{q}$ as in (14).
Then all solutions of the diffusive cellular network (4) and (5) are ultimately bounded and there exists a positive $\bar{\lambda}$ such that if $\lambda^{\prime}>\bar{\lambda}$ the set $\operatorname{ker}\left(\mathbf{I}_{k n}-\boldsymbol{\Pi} \otimes \mathbf{I}_{n}\right)$ contains a globally asymptotically stable compact subset.

## 3. SETUP

An experimental setup consisting of a network of 4 Chua circuits is used, shown in figure 1. Consider


Fig. 1. Schematic layout of four symmetrically coupled Chua circuits.
the well known system description of a single Chua circuit (Matsumoto, 1984)

$$
\begin{align*}
C_{1} \dot{x}_{i 1} & =G\left(x_{i 2}-x_{i 1}\right)-f\left(x_{i 1}\right) \\
C_{2} \dot{x}_{i 2} & =G\left(x_{i 1}-x_{i 2}\right)+x_{i 3}  \tag{15}\\
L \dot{x}_{i 3} & =-x_{i 2}-R_{0} x_{i 3}
\end{align*}
$$

with $G=\frac{1}{R}$ and the function $f\left(x_{i 1}\right)$ defined as

$$
\begin{equation*}
G_{b} x_{i 1}+\frac{1}{2}\left(G_{a}-G_{b}\right)\left(\left|x_{i 1}+B_{p}\right|-\left|x_{i 1}-B_{p}\right|\right) \tag{16}
\end{equation*}
$$

In these equations the variables $x_{i 1}$ and $x_{i 2}$ are the voltages across the capacitors, $C_{1}$ and $C_{2}, x_{i 3}$ is the current flowing through the inductor $L$, which has an internal resistance $R_{0} . G_{a}$ and $G_{b}$ are the conductances of the piecewise characteristic for $\left|x_{i 1}\right|<B_{p}$ and $\left|x_{i 1}\right| \geq B_{p}$ respectively. $B_{p}$ is voltage of the breakpoint. Measurements of $x_{i 1}$ and $x_{i 2}$ are available. For the coupling between systems the matrices $\mathbf{B}$ and $\mathbf{C}$ are as follows

$$
\mathbf{B}=\left[\begin{array}{lll}
1 & 0 & 0
\end{array}\right]^{T} \quad \mathbf{C}=\left[\begin{array}{lll}
1 & 0 & 0 \tag{17}
\end{array}\right]
$$

The coupling strength between systems is controlled by four variable resistors. The nonlinear resistor, $N_{r}$, in the circuits is build with
operational amplifiers (AD712JN) as described in (Kennedy, 1992). The nominal values of the components can be found in table 1, however due to tolerances of the components each circuit is slightly different. Therefore synchronization in the sense that $\left|\mathbf{x}_{i}(t)-\mathbf{x}_{j}(t)\right|=0$ is not possible and practical synchronization is defined as $\left|\mathbf{x}_{i}(t)-\mathbf{x}_{j}(t)\right| \leq \delta$, for some fixed $\delta>0$.

Table 1. Nominal values for each circuit.

| Component | Value |
| :--- | ---: |
| $C_{1}$ | $10[\mathrm{nF}]$ |
| $C_{2}$ | $100[\mathrm{nF}]$ |
| $L$ | $22[\mathrm{mH}]$ |
| $R_{0}$ | $22[\Omega]$ |
| $R$ | $1.5-2.0[\mathrm{k} \Omega]$ |
| $G_{a}$ | $-0.758[\mathrm{mS}]$ |
| $G_{b}$ | $-0.409[\mathrm{mS}]$ |
| $B_{p}$ | $1.75[\mathrm{~V}]$ |

## 4. TWO SYSTEMS

Before synchronization of four systems is considered, the threshold value for synchronization of two circuits is determined, i.e., the minimal value $K$ such that practical synchronization occurs. The two circuits are diffusely coupled with a variable resistor, $R_{c}$, which gives the coupling constant $K$ as $\frac{1}{R_{c}}$ and a coupling matrix

$$
\boldsymbol{\Gamma}_{1}=\left[\begin{array}{cc}
K & -K  \tag{18}\\
-K & K
\end{array}\right]
$$

The variable resistor $R$ is set to $1775[\Omega]$ on both circuits, so the circuits operate on the double scroll attractor. Synchronization is visualized by the phase portrait of $x_{11}$ and $x_{21}$, shown in figure 2(b). The value of $R_{c}$ for the synchronization threshold is around $3400[\Omega]$ with $\delta=0.15[V]$.
Remark: If the value $R_{c}$ is increased, desynchronization occurs and at about $R_{c}=10[\mathrm{k} \Omega]$ the trajectories are no longer bounded. A possible explanation for this phenomena is the following. A single circuit (15), with $\mathbf{u}=0$, can have unbounded trajectories. No storage function $V(\mathbf{x})$ can be found such that inequality (8) is satisfied to prove semipassivity for system (15). Therefore it is not guaranteed that the solutions are ultimately bounded. This may cause the bursting phenomenon above $10[\mathrm{k} \Omega]$. For $R_{c}=30[\mathrm{k} \Omega]$ and above the current through the coupling resistor becomes negligible such that both circuits operate as free systems and the trajectories of both are bounded by their attractors again.

## 5. THREE SYSTEMS

When the network is expanded by adding a circuit, see figure 3 , the following coupling matrix $\Gamma$ is obtained


Fig. 2. Experimental synchronization for $R_{c}=$ 3430 [ $\Omega$ ].

$$
\boldsymbol{\Gamma}_{2}=\left[\begin{array}{ccc}
2 K & -K & -K  \tag{19}\\
-K & 2 K & -K \\
-K & -K & 2 K
\end{array}\right] .
$$

The coupling constant needed to globally syn-


Fig. 3. Layout of three coupled systems.
chronize this structure can be estimated using the conjecture stated in (Wu and Chua, 1996):

$$
\begin{equation*}
\mu_{1} \alpha_{1}=\mu_{2} \alpha_{2} \tag{20}
\end{equation*}
$$

where $\mu_{i}, i=1,2$ is the smallest nonzero eigenvalue of the coupling matrix $\boldsymbol{\Gamma}_{i}$ and $\alpha_{i}$ the coupling coefficient. Although it has been pointed out in (Pecora, 1998) that this conjecture is in general wrong, it holds in this particular case.
The synchronization threshold for three systems, using (20), requires a resistor value of $5100[\Omega]$. The threshold value in experiments is found to be 4950 [ $\Omega$ ], confirming that three systems, coupled in a ring structure, synchronize with a lower coupling constant $K$.

## 6. FOUR SYSTEMS

Four systems are symmetrically coupled in a ring structure with two coupling constants $K_{0}$ and $K_{1}$. With the proposed coupling, as shown in figure 4,


Fig. 4. Layout of four coupled systems.
the coupling matrix $\boldsymbol{\Gamma}$ can be written as follows

$$
\boldsymbol{\Gamma}_{3}=\left[\begin{array}{cccc}
K_{0}+K_{1} & -K_{0} & 0 & -K_{1}  \tag{21}\\
-K_{0} & K_{0}+K_{1} & -K_{1} & 0 \\
0 & -K_{1} & K_{0}+K_{1} & -K_{0} \\
-K_{1} & 0 & -K_{0} & K_{0}+K_{1}
\end{array}\right]
$$

If $K_{0}=K_{1}=K$ the smallest nonzero eigenvalue of $\boldsymbol{\Gamma}_{3}, \mu_{3}$, is equal to the smallest eigenvalue of $\boldsymbol{\Gamma}_{1}$, and the ring structure should synchronize with $R_{c}$ around $3400[\Omega]$. In figure 5 global synchronization is shown for $R_{c}=3200[\Omega]$ again with $\delta=$ 0.15 [ V$]$.

With the symmetric coupling matrix (21) there are four permutation matrices $\boldsymbol{\Pi}$ commuting with $\Gamma_{3}$

$$
\begin{array}{ll}
\boldsymbol{\Pi}_{1}=\left[\begin{array}{ll}
\mathbf{A} & \mathbf{O} \\
\mathbf{O} & \mathbf{A}
\end{array}\right], & \boldsymbol{\Pi}_{2}=\left[\begin{array}{ll}
\mathbf{O} & \mathbf{I}_{2} \\
\mathbf{I}_{2} & \mathbf{O}
\end{array}\right] \\
\boldsymbol{\Pi}_{3}=\left[\begin{array}{ll}
\mathbf{O} & \mathbf{A} \\
\mathbf{A} & \mathbf{O}
\end{array}\right], & \boldsymbol{\Pi}_{4}=\mathbf{I}_{4} \tag{22}
\end{array}
$$

where

$$
\mathbf{A}=\left[\begin{array}{ll}
0 & 1  \tag{23}\\
1 & 0
\end{array}\right], \quad \mathbf{O}=\left[\begin{array}{ll}
0 & 0 \\
0 & 0
\end{array}\right]
$$

Three linear invariant manifolds associated with (22) exist independently of systems (1) and are given by

$$
\begin{align*}
& \mathcal{A}_{1}=\mathbf{x} \in \mathbb{R}^{4 n}: \mathbf{x}_{1}=\mathbf{x}_{2}, \mathbf{x}_{3}=\mathbf{x}_{4}  \tag{24}\\
& \mathcal{A}_{2}=\mathbf{x} \in \mathbb{R}^{4 n}: \mathbf{x}_{1}=\mathbf{x}_{3}, \mathbf{x}_{2}=\mathbf{x}_{4}  \tag{25}\\
& \mathcal{A}_{3}=\mathbf{x} \in \mathbb{R}^{4 n}: \mathbf{x}_{1}=\mathbf{x}_{4}, \mathbf{x}_{2}=\mathbf{x}_{3} . \tag{26}
\end{align*}
$$

The intersection of any two of these manifolds describes the full synchronization manifold $\left(\mathrm{x}_{1}=\mathrm{x}_{2}=\mathrm{x}_{3}=\mathrm{x}_{4}\right)$. There are two possible ways to global synchronization depending on the ratio $K_{0}$ and $K_{1}$

$$
\begin{align*}
& \mathcal{A}_{1} \rightarrow \mathcal{A}_{1} \cap \mathcal{A}_{2}  \tag{27}\\
& \mathcal{A}_{3} \rightarrow \mathcal{A}_{3} \cap \mathcal{A}_{2} \tag{28}
\end{align*}
$$

Theorem 1, to prove stability of these manifolds, depends on two conditions. It is already pointed out that the first condition is not satisfied, since system (15) is not semipassive. However on an experimental setup the solutions are normally bounded by the attractor.
With $x_{1}$ chosen as the external signal in (14) and $\mathbf{z}=\left[\begin{array}{ll}x_{2} & x_{3}\end{array}\right]^{T}$ and the parameter values in table 1 , it is possible to find a matrix $\mathbf{P}$ such that (12) is satisfied and therefore system (15) is convergent. Hence it is expected that these manifolds are locally stable on the Chua circuits as long as the solutions remain bounded by the attractor. This is confirmed as shown in figures 6 and 7 . In figure 6 it can be seen that the circuits one and two and also three and four are synchronized with a coupling constant $K_{0}=\frac{1}{R_{c 0}}$. Circuits two
and three, coupled with $K_{1}=\frac{1}{R_{c 1}}$, are, as well as one and four, not synchronized. This corresponds to manifold $\mathcal{A}_{1}(24)$. In figure 7 the situation corresponding to manifold $\mathcal{A}_{3}(26)$ is depicted. These manifolds are robust to parameter variation of the variable resistors $R$ of the circuits. However if the coupling resistances are increased the same phenomena with two coupled systems occurs, i.e., the trajectories are no longer bounded. And again above a second threshold the four circuits operate as free systems. All these phenomena are summarized in a stability diagram shown in figure 8 .


Fig. 5. Phase portraits for global synchronization with $R_{c 0}=R_{c 1}=3200[\Omega]$.


Fig. 6. Phase portraits for partial synchronization with $R_{c 0}=3200$ and $R_{c 1}=9300$.

This stability diagram can also be obtained by numerically integrating four Chua systems (15), taking the tolerances of the components into account. The presented experimental results are qualitative comparable with numerical simulations. As an illustrative numerical example partial synchronization is considered. For the individual systems the capacitors and variable resistor are chosen as in table 2, while the other parameters of system (15) are the same as in table 1 . The coupling constants $K_{0}$ and $K_{1}$ are $\frac{1}{8500}$ and $\frac{1}{3200}$ respectively. In figure 9 the error signals $x_{i 1}-x_{j 1}, i, j=1,2,3,4$ are


Fig. 7. Phase portraits for partial synchronization with $R_{c 0}=9300$ and $R_{c 1}=3200$.


Fig. 8. Stability diagram.
shown. After the transients are vanished systems two and three synchronize as well as one and four. This simulation result matches with the experimental partial synchronization as shown in figure 7.

Table 2. Parameter values for individual circuits.

| System | Component | Value |
| ---: | :--- | ---: |
| 1 | $C_{1}$ | $10.90[\mathrm{nF}]$ |
|  | $C_{2}$ | $97.93[\mathrm{nF}]$ |
|  | $R$ | $1775[\Omega]$ |
| 2 | $C_{1}$ | $10.80[\mathrm{nF}]$ |
|  | $C_{2}$ | $99.60[\mathrm{nF}]$ |
|  | $R$ | $1778[\Omega]$ |
| 3 | $C_{1}$ | $10.98[\mathrm{nF}]$ |
|  | $C_{2}$ | $101.90[\mathrm{nF}]$ |
|  | $R$ | $1770[\Omega]$ |
| 4 | $C_{1}$ | $10.65[\mathrm{nF}]$ |
|  | $C_{2}$ | $100.50[\mathrm{nF}]$ |
|  | $R$ | $1780[\Omega]$ |



Fig. 9. Simulation synchronization errors for $R_{c 0}=8500$ and $R_{c 1}=3200$.

## 7. CONCLUSIONS

In this paper experimental partial synchronization of diffusively coupled Chua circuits is presented. With the experimental setup it is impossible to achieve a zero synchronization error due to tolerances of the electrical components. Therefore a form of practical synchronization is introduced to be able to specify synchronization of two systems. Besides global synchronization of two circuits, a bursting phenomena is observed if the diffusive coupling between two systems is above a certain threshold value. At this point the trajectories are no longer bounded by the double scroll attractor. In the case where four circuits are symmetrically coupled it shown that partial synchronization is possible. The stability of the linear invariant manifolds, describing this partial synchronization, can not be proven globally. However the manifolds are locally stable if the solutions remain bounded by the double scroll attractors.

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### 5.2.2 Additional simulations

As an addition to this article some numerical results are included to complement the presented experimental results. These simulations are carried out for the situation of the four diffusively coupled systems.
The stability diagram, see fig. 8 in the article, is used as a guideline for the simulations. We consider the second partial synchronized and the full synchronized situation. Further the different situations for the unsynchronized state are given.
In figure 5.14 the second partial synchronization is shown. This corresponds to the situation that $\mathbf{x}_{1}=\mathbf{x}_{2} \neq \mathbf{x}_{3}=\mathbf{x}_{4}$. Full synchronization can be seen in figure 5.15.
In the case of weak coupling three situations exist. If the coupling resistances are increased the circuits first desynchronize, then the trajectories grow unbounded and finally the resistance is so high that the current is negligible and the circuits operate as it are free systems.


Figure 5.14: Synchronization error for $R_{c 0}=3200[\Omega]$ and $R_{c 1}=8500[\Omega]$.


Figure 5.15: Synchronization error for $R_{c 0}=R_{c 1}=3200[\Omega]$.


Figure 5.16: Synchronization error for $R_{c 0}=R_{c 1}=8500[\Omega]$.


Figure 5.17: Synchronization error for $R_{c 0}=R_{c 1}=20000[\Omega]$.


Figure 5.18: Synchronization error for $R_{c 0}=R_{c 1}=70000[\Omega]$.

## Chapter 6

## Conclusions and recommendations

As final part of this thesis conclusions are drawn and recommendations for future research are presented.

### 6.1 Conclusions

In this thesis the complex behavior of multiple hybrid chaotic systems has been investigated. This objective has been achieved by analyzing multiple Chua circuits. The objective has been divided into three main aspects and conclusions will be drawn accordingly.

### 6.1.1 Analysis of a single Chua circuit

The Chua circuit is a complex system capable of generating bifurcation and chaos phenomena. The nonlinearity of the circuit is given by a piece-wise linear characteristic, which consists of three linear parts. The dynamic behavior is analyzed by looking at the stability properties of the equilibria and a 2D Poincaré map is derived to explain the bifurcations. This complex Poincaré map is approximated by a simpler 1D Poincaré map, which can be calculated for any given parameter set. The qualitative behavior of the Chua circuit can be visualized by iterating this 1D Poincaré map, without actual integration of the differential equations. With the 1D Poincaré map the single scroll and 'birth' and 'death' of the double scroll trajectories are explained.
To classify a system as chaotic it has to satisfy certain properties. We showed that these conditions can be fulfilled and therefore we can conclude that a Chua circuit is a chaotic system.

### 6.1.2 Development and comparison with experimental results

To compare the obtained numerical results with experimental results a new experimental setup has been designed. This resulted in a small portable setup with some additional advantages to previous designs. The circuit can be divided into two subsystems to measure the nonlinear resistor or a resonance circuit, consisting of the inductor, separately. Further data acquisition devices or other circuits can be easily connected.

To identify the parameters of the experimental setup two identification methods are used. First the parameters are determined by experiments of the separate subsystems and secondly three different Kalman filters are developed and compared.
With the obtained parameters numerical simulations are carried out to compare the results with experiments. However if the obtained parameters from measurement identification or the Kalman filters are used in a numerical simulation, without using measurement data, the solution always differs from experimental data. There are two important reasons for this.
Although all the experimental trajectories are qualitatively comparable with numerical results, there is a quantitative difference. This difference is also present if we compare a measured single
limit cycle with a filtered result obtained by a Kalman filter, which indicates that the used model in the Kalman filter is not correct.
This is caused by the assumptions we made on the location of the breakpoints and the fact the dynamic behavior of the electric components is neglected.
Secondly if we are able to develop an exact model, it is still impossible to obtain the same trajectory when the circuit generates a chaotic time signal. We can only measure with a finite precision and a numerical solver causes unwanted rounding and truncation errors. Due to the sensitivity of initial conditions the trajectories will always diverge from each other.

### 6.1.3 Synchronization of multiple circuits

As last part of this thesis we considered two synchronization methods, master-slave synchronization and mutual synchronization.

In the case of master-slave synchronization two Chua circuits, operating at the double scroll attractor, are unidirectionally coupled. This is achieved by applying a negative feedback to the slave system. Two different feedback signals are considered. We used either the voltage across the capacitor in parallel with the nonlinear resistor or the voltage across the capacitor in parallel with the inductor to create the feedback signal, state one or two respectively.
To be able to specify asymptotic synchronization of non-identical systems we introduced a form of practical synchronization. This is done because one of the drawbacks of the experimental setup is that it is impossible to achieve a zero synchronization error due to tolerances in presumably identical electrical components of the circuits.
It followed from both simulations and experiments that synchronization is possible. A smaller synchronization error is obtained if the coupling strength is increased and the error is also smaller if the first state is used to compose the feedback signal.
Besides synchronization we also observed intermittency when the second state is used in the feedback signal.

In the case of mutual synchronization the circuits are diffusively coupled using the first state. First synchronization of two systems is considered. Again by increasing the coupling strength the synchronization error decreases. However if the coupling strength is decreased the systems starts to desynchronize and above a certain threshold value a bursting phenomenon is observed and the trajectories are no longer bounded by the attractor. If a second threshold is crossed, the coupling is too weak to synchronize the circuits and they operate as free systems.
Secondly we looked at synchronization in a network of four symmetrically coupled Chua circuits. In this case there are three linear invariant manifolds, corresponding to a partial synchronized situation. Two of these manifolds can be made locally stable if the trajectories remain bounded by the double scroll attractors. Global stability is not obtained because a Chua circuit is not (semi)passive, which can result in unbounded trajectories. However on the experimental setup the solutions are normally bounded by the attractor. Therefore by altering the coupling strength one of the manifolds can be made locally stable and as a result partial synchronization is experimentally observed. These experimental observations are confirmed by performing several simulations.

To achieve synchronization of chaotic dynamical systems several aspects play a role. The way of coupling and the coupling strength are two important factors for experimental synchronization of non-identical systems. The systems' input-output relations can be used to prove (semi)passivity of the system, which guarantees bounded solutions.
If a dynamical system, driven by a bounded external signal, is convergent the solutions will forget their specific initial conditions. This indicates that if convergent systems are synchronized they remain synchronized.
For the Chua circuit we proved that the using the first state the system is convergent, while using the second state we can not prove this. This might explain the intermittency when the second state is used in the feedback signal.

### 6.2 Recommendations

Besides a better understanding of the very complex behavior of a Chua circuit several observations have been made in this thesis leading to a number of recommendations for further research.

- We concluded that it is impossible to obtain the same numerical trajectory when a dynamical system generates a chaotic time signal if no measurement information is used. To use measurement information a robust observer can be used to correct a numerical solution from time to time. This can also be viewed as a form of synchronization (Nijmeijer and Mareels, 1997).
- Experimental setup: At the current setup a fixed voltage supply of $\pm 15 \mathrm{~V}$ is used, a variable voltage supply could be used to place the locations of the breakpoints more symmetrically. The possible frequency dependency of the nonlinear resistor should be investigated further. If this effect is caused by the operational amplifiers a different realization of the nonlinear resistor may be desirable.
- Identification: First of all one should consider if the used numerical model is accurate enough to fully describe the electrical dynamics.
One of the drawbacks of using a Kalman filter is the initialization ( $P, R, Q$ matrices) of the filter and this needs more attention. One of the problems with the current model is that it is ill-conditioned and this can lead to instability of the filter. The estimation of the states and parameters might improve if a dimensionless model can be implemented.
A filter that is also capable of estimating the location of the breakpoints can improve the current results, e.g. (Santoboni and Nijmeijer, 2001)
- Synchronization: The bursting phenomenon, which occurs if the coupling strength crosses a certain threshold value, should be investigated further. A possible technique that can be used is the examination of transverse Lyapunov exponents of the synchronization manifold (Heagy et al., 1994; Pecora, 1998).
The possibility of real-time synchronization / identification between an experimental setup and a numerical model may be realized using a TUeDacs in combination with a TUeDAX Linux Live CD.


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## Appendix A

## Reference frame

In this appendix the explicit coordinates of the points $A, B, C, D, E$ and $F$ in the $D_{0}$ and $D_{1}$ units, see figure A.1, are given in terms of the real and complex eigenvalues $\sigma_{0}=\frac{\tilde{\sigma}_{0}}{\tilde{\omega}_{0}}, \gamma_{0}=$ $\frac{\tilde{\gamma}_{0}}{\tilde{\omega}_{0}}, \sigma_{1}=\frac{\tilde{\sigma}_{1}}{\tilde{\omega}_{1}}, \gamma_{1}=\frac{\tilde{\gamma}_{1}}{\tilde{\omega}_{1}}$. These coordinates are used to compute the approximated 1D Poincaré map.


Figure A.1: Original system and their images in the $D_{0}$ and $D_{1}$ units of the transformed system.

## A. $1 D_{0}$ unit

$$
\begin{align*}
A_{0} & =\left[\begin{array}{lll}
1 & p_{0} & 0
\end{array}\right]  \tag{A.1}\\
B_{0} & =\left[\frac{\gamma_{0}\left(\gamma_{0}-\sigma_{0}-p_{0}\right)}{Q_{0}} \frac{\gamma_{0}\left[1-p_{0}\left(\sigma_{0}-\gamma_{0}\right)\right]}{Q_{0}} 1-\frac{\gamma_{0}\left(\gamma_{0}-\sigma_{0}-p_{0}\right)}{Q_{0}}\right]  \tag{A.2}\\
C_{0} & =\left[\begin{array}{lll}
0 & 0 & 1
\end{array}\right]  \tag{A.3}\\
E_{0} & =\left[\begin{array}{lll}
1 & \sigma_{0} & 0
\end{array}\right]  \tag{A.4}\\
F_{0} & =\left[\frac{\gamma_{0}\left(\gamma_{0}-2 \sigma_{0}\right)}{Q_{0}} \frac{\gamma_{0}\left[1-\sigma_{0}\left(\sigma_{0}-\gamma_{0}\right)\right]}{Q_{0}} \frac{\left(\sigma_{0}^{2}+1\right)}{Q_{0}}\right]  \tag{A.5}\\
p_{0} & =\sigma_{0}+\frac{k_{0}}{\gamma_{0}}\left(\sigma_{0}^{2}+1\right)  \tag{A.6}\\
Q_{0} & =\left(\sigma_{0}-\gamma_{0}\right)^{2}+1  \tag{A.7}\\
k_{0} & =-\frac{\tilde{\gamma}_{0}}{\tilde{\gamma}_{1}} \tag{A.8}
\end{align*}
$$

## A. $2 D_{1}$ unit

$$
\begin{align*}
A_{1} & =\left[\begin{array}{lll}
1 & p_{1} & 0
\end{array}\right]  \tag{A.9}\\
B_{1} & =\left[\begin{array}{lll}
1 & \sigma_{1} & 0
\end{array}\right]  \tag{A.10}\\
D_{1} & =\left[\begin{array}{lll}
0 & 0 & 1
\end{array}\right]  \tag{A.11}\\
E_{1} & =\left[\frac{\gamma_{1}\left(\gamma_{1}-\sigma_{1}-p_{1}\right)}{Q_{1}} \frac{\gamma_{1}\left[1-p_{1}\left(\sigma_{1}-\gamma_{1}\right)\right]}{Q_{1}} 1-\frac{\gamma_{1}\left(\gamma_{1}-\sigma_{1}-p_{1}\right)}{Q_{1}}\right]  \tag{A.12}\\
F_{1} & =\left[\frac{\gamma_{1}\left(\gamma_{1}-2 \sigma_{1}\right)}{Q_{1}} \frac{\gamma_{1}\left[1-\sigma_{1}\left(\sigma_{1}-\gamma_{1}\right)\right]}{Q_{1}} \frac{\left(\sigma_{1}^{2}+1\right)}{Q_{1}}\right]  \tag{A.13}\\
p_{1} & =\sigma_{1}+\frac{k_{1}}{\gamma_{1}}\left(\sigma_{1}^{2}+1\right)  \tag{A.14}\\
Q_{1} & =\left(\sigma_{1}-\gamma_{1}\right)^{2}+1  \tag{A.15}\\
k_{1} & =\frac{1}{k_{0}} \tag{A.16}
\end{align*}
$$

## A. 3 Connection map

The connection map $\Phi$ can be used to transform coordinates from the $D_{0}$ unit into the $D_{1}$ unit. Since $z_{i}=1-x_{i}$, we only need a relation for the $x$ and $y$ coordinates

$$
\left[\begin{array}{l}
x_{1}  \tag{A.17}\\
y_{1}
\end{array}\right]=\Phi\left[\begin{array}{l}
x_{0} \\
y_{0}
\end{array}\right]
$$

which can be rewritten as, cf. (Chua et al., 1986),

$$
\left[\begin{array}{c}
x_{1}-1  \tag{A.18}\\
y_{1}-p_{1}
\end{array}\right]=L\left[\begin{array}{c}
x_{0}-1 \\
y_{0}-p_{0}
\end{array}\right]
$$

where $L$ is given by
$L=\frac{\left(\sigma_{1}^{2}+1\right) k_{1}}{\left(\sigma_{0}^{2}+1\right)\left(k_{0}+1\right) Q_{1} \gamma_{1}}\left[\begin{array}{cc}-\gamma_{1}\left(k_{0}+1\right)\left[Q_{0}+\gamma_{0}\left(\sigma_{0}-\gamma_{0}\right)\left(k_{1}+1\right)\right] & \gamma_{0} \gamma_{1}\left(k_{0}+1\right)\left(k_{1}+1\right) \\ -\gamma_{0}\left(k_{1}+1\right)\left(\sigma_{0}-\gamma_{0}\right)\left[\sigma_{1}\left(\sigma_{1}-\gamma_{1}\right)+1\right] & \gamma_{0}\left(k_{1}+1\right)\left[Q_{1}+\gamma_{1}\left(\sigma_{1}-\gamma_{1}\right)\left(k_{0}+1\right)\right] \\ -\gamma_{1}\left(k_{0}+1\right)\left(\sigma_{1}-\gamma_{1}\right)\left[\sigma_{0}\left(\sigma_{0}-\gamma_{0}\right)+1\right] & \end{array}\right]$

## Appendix B

## Component list and layout

## B. 1 Component list

Table B.1: Component list.

| Component | Value | Description |
| :--- | ---: | :--- |
| $R_{1}$ | $220[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{2}$ | $220[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{3}$ | $2200[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{4}$ | $22000[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{5}$ | $22000[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{6}$ | $3300[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{7}$ | $22[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{8}$ | $100000[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{9}$ | $1500[\Omega]$ | $\frac{1}{4} W$ resistor |
| $R_{10}$ | $0-500[\Omega]$ | $\frac{1}{4} W$ resistor |
| $C_{1}$ | $10[n F]$ | capacitor |
| $C_{2}$ | $100[n F]$ | capacitor |
| $C_{3}$ | $10[n F]$ | capacitor |
| $C_{4}$ | $1[\mu F]$ | tantalum capacitor |
| $C_{5}$ | $1[\mu F]$ | tantalum capacitor |
| $C_{6}$ | $1[\mu F]$ | tantalum capacitor |
| $C_{7}$ | $1[\mu F]$ | tantalum capacitor |
| $C_{8}$ | $1[\mu F]$ | tantalum capacitor |
| $C_{9}$ | $1[\mu F]$ | tantalum capacitor |
| $C_{10}$ | $1[\mu F]$ | tantalum capacitor |
| $C_{11}$ | $1[\mu F]$ | tantalum capacitor |
| $U_{1}, U_{2}, U_{3}, U_{4}$ |  | AD712JN operational amplifier |
| $V 1 a, V 2 a$ |  | BNC connector |
| $V 1 b, V 2 b$ |  | pin |
| $X 2 a, X 2 b$ |  | pin+jumper |
| $X 3 a, X 3 b$ |  | pin+jumper |

## B. 2 Layout



Figure B.1: Electrical layout

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