

Author / Eingereicht von **Felix Dreßler** k12105003

Submission / Angefertigt am Institute of Numerical Analysis

Thesis Supervisor / First Supervisor / BeurteilerIn / ErstbeurteilerIn / ErstbetreuerIn DI Dr. Andreas SchafeIner

April 2025

# **Circuit Modelling**



Bachelor Thesis

to obtain the academic degree of

Bachelor of Science

in the Bachelor's Program

Technische Mathematik

JOHANNES KEPLER UNIVERSITY LINZ

Altenbergerstraße 69 4040 Linz, Austria www.jku.at DVR 0093696

# Kurzfassung

Kapitel zwei gibt Überblick über die Physikalischen Zusammenhänge und deren Mathematische Beschreibung. Anschließend wird ein Mathematisches Modell nach der Methode der "Modified Nodal Analysis" erstellt. Im dritten Kapitel analysieren wir die Systeme, welche durch diese Analyse entstehen, wobei das vierte Kapitel eine wichtige Eigenschaft dieser Systeme genauer untersucht. Das letzte Kapitel beschäftigt sich mit der Numerischen Lösung dieser Systeme.

→ je etwas mehr Text zu den Kapiteln (Struktur der Kapitel)

# **Abstract**

Abstract in English.

# **Contents**

| 1 | Intro | oduction                                     | 1  |
|---|-------|--|----|
| 2 | For   | mulating a Mathematical Model                | 2  |
|   | 2.1   | Network Topology                             | 2  |
|   | 2.2   | Energy Conservation Laws                     | 6  |
|   | 2.3   | Electrical Components and their relations    | 7  |
|   | 2.4   | Modified Nodal Analysis - MNA                | 10 |
| 3 | Diff  | erential Algebraic Equations                 | 14 |
|   | 3.1   | Types of DAEs                                | 14 |
|   |       | 3.1.1 Weierstraß-Kronecker Normalform        | 16 |
|   | 3.2   | Index of a Differential Algebraic Equation   | 21 |
|   | 3.3   | Consistent Initial Values                    | 23 |
| 4 | Inde  | ex Analysis of the Modified Nodal Analysis   | 26 |
|   | 4.1   | General Index analysis                       | 27 |
|   | 4.2   | Topological Conditions                       | 28 |
| 5 | Nun   | nerical Solutions                            | 34 |
|   | 5.1   | Single-Step-Methods                          | 35 |
|   |       | 5.1.1 Consistency, Stability and Convergence | 35 |
|   |       | 5.1.2 further stability properties           | 38 |
|   | 5.2   | Multistep-Methods                            | 38 |
|   |       | 5.2.1 Consistency, Stability and Convergence | 39 |
|   |       | 5.2.2 further stability properties           | 41 |
|   | 5.3   | Implicit linear multi-step formulas          | 43 |
|   |       | 5.3.1 BDF-k methods                          | 43 |
|   |       | 5.3.2 trapezoidal rule                       | 45 |
|   | 5.4   | Numerical Examples                           | 46 |

# **List of Figures**

| charging capacitor with series resistor and voltage source                              | 3   |
|---|---|
| capacitor and inductor oscillator   | 5   |
|   | 6   |
|   | 8   |
|   | 8   |
|   | 8   |
|   | 9   |
| current source symbol   | 9   |
| charging capacitor with series resistor and voltage source (note the difference to 2.1) | 30  |
|   | 31  |
| current source with resistor and capacitor  | 33  |
| approximation of a function using numerical methods                                     | 35  |
| $\widehat{A(\alpha)}$ stability for $\alpha = \frac{\pi}{4}$                            | 42  |
| stability regions of BDF-schemes  | 45  |
| charging capacitor with series resistor and voltage source                              | 46  |
| Exact solution for example 1  | 46  |
| LC-circuit  | 48  |
| Exact solution for example 2  | 48  |
| Current source with capacitor and resistor  | 49  |
| Exact solution for example 3  | 49  |
|   | charging capacitor with series resistor and voltage source (note the difference to 2.1) |

# **List of Tables**

| 5.1 | Resulting errors for the BDF-k methods and ther trapezoidal rule | 47 |
|-----|--|----|
| 5.2 | Resulting errors for the BDF-k methods and ther trapezoidal rule | 48 |
| 5.3 | Resulting errors for the BDF-k methods and the trapezoidal rule  | 49 |

## 1 Introduction

This chapter should include information about what circuit modelling wants to achieve as well as giving an overview of what this bachelor-thesis is about.

What is this thesis about?

Modelling and numerically solving systems that arrise from electrical circuits with RLC elements. Furthermore it will briefly discuss on expanding this baseline with more complicated electrical components. What is the goal of this thesis?

The goal of this thesis is to give insight into the state of the art of circuit modelling. It aims to elaborate on the underlying concepts of MNA as well as on the most commonly used numerical methods.

based on DAE lecture and [2]

To accurately represent a physical system in a mathematical model we first have to think about **how** to best formulate a mathematical model that formulates this system precise enough. This chapter aims to elaborate on the Modified Nodal Analysis (MNA) of an electrical circuit, which is the common choice for these kinds of models. For this we have to introduce the notions of network topology as well as some basic electrical components and physical laws. Most of the following notions such as voltage, current and electrical potentials can depend on time. For better readability we will leave out the time argument t in some cases.

## 2.1 Network Topology

An electrical circuit is usually considered as a graph  $(\mathcal{N}, \mathcal{E})$  where  $\mathcal{N} = (n_0, n_1, n_2, ..., n_k)$  denotes nodes and  $\mathcal{E} = \{e_j : j = 1, ..., l\}$  is the set of edges where  $|\mathcal{N}| = k$  is the number of nodes and  $|\mathcal{E}| = l$  the number of edges. If for some i and some j the nodes  $n_i$  and  $n_j$  are connected, then there exists an edge connecting those two nodes. We can store this information in an *incidence matrix*  $\tilde{A} = (\tilde{a}_{ij}) \in \mathbb{R}^{k \times l}$  which is defined by

$$\tilde{a}_{ij} = \begin{cases} 1 & \text{edge } j \text{ starts at node } i, \\ -1 & \text{edge } j \text{ ends at node } i, \\ 0 & \text{else.} \end{cases}$$

We call  $u = (u_0, u_1, u_2, ...)$  the corresponding electrical potentials (or just potentials) to the nodes  $\mathcal{N}$ . The difference between the potentials at two connected nodes is called the voltage at the respective edge. To fix the absolute values of these potentials we have

to set one node to a fixed value. We will do that by "grounding" the node  $n_0$ , which means we set the potential  $u_0 = 0$ . The grounding of a node allows us to remove the corresponding row from the incidence matrix to get the *reduced incidence matrix A*. The vector  $v = (v_{ij})_{ij}$  represents the voltages at the edges. For some i and some j the voltage at edge ij is  $v_{ij} = u_i - u_j$ .

We will later see, that the components of an electrical circuit, which will be installed along the edges, describe a relationship between the current along the edge the corresponding voltage. Thus a current vector  $i = (i_1, i_2, i_3, ...)$  containing the currents along the edges is required.

### **Example 1.1** (Incidence Matrix)

As an example we consider the charging of a capacitor. The circuit is given in Figure 2.1.

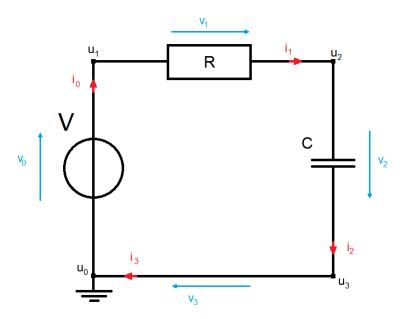


Figure 2.1: charging capacitor with series resistor and voltage source

With the node-potentials, the voltages and the currents collected in the vectors

$$u = \begin{pmatrix} u_0 \\ u_1 \\ u_2 \\ u_3 \end{pmatrix}, \quad v = \begin{pmatrix} v_0 \\ v_1 \\ v_2 \\ v_3 \end{pmatrix}, \quad i = \begin{pmatrix} i_0 \\ i_1 \\ i_2 \\ i_3 \end{pmatrix},$$

the incidence matrix of this circuit has the form

$$\tilde{A} = \begin{pmatrix} 1 & 0 & 0 & -1 \\ -1 & 1 & 0 & 0 \\ 0 & -1 & 1 & 0 \\ 0 & 0 & -1 & 1 \end{pmatrix}.$$

The rows of this matric correspond to the nodes of the circuit and the columns correspond to the edges and tell us if the edge either starts (1) or end (-1) at a node. By grounding node  $u_0$  this matrix reduces to

$$A = \begin{pmatrix} -1 & 1 & 0 & 0 \\ 0 & -1 & 1 & 0 \\ 0 & 0 & -1 & 1 \end{pmatrix}.$$

In this circuit edge 3 is not populated with any components. This means that the voltage along this edge does not change, or in terms of potentials

$$u_0 - u_3 = 0.$$

Thus we can consider this circuit with node 0 and node 3 merged. This leads to a slightly different incidence and reduced incidence matrix, i.e.

$$\tilde{A} = \begin{pmatrix} 1 & 0 & -1 \\ -1 & 1 & 0 \\ 0 & -1 & 1 \end{pmatrix}$$
 and  $A = \begin{pmatrix} -1 & 1 & 0 \\ 0 & -1 & 1 \end{pmatrix}$ , respectively.

### Example 2.1 (Incidence Matrix)

Another example is the a circuit that only contains a capacitor and an inductor, see Figure 2.2. This circuit will oscillate, if given appropriate initial conditions.

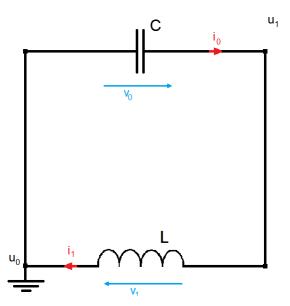


Figure 2.2: capacitor and inductor oscillator

We again collect the node-potentials, the voltages and the currents in the vectors

$$u = \begin{pmatrix} u_0 \\ u_1 \end{pmatrix}, \quad v = \begin{pmatrix} v_0 \\ v_1 \end{pmatrix}, \quad i = \begin{pmatrix} i_0 \\ i_1 \end{pmatrix}.$$

The incidence matrix of this circuit has the form

$$\tilde{A} = \begin{pmatrix} 1 & -1 \\ -1 & 1 \end{pmatrix}.$$

By grounding node  $u_0$  this can again be reduced, we obtain the reduced incidence matrix

$$A = \begin{pmatrix} -1 & 1 \end{pmatrix}$$
.

### Example 3.1 (Incidence Matrix)

As a third example we consider a similar circuit to the circuit presented in Example 1.1, the only difference is that we connected the resistor parallel to the capacitor instead of connecting it in series.

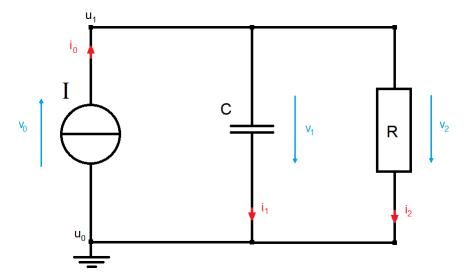


Figure 2.3: Capacitor with resistor and voltage source

We collect the node-potentials, the voltages and the currents in the vectors

$$u = \begin{pmatrix} u_0 \\ u_1 \end{pmatrix}, \quad v = \begin{pmatrix} v_0 \\ v_1 \\ v_2 \end{pmatrix}, \quad i = \begin{pmatrix} i_0 \\ i_1 \\ i_2 \end{pmatrix}.$$

The incidence matrix of this circuit has the form

$$\tilde{A} = \begin{pmatrix} 1 & -1 & -1 \\ -1 & 1 & 1 \end{pmatrix},$$

this leads to the reduced incidence matrix

$$A = \begin{pmatrix} -1 & 1 & 1 \end{pmatrix}.$$

### 2.2 Energy Conservation Laws

To fully fix all the variables that arise in the model of an electrical circuit we will need some *conservation laws*. In this context, conservation laws are a description of physical properties that can be observed in electrical circuits. This means that they have to be fulfilled non the

less, thus including them into the system also makes sense from a physical perspective. They were proposed by Kirchhoff in 1845, see [6].

### • Kirchhoff's voltage law (KVL):

The sum of voltages along each loop of the network must equal to zero. Using the incidence matrix *A* this law can be formulated as

$$A^{\top}u = v. \tag{2.1}$$

### • Kirchhoff's current law (KCL):

For any node, the sum of currents flowing into the node is equal to the sum of currents flowing out of the node. Using the incidence matrix *A* again, this law can be formulated as

$$Ai = 0. (2.2)$$

### 2.3 Electrical Components and their relations

Electrical components are described by equations relating their edge voltage v to their edge current i. We will mainly focus on so-called RLC-networks (short for Resistor, Inductor, Capacitor), which consist of resistors, capacitors, inductors, voltage sources and current sources. Diodes and transistors as well as other electrical components can be described in a similar way, although these lead to a more difficult analysis of the system, see e.g. [2] and [1]. Using the notation introduced in Section 2.1, the relations read as:

#### • Resistor

Resistors "resist" the flow of current, which causes voltage to drop. This behaviour is described by the *resistance*  $R \in \mathbb{R}^+ := \{x \in \mathbb{R} : x > 0\}$  which is given in *Ohm* ( $\Omega$ ) and its reciprocal, the *conductance*  $G \in \mathbb{R}^+$ , which is given in *Siemens* ( $S = \frac{1}{\Omega}$ ).

$$v = R i \quad \text{or} \quad i = G u. \tag{2.3}$$



Figure 2.4: resistor symbol

### Capacitor

Capacitors "store" electrical energy by accumulating electrical charge. Their characteristic equations can be described directly using the stored charge  $Q \in \mathbb{R}_0^+ := \{x \in \mathbb{R} : x \geq 0\}$  or indirectly using the change in charge, which is nothing other than the current *I*. The *capacitance*  $C \in \mathbb{R}^+$  is given in *Farads* (*F*).

$$Q = C v$$
 and by derivation in t  $I = C \frac{d}{dt}v = C v'$ . (2.4)



Figure 2.5: capacitor symbol

#### • Inductor (Coil)

An electric current flowing through a conductor generates a magnetic field  $\Phi \in \mathbb{R}^{n \times n}$  surrounding it. This magnetic field causes a voltage drop dependant on the change in current. The *inductance*  $L \in \mathbb{R}^+$  is given in *Henry* (H).

$$\Phi = L i$$
 and by derivation in t  $v = L i'$ . (2.5)



Figure 2.6: inductor symbol

### • Voltage Source

A voltage source supplies the system with a voltage. It can either supply varying

amounts of voltage (with the special case of alternating current AC) or a fixed amount of voltage. The unit of voltage is *Volts* (*V*).

$$v = v_{src} (2.6)$$



Figure 2.7: voltage source symbol

#### • Current Source

A current source supplies the system with current. It can either supply varying amounts of current (with the special case of alternating current AC) or a fixed amount of current. The unit of current is Ampere(A).

$$i = i_{src} \tag{2.7}$$

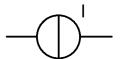


Figure 2.8: current source symbol

The resistance *R*, conductance *G*, capacitance *C* and inductance *L* are respectively described as a scalar constant that relates the edge-current to the edge-voltage. If there are more components of the same kind in one circuit their corresponding constants will usually be collected into a matrix which is also called the same and also denoted by the same letter, respectively. These matrices are then positive definite diagonal matrices.

### 2.4 Modified Nodal Analysis - MNA

[2] and [5]

To analyse the network further we will rearrange the columns of the reduced incidence matrix *A* such that it has the block form

$$A = (A_R A_C A_L A_V A_I)$$

where  $A_R$ ,  $A_C$ ,  $A_L$ ,  $A_V$  and  $A_I$  include the columns that are related to the resistors, capacitors, coils, voltage sources and current sources, respectively.

To mathematically describe the circuit we will use *modified nodal analysis* (or short MNA), see the paper by Chung-Wen Ho, Albert E. Ruehli and Pierce A. Brennan, [4]. MNA uses the node voltages as well as the currents of the coils and the voltage sources as unknowns. It is based on the conservation laws (2.2) and (2.1) as well as on the voltage-current relations of the electrical components, discussed in Section 2.3. The voltages can be represented using the node-potentials

$$v = A^{\top}u$$

The vector v can thus be rearranged into  $v = (v_R, v_C, v_L, v_{src}, v_I)$ . In a similar way we also rearrange the current vector into  $i = (i_R, i_C, i_L, i_V, i_s rc)$ . Using the sorted incidence matrix blocks we can rewrite the resistor current relation as

$$i_R = G v_R = G A_R^\top u.$$

Analogously, we rewrite the capacitor relation as

$$i_C = C v_C' = C A_C^\top u'.$$

Kirchhoffs current law (2.2) gives us that

$$A_C i_C + A_R i_R + A_L i_L + A_V i_V = -A_I i_{src}$$
.

We plug in the component relations derived above and obtain

$$A_{\mathcal{C}}CA_{\mathcal{C}}^{\top}u' + A_{\mathcal{R}}GA_{\mathcal{R}}^{\top}u + A_{\mathcal{L}}i_{\mathcal{L}} + A_{\mathcal{V}}i_{\mathcal{V}} = -A_{\mathcal{I}}i_{src}.$$

Combining this with the component law for inductors (2.5) and the potential-voltage relation for voltage sources (2.6) we finally get the modified nodal analysis equations

$$\begin{split} A_C C A_C^\top u' + A_R G A_R^\top u + A_L i_L + A_V i_V &= -A_I i_{src}, \\ L i_L' - A_L^\top u &= 0, \\ -A_V^\top u &= -v_{src}. \end{split}$$

In matrix form these read

$$\begin{pmatrix} A_{C}CA_{C}^{\top} & 0 & 0 \\ 0 & L & 0 \\ 0 & 0 & 0 \end{pmatrix} * \begin{pmatrix} u' \\ i'_{L} \\ i'_{V} \end{pmatrix} + \begin{pmatrix} A_{R}GA_{R}^{\top} & A_{L} & A_{V} \\ -A_{L}^{\top} & 0 & 0 \\ -A_{V}^{\top} & 0 & 0 \end{pmatrix} * \begin{pmatrix} u \\ i_{L} \\ i_{V} \end{pmatrix} = \begin{pmatrix} -A_{I}i_{src} \\ 0 \\ -v_{src} \end{pmatrix}, \quad (2.8)$$

where the diagonal matrices C, G and L contain the capacities, conductivities and inductivities.

The resulting systems are *stiff* systems. This means that for their numerical solution special care has to be put into which methods are suitable for solving these systems in a stable manner. For analyzing nonlinear components, the charge-flux oriented formulation might be useful for numerical considerations, for this see [2].

#### **Example 1.2** (Modified Nodal Analysis System)

We again consider the charging of a capacitor. The circuit is again given as in Example 1.1, where we have already derived the reduced incidence matrix

$$A = \begin{pmatrix} -1 & 1 & 0 \\ 0 & -1 & 1 \end{pmatrix}.$$

This matrix can be split into the three submatrices  $A_V$ ,  $A_R$  and  $A_C$  containing the columns of the matrix A. The matrices  $A_L$  and  $A_V$  are empty in this case. The diagonal matrices containing the components constants are

$$C = (C), \qquad L = (), \qquad G = (\frac{1}{R}).$$

Plugging this into the formula (2.8) we obtain the system

$$\begin{pmatrix} 0 & 0 & 0 \\ 0 & C & 0 \\ 0 & 0 & 0 \end{pmatrix} * \begin{pmatrix} u_1' \\ u_2' \\ i_0' \end{pmatrix} + \begin{pmatrix} \frac{1}{R} & -\frac{1}{R} & -1 \\ -\frac{1}{R} & \frac{1}{R} & 0 \\ 1 & 0 & 0 \end{pmatrix} * \begin{pmatrix} u_1 \\ u_2 \\ i_0 \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \\ -v_{src} \end{pmatrix}.$$
 (2.9)

Note that because L is empty, the system does not contain any related equations to inductances.

### Example 2.2 (Modified Nodal Analysis System)

Similarly we also derived the incidence Matrix for this system in Example 2.1.

$$A = \begin{pmatrix} -1 & 1 \end{pmatrix}$$
.

This matrix can be split into the two submatrices  $A_C$  and  $A_L$  containing the columns of the matrix A. The matrices  $A_R$ ,  $A_V$  and  $A_I$  are empty in this case. The diagonal matrices containing the components constants are

$$C = (C), \qquad L = (L), \qquad G = ().$$

Plugging this into the formula (2.8) we obtain the system

$$\begin{pmatrix} C & 0 \\ 0 & L \end{pmatrix} * \begin{pmatrix} u_1' \\ i_L' \end{pmatrix} + \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} * \begin{pmatrix} u_1 \\ i_L \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix}.$$

### Example 3.2 (Modified Nodal Analysis System)

Consider the circuit of Example 3.1. We have already found the incidence matrix for this circuit.

$$A = \begin{pmatrix} -1 & 1 & 1 \end{pmatrix}.$$

This matrix can be split into the three submatrices  $A_I$ ,  $A_R$  and  $A_C$  containing the columns of the matrix A. The matrices  $A_L$  and  $A_V$  are empty in this case. The diagonal matrices containing the components constants are

$$C = (C), \qquad L = (), \qquad G = (\frac{1}{R}).$$

Plugging this into the formula (2.8) we obtain the system

$$\begin{pmatrix} C & 0 \\ 0 & 0 \end{pmatrix} * \begin{pmatrix} u_1' \\ i_V' \end{pmatrix} + \begin{pmatrix} \frac{1}{R} & -1 \\ 1 & 0 \end{pmatrix} * \begin{pmatrix} u_1 \\ i_V \end{pmatrix} = \begin{pmatrix} 0 \\ -v_{src} \end{pmatrix}. \tag{2.10}$$

von [5] und DAE lecture

The resulting MNA of electrical circuits leads to a special form of differential equations, namely *differential algebraic equations* (DAE). Additionally to differential equations they also contain algebraic equations or are equivalent to such a system. They are distinct from ordinary differential equations in that the Jacobian matrix for a DAE system is singular. To understand the general solvability of those systems and to find good numerical methods we first have to take a look at the theory of such DAEs.

In the most general form a DAE can be written as: Find  $y : \mathbb{R} \to \mathbb{R}^n$  such that

$$F(t, y(t), y'(t)) = 0, \qquad \forall t \in I$$
(3.1)

with  $F: \mathbb{R} \times \mathbb{R}^n \times \mathbb{R}^n \to \mathbb{R}^n$  sufficiently smooth and I the time-interval. As this is usually too general, one has to consider more specific types of DAEs. This chapter focuses on giving a brief overview of the different types of differential algebraic equations. We will discuss the specific form of linear systems with constant coefficients in more detail, as those are the systems that arise from modified nodal analysis of a RLC-system. Additionally this chapter also discusses the notion of the index of an MNA-system.

## 3.1 Types of DAEs

Commonly systems of either of the three types below are considered when talking about DAEs.

### • Linear systems with constant coefficcients

are systems of the form: find y such that

$$Ay'(t) + By(t) = f(t), \tag{3.2}$$

with  $A, B \in \mathbb{R}^{n \times n}$ , A singular, B regular and  $f : \mathbb{R} \to \mathbb{R}^n$  a function in time. If A is regular this equation can be transformed into an ordinary differental equation, since

$$y'(t) + A^{-1}By(t) = A^{-1}f(t)$$

Thus it is reasonable to assume *A* singular.

• **Linear time dependent systems** are systems of the form: find *y* such that

$$A(t)y'(t) + B(t)y(t) = f(t),$$

with  $A, B : \mathbb{R} \to \mathbb{R}^{n \times n}$ ,  $f : \mathbb{R} \to \mathbb{R}^n$  functions, such that for every  $t \in \mathbb{R}$  the matrix A(t) is singular and the matrix B(t) regular.

### • Structured (non-linear) systems

are semi-explicit systems of the form: find (y, z) such that

$$y'(t) = f(t, y(t), z(t)),$$
 (3.3)

$$0 = g(t, y(t), z(t)), (3.4)$$

with  $f : \mathbb{R} \to \mathbb{R}^n$  and  $g : \mathbb{R} \to \mathbb{R}^d$  functions.

For our analysis of electrical networks we will focus on linear systems with constant coefficients. The other systems could for example occur when we also consider time-dependant component relations e.g. Temperature. A splitting of the variable y is usually common, where we differ between the *differential variable* and the *algebraic variable*. A differential variable in this sense is the part of y, which derivatives also appear in the differential algebraic equation. The algebraic variable on the other hand only appears in algebraic equations.

#### 3.1.1 Weierstraß-Kronecker Normalform

To determine the solvability of a linear system with constant coefficients (3.2) we first need to introduce a normalform for the system, the *Weierstraß-Kronecker normalform*. This normalform is dependent on the family  $\{A, B\} := \{\mu A + B | \mu \in \mathbb{R}\}$ , which is called the *matrix pencil* of the DAE.

**Definition 1.** The matrix pencil  $\{A, B\}$  is called regular if there exists some  $c \in \mathbb{R}$ , such that (cA + B) is regular (i.e.  $det(cA + B) \neq 0$ ), otherwise it is called singular.

**Theorem 1** (Jordan Normalform [5, Theorem 13.2.1]). For every matrix  $Q \in \mathbb{R}^{n \times n}$  there exists a regular matrix  $T \in \mathbb{C}^{n \times n}$ , such that

$$T^{-1}QT = J = diag(J_1, ..., J_r) \quad with \quad J_i = \begin{pmatrix} \lambda_i & 1 & & 0 \\ 0 & \lambda_i & \ddots & \vdots \\ & \ddots & \ddots & 1 \\ 0 & \dots & 0 & \lambda_i \end{pmatrix} \in \mathbb{C}^{m_i \times m_i}$$

and  $n = m_1 + ... + m_r$ .

The matrix J is called Jordan Normalform of Q, the matrices  $J_i$  are called Jordan Blocks, where  $\lambda_i$  are corresponding the eigenvalues of Q. The matrix J is uniquely determined by Q except for the arrangement of the diagonal blocks. If Q possesses only real eigenvalues, then T can also be choosen from the reals.

In order to separate the equation (3.2) into an algebraic and a differential part, we transform the matrices A and B according to the following theorem.

**Theorem 2** ([5, Satz 13.2.2]). *Let*  $\{A, B\}$  *be a regular matrix pencil. There exist*  $P, Q \in \mathbb{C}^{n \times n}$  *such that* 

$$PAQ = \begin{pmatrix} I_d & 0 \\ 0 & N \end{pmatrix}$$
,  $PBQ = \begin{pmatrix} R & 0 \\ 0 & I_{n-d} \end{pmatrix}$ 

where

$$N = diag(N_1,...,N_r) \quad with \quad N_i = egin{pmatrix} 0 & 1 & & 0 & \\ & \ddots & \ddots & \\ & & & 0 & 1 \\ 0 & & & 0 \end{pmatrix} \in \mathbb{R}^{n_i \times n_i}$$

and R has Jordan Normalform. By  $I_k$  we denote the identity matrix of size  $k \times k$ .

*Proof.* Because  $\{A, B\}$  is a regular matrix pencil, there exists  $c \in \mathbb{R}$  such that (cA + B) is regular. Set

$$\hat{A} := (cA + B)^{-1}A, \quad \hat{B} := (cA + B)^{-1}B.$$

Considering

$$(cA + B)^{-1}(cA + B) = I \implies (cA + B)^{-1}B + c(cA + B)^{-1}A = I$$

we get that

$$\hat{B} = I - c\hat{A}$$
.

Let  $J_{\hat{A}}$  be the Jordan Normalform of  $\hat{A}$ , this means that there exists a regular matrix  $T_1$  such that

$$T_1^{-1}AT_1 = J_{\hat{A}} = \begin{pmatrix} W & 0 \\ 0 & \tilde{N} \end{pmatrix}.$$

The matrix W contains the Jordanblocks with Eigenvalues which are nonzero, the matrix  $\tilde{N}$  contains the Jordan blocks with Eigenvalues equal to zero, thus  $\tilde{N}$  is *nilpotent*. The Jordan Normalform  $J_{\hat{B}}$  of  $\hat{B}$  is given by

$$T_1^{-1}\hat{B}T_1 = J_{\hat{B}} = \begin{pmatrix} I - cW & 0\\ 0 & I - c\tilde{N} \end{pmatrix}.$$

The following two transformations will allow us to get the desired structure. First we will transform  $J_{\hat{A}}$  with

$$T_2 := \begin{pmatrix} W & 0 \\ 0 & I - c\tilde{N} \end{pmatrix}$$

in

$$T_2^{-1}J_{\hat{A}} = \begin{pmatrix} I & 0 \\ 0 & (I - c\tilde{N})^{-1}\tilde{N} \end{pmatrix}$$

and  $J_{\hat{B}}$  in

$$T_2^{-1}J_{\hat{B}} = \begin{pmatrix} W^{-1} - cI & 0\\ 0 & I \end{pmatrix}.$$

Let now R be the Jordan Normalform of  $(W^{-1}-cI)$  and N be the Normalform of  $(I-c\tilde{N})^{-1}\tilde{N}$ , this means

$$T_W^{-1}(W^{-1} - cI)T_W = R$$
 and  $T_{\tilde{N}}^{-1}(I - c\tilde{N})^{-1}\tilde{N}T_{\tilde{N}} = N$ 

Considering this definition together with the Neumann-series of  $(I-c\tilde{N})^{-1}$  we obtain

$$\tilde{N}(I - c\tilde{N})^{-1} = \tilde{N}(c\sum_{i=0}^{\infty} \tilde{N}^i) = \tilde{N}(c\sum_{i=0}^{k-1} \tilde{N}^i) = c\sum_{i=0}^{k-1} \tilde{N}^{i+1} = (c\sum_{i=0}^{k-1} \tilde{N}^i)\tilde{N}$$

where we used that  $\tilde{N}$  is nilpotent with nilpotency index k. This shows that  $\tilde{N}$  and  $(I - c\tilde{N})^{-1}$  commute.

From this we can conclude that

$$N^k = [T_{\tilde{N}}^{-1}(I - c\tilde{N})^{-1}\tilde{N}T_{\tilde{N}}]^k = T_{\tilde{N}}^{-1}[(I - c\tilde{N})^{-1}\tilde{N}]^kT_{\tilde{N}} = T_{\tilde{N}}^{-1}(I - c\tilde{N})^{-k}\underbrace{\tilde{N}^k}_{-0}T_{\tilde{N}} = 0$$

We used the commutativity in the third step here. The nilpotent matrix N thus also has the nilpotency index k. A transformation with

$$T_3 := \begin{pmatrix} T_W & 0 \\ 0 & T_{\tilde{N}} \end{pmatrix}$$

transforms  $T_2^{-1}J_{\hat{A}}$  into the Jordan Normalform

$$J_{\tilde{A}} := T_3^{-1} T_2^{-1} J_{\hat{A}} T_3 = T_3^{-1} T_2^{-1} T_1^{-1} \hat{A} T_1 T_3 = \begin{pmatrix} I & 0 \\ 0 & N \end{pmatrix}$$

and  $T_2^{-1}J_{\hat{B}}$  into

$$J_{\tilde{B}} := T_3^{-1} T_2^{-1} J_{\hat{B}} T_3 = T_3^{-1} T_2^{-1} T_1^{-1} \hat{B} T_1 T_3 = \begin{pmatrix} R & 0 \\ 0 & I \end{pmatrix}.$$

Now set

$$P := T_3^{-1} T_2^{-1} T_1^{-1} (cA + B)^{-1}$$
 and  $Q = T_1 T_3$ 

to get the statement.

Using the findings above we are able to transform the initial DAE (3.2) using the matrix P from Theorem 2. By multiplying with P from the left, we obtain

$$PAy'(t) + PBy(t) = Pf(t).$$

Setting

$$y(t) = Q \begin{pmatrix} u(t) \\ v(t) \end{pmatrix}, \quad Pf(t) = \begin{pmatrix} s(t) \\ q(t) \end{pmatrix},$$

with u(t),  $s(t) : \mathbb{R} \to \mathbb{R}^d$  and q(t),  $v(t) : \mathbb{R} \to \mathbb{R}^{n-d}$ .

We get a system of the form

$$u'(t) + Ru(t) = s(t),$$
  
 $Nv'(t) + v(t) = q(t),$  (3.5)

where 
$$PAQ = \begin{pmatrix} I \\ N \end{pmatrix}$$
 and  $PBQ = \begin{pmatrix} R \\ I \end{pmatrix}$ .

We call this system the Weierstraß-Kronecker normalform of the DAE. The first equation of (3.5) is a first order ordinary differential equation and posesses a unique solution u(t) in  $[t_0, t_l]$  for any initial values  $u_0 \in \mathbb{R}^d$ . Now let  $q(t) \in C^{k-1}([t_0, t_l])$ , where  $C^{k-1}(I)$  denotes the space of all functions with domain I that are k-1 times differentiable. We differentiate the second equation in (3.5) and obtain

$$v(t) = q(t) - Nv'(t) = q(t) - N(\underline{q(t) - Nv'(t)})' = q - Nq' + N^{2}v''$$

$$= q - Nq' + N^{2}(q - Nv')'' = q - Nq' + N^{2}q'' - N^{3}v'''$$

$$\vdots$$

$$= q - Nq' + \dots + (-1)^{k-1}N^{k-1} \underbrace{\frac{d^{k}}{dt^{k}}q}_{:=q^{(k-1)}} + (-1)\underbrace{N^{k}v^{(k)}}_{=0}$$

$$= \sum_{i=0}^{k-1} (-1)^{i}N^{i}q^{(i)}(t)$$
(3.6)

where k is the nilpotency index of N. Hence we have an explicit form of the solution of v(t). The Kronecker index k thus tells us how many differentiations are required to obtain an ordinary differential equation. Note that this way of splitting the DAE allows us to differ between algebraic and differential variables.

We can also give a general definition of a DAE in Weierstraß-Kronecker normalform.

**Definition 2** ([5, Definition 13.2.4]). A linear differential equation (3.2) is said to be in Weierstraß-Kronecker normalform if

$$A = \begin{pmatrix} I & 0 \\ 0 & N \end{pmatrix}, \quad B = \begin{pmatrix} R & 0 \\ 0 & I \end{pmatrix},$$

where N is a nilpotent Jordan-block matrix.

**Definition 3.** The nilpotency index k of the matrix N from the Weierstraß-Kronecker Normalform of a matrix pencil  $\{A, B\}$  with A singular is called the Kronecker-Index of  $\{A, B\}$ , which we denote by  $\operatorname{ind}\{A, B\}$ . Note that for A regular we set  $\operatorname{ind}\{A, B\} = 0$ .

**Lemma 1.** The Kronecker-Index ind  $\{A, B\}$  is independent of the choice of the matrices P and Q.

*Proof.* See Lemma 13.2.1 in [5, p.403].

### 3.2 Index of a Differential Algebraic Equation

The index of a DAE gives us insight about its numerical properties and in general about the solvability where the intuition is "The higher the index, the harder to solve".

We will consider two types of index concepts, the differentiation index and the perturbation index.

Since numerical differentiation is an unstable procedure the differentiation index aims to give a measure for the numerical problems to be expected when solving such systems.

**Definition 4** (differentiation index [5, Definition 13.3.1]). *Consider the differential algebraic equation* (3.1) *to be uniquely locally solvable and F sufficiently smooth. For a given*  $m \in \mathbb{N}$  *consider the equations* 

$$F(t, y, y') = 0,$$

$$\frac{dF(t, y, y')}{dt} = 0,$$

$$\vdots$$

$$\frac{d^m F(t, y, y')}{dt^m} = 0.$$

The smallest natural number m for which the above system results in an explicit system of ordinary differential equations (ODEs), i.e. it has the form

$$y' = \phi(t, y),$$

is called differentiation index.

In the previous chapter we have already discussed, that for a DAE with constant coefficients (3.2) and a regular matrix pencil  $\{A, B\}$  we need  $k = ind\{A, B\}$  differentiations to receive an ordinary differential equation. Hence for a DAE with constant coefficients the Kronecker index k is equal to the differentiation index.

**Definition 5** (perturbation index [5, Definition 13.3.3]). Let y(t) be the exact solution of (3.1) and  $\tilde{y}(t)$  be the solution of the perturbed system  $F(t, \tilde{y}, \tilde{y}') = \delta(t)$ . The smallest number  $k \in \mathbb{N}$  such that

$$\|y(t) - \tilde{y}(t)\| \le C \left( \|y(t_0) - \tilde{y}(t_0)\| + \sum_{j=0}^k \max_{t_0 \le \xi \le T} \left\| \int_{t_0}^{\xi} \frac{\mathrm{d}^j \delta}{\mathrm{d}\tau^j}(\tau) d\tau \right\| \right)$$

for all  $\tilde{y}(t)$ , is called the **perturbation index** of this system.

Considering the case of a perturbed linear DAE with constant coefficients and a regular matrix pencil  $\{A, B\}$ 

$$A\tilde{y}'(t) + B\tilde{y}(t) = -\delta(t).$$

Then for the difference of the unperturbed system with the perturbed system we obtain

$$A(y'(t) - \tilde{y}'(t)) + B(y'(t) - \tilde{y}(t)) = \delta(t)$$

Due to 2 we can transform this system into a system of the form (3.5)

$$u'(t) - \tilde{u}'(t) + \hat{R}(u(t) - \tilde{u}(t)) = \delta_1(t),$$
  
$$N(v'(t) - \tilde{v}'(t)) = \delta_2(t).$$

Where u(t) holds the first d entries of  $y(t) : \mathbb{R} \to \mathbb{R}^n$  and v(t) holds the remaining n - d entries. According to (3.6), the solution of the algebraic variable has the form

$$v(t) - \tilde{v}(t) = \sum_{i=0}^{k-1} (-1)^i N^i \delta_2^{(i)}(t)$$

Thus the perturbation index is the same as the Kronecker-index  $ind\{A, B\}$ , see [5]. More over we observe that the error depends on the derivatives of the pertubation.

### 3.3 Consistent Initial Values

Especially for numerical considerations the calculation of initial values that are consistent with the problem is important. This leads us to differentiate between three cases. Namely we need to differ with respect to the index of the DAE. As we will later see, only indices less or equal to 2 are important for our situation, this is why we will consider them especially. For this we also refer to [5, chapter 13.3.1]

from imp funkionensatz usw

**index**  $\nu = 0$ . We have already discussed, that this case leads to an ODE, and Picard-Lindelöff theorem [5, Satz 1.2.1] tells us that every initial condition is consistent for ordinary differential equations.

### Example 2.3 (Consistent initial vaues)

We will later see that this system is of index  $\nu=0$ . This means that we can prescribe arbitrary initial values. As the initial values still have a physical interpretation, still not every consistent initial value also makes sense! We will assume, that at time t=0 the capacitor holds charge, i.e.  $u_1(0)=1$  and that there is no energy stored in the inductance  $i_L(0)=0$ .

**Case: Index** v = 1. By rewriting our system into the form

$$y'(t) = f(t, y(t), z(t)),$$
  
 $0 = g(t, y(t), z(t)).$ 

we are able to give conditions for consistent initial values. Namely  $y_0$  and  $z_0$  are consistent initial values for this system, if  $g(t_0, y_0, z_0) = 0$  holds.

#### **Example 1.3** (Consistent initial vaues)

We will later find, that this system is of index v=1. Due to the fact that  $u_1$  and  $i_V$  are only algebraically constrained, we have to determine their initial values by solving for them at the point t=0. On the other hand  $u_2$  and its derivative both appear, thus we can simply prescribe an initial value for  $u_2$ . Again prescribing an arbitrary initial value might make sense from a mathematical perspective, since we might only be interested in the existence of a solution. But due to the physical interpretation of  $u_2$  as being related to the charge that might already be stored in the capacitor, it makes sense to consider the capacitor to have no initial charge, i.e.  $u_2(0)=0$ . The third line of the MNA-system (2.9) together with setting  $v_{src}=sin(\pi t)$  immediately gives us, that

$$u_1(0) = -\sin(\pi * 0) = 0.$$

From line one follows, that

$$u_1(0) - u_2(0) - i_V(0) = 0$$
,

thus  $i_V(0) = 0$  is a consistent initial value for  $i_V$ .

**Case:** Index  $\nu = 2$ . For index-2 systems we rewrite our system into

$$y' = f(t, y(t), z(t)),$$
  
$$0 = g(t, y(t)).$$

Consistent initial values  $y_0$ ,  $z_0$  for this case not only have to fulfill  $g(t_0, y_0) = 0$  but also the *hidden constraint*  $g_t(t_0, y_0) + g_y(t_0, y_0) f(t_0, y_0, z_0)$ . By  $g_t$  and  $g_y$  we denote the derivative of g with respect to t or y, respectively.

### Example 3.3 (Consistent initial vaues)

Again we will soon investigate the index of this system, it will turn out to be  $\nu=2$ . Considering (2.10), we can define

$$f(t, u_1(t), i_v(t)) = i_v(t) - u_1(t),$$
  
 $g(t, u_1(t)) = u_1(t) + v_{src}(t).$ 

Computing the first condition gives together with the choice  $v_{src} = sin(\pi t)$ 

$$u_1(0) = -v_{src}(0) = -\sin(\pi * 0) = 0.$$

While the second condition leads to

$$0 = v'_{src}(0) + 1 * (i_V(0) - u_1(0)) = \pi cos(\pi * 0) + i_v(0) - 0 = \pi + i_V(0).$$

Thus we obtain  $i_V(0) = -\pi$ .

[2]

In the previous chapter we have seen two different kinds of index concepts for differential algebraic equations. We have also seen that these two, even though they describe rather different structural aspects of the equation, are identical for our use-cases. This leads to the question, what indices we expect obtain for our MNA system.

This chapter answers this question for the case of linear RLC-circuits. As already discussed in Section ??, RLC circuits contain only resistors, inductors and capacities as well as current and voltage sources (RLC-circuits). Linear means that all circuit components are described by linear positive functions. Thus the matrices C, R, and L are positive definite and symmetric.

Recall that we consider the equations resulting from the analysis in Section 2.4. These equations are of the form (3.2): find y such that

$$Ay'(t) + By(t) = f(t).$$

In particular we obtain the equation

$$\begin{pmatrix} A_{C}CA_{C}^{\top} & 0 & 0 \\ 0 & L & 0 \\ 0 & 0 & 0 \end{pmatrix} * \begin{pmatrix} u' \\ i'_{L} \\ i'_{V} \end{pmatrix} + \begin{pmatrix} A_{R}GA_{R}^{\top} & A_{L} & A_{V} \\ -A_{L}^{\top} & 0 & 0 \\ -A_{V}^{\top} & 0 & 0 \end{pmatrix} * \begin{pmatrix} u \\ i_{L} \\ i_{V} \end{pmatrix} = \begin{pmatrix} -A_{I}i_{src} \\ 0 \\ -v_{src} \end{pmatrix}.$$

### 4.1 General Index analysis

Assuming the system only contains linear elements, or is linearized at an operating point in order to investigate the system behaviour, then the corresponding network equation represents a DAE with constant coefficients (3.2). We collect all unknowns in  $x(t) = (u(t), i_L(t), i_V(t))^{\top}$ . The structure of the system is reliant on the matrix B, thus we consider two cases.

**ODE-case:** The matrix B in (3.2) is regular. This is the case iff the circuit contains no voltage sources and every node has a path to ground, that leads through at least one capacitor. Then the system represents a linear-implicit system of ODEs and can be transformed into the explicit ODE system

$$x' = B^{-1}(-Ax + f(t)).$$

Thus we obtain an index of 0.

**DAE-case:** The matrix B in (3.2) is singular. This is the interesting case which we will analyse in the following.

For a singular matrix *B* we have already obtained a representation of the form

$$u'(t) + Ru(t) = s(t),$$
  
 
$$Nv'(t) + v(t) = q(t)$$

in the previous section, see (3.5). We now consider the nilpotency index  $\nu$  of the matrix N, since we already know that this index correlates with the differentiation index as well as the perturbation index. We consider two cases:

### 1. The nilpotency index is 1

Because N is nilpotent with nipotency index  $\nu = 1$ , it holds that  $N^1 = 0$ , thus the system transforms to

$$u'(t) + Ru(t) = s(t),$$
  
$$v(t) = q(t).$$

This means that the algebraic variables are given explicitly. Thus the system is in ODE form.

### 2. The nilpotency index is $\geq 2$

This case is the situation we have discussed in section 3.1.1, where we obtained the representation (3.6)

$$v(t) = \sum_{i=0}^{\nu-1} (-1)^i N^i q^{(i)}(t). \tag{4.1}$$

By differentiating the above equation one more time, we obtain

$$v'(t) = \sum_{i=0}^{\nu-1} (-1)^i N^i q^{(i+1)}(t)$$
(4.2)

with the assumption that q is  $\nu$  times differentiable.

We observe that in both cases the solution v has to fulfill an algebraic constraint. In the case of index v=1 this equation is given explicitly, while for index  $v\geq 2$  it is given implicitly by (4.2).

This system is also sensitive to perturbations. Note that the solution formula (4.1) for v also contains derivatives of the input q. Thus small noise in the input can have arbitrarily large derivatives, which can influence the solution.

## 4.2 Topological Conditions

From analyzing the MNA some conditions to the circuit topology can be obtained. We will be considering the impact of special arrangements of components on the index of

the system. In [9] as well as in [8] and [7] they present results about the index of MNA equations.

**Theorem 3** (Index conditions [8, Theorem 2.2.1]). Let the matrices of the capacitances, inductances and resistances be positive definite.

If

$$ker([A_R, A_C, A_V, A_L]^\top) = \{0\} \quad and \quad ker(A_V) = \{0\}$$
 (4.3)

holds, then the MNA (2.8) leads to a system with index  $v \leq 2$ .

• If additionally

$$ker([A_R, A_C, A_V]^\top) = \{0\} \quad and \quad ker([A_C, A_V]) = \{0\}$$
 (4.4)

*holds, then the system is of index*  $v \leq 1$ 

• *If further* 

$$ker(A_C^{\top}) = \{0\} \quad and \quad dim(v_{src}) = 0$$
 (4.5)

*holds, then the system has index*  $\nu = 0$ .

By ker(A) we denote the nullspace of the matrix A and by dim(y) the size of the vector y. As (4.3) - (4.5) are just conditions concerning the topology of the circuit they can also be interpreted as follows:

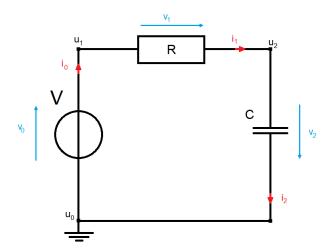
check again, maybe need to add "only contains" somwhere - maybe at index 1 condition

- Condition (4.3) can be interpreted, as the circuit neither containing loops of voltage sources nor cutsets of current sources.
- Condition (4.4) can be interpreted, as the circuit containing neither loops of capacitors and/or voltage sources nor cutsets of inductors and/or current sources.
- Condition (4.5) can be interpreted, as every node in the circuit being connected to the reference node (ground) through a path containing only the capacitors.

We will apply those results to some examples:

### Example 1.4 (Index analysis)

As our first example we consider again the charging of a capacitor with a series resistor as in **Example 1**.



**Figure 4.1:** charging capacitor with series resistor and voltage source (note the difference to 2.1)

Because this circuit only contains one loop, which has a voltage source as well as a capacitor and a resistor we expect the resulting MNA-system to be of index  $\nu = 1$ .

The incidence matrix for this system is

$$A = [A_V, A_R, A_C] = \begin{pmatrix} -1 & 1 & 0 \\ 0 & -1 & 1 \end{pmatrix}.$$

Thus condition (4.3) gives us

$$ker([A_R, A_C, A_V, A_L]^\top) = ker\begin{pmatrix} 1 & -1 \\ 0 & 1 \\ -1 & 0 \end{pmatrix} = \{0\} \text{ and } ker(A_V) = ker\begin{pmatrix} -1 \\ 0 \end{pmatrix} = \{0\}.$$

# 4 Index Analysis of the Modified Nodal Analysis

This condition is fulfilled, thus the index has to be smaller or equal to 2. The condition (4.4) results in

$$ker([A_R, A_C, A_V]^\top) = ker\begin{pmatrix} 1 & -1 \\ 0 & 1 \\ -1 & 0 \end{pmatrix} = \{0\} \text{ and } ker([A_C, A_V]) = ker\begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix} = \{0\}.$$

Thus the index also has to be smaller or equal to 1. Finally condition (4.5) results in

$$ker(A_C^\top) = ker(0 \ 1) = \{0\} \ and \ dim(v_src) \neq 0.$$

This means that this system has index  $\nu = 1$ .

# Example 2.4 (Index analysis)

This example depicts an energy conserving system. The electric energy stored in the capacitor is converted into the the magnetic energy of the coil and vice versa.

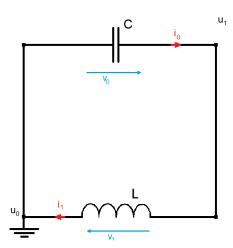


Figure 4.2: LC circuit

Because this circuit only contains one capacitor and one inductance in a singular loop and also has no sources, we expect the resulting MNA-system to be of index  $\nu=0$ . This means that the system is an ODE!

## 4 Index Analysis of the Modified Nodal Analysis

The incidence matrix for this system is

$$A = [A_L A_C] = \begin{pmatrix} -1 & 1 \end{pmatrix}.$$

Thus condition (4.3) gives us

$$ker([A_R, A_C, A_V, A_L]^\top) = ker\begin{pmatrix} -1\\1 \end{pmatrix} = \{0\} \text{ and } ker(A_V) = \{0\}.$$

This condition is fulfilled, thus the index has to be smaller or equal to 2. Condition (4.4) results in

$$ker([A_R, A_C, A_V]^{\top}) = ker(-1) = \{0\} \text{ and } ker([A_C, A_V]) = ker(-1) = \{0\}.$$

This condition is also fulfilled, thus the index has to be smaller or equal to 1. Finally condition (4.5) results in

$$ker(A_C^\top) = ker(-1) = \{0\}$$
 and  $dim(v_src) = 0$ .

Thus it is fulfilled and we obtain index 0.

If we have consider the MNA directly, we also see that this system is of index  $\nu = 0$ , because the MNA has the form

$$Cu_1'-i_L=0,$$

$$Lu_L' + u_1 = 0$$

which is an ODE.

## Example 3.4 (Index analysis)

Finally, as the third example seems very similar to the first one, we might also expect them to have the same index. But this is indeed not the case as the following analysis shows.

# 4 Index Analysis of the Modified Nodal Analysis

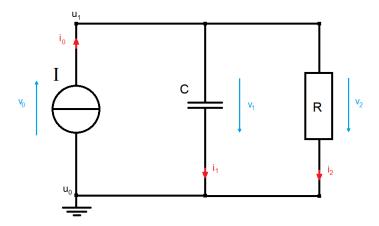


Figure 4.3: current source with resistor and capacitor

This circuit contains two loops, the first with only a voltage source and a capacitor and the second with a capacitor and a resistor. Because of the first loop, we expect this circuit to omit a system of index  $\nu = 2$ .

The incidence matrix for this system is

$$A = \begin{bmatrix} A_V & A_C & A_R \end{bmatrix} = \begin{pmatrix} -1 & 1 & 1 \end{pmatrix}.$$

Thus condition (4.3) gives us (first ker is indeed not 0, no? - would need to not have transposed - maybe because  $A_L$  is empty)

$$ker([A_R, A_C, A_V, A_L]^\top) = ker\begin{pmatrix} 1\\1\\-1 \end{pmatrix} = \{0\} \text{ and } ker(A_V) = ker(-1) = \{0\}$$

This condition is fulfilled, thus the index has to be smaller or equal to 2. Condition (4.4) results in

$$ker([A_R, A_C, A_V]^\top) = ker \begin{pmatrix} 1\\1\\-1 \end{pmatrix} = \{0\} \text{ and } ker([A_C, A_V]) = ker \begin{pmatrix} 1\\-1 \end{pmatrix} \neq \{0\}.$$

Thus this condition is not fulfilled, which means that the system indeed has index  $\nu = 2$ .

As we are interested in finding real valued solutions of our systems we also have to look into solving them numerically. This chapter focuses on the numerical solution of the mentioned systems.

We will first focus on methods used to solve a more general initial value problem of the form:

Find y, such that

$$y'(t) = f(t, y), \quad t \in [t_0, t_1],$$
 (5.1)

$$y(t_0) = y_0. (5.2)$$

For this we presume that the function f(t, y) is continuous and Lipschitz, we can apply the theorem of *Picard-Lindelöf* [5, Satz 1.2.1], which states, that for every initial value  $y_0$ , the initial value problem (5.1) is uniquely solvable in  $[t_0, t_l]$ .

In order to obtain an approximation to this unique solution we have to discretize. We first divide the time-intervall  $[t_0, t_l]$  into

$$t_0 < t_1 < \dots < t_N \le t_1$$

and consider approximations  $y_m \approx y(t_m)$  for m = 1, ..., N, see 5.1. We call this a *time-grid* and we call the difference  $h_{ij} = t_j - t_i$  with  $0 \le i < j \le N$  the *step size* from  $t_i$  to  $t_j$ . We will only consider equidistant grids, i.e.  $h_{ij} = h$  for all  $0 \le i, j \le N$ . A function defined on such a grid is called a *grid-function*.

The grid function  $y_h : T \to A$  with the time-grid  $T = t_0, ..., t_N$  and the approximated values  $A = y_0, ..., y_N$  is called the *numerical approximation* (of y on the time grid).



Figure 5.1: approximation of a function using numerical methods

# 5.1 Single-Step-Methods

The first class of numerical methods we will consider are single-step methods. These methods use the previous approximated value  $y_j$  and (for implicit methods) also the current approximated value  $y_{j+1}$  to determine the current value  $y_{j+1}$  through a *procedural function*.

**Definition 6.** A numerical method to approximate a differential equation 5.1 on a time-grid  $t_0, ..., t_1$  with the intermediate values  $y_0, ..., y_1$  is called a single-step method if it is of the form

$$y_{j+1} = y_j + h_j \phi(t_j, y_j, y_{j+1}, h_j). \tag{5.3}$$

We call  $\phi$  the procedural function. If  $\phi$  does not depend on  $y_{j+1}$ , then the method is called explicit, otherwise it is called implicit.

# 5.1.1 Consistency, Stability and Convergence

In order to compare different single-step methods we have to define some notions to compare their quality. This leads to the definition of the error of the method, its consistency and its convergence. We begin with the definition of the error.

**Definition 7.** Let  $y_{m+1}$  be the result of one step of a single step method (5.3) with the exact start-vector  $y_m = y(t_m)$  then

$$\delta_{m+1} = \delta(t_m + h) = y(t_{m+1}) - \tilde{y}_{m+1}, \quad m = 0, ..., N - 1$$
 (5.4)

is called the local discretization error of the single step method at the point  $t_{m+1}$ .

The local error quantifies the error of every step of the method with respect to the exact solution. In most applications the exact solution is not known. Next we consider the consistency.

**Definition 8.** A single-step method is called consistent if for all initial value problems (5.1)

$$\lim_{h \to 0} \frac{\|\delta(t+h)\|}{h} = 0 \quad \text{for} \quad t_0 \le t \le t_1$$
 (5.5)

holds.

It is called consistent of order p, if for a sufficiently smooth function f

$$\|\delta(t+h)\| \le Ch^{p+1}$$
 for all  $h \in (0,H]$  and  $t_0 \le t \le t_l - h$  (5.6)

holds with C independent of h.

Consistency aims to give insight in how similar the problem (5.3) that the numerical method solves is to the real problem (5.1). Finally we consider convergence of a single step method 6.

**Definition 9.** A single-step method is called convergent, if for all initial value problems (5.1) for the global discretization error

$$e_m = y(t_m) - y_m$$

holds that

$$\max_{m} \|e_m\| \to 0$$
 for  $h_{max} \to 0$ .

The single-step method is called to have the convergence order p, if

$$\max_{m} \|e_m\| \le Ch_{max}^p$$
 for  $h_{max} \in (0, H]$  with  $t_0 \le t_m \le t_1$ 

with the constant C not dependent on the step size h.

As the name suggestes convergence tries to quantify how far off a numerical solution is from the real solution of a system. A very interesting result follows if we additionally require the single-step method to be stable.

**Definition 10.** A single-step method is called (discretely) stable if for grid-functions  $y_h$  and  $\tilde{y}_h$  with

$$y_{i+1} = y_i + h\phi(t_i, y_i),$$
 (5.7)

$$\tilde{y}_{i+1} = \tilde{y}_i + h[\phi(t_i, \tilde{y}_i) + \theta_i], \tag{5.8}$$

and perturbations  $\theta_i = \theta_h(t_i)$  of the right side as well as a bounded perturbation in the initial-values  $y_0 - \tilde{y}_0$  the error is bounded by

$$||y_h - \tilde{y}_h||_{\infty,h} \le C(||y_0 - \tilde{y}_0||_{l^2} + ||\theta_h||_{\infty,h})$$

with a constant C which is not dependent on h. The norm  $\|.\|_{\infty,h}$  denotes the maximum norm over the time-grid, i.e. for a function  $b: T = t_0, ..., t_N \to \mathbb{R}^d$  we have  $\|b\|_{\infty,h} = \max_{t \in T} \|b(t)\|$ ,  $\|b\|$  is the euclidean norm.

For single-step methods which are consistent and stable we obtain the following convergence theorem.

**Theorem 4** (Lax-Richtmyer). A consistent (with order p) and discretely stable single-step method is convergent (with order p).

This theorem is due to Lax and Richtmyer. The converse of this statement is also true.

# 5.1.2 further stability properties

from numpdgl skript

In this section we consider as a model problem the Dahlquist equation, i.e. find *y* such that

$$y' = \lambda y, \quad t > 0 \tag{5.9}$$

$$y(0) = y_0 (5.10)$$

with  $\lambda \in \mathbb{C}$  and  $u_0$  fixed.

**Definition 11.** 1. If a single-step method can be written in the form

$$y_{i+1} = R(z) y_i, \quad z := h\lambda \tag{5.11}$$

*then we call*  $R : \mathbb{C} \to \mathbb{C}$  *the* stability function *of the single-step method.* 

2. The set

$$S := \{ z \in \mathbb{C} : |R(z)| \le 1 \} \tag{5.12}$$

is called the region of stability of the method.

- 3. A single-step method is called
  - 0-stable, if  $0 \in S$ .
  - A-stable, if  $\mathbb{C}^- \subset S$ .
  - L-stable, if  $R(z) \to 0$  for  $Re(z) \to -\infty$ .

# 5.2 Multistep-Methods

based on chapter 4 of book num gew dgl steif nichtsteif

Linear multistep methods (LMSM) use approximations  $y_{m+l}$  along the gridpoints  $t_{m+l}$ , l = 0, 1, ..., k-1 to calculate the new approximation  $y_{m+k}$  at  $t_{m+k}$ . We will first discuss topics

related to the order of the methods depending on its parameters, stability and convergence.

**Definition 12.** For given  $\alpha_0, ..., \alpha_k$  and  $\beta_0, ..., \beta_k$  the iteration rule

$$\sum_{l=0}^{k} \alpha_l y_{m+l} = h \sum_{l=0}^{k} \beta_l f(t_{m+l}, y_{m+l}), \quad m = 0, 1, ..., N - k$$
(5.13)

is called a linear multistep method (linear k-step method). It is always assumed that  $\alpha_k \neq 0$  and  $|\alpha_0| + |\beta_k| > 0$ . If  $\beta_k = 0$  holds, then the method is called explicit, otherwise implicit.

A linear multi-step method consists of two parts:

- 1. In the *starting-phase* approximations  $y_1, ..., y_{k-1}$  for the first k-1 gridpoints  $t_l = t_0 + th$ , l = 1, ..., k-1 have to be calculated. For example we can use as single-step method or a multi-step method with fewer steps.
- 2. In the *run-phase* the multi-step formula is used to determine new approximations  $u_{m+k}$  for the gridpoint  $t_{m+k}$

For theoretical analysis of the multi-step methods we consider the generating polynomials of a multi-step method 12.

$$\rho(x) := \sum_{l=0}^{k} \alpha_l x^l \quad \text{and} \quad \sigma(x) := \sum_{l=0}^{k} \beta_l x^l$$
 (5.14)

# 5.2.1 Consistency, Stability and Convergence

[5]

Again we need to define the properties consistency, stability and convergence for this method.

**Definition 13.** Let  $y_{m+k}$  be the result of one step of the multi-step method (5.13) with the start-values given as the evaluations of the exact solution  $y_{m+l} = y(t_{m+l})$  at  $0 \le l < k$ . This means

$$\alpha_k \tilde{u}_{m+k} = \sum_{l=0}^{k-1} (h\beta_l f(t_{m+l}, y(t_{m+l})) - \alpha_l y(t_{m+l})) + h\beta_k f(t_{m+k}, y_{m+k}).$$

Then

$$\delta_{m+k} = \delta(t_{m+k}) = y(t_{m+k}) - y_{m+k}, \quad m = 0, 1, ..., N - k$$

is called the local discretization error (local error) of the linear multi-step method, see Def. 5.13 at the point  $t_{m+k}$ .

For k = 1 this definition agrees with the definition of the local discretization error for single-step methods 7. If we assign the difference operator

$$L[y(t), h] = \sum_{l=0}^{k} (\alpha_l y(t+lh) - h\beta_l y'(t+lh))$$
 (5.15)

to the linear mutli-step method, we gain the following definition for consistency of such a methods.

**Definition 14.** A linear multi-step method is called consistent, if for all functions  $y(t) \in C^2([t_0, t_l])$ 

$$\lim_{h\to 0} \frac{1}{h} L[y(t), h] = 0$$

*holds.* It has the consistency order p, if for all functions  $y(t) \in C^{p+1}[t_0, t_l]$ 

$$L[y(t), h] = \mathcal{O}(h^{p+1})$$
 for  $h \to 0$ 

holds.

From the generating polynomials (5.14) we can derive simple consistency conditions, i.e.

$$\rho(1) = 0$$
 and  $\rho'(1) = \sigma(1)$ .

Next we also need convergence of such methods.

**Definition 15.** We say that a linear multi-step method is convergent if for a solution y of the problem a solution vector created by an LMSM  $y_i$  for  $j \in 0, ..., k$  we have that

$$\lim_{h\to\infty} \max_{0\leq j\leq k} ||y(t_j)-y_j||=0.$$

The discrete stability is also very similar to the single-step methods.

**Definition 16.** A linear multi-step method is called (discretely) stable, if for solutions  $y_h$  and  $\tilde{y}_h$  of

$$\sum_{l=0}^{k} \alpha_l y_{m+l} = h \sum_{l=0}^{k} \beta_l f(t_{m+l}, y_{m+l}), \tag{5.16}$$

$$\sum_{l=0}^{k} \alpha_{l} \tilde{y}_{m+l} = h \sum_{l=0}^{k} \beta_{l} f(t_{m+l}, \tilde{y}_{m+l}) + h \theta_{n}$$
(5.17)

and bounded initial values  $y_j - \tilde{y}_j$  for  $j \in 0, ..., k$  we have that

$$\max_{t_0 \le t_n \le T} ||y_n - \tilde{y}_n|| \le C \sum_{j=0}^{k-1} ||y_j - \tilde{y}_j|| + \max_{t_0 \le t_n \le T} ||\theta_n||.$$

# 5.2.2 further stability properties

In this section we consider again the Dahlquist test problem as a model problem, find y such that

$$y' = \lambda y, \quad t > 0 \tag{5.18}$$

$$y(0) = y_0 (5.19)$$

with  $\lambda \in \mathbb{C}$  and  $y_0$  fixed.

Thus the resulting linear multistep method is of the form

$$\sum_{l=0}^{k} \alpha_l y_{n+l} = h \sum_{l=0}^{k} \beta_l \lambda y_{n+l}$$

$$\iff \sum_{l=0}^{k} [\alpha_l - h \beta_l \lambda] y_{n+l}$$

Using this we define the following important stability notions.

**Definition 17.** 1. The set

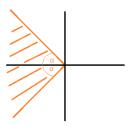
$$S := \{ z \in \mathbb{C} : \rho(\xi) - z\sigma(\xi) = 0 \implies \xi \in \mathbb{C} \text{ and } |\xi| \le 1.$$

$$If \xi \text{ has multiplicity greater than 1, then } |\xi| < 1 \}$$

$$(5.20)$$

is called the region of stability of the method.

- 2. A linear multistep method is called
  - 0-stable, if  $0 \in S$ .
  - stable in the point  $z \in \mathbb{C}$ , if  $z \in S$ .
  - $A(\alpha)$ -stable, if it is stable in all z that lie within the set  $\{z \in \mathbb{C}^- : |arg(z) \pi| \le \alpha\}$  for  $\alpha \in (0, \frac{\pi}{2})$ .



**Figure 5.2:**  $A(\alpha)$  stability for  $\alpha = \frac{\pi}{4}$ 

**Theorem 5** ([5, Satz 4.2.10]). Let f(t, y) be sufficiently smooth and the linear multi-step method be zero-stable and consistent of order p, then it is also convergent of order p.

# 5.3 Implicit linear multi-step formulas

These kinds of multi-step methods are commonly used to numerically solve the systems obtained using modified nodal analysis.

The approach to solve those systems numerically can be split into two main steps:

#### 1. Start-Phase

Here we calculate the first k-1 steps using a method of lower order than the method we want to use ultimately. To achieve the desired convergence rates it is also important to compute these initial steps appropriately, i.e. the methods with which these initial steps are calculated have to also be of the desired order with respect to the step size h. (In applications this can mean that we calculate the initial steps with methods of lower order, but compensate by using a smaller step-size)

### 2. Run-Phase

We successively apply the multi-step method to approximate at new points of the time-grid.

In the following two subchapters we will give a short overview over the two most commonly used methods for solving these kinds of systems.

#### 5.3.1 BDF-k methods

chapter 9.2 numerik book wikipedia

The most commonly used numerical methods for solving the systems that arise in electrical circuits are the BDF-scheme and the trapezoidal rule.

We will not give a deeper look into their construction but will only state their properties. BDF schemes are appealing because they save function evaluations as much as possible, since this is very costly in circuit simulation.

The *backward differentiation formula* (*BDF*) is a family of implicit linear multistep methods. They have the general form

$$\sum_{k=0}^{s} \alpha_k y_{n+k} = h\beta f(t_{n+s}, y_{n+s})$$
(5.21)

Since we are interested in the unknown  $y_{n+s}$  which is used to evaluate f, this method is implicit. The coefficients  $\alpha_k$  and  $\beta$  are chosen, such that the method achieves the best possible order s.

The BDF or BDF-k formulas for k = 1, ..., 3 have the following form

$$k = 1 : hf_{m+1} = y_{m+1} - y_m$$

$$k = 2 : hf_{m+2} = \frac{1}{2}(3y_{m+2} - 4y_{m+1} + y_m)$$

$$k = 3 : hf_{m+3} = \frac{1}{6}(11y_{m+3} - 18y_{m+2} + 9y_{m+1} - 2y_m)$$

Methods with k > 6 are not zero-stable. Indeed in practice, methods with order greater than 3 are rarely used because of their stability properties. The stability-region shrinks with increasing order.

BDF schemes have consistency order p = k. They are methods for solving stiff equations, thus their stability is indicated by their region of absolute stability. Unfortunately not all BDF-schemes are A-stable, but their stability region still contains a large part of the complex left half-plane, see Figure 5.3 They are the most efficient linear multistep methods of this kind.

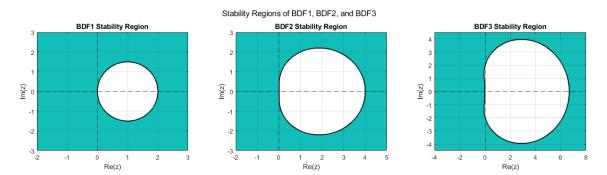


Figure 5.3: stability regions of BDF-schemes

The first timestep is always performed by BDF1 (implicit Euler scheme) as a starting procedure. sure? not really

The Order of consistency of this family of multi-step methods is given by the following theorem

**Theorem 6** ([5, Satz 9.2.1]). *The BDF-k methods have consistency order* p = k.

Due to Theorem 5 and the fact that BDF-k methods are 0-stable, they are also convergent with order k.

## 5.3.2 trapezoidal rule

The trapezoidal rule is a natural alternative to BDF2 since it is A-stable and a linear multistep method of order 2. It also has the smallest leading error coefficient. ([2])

It works by approximating the region under the graph of the function f(x) as a trapezoid, hence the name. It follows that

$$\int_{a}^{b} f(x)dx \approx (b-a)\frac{1}{2}(f(a)+f(b))$$

This procedure is repeated for small subsections of the interval [a, b]. Thus we obtain the iteration formula

$$u_h(t+h) = u_h(t) + \frac{h}{2}[f(t,u_h(t)) + f(t+h,u_h(t+h))].$$

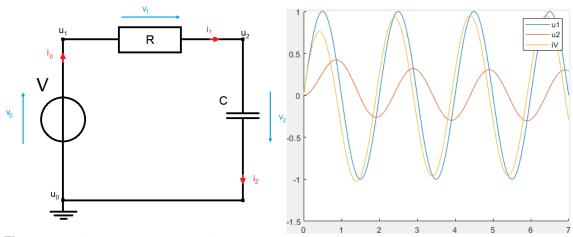
Because  $u_h(t+h)$  appears in f again we see that this is an implicit method. The trapezoidal rule has convergence order p=2 for our applications. For details see [3], where they discuss the more general form of the trapezoidal rule, namely the Lobatto III A methods amongst other methods.

# 5.4 Numerical Examples

In this section we will apply the Trapezoidal rule as well as the BDF-k schemes discussed in Section 5.3.1 and Section 5.3.2 to our three examples. We will also discuss the occurring errors in the solution and compute appropriate consistant inital values for our problems.

# **Example 1.5** (Numerical computation)

We again consider the charging of a capacitor, depicted in the schematics in Figure 5.4.



**Figure 5.4:** charging capacitor with series resistor and voltage source

**Figure 5.5:** Exact solution for example 1.

For testing we set the resistance to R=1 and the capacitance to C=1. In our case we let the voltage source supply a voltage of the form  $v_{src}=sin(\pi t)$ . The resulting MNA-system for this example reads

$$\begin{pmatrix} 0 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 0 \end{pmatrix} * \begin{pmatrix} u' \\ u'_2 \\ i'_0 \end{pmatrix} + \begin{pmatrix} 1 & 1 & -1 \\ 1 & 1 & 0 \\ 1 & 0 & 0 \end{pmatrix} * \begin{pmatrix} u_1 \\ u_2 \\ i_0 \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \\ -sin(\pi t) \end{pmatrix}.$$
 (5.22)

We expect the resulting potentials and currents to be of the form seen in Figure 5.5.

The resulting errors for our numerical tests are displayed in table 5.1, we omit  $u_1$  since it just equals  $-v_{src}$  and thus produces no interesting error results. The reported errors are computed using  $err = \max_{0 \le n \le N} |a_n|t_n - \tilde{y}(t_n)|$  where  $y_n(t_n)$  denotes the numerical approximation and  $\tilde{y}(t_n)$  denoted the exact solution at the time-step  $t_n$ .

| h     | k = 1                  |                        | k = 2                  |                        | k = 3                  |                        | trapezoidal            |                        |
|-------|------------------------|------------------------|------------------------|------------------------|------------------------|------------------------|------------------------|------------------------|
|       | u2                     | iV                     | u2                     | iV                     | u2                     | iV                     | u2                     | iV                     |
| 0.1   | $4.620 \times 10^{-2}$ | $4.620 \times 10^{-2}$ | $9.567 \times 10^{-3}$ | $9.567 \times 10^{-3}$ | $3.057 \times 10^{-3}$ | $3.057 \times 10^{-3}$ | $3.344 \times 10^{-3}$ | $3.344 \times 10^{-3}$ |
| 0.05  | $2.339 \times 10^{-2}$ | $2.339 \times 10^{-2}$ | $2.454 \times 10^{-3}$ | $2.454 \times 10^{-3}$ | $6.083 \times 10^{-4}$ | $6.083 \times 10^{-4}$ | $8.367 \times 10^{-4}$ | $8.367 \times 10^{-4}$ |
| 0.025 | $1.178 \times 10^{-2}$ | $1.178 \times 10^{-2}$ | $6.264 \times 10^{-4}$ | $6.264 \times 10^{-4}$ | $1.672 \times 10^{-4}$ | $1.672 \times 10^{-4}$ | $2.092 \times 10^{-4}$ | $2.092 \times 10^{-4}$ |

**Table 5.1:** Resulting errors for the BDF-k methods and ther trapezoidal rule.

The table 5.1 confirms our results from Section 5.3.1 and 5.3.2, that the BDF-k schemes have convergence order of p = k and the trapezoidal rule has convergence order p = 2.

# Example 2.5 (Numerical computation)

For this example we again want to consider the LC-circuit given in Figure 5.6.

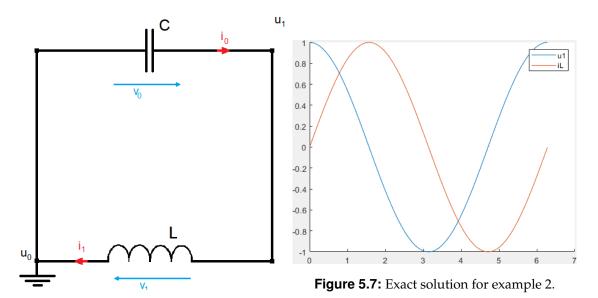


Figure 5.6: LC-circuit

$$\begin{pmatrix} 1 & 0 \end{pmatrix} \begin{pmatrix} u' \end{pmatrix} \begin{pmatrix} 0 & 1 \end{pmatrix} \begin{pmatrix} u_1 \end{pmatrix} \begin{pmatrix} 0 \end{pmatrix}$$

$$\begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} * \begin{pmatrix} u_1' \\ i_L' \end{pmatrix} + \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} * \begin{pmatrix} u_1 \\ i_L \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix}.$$

When applying the BDF-k methods ?? and the Trapezoidal rule 5.3.2 to this problem, we can observe the maximal error  $err = \max_{0 \le t \le T} \|y(t) - \tilde{y}(t)\|$ . This error is displayed in table 5.2.

We again set L = 1 and C = 1. The resulting MNA-system for this example then reads

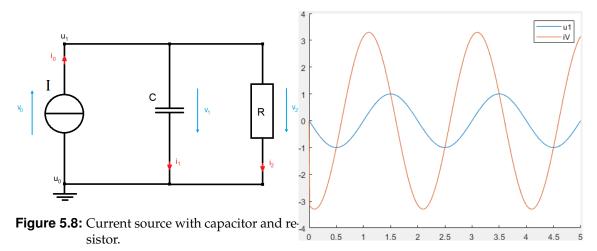
| h     | k = 1                  |                        | k = 2                  |                        | k = 3                  |                        | trapezoidal            |                        |
|-------|------------------------|------------------------|------------------------|------------------------|------------------------|------------------------|------------------------|------------------------|
|       | u1                     | iL                     | u1                     | iL                     | u1                     | iL                     | u1                     | iL                     |
| 0.1   | $7.145 \times 10^{-1}$ | $6.905 \times 10^{-1}$ | $7.763 \times 10^{-2}$ | $8.060 \times 10^{-2}$ | $5.395 \times 10^{-3}$ | $5.180 \times 10^{-3}$ | $1.963 \times 10^{-2}$ | $2.087 \times 10^{-2}$ |
| 0.05  | $4.659 \times 10^{-1}$ | $4.448 \times 10^{-1}$ | $1.964 \times 10^{-2}$ | $2.066 \times 10^{-2}$ | $5.938 \times 10^{-4}$ | $5.579 \times 10^{-4}$ | $4.912 \times 10^{-3}$ | $5.224 \times 10^{-3}$ |
| 0.025 | $2.695 \times 10^{-1}$ | $2.551 \times 10^{-1}$ | $4.924 \times 10^{-3}$ | $5.216 \times 10^{-3}$ | $5.773 \times 10^{-5}$ | $4.740 \times 10^{-5}$ | $1.228 \times 10^{-3}$ | $1.308 \times 10^{-3}$ |

**Table 5.2:** Resulting errors for the BDF-k methods and ther trapezoidal rule.

once again we may see, that the errors reflect the predicted convergence rates.

# Example 3.5 (Numerical computation)

In our last example we consider again the circuit given in Figure 5.8.



**Figure 5.9:** Exact solution for example 3.

Setting R=1, C=1 and  $v_{src}=sin(\pi t)$  the resulting system has the form

$$\begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix} * \begin{pmatrix} u_1' \\ i_V' \end{pmatrix} + \begin{pmatrix} 1 & -1 \\ 1 & 0 \end{pmatrix} * \begin{pmatrix} u_1 \\ i_V \end{pmatrix} = \begin{pmatrix} 0 \\ -sin(\pi t) \end{pmatrix}.$$

With this we expect the potential and current to be given as illustrated in Figure 5.9.

The numerical methods result in errors displayed in Table 5.3. Due to the fact, that  $u_1$  is again the same as  $-v_{src}$ , it is not of interest for our comparison of the errors. We omit  $u_1$ .

| h     | k = 1                  | k = 2                  | k = 3                  | trapezoidal            |  |
|-------|------------------------|------------------------|------------------------|------------------------|--|
|       | iV                     | iV                     | iV                     | iV                     |  |
| 0.1   | $4.894 \times 10^{-1}$ | $1.023 \times 10^{-1}$ | $2.530 \times 10^{-2}$ | $5.219 \times 10^{-2}$ |  |
| 0.05  | $2.462 \times 10^{-1}$ | $2.577 \times 10^{-2}$ | $6.426 \times 10^{-3}$ | $1.295 \times 10^{-2}$ |  |
| 0.025 | $1.233 \times 10^{-1}$ | $6.456 \times 10^{-3}$ | $1.613 \times 10^{-3}$ | $3.232\times10^{-3}$   |  |

**Table 5.3:** Resulting errors for the BDF-k methods and the trapezoidal rule.

The errors lie again within expected convergence orders.

# **Bibliography**

- [1] D. Estévez Schwarz and C. Tischendorf. "Structural analysis of electric circuits and consequences for MNA". In: *International Journal of Circuit Theory and Applications* 28.2 (2000), pp. 131–162. DOI: https://doi.org/10.1002/(SICI)1097-007X(200003/04) 28:2<131::AID-CTA100>3.0.CO;2-W.
- [2] M. Günther, U. Feldmann, and E.J.W. Maten, ter. *Modelling and discretization of circuit problems*. English. CASA-report. Technische Universiteit Eindhoven, 2005.
- [3] Ernst Hairer. The numerical solution of differential-algebraic systems by Runge-Kutta methods. eng. (AT-OBV)AC00001636 1409. Springer, 1989. ISBN: 0387518606. URL: http://media.obvsg.at/AC00065128-1001.
- [4] Chung-Wen Ho, A. Ruehli, and P. Brennan. "The modified nodal approach to network analysis". In: *IEEE Transactions on Circuits and Systems* 22.6 (1975), pp. 504–509. DOI: 10.1109/TCS.1975.1084079.
- [5] Helmut Podhaisky Karl Strehmel Rüdiger Weiner. *Numerik gewöhnlicher Differential-gleichungen Nichtsteife, steife und differential-algebraische Gleichungen*. Vieweg+Teubner Verlag Wiesbaden, 2012. DOI: https://doi.org/10.1007/978-3-8348-2263-5.
- [6] Studiosus Kirchhoff. "Ueber den Durchgang eines elektrischen Stromes durch eine Ebene, insbesondere durch eine kreisförmige". In: *Annalen der Physik* 140.4 (1845), pp. 497–514. DOI: https://doi.org/10.1002/andp.18451400402. eprint: https://onlinelibrary.wiley.com/doi/pdf/10.1002/andp.18451400402. URL: https://onlinelibrary.wiley.com/doi/abs/10.1002/andp.18451400402.
- [7] Timo Reis. "Mathematical Modeling and Analysis of Nonlinear Time-Invariant RLC Circuits". In: *Large-Scale Networks in Engineering and Life Sciences*. Ed. by Peter Benner et al. Cham: Springer International Publishing, 2014, pp. 125–198. ISBN: 978-3-319-08437-4. DOI: 10.1007/978-3-319-08437-4\_2. URL: https://doi.org/10.1007/978-3-319-08437-4\_2.

# **Bibliography**

- [8] Vsevolod Shashkov. "On structure preserving simulations in nonlinear electromagnetics, electric circuits, and efficient treatment of systems with memory". en. PhD thesis. Darmstadt: Technische Universität Darmstadt, June 2024, 133 Seiten. DOI: https://doi.org/10.26083/tuprints-00027452. URL: http://tuprints.ulb.tu-darmstadt.de/27452/.
- [9] Caren Tischendorf. *Topological index calculation of DAEs in circuit simulation*. Humboldt-Universität zu Berlin, Mathematisch-Naturwissenschaftliche Fakultät II, Institut für Mathematik, 2005. DOI: http://dx.doi.org/10.18452/2717.