

DEPARTMENT OF INFORMATICS

TECHNISCHE UNIVERSITÄT MÜNCHEN

Bachelor's Thesis in Informatics: Games Engineering

**Tsunami simulation in the ExaHyPE  
framework**

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## **Tsunami simulation in the ExaHyPE framework**

### **Tsunamisimulation innerhalb des ExaHyPE-Frameworks**

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I confirm that this bachelor's thesis in informatics: games engineering is my own work and I have documented all sources and material used.

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Lisa Scheller

## Acknowledgments

# Abstract

## List of Abbreviations

ODE ordinary differential equation

PDE partial differential equation

## Nomenclature

$\Delta x$  Cell size,  $x_{i+1/2} - x_{i-1/2}$

$f'(q)$  Jacobi matrix of  $f(q)$

$q$  vector of unknowns

$q_x$  partial derivative  $\frac{\delta q}{\delta x}$

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# 1 Introduction

## 1.1 Section

### 1.1.1 Subsection

## 2 The shallow water equations

### 2.1 Hyperbolic partial differential equations

#### 2.1.1 Definition

Partial differential equations (PDEs) are used to describe a multitude of problems. In contrast to ordinary differential equations (ODEs) PDEs contain partial derivatives, whereas ODEs only contain derivatives of one variable. There are a multitude of phenomenons that can be described by PDEs, ranging from fluid mechanics and traffic simulation over thermodynamic problems to quantum mechanics. In this bachelor's thesis we restrict ourselves to a certain subclass of PDEs, the so-called hyperbolic PDEs. The core equations of this thesis, the shallow water equations, which will be introduced in detail in the next section, as well as the standard advection equation, are hyperbolic PDEs.

**Definition 2.1.1.** "A linear system of the form

$$q_t + Aq_x = 0$$

is called hyperbolic if the  $m \times m$  matrix  $A$  is diagonalizable with real eigenvalues." [Lev04, p. 31]

This is a one-dimensional, homogenous, first-order system of PDEs.  $q$  is a vector from  $\mathbb{R}^m$  that contains the unknowns of the equations and  $A$  a  $m \times m$  matrix. The restriction that the matrix has to be diagonalizable means that it must be possible to describe  $A$  in the following way:

$$A = S \cdot D \cdot S^{-1} \tag{2.1}$$

where  $D$  is a diagonal matrix, and  $S$  is a basis transformation onto the basis defined by all eigenvectors of  $A$ .

Other types of PDEs are elliptic or parabolic PDEs.

### 2.1.2 Classification

There are several sub-groups of hyperbolic PDEs, in regard to the eigenvalues of  $A$  as well as in regard to the structure of the equations.

#### Eigenvalues

Since a symmetric matrix is always diagonalizable with real eigenvalues, all symmetric matrices can be used to describe a hyperbolic problem [Lev04]. These equations are called symmetric hyperbolic.

If the matrix has distinct eigenvalues and is diagonalizable, the system of PDEs is called strictly hyperbolic. [Lev04]

In contrast, a matrix that is not diagonalizable but has real eigenvalues is called weakly hyperbolic.[Lev04]

#### Structure of equations

In the simplest case of a hyperbolic PDE  $A$  is not a matrix, but a constant scalar. Equation (??) simplifies to

$$q_t + a q_x = 0 \quad (2.2)$$

and  $a$  has to be real for the system to be hyperbolic.

If the scalar  $a$  or the matrix  $A$  are time-dependent or depend on another value (e.g.  $a(x, t)$ ,  $A(x, p)$ ), they have to fulfill the criteria for a system of hyperbolic PDEs for each value of  $t$  (or in the examples also of  $x$  and  $p$ ).

**Conservation laws** A special group of hyperbolic PDEs are called conservation laws. In the simplest case (1 dimension) they have the following structure:

$$q(x, t)_t + f(q(x, t))_x = 0 \quad (2.3)$$

or in the quasilinear form

$$q_t + f'(q) q_x = 0 \quad (2.4)$$

[Lev04]  $f(q)$  can be linear as well as non-linear in regard to  $q$ .

**The integral formulation** This formulation can also be used to describe a conservation law, in fact it is the more fundamental form to describe one because it also holds across discontinuities.

The integral formulation of a hyperbolic conservation law looks as follows:

$$\frac{d}{dt} \int_{x_1}^{x_2} q(x, t) dx = F(x_2, t) - F(x_1, t) \quad (2.5)$$

$f(x, t)$  represent fluxes going either into or out from the element. The differential formulation ((2.3)), which is used mostly, can be derived from the integral formulation under the assumption that  $q$  is smooth. [Lev04]

**Weak solutions** However, to be able to treat naturally occurring discontinuities, we need to extend the native definition of a solution by the weak formulation. This leads (detailed description [Lev04, p.215]) to the following equation:

$$\int_0^\infty \int_{-\infty}^\infty [q\Phi_t + f(q)\Phi_x] dx dt = - \int_0^\infty q(x, 0)\Phi(x, 0) dx \quad (2.6)$$

**Definition 2.1.2.** "The function  $q(x, t)$  is a weak solution of the conservation law  $q_t + f(q)_x = 0$  with given initial data  $q(x, 0)$  if (2.6) holds for all functions  $\phi$  in  $C_0^1$ ."

[Lev04, p.215] Here  $C_0^1$  describes all functions with compact support that are continuously differentiable.

### 2.1.3 Examples

Only very few hyperbolic PDEs are easy to solve. Their complexity ranges from one-dimensional, linear equations to multidimensional nonlinear PDEs. For most of these problems, especially if they are nonlinear or multidimensional, no analytic solution are known. One example of a very basic hyperbolic PDE is the one-dimensional advection equation with constant speed:

$$q_t + uq_x = 0 \quad (2.7)$$

where  $u$  is the constant velocity. This equation is solved by

$$q(x, t) = q_0(x - ut) \quad (2.8)$$

where  $q_0(x) = q(x, 0)$  are the known initial conditions.

Even though this equation is very simple, its structure still resembles that of more complex problems. So its understanding plays a role at solving more difficult problems.[Lev04]

Another, more complex example are the Maxwell equations

$$\begin{aligned} E_t &= \nabla \times B \\ B_t &= -\nabla \times E \end{aligned} \quad (2.9)$$

([Stu13]), which are used to describe electromagnetic fields.

In general, most problems which are solved by waves or some superposition of waves are described by hyperbolic PDEs. Of course this also includes the description of water waves.

### 2.1.4 Numerical difficulties

The main problem with solving hyperbolic PDEs is to correctly handle the discontinuities that may arise in the solutions. While weak solutions can be analytically found with the criteria listed above, not all weak solutions solve the problem physically correct [Lev04, p.217]. To dismiss those incorrect solutions, the numerical methods often make use of entropy conditions. Also it is usually not possible to simply use the differential formulation and assume it will be sufficiently smooth. Another problem is that the solutions often can be described as several waves of different types. The composition of types depends on the problem and has to be computed first. If the waves interact with each other, it complicates the solution further.

## 2.2 The shallow water equations

This thesis focuses on solving the shallow water equations, which describe water height and momentum of a fluid, in one dimension. For the shallow water equations to be applicable the vertical scale has to be much smaller than the horizontal scale. This is true for ocean waves (period, small wavelength) and even tsunamis (caused by singular event, long wavelength, small amplitude while in deep water, high altitude in shallow water), which are also described by the shallow water equations.

### 2.2.1 Equations in one dimension

The one-dimensional shallow water equations consist of two equations. The first one is:

$$h_t + (uh)_x = 0 \quad (2.10)$$

and the second one is:

$$(hu)_t + (hu^2 + \frac{1}{2}gh^2)_x = 0 \quad (2.11)$$

If  $q(\vec{x}, t) = \begin{pmatrix} q_1 \\ q_2 \end{pmatrix} = \begin{pmatrix} h \\ hu \end{pmatrix}$  is introduced and  $f(q) = \begin{bmatrix} hu \\ hu^2 + \frac{1}{2}gh^2 \end{bmatrix} = \begin{bmatrix} q_2 \\ (q_2)^2/q_1 + \frac{1}{2}gq_1 \end{bmatrix}$  is defined, then the equations can be written as

$$q_t + f(q)_x = 0 \quad (2.12)$$

For smooth solutions, the formulation

$$q_t + f'(q)q_x = 0 \quad (2.13)$$

is also valid.[Lev04]

### 2.2.2 Properties

Since the eigenvalues and eigenvectors of  $f'(q)$  are defining the entire system of equations, they are also basic properties of the problem. They have to be known to be able to compute a numeric or exact solution to any problem. They are given by:

$$\lambda_1 = u - \sqrt{gh} \quad (2.14)$$

and

$$\lambda_2 = u + \sqrt{gh} \quad (2.15)$$

This leads to the eigenvectors

$$\vec{v}_1 = \begin{pmatrix} 1 \\ \lambda_1 \end{pmatrix} \quad (2.16)$$

and

$$\vec{v}_2 = \begin{pmatrix} 1 \\ \lambda_2 \end{pmatrix} \quad (2.17)$$

[Lev04].

If  $h$  is zero, both eigenvalues and thus the eigenvectors have the same value. Because the eigenvalues (and thus the eigenvectors) depend on  $h$  and  $u$ , which are the unknowns of the equation system and can be different for every point in the domain, every of these points may have different eigenvalues. It is not sufficient to compute one set of eigenvalues at the beginning and then never change it again, like with the constant velocity advection equation. The eigenvalues have to be computed for each point or cell in the domain and for each timestep, otherwise the numerical solution disregards the distribution in space as well as the time-development.

### 2.2.3 Example problems

We will look at two well known problems for the Shallow Water Equations. The first one has continuous initial conditions and is the already described gaussian water hump. The second one is the dam break problem, a so called Riemann problem with initial velocities zero and a difference in initial height.

## Riemann problems

**Definition 2.2.1.** A problem with initial conditions of the form

$$q(x, 0) = \begin{cases} q_l & x \leq 0 \\ q_r & 0 \leq x \end{cases}$$

is called a Riemann problem[Lev04].

The main characteristic of a Riemann problem is the discontinuity. To solve such a problem analytically, it has to be fragmented into several states. Two of them are the outer left and right states, where the solution takes the value  $q_l$  or  $q_r$  respectively. The middle states depend on the exact composition of the problem, the possibilities for the cases viewed here are either 2 shock waves, 2 rarefaction waves or one of each. Also a contact discontinuity often exists. In every case a middle value  $q_m$  has to be computed. The exact way to solve a Riemann problem as well as an explanation of the different types of waves will be given in chapter 3. The analytical solution for a dam-break problem will also be computed.

### 2.2.4 Two-dimensional equations

In two dimension the shallow water equations can be written as

$$\begin{aligned} h_t + (hu)_x + (hv)_y &= 0 \\ (hu)_t + (hu^2 + \frac{1}{2}gh^2)_x + (huv)_y &= 0 \\ (hv)_t + (huv)_x + (hv^2 + \frac{1}{2}gh^2)_y &= 0 \end{aligned} \tag{2.18}$$

[Lev04] This differs from the one-dimensional form by the addition of a additional dimension  $y$ . Also the velocity is now described by a vector,  $\vec{u} = \begin{pmatrix} u \\ v \end{pmatrix}$ . The one-dimensional equations can be derived from the two-dimensional case by setting one of the velocities to zero. The two-dimensional equations can be solved using similar techniques as for the one-dimensional case, however for the sake of complexity this thesis will only cover the one-dimensional equations.

## 3 Methods for solving hyperbolic partial differential equations

There are a multitude of numerical methods to solve PDEs. However, some of them are more suited to the task of solving hyperbolic PDEs than others. As pointed out earlier some methods become instable along discontinuities. So it's important to choose a method that handles these discontinuities well if a hyperbolic PDE has to be solved numerically.

### 3.1 Finite Differences

#### 3.1.1 General description

The finite differences method is one of the better known numerical schemes. It approximates derivatives by replacing them with a difference quotient of the form

$$\frac{u(x+h) - u(x)}{h} \quad (3.1)$$

and solving the resulting equation for  $u(x+h)$ . The result can then be integrated by using an integration scheme such as explicit Euler etc. The domain is divided into intervals of length  $h$ , for which the equation is evaluated. Higher accuracy can be achieved by choosing a smaller  $h$  or using a higher order integration scheme such as Runge-Kutta, however this will also impact the definition of the difference quotient.

The computational domain is divided so that the method can be evaluated at certain pre-defined grid points.

#### 3.1.2 Advantages and disadvantages

**Advantages** The method is easy to understand and since only one value is computed per interval it can also be very efficient[JSH08]. Also it is possible to modify the scheme to get high order accuracy.



**Disadvantages** However since the finite difference becomes quite nonsensical at a discontinuity it is not well suited for solving hyperbolic PDEs. Also it doesn't simply allow more complex geometry because it only supports equidistant intervals in its standard form [JSH08]. While schemes with non-equidistant intervals exist, they are quite convoluted. This becomes even more of a problem for solving higher dimensional problems.

### 3.1.3 Example

## 3.2 Finite Volumes

### 3.2.1 General description

The main difference between the method of finite volumes and the method of finite differences is that the first is based on the integral form, whereas the latter is based on the differential form. Using the integral form makes it possible to handle discontinuities. The computational domain is divided in cells, and for each cell the integral over this cell is approximated. For each cell, the cell average is computed, which is the quotient of the approximated integral and the cell volume [Lev04]. This means that each cell stores only this average and not an approximation polynomial or an interpolation polynomial.

To evaluate the change over time now also the fluxes from the neighbor cells have to be considered. There are several possibilities to compute those fluxes, some of which will be explained later.

From now on, the description will be restricted to only one dimension.  
The basic idea is to describe the  $i$ th grid cell as

$$C_i = (x_{i-1/2}, x_{i+1/2}) \quad (3.2)$$

[Lev04] The cell average of the  $i$ th cell in the  $n$ th iteration can now be approximated as:

$$Q_i^n = \frac{1}{\Delta x} \int_{x_{i-1/2}}^{x_{i+1/2}} q(x, t_n) dx = \frac{1}{\Delta x} \int_{C_i} q(x, t_n) dx \quad (3.3)$$

[Lev04], where  $\Delta x$  is the cell size.

If this is applied to the integral formulation of a conservation law, it leads to the following equation:

$$\frac{d}{dt} \int_{C_i} q(x, t) dx = f(q(x_{i-1/2}, t)) - f(q(x_{i+1/2}, t)) \quad (3.4)$$

Since the goal is to propagate the degrees of freedom in time and thus compute the values  $Q_i^{n+1}$  at time  $t^{n+1}$ , the above equation has to be time-integrated. If the result of this is then divided by  $\Delta x$  to get the cell average it leads to an approximation of exactly this.

$$\frac{1}{\Delta x} \int_{C_i} q(x, t^{n+1}) dx = \frac{1}{\Delta x} \int_{C_i} q(x, t^n) dx - \frac{1}{\Delta x} \left[ \int_{t^n}^{t^{n+1}} f(q(x_{i+1/2}, t)) dt - \int_{t^n}^{t^{n+1}} f(q(x_{i-1/2}, t)) dt \right] \quad (3.5)$$

[Lev04]. While this equation would compute exactly what we want it to compute, it can't normally be used because the time integrals on the right side are usually not known. Usually the time integrals over  $\Delta t$  are approximated by some numerical flux in the time interval  $\Delta t$ , which leads to schemes of the following form:

$$Q_i^{n+1} = Q_i^n - \frac{\Delta t}{\Delta x} (F_{i+1/2}^n - F_{i-1/2}^n) \quad (3.6)$$

$F$  is a approximation of the time integrals in (3.5). The choice of this flux approximation heavily impacts the quality of the approximated solution.  $F_{i-1/2} = F((Q_{i-1}, Q_i))$  is called a numerical flux function and it can now be seen that the approximation of  $Q_i^{n+1}$  depends on  $Q_{i-1}$ ,  $Q_i$  and  $Q_{i+1}$ , a so called three point stencil [Lev04].

### 3.2.2 Advantages and disadvantages

**Advantages** One of the main benefits of the method of finite volumes is that it is able to handle problems with discontinuities. Also it is possible to use varying cell sizes, which is useful for more complex geometries. Since the solution depends only on the initial conditions and the fluxes one of the main tasks is to use the flux functions most well-suited for the specific problem. This depends on required accuracy, stability and also on the properties of the problem itself.

**Disadvantages** Since only the cell average is stored for each cell, the approximation inside of one cell is quite rudimentary. It would be more accurate to give the cell more information, e.g. a value at the left boundary and the right boundary of the cell as well as a interpolation function to link those values. The inaccuracy of the cell average can of course be lessened by making the cells smaller, but it is still not as accurate as for example the discontinuous Galerkin method, which will be covered later.

Also if a higher order of accuracy shall be reached, the restrictions on the grid come back again, as described in [JSH08]. This means the method loses one of it's main advantages. If these disadvantages are a problem for the specific problem that should

be solved, and a finite difference scheme is also not the right solution, it might be necessary to take a look at finite elements methods.

### 3.2.3 The CFL condition

The stability of a finite volume method depends heavily on the chosen timestep. If it is too big, the changed information has travelled too far and so the solution depends not only on the neighboring cells but on their neighbors etc. as well. Because those are not part of the numerical scheme the results will become increasingly wrong with time. So a condition to determine the timestep size is necessary.

The Courant-Friedrichs-Lewy-condition (CFL-condition) is such a condition.

**Definition 3.2.1.** "A numerical method can be convergent only if its numerical domain of dependence contains the true domain of dependence of the PDE, at least in the limit as  $\Delta t$  and  $\Delta x$  go to zero"[Lev04]

The CFL condition is a necessary condition for stability and convergence. It leads to the conclusion that the Courant number

$$v = \frac{\Delta t}{\Delta x} \max |\lambda^p| \quad (3.7)$$

, where the  $\lambda^p$  are the wave speeds of the system of equations, has to stay below a certain bound, usually 1 for the finite volume methods we use. From this equation the biggest timestep that still leads to a stable solution can be computed by solving it for  $\Delta t$ . Depending on the problem (for problems with variable wavespeeds) it is not possible to use a stable timestep, since it would become bigger than the maximum allowed timestep over time and so the timestep for the next time iteration has to be computed from the wave speeds at the respective time.

### 3.2.4 Flux approximations

The choice of the flux function is a major part of the finite volumes scheme. The following section will give short overview over some flux functions. Most of these fluxes were also implemented in the program.

**Unstable Flux** This method follows a very simple approach. The flux function is approximated by

$$F_{i-1/2}^n = \frac{1}{2} [f(Q_{i-1}^n) + f(Q_i^n)] \quad (3.8)$$

This would lead to

$$Q_i^{n+1} = Q_i^n - \frac{\Delta t}{2\Delta x} [f(Q_{i+1}^n) - f(Q_{i-1}^n)] \quad (3.9)$$

However, this approach is unusable because it is unstable for hyperbolic problems and all timestep sizes [Lev04].

**Lax-Friedrichs flux** This flux looks very similar to the unstable flux, but it is different enough that it leads to a stable result. It has the advantage that it is still a very simple method that is easy to implement. The flux function is

$$F_{i-1/2} = \frac{1}{2}[f(Q_{i-1}^n) + f(Q_i^n)] - \frac{\Delta x}{2\Delta t}(Q_i^n - Q_{i-1}^n) \quad (3.10)$$

This can then be used in (3.6) to compute the value of  $Q_i$  at the next timestep. While the Lax-Friedrichs flux leads to a stable scheme the order of accuracy still isn't too good, as will be outlined later when the results of the Shallow Water Equations computed with this flux will be compared to the analytic solution.

**Richtmyer Two-Step Lax-Wendroff Method** A better order of accuracy can be achieved by evaluating  $q$  at time  $t + \frac{1}{2}\Delta t$  and using this value to compute the fluxes at time  $t + \Delta t$ . However, this method is less performant than the one-step methods, because for each timestep two values of  $Q$  have to be computed.

**The local Lax-Friedrichs-Flux** The main difference between the local Lax-Friedrichs method is that the factor  $\frac{\Delta x}{\Delta t}$  in the computation of the fluxes is now replaced with another value. In particular this value is not fixed, but changes for every cell and every time. The method for determining the fluxes becomes

$$F_{i-1/2} = \frac{1}{2}(f(Q_{i-1}) + f(Q_i) - a_{i-1/2}(Q_i - Q_{i-1})) \quad (3.11)$$

and the factor  $a$  is determined by

$$a_{i-1/2} = \max(|f(Q_{i-1})|, |f(Q_i)|) \quad (3.12)$$

For the shallow water equations this leads to the biggest eigenvalue (which are  $u \pm \sqrt{gh}$ ) at both indexes ( $i$  and  $i-1$  for the left side and  $i$  and  $i+1$  for the right side). Using the local Lax-Friedrichs flux leads to better results especially concerning rarefaction waves [Lev04]. This method is also known as Rusanov's method. The name "local Lax-Friedrichs method" is based on the similarity to the Lax Friedrichs method and the choice of different factors for each cell, which makes the method local.

## Upwind methods

### 3.3 Finite Elements

#### 3.3.1 General description

The method of finite elements divides the domain in several elements. Even if the general geometry is very complex, the geometry of the elements is a lot simpler and makes it possible to approximate the solution inside of these elements. The solution inside the elements is approximated by local basis functions [JSH08]. The approximated solution is easier to compute for a lower degree of these basis functions, but the accuracy increases with the degree of those functions. A popular choice for the basis functions are the Lagrange polynomials,

$$\ell_i^k(x) = \frac{x - x^{k+1-i}}{x^{k+1} - x^{k+1-i}} \quad (3.13)$$

It is of course possible to choose other interpolation polynomials.

The global formulation is retrieved by choosing a set of test functions that are orthogonal to the residual of the problem. For a Galerkin scheme the test functions span the same spaces as the basis functions. Other choices of test functions are possible. For the standard advection equation this leads to a scheme of the form

$$M \frac{du_h}{dt} + S f_h = 0 \quad (3.14)$$

where

$$M_{ij} = \int_{\Omega} N^i(x) N^j(x) dx \quad (3.15)$$

is the globally defined mass matrices ( $N^i$  are the basis functions) and

$$S_{ij} = \int_{\Omega} N^i \frac{dN^j}{dx} dx \quad (3.16)$$

is the globally defined stiffness matrix [JSH08].

#### 3.3.2 Advantages and disadvantages

**Advantages** This approach has the advantage that it is possible to choose different orders of approximation in different cells [JSH08] and so adapt to the specifications of the problem. Also the choice of the elements size and geometry is very flexible.

**Disadvantages** However, the above choice of basis functions and the requirements they must fulfill means that the matrix  $M$  has to be inverted and the scheme becomes implicit [JSH08]. Also the choice of symmetric basis functions is not helpful for problems which prefer a certain direction, e.g. wave problems [JSH08].

## 3.4 Discontinuous Galerkin methods

### 3.4.1 General description

As the previous sections show, both finite volume methods and finite element methods have some strenghts, but also some weaknesses. The discontinuous Galerkin method is a method that tries to combine these strenghts and eliminate the weaknesses as good as possible.

It uses the same elements as a finite elements method, but the inner nodes are duplicated [JSH08]. The cells/elements are then (in the one-dimensional case) limited to the left by  $u_i$  and to the right by  $u_{i+1}$ . For example the first cell stretches from  $u_0$  to  $u_1$ , while the second cell goes from  $u_1$  to  $u_2$ . This shows why the inner nodes have to be duplicated: The value at the right end of cell 1 doesn't have to be the same value as the left end of cell 2. This means the solution can be discontinuos at the cell interfaces. Inside an element the solution is approximated as

$$u_h^k(x) = \sum_{i=0}^1 u^{k+i} \ell_i^k(x) \quad (3.17)$$

as well as the flux, as explained in [JSH08].

To determine the test functions, the local residual (for a scalar conservation law)

$$R_h(x, t) = \frac{\delta u_h^k}{\delta t} + \frac{\delta f_h^k}{\delta x} \quad (3.18)$$

has to be orthogonal to the test functions, this means the test functions have to fulfill the following equation:

$$\int_{D^k} R_h(x, t) \ell_j^k(x) dx = 0 \quad (3.19)$$

Applying Gauss' theorem and introducing a numerical flux  $f^*$  leads to the equation

$$\int_{D^k} R_h(x, t) \ell_j^k(x) dx = [(f_h^k - f^*) \ell_j^k]_{x^k}^{x^{k+1}} \quad (3.20)$$

Like with finite volumes the choice of the numerical flux is an important part of the method. It is also possible to convert this equation so that it mimics the matrix formulation of finite element methods. The mass and stiffness matrixes are computed analogous to this.

### 3.4.2 Nodal and modal formulation

There a two ways to express the solution when using a discontinuous Galerkin method. These two ways are called the nodal and the modal form.

**Nodal form** The nodal form uses the values at the grid points and interpolates the function between them by using some form of interpolation polynomial [JSH08]:

$$u_h^k = \sum_{i=1}^{N_p} u_h^n(x_i^k, t) \ell_i^k(x) \quad (3.21)$$

In this case the interpolation polynomials are the Lagrange polynomials of order  $N_p + 1$ .

"The global solution  $u(x, t)$  is then assumed to be approximated by the piecewise  $N$ -th order polynomial approximation  $u_h(x, t)$ ,

$$u(x, t) \simeq u_h(x, t) = \oplus_{k=1}^K u_h^k(x, t) \quad (3.22)$$

defined as the direct sum of the  $K$  local polynomial solutions  $u_h^k(x, t)$ ." [JSH08, p.21]  
Since the residual (here for the constant coefficient advection equation)

$$R_h = \frac{\delta u_h}{\delta t} + \frac{\delta a u_h}{\delta x} \quad (3.23)$$

has to be orthogonal to all test functions

$$\phi_h^k(x) = \sum_{n=1}^{N_p} \Phi_n^k \psi_n(x) \quad (3.24)$$

the following statement has to be fulfilled for all  $n \in [0, N_p]$ :

$$\int_{D^k} R_h(x, t) \psi_n(x) dx = 0 \quad (3.25)$$

However since this is only a local statement, it doesn't help in finding the global solution. Spacial integration and the choice of a sensible numerical flux  $(au_h)^*$  (see [JSH08, p.22] leads to the weak form

$$\int_{D^k} \left( \frac{\delta u_h^k}{\delta t} \psi_n - a u_h^k \frac{\delta \psi_n}{\delta x} \right) dx = - \int_{\delta D^k} \hat{n} \cdot (au_h)^* \psi_n dx \quad (3.26)$$

and the strong form

$$\int_{D^k} R_h(x, t) \psi_n(x) dx = \int_{\delta D^k} \hat{n} \cdot (au_h^k - (au_h)^*) \psi_n dx \quad (3.27)$$

where  $\hat{n}$  is the outward pointing normal (+1 and -1 for one dimension) [JSH08] and  $(au_h)^*$  is a chosen numerical flux function.

For a mass matrix

$$\hat{M}_{ij}^k = \int_{D^k} \ell_i \ell_j dx \quad (3.28)$$

and stiffness matrix

$$\hat{S}_{ij}^k = \int_{D^k} \ell_i \frac{d\ell_j}{dx} dx \quad (3.29)$$

the weak form can be written as

$$\hat{M}^k \frac{d}{dt} u_h^k - (\hat{S})^T a u_h^k = -(a u_h)^* \psi(x_r^k) + (a u_h)^* \psi(x_l^k) \quad (3.30)$$

and the strong form as

$$\hat{M}^k \frac{d}{dt} u_h^k + \hat{S}^k a u_h^k = (a u_h^k - (a u_h)^*) \psi(x_r^k) - (a u_h^k - (a u_h)^*) \psi(x_l^k) \quad (3.31)$$

[JSH08]

**Modal form** The modal form uses a local polynomial basis [JSH08]:

$$u_h^k(x, t) = \sum_{n=1}^{N_p} \hat{u}_n^k(t) \psi_n(x) \quad (3.32)$$

where  $\psi_n(x)$  is the polynomial basis function and  $N_p$  is the polynomial order plus one. Following the same approach as above to determine the weak and strong forms for test functions that have the same space as the solution space (called a Galerkin approach) [JSH08, p.23] leads to

$$\hat{M}^k \frac{d}{dt} \hat{u}_h^k - (\hat{S})^T a \hat{u}_h^k = -(a u_h)^* \psi(x_r^k) + (a u_h)^* \psi(x_l^k) \quad (3.33)$$

with local mass matrix

$$\hat{M}_{ij}^k = \int_{D^k} \psi_i \psi_j dx \quad (3.34)$$

and local stiffness matrix

$$\hat{S}_{ij}^k = \int_{D^k} \psi_i \frac{d\psi_j}{dx} dx \quad (3.35)$$

The strong form is given by

$$\hat{M}^k \frac{d}{dt} \hat{u}_h^k + \hat{S}^k a \hat{u}_h^k = (a u_h^k - (a u_h)^*) \psi(x_r^k) - (a u_h^k - (a u_h)^*) \psi(x_l^k) \quad (3.36)$$



### 3.4.3 Handling more complex equations

The explanations until now only used a very basic linear flux function  $au$  and only one unknown. If we assume that now the flux is given by some flux function  $f(u)$ , we can use the same equations ((3.27) and (3.26)) for the weak and strong formulation, only that now  $au_h^k$  gets replaced by  $f_h^k(u_h^k)$  (which is also defined as a sum of coefficients and basis functions) and  $(au_h)^*$  by  $f^*$ . So they become

$$\int_{D^k} \left( \frac{\delta u_h^k}{\delta t} \psi_n - f_h^k(u_h^k) \frac{\delta \psi_n}{\delta x} \right) dx = - \int_{\delta D^k} \hat{n} \cdot f^* \psi_n dx \quad (3.37)$$

and

$$\int_{D^k} \left( \frac{\delta u_h^k}{\delta t} + \frac{\delta f_h^k(u_h^k)}{\delta x} \right) \psi_h^k dx = \int_{\delta D^k} \hat{n} \cdot (f_h^k(u_h^k) - f^*) \psi_h^k dx \quad (3.38)$$

[JSH08, p.31] The generalization into more than one dimension leaves this structure intact, the scalars for  $u$ ,  $\hat{n}$  and  $f$  get replaced by their vector equivalents and the space derivative gets replaced by the gradient (weak form) or the divergence (strong form).

### 3.4.4 Advantages and disadvantages

**Advantages** The discontinuous Galerkin method makes it possible to achieve higher order accuracy even on unstructured grids. Also it can make use of different flux functions, making it possible to choose a flux tailored to the specific problem that shall be solved [JSH08]. Choosing the right flux makes it also possible to reduce oscillations in comparison with finite element methods [JSH08].

**Disadvantages** The computational work is a lot higher than in the simpler methods. This is caused by the doubling of the nodes but also by the more complex numerical computation. Especially with high-order interpolation and thus bigger matrices to compute this is an issue. Also, of course, for problems that don't need the advantages of the discontinuous Galerkin method and are not hit especially hard by the problems of other methods, this computational effort is often not worth it.

### 3.4.5 Comparison between the methods

The following figure gives an overview over the strenghts and weaknesses of each of the methods we analyzed so far. It can be seen that the discontinuous Galerkin method is a good combination of the strengths of finite volume methods and finite element methods.

Figure 3.1: Comparison of the advantages and disadvantages of different methods to solve hyperbolic PDEs, source [JSH08]

	Complex geometries	High-order accuracy and $hp$ -adaptivity	Explicit semi-discrete form	Conservation laws	Elliptic problems
FDM	×	✓	✓	✓	✓
FVM	✓	×	✓	✓	(✓)
FEM	✓	✓	×	(✓)	✓
DG-FEM	✓	✓	✓	✓	(✓)

### 3.5 Overview over the integration methods

The following section aims to give a short overview over the most common time integration methods and their usage without going too much into detail.

#### 3.5.1 Explicit Euler

For a given  $\frac{d}{dt}x = f(t, x)$  with initial value  $x(t_0) = x_0$  the explicit Euler step is

$$x_{n+1} = x_n + hf(t_n, x_n) \quad (3.39)$$

where  $h$  is the timestep size.

While this method is simple to use, it doesn't suit all kinds of numerical problems. Especially problems with high stiffness prefer an implicit method. Also more accurate schemes exist. When using an explicit euler step the time step size is also limited by the CFL-condition (has to be lower than 1).

#### 3.5.2 Implicit Euler

Used especially for problems of high stiffness, the implicit euler step is defined as

$$x_{n+1} = x_n + hf(x_{n+1}, t_{n+1}) \quad (3.40)$$

The new computation of  $f$  for all time steps needs more computational effort, but since the implicit euler method is not restricted in its timestep size by the CFL condition, this may be worth it for some problems.

### 3.5.3 Heun's method

To get a higher rate of convergence (the euler methods have 1, Heun's method has 2) a temporary value  $\tilde{x}_{k+1}$  is computed. The complete method is then given by

$$\begin{aligned}\tilde{x}_{n+1} &= x_n + hf(x_n, t_n) \\ x_{n+1} &= x_n + \frac{1}{2}h(f(t_n, x_n) + f(t_{n+1}, \tilde{x}_{n+1}))\end{aligned}$$

The disadvantage of this method is that for each timestep two values have to be computed.

### 3.5.4 Runge-Kutta method

Runge-Kutta methods are a family of methods with varying number of middle values per timestep and explicit as well as implicit forms. The most common Runge-Kutta method has 4 middle values and is explicit. It has the form

$$x_{k+1} = x_k + h \cdot \frac{1}{6}(k_1 + 2k_2 + 2k_3 + k_4) \quad (3.41)$$

with

$$\begin{aligned}k_1 &= f(x_n, t_n) \\ k_2 &= f(x_n + \frac{h}{2}k_1, t_n + \frac{h}{2}) \\ k_3 &= f(x_n + \frac{h}{2}k_2, t_n + \frac{h}{2}) \\ k_4 &= f(x_n + hk_3, t_n + h)\end{aligned}$$

With this method a rate of convergence of 4 is possible, but at the cost of having to compute 4 extra values per time step.

## 3.6 Analytical solution of the dam-break problem

To evaluate how well the numerical methods used for solving the dam-break problem in the implementation work and how they compare to each other, we need to know the exact solution for this problem. The following section will explain how to solve a given dam-break problem using the shallow water equations. The specific solution for the problem used in the implementation will be computed based on this knowledge in the next chapter.

### 3.6.1 Introduction

As defined earlier, a dam-break problem is a special set of initial conditions, where the initial velocities are zero over the whole domain and the height is piecewise constant, with a left and a right starting height, as in (2.2.1).

The solution can be divided into several areas with identical behaviour. Two of those sections are the areas where the water has still the height before the waves reach it (so the smaller of  $h_l$  and  $h_r$  and another behind the waves where the height is still the bigger one. What happens between those boundaries depends on the properties of the problem, more precisely on which types of waves occur in which order. For the dam-break problem as defined above, the solution always contains one shock wave and one rarefaction wave [Lev04, p.260].

**Shock waves** Physically a shock wave is a wave travelling faster than the speed of sound of the fluid it travels in. It is characterized by a sudden change of value of the examined variables. A well known example is the sonic boom observed when aircraft travel faster than the speed of sound in air. A shock wave also occurs at the end of a traffic jam, where cars' velocity changes abruptly from some positive value to zero. In the case of the shallow water equations, shock waves are discontinuities in height and/or velocity that travel through the water.

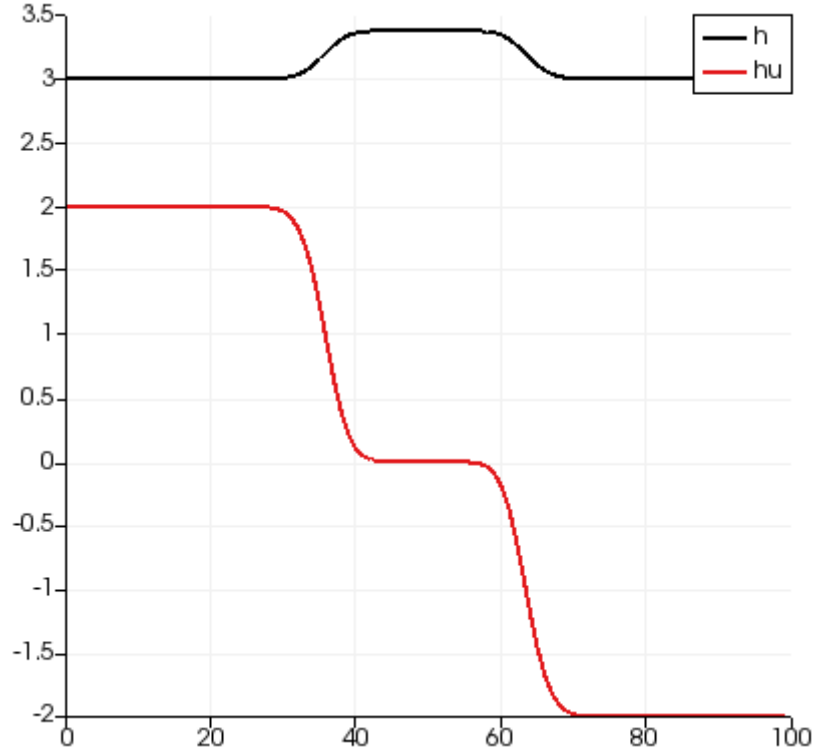
**Rarefaction waves** Rarefaction waves are defined by a gradual increase or decline of (in our case) height and/or velocity, in contrast to the sudden changes associated with shock waves. Classically, they occur behind a travelling shock wave. They usually widen and spread out over time.

### 3.6.2 Different types of wave composition

**Two shocks** If the initial height is uniform over the whole domain and the velocity is piecewise constant with  $u_r = -u_l$ , the solution consists of two shock waves. The approximated shock waves can be seen, especially in  $hu$ . The left shock travels to the left and the right shock to the right.

**Two rarefaction waves** While the example for two shock waves can be explained as two bodies of water colliding, two bodies of water that travel away from each other lead to a solution consisting of two rarefaction waves. To achieve this, the sign of the velocities used in (??) have to be reversed. Since the graphics were made from a numerical approximation using a finite volume scheme the difference between the shocks and rarefaction waves is not as strong as in an exact solution.

Figure 3.2: Paraview visualization of a numerical approximation to a problem with  $u_l = 2$ ,  $u_r = -2$  and  $h_0 = 3$  at time  $t \approx 2.65$  for a computational domain of size 100

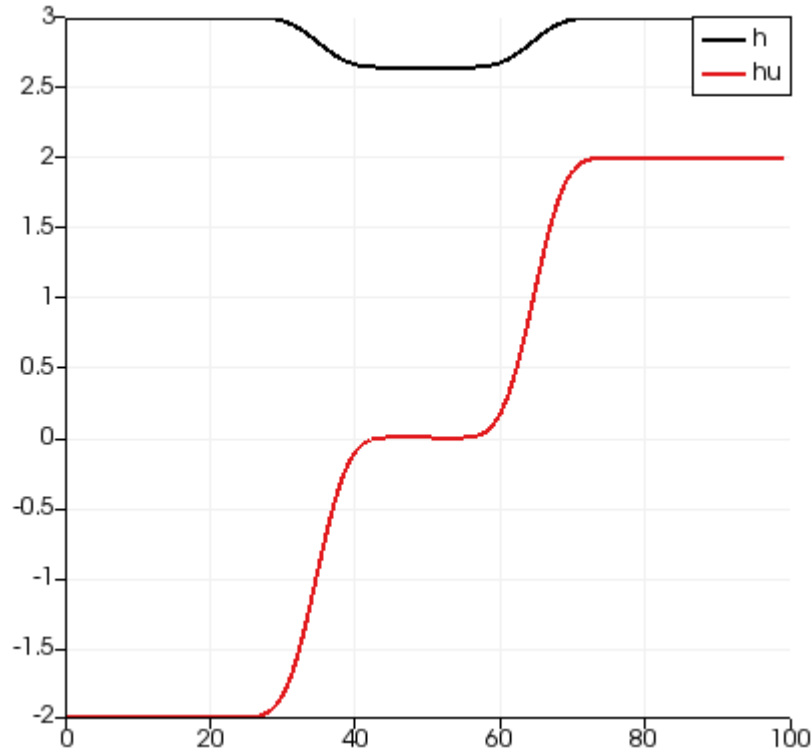


**One rarefaction wave and one shock** The problem that is used in the implementation section is solved by a combination of a shock and a rarefaction wave. The shock travels to the right and the abrupt change in height is clearly visible.

### 3.6.3 Solving the problem

In all types of wave composition the solution can be split in five sections. The first is, as already discussed earlier, the area behind the wave. In ?? this is the left section with  $h = h_l$ . There is also an area in front of the wave, in ?? this is the section with  $h = h_r$ . Between those two areas the two waves impact the variables. Each of them defines the behaviour in one section. Between the waves there is some middle state with properties that depend on the initial problem.

Figure 3.3: Paraview visualization of a numerical approximation to a problem with  $u_l = -2$ ,  $u_r = 2$  and  $h_0 = 3$  at time  $t \approx 2.65$  for a computational domain of size 100



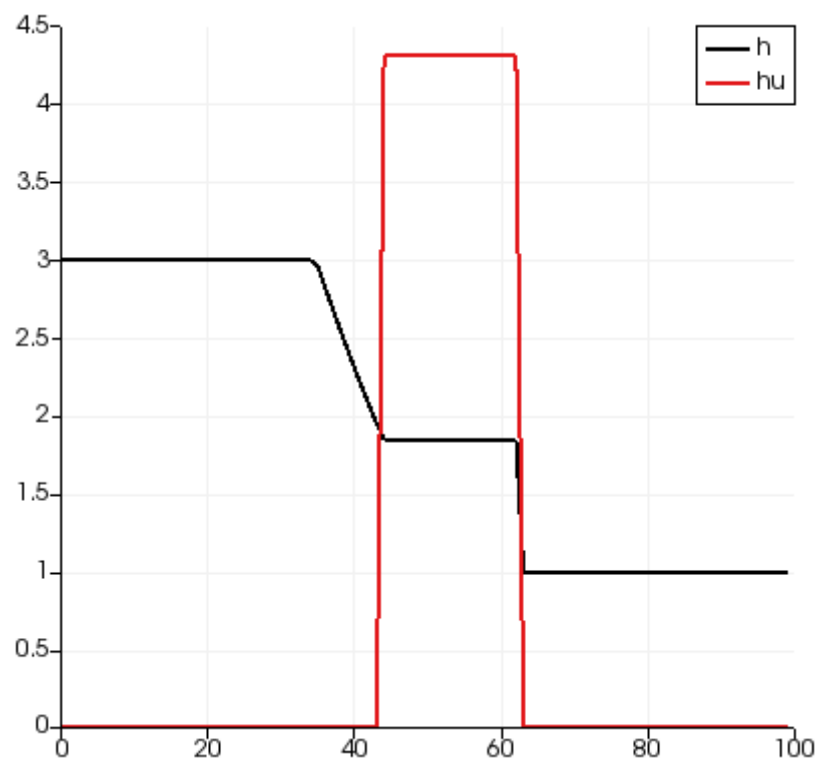
To find an exact solution, the first step is to find out what types of waves occur in the given scenario. To do this, an entropy condition like the Lax condition, which gives requirements for a shock wave, can be used [Lev04, p.285]. Also a way to define which  $x$  value corresponds to which part of the solution is needed. The third step is then to compute the middle state  $q_m$  between the two waves.

**Getting the middle state**

**Determining the boundaries of each section**

**Putting the solution together**

Figure 3.4: Paraview visualization of the exact solution to a problem with  $u_l = u_r = 0$ ,  $h_l = 3$  and  $h_r = 1$  at time  $t \approx 2.65$  for a computational domain of size 100



## 4 Implementation

### 4.1 Computed scenarios

Two types of initial conditions were used, one of them continuous and one discontinuous. The same initial conditions were used for the finite volumes solver as well as for the discontinuous galerkin solver to be able to compare the results.

#### 4.1.1 Continuous initial conditions - gaussian distribution

As an example of continuous initial conditions a water hump with gaussian height distribution was used. The water height (for the advection equation it might also be some density or similar) is described by a gaussian distribution

$$f(x) = \frac{1}{\sqrt{2\pi\sigma^2}} e^{-\frac{(x-\mu)^2}{2\sigma^2}} \quad (4.1)$$

with  $\sigma = 10$ ,  $\mu = 0$  and  $x = pos - x_0$ , where  $x_0$  is the middle of the computational domain. So the height distribution is given by

$$h(x) = \frac{1}{\sqrt{2\pi 100}} e^{-\frac{x^2}{200}} \quad (4.2)$$

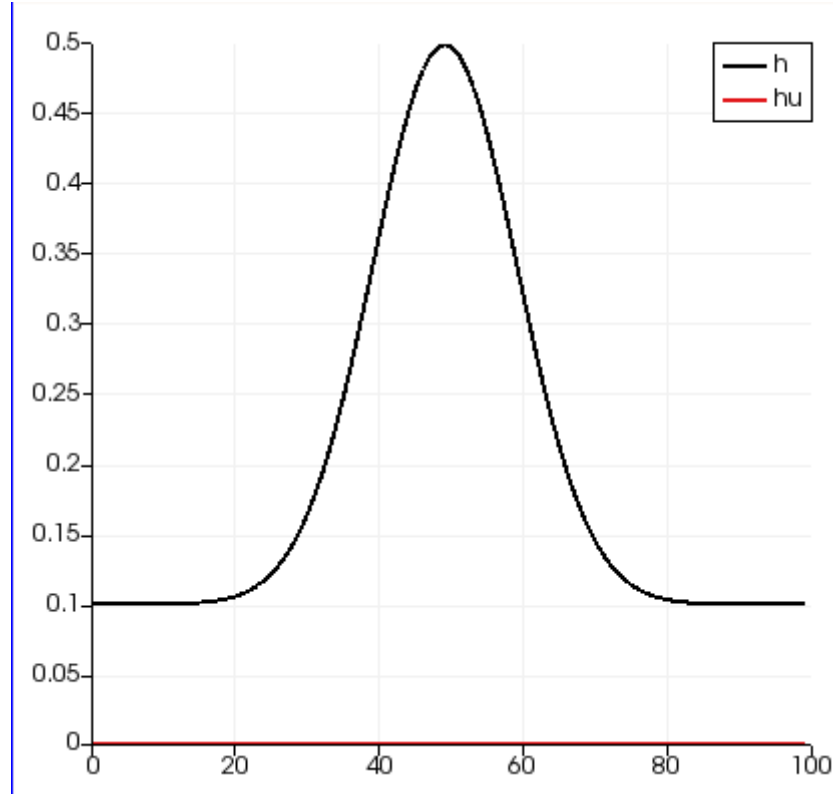
This height is then multiplied by 10 and added to a basic height of 0.1 to prevent the water height becoming zero and disturbing the calculations by using negative values, which might lead to uncontrolled oscillations. The water hump then has its maximum at the middle of the domain with a corresponding height of  $\approx 0.5$ .

**Boundary conditions** For the advection equation periodic boundary conditions were used. Since the wave reaches its original position again, it is possible to use this fact to compute the error over time (for the time where the maximum reaches the middle for the  $n$ -th time) by simply comparing it with the starting conditions.

For the shallow water equations the boundary elements were simply estimated by their direct neighbour and thus the boundary conditions were not periodic.



Figure 4.1: Paraview visualization of the initial conditions following a gaussian distribution for a computational domain of size 100



#### 4.1.2 Discontinuous initial conditions - dam-break problem

To test the numerical solution's ability to handle discontinuities a discontinuous test case was used. This was a classical dam-break problem with initial values

$$h(x,0) = \begin{cases} 3 & x \leq x_0 \\ 1 & x_0 \leq x \end{cases} \quad (4.3)$$

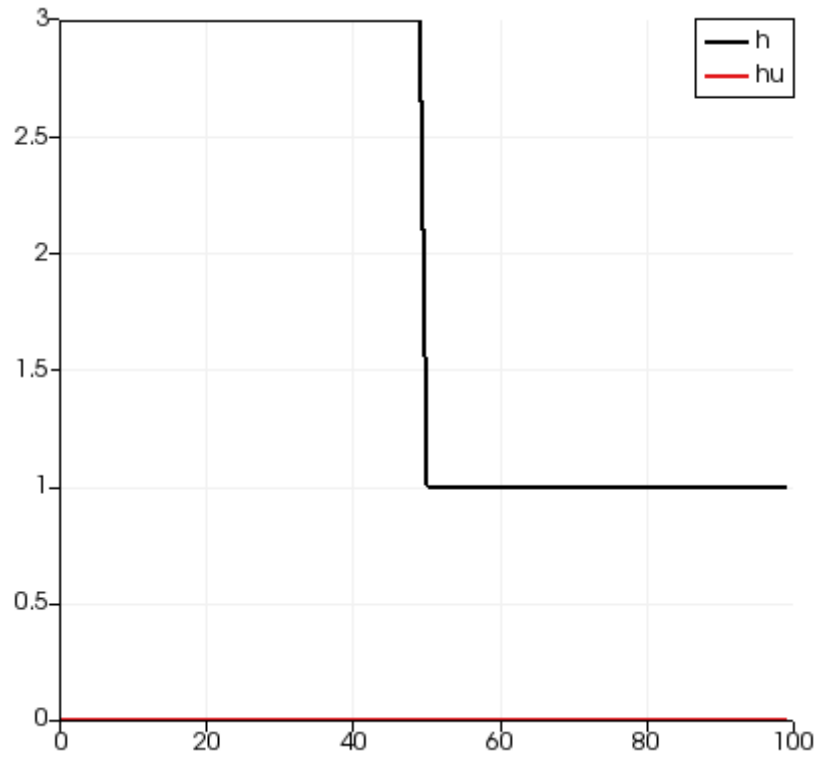
$$hu(x,0) = 0 \quad (4.4)$$

with  $x_0$  defining the middle of the computational domain.

This defines a non-moving area of water, where the left half has height 3 and the left half has height 1.

For this dam-break problem an analytic solution can be found, as given in the next section. This can then be used to compute an error to test the accuracy of the used numerical methods.

Figure 4.2: Paraview visualization of the initial conditions of the used dam-break problem for a computational domain of size 100



**Boundary conditions** Like for the shallow water solution for the gaussian bump, non-periodic boundary conditions were used that simply extrapolate from their direct neighbour.

**4.1.3 Exact solution of the used dam-break problem**

**4.2 Details of implementation**

**4.3 Comparison of finite volume and discontinuos Galerkin**

**4.3.1 Advection equation**

**4.3.2 Shallow Water Equations**

**4.4 Error estimation**

## 5 Conclusion

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