

# **Wave Propagation in Fluids**

*Models and Numerical Techniques*

*Second Edition*

**Vincent Guinot**

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# Table of Contents

## Introduction

## Chapter 1. Scalar Hyperbolic Conservation Laws in One Dimension of Space

- 1.1. Definitions
- 1.2. Determination of the solution
- 1.3. A linear law: the advection equation
- 1.4. A convex law: the inviscid Burgers equation
- 1.5. Another convex law: the kinematic wave for free-surface hydraulics
- 1.6. A non-convex conservation law: the Buckley-Leverett equation
- 1.7. Advection with adsorption/desorption
- 1.8. Summary of Chapter 1

## Chapter 2. Hyperbolic Systems of Conservation Laws in One Dimension of Space

- 2.1. Definitions
- 2.2. Determination of the solution
- 2.3. A particular case: compressible flows
- 2.4. A linear  $2 \times 2$  system: the water hammer equations
- 2.5. A nonlinear  $2 \times 2$  system: the Saint Venant equations

- [2.6. A nonlinear 3×3 system: the Euler equations](#)
- [2.7. Summary of Chapter 2](#)

## **Chapter 3. Weak Solutions and their Properties**

- [3.1. Appearance of discontinuous solutions](#)
- [3.2. Classification of waves](#)
- [3.3. Simple waves](#)
- [3.4. Weak solutions and their properties](#)
- [3.5. Summary](#)

## **Chapter 4. The Riemann Problem**

- [4.1. Definitions – solution properties](#)
- [4.2. Solution for scalar conservation laws](#)
- [4.3. Solution for hyperbolic systems of conservation laws](#)
- [4.4. Summary](#)

## **Chapter 5. Multidimensional Hyperbolic Systems**

- [5.1. Definitions](#)
- [5.2. Derivation from conservation principles](#)
- [5.3. Solution properties](#)
- [5.4. Application: the two-dimensional shallow water equations](#)
- [5.5. Summary](#)

## **Chapter 6. Finite Difference Methods for Hyperbolic Systems**

- [6.1. Discretization of time and space](#)
- [6.2. The method of characteristics \(MOC\)](#)
- [6.3. Upwind schemes for scalar laws](#)
- [6.4. The Preissmann scheme](#)
- [6.5. Centered schemes](#)
- [6.6. TVD schemes](#)
- [6.7. The flux splitting technique](#)
- [6.8. Conservative discretizations: Roe's matrix](#)
- [6.9. Multidimensional problems](#)
- [6.10. Summary](#)
  - [6.10.1. What you should remember](#)
  - [6.10.2. Application exercises](#)

## **Chapter 7. Finite Volume Methods for Hyperbolic Systems**

- [7.1. Principle](#)
- [7.2. Godunov's scheme](#)
- [7.3. Higher-order Godunov-type schemes](#)
- [7.4. EVR approach](#)
- [7.5. Summary](#)

## **Chapter 8. Finite Element Methods for Hyperbolic Systems**

- [8.1. Principle for one-dimensional scalar laws](#)
- [8.2. One-dimensional hyperbolic systems](#)
- [8.3. Extension to multidimensional problems](#)
- [8.4. Discontinuous Galerkin techniques](#)
- [8.5. Application examples](#)
- [8.6. Summary](#)

## **Chapter 9. Treatment of Source Terms**

[9.1. Introduction](#)

[9.2. Problem position](#)

[9.3. Source term upwinding techniques](#)

[9.4. The quasi-steady wave algorithm](#)

[9.5. Balancing techniques](#)

[9.6. Computational example](#)

[9.7. Summary](#)

## **Chapter 10. Sensitivity Equations for Hyperbolic Systems**

[10.1. Introduction](#)

[10.2. Forward sensitivity equations for scalar laws](#)

[10.3. Forward sensitivity equations for hyperbolic systems](#)

[10.4. Adjoint sensitivity equations](#)

[10.5. Finite volume solution of the forward sensitivity equations](#)

[10.6. Summary](#)

## **Chapter 11. Modeling in Practice**

[11.1. Modeling software](#)

[11.2. Mesh quality](#)

[11.3. Boundary conditions](#)

[11.4. Numerical parameters](#)

[11.5. Simplifications in the governing equations](#)

[11.6. Numerical solution assessment](#)

[11.7. Getting started with a simulation package](#)

## **Appendix A. Linear Algebra**

A.1. Definitions

A.2. Operations on matrices and vectors

A.3. Differential operations using matrices and vectors

A.4. Eigenvalues, eigenvectors

## **Appendix B. Numerical Analysis**

B.1. Consistency

B.2. Stability

B.3. Convergence

## **Appendix C. Approximate Riemann Solvers**

C.1. The HLL and HLLC solvers

C.2. Roe's solver

C.3. The Lax-Friedrichs solver

C.4. Approximate-state solvers

## **Appendix D. Summary of the Formulae**

## **Bibliography**

## **Index**



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# Introduction

## **What is wave propagation?**

In a kitchen or in a bathroom, the number of times we turn a tap every day is countless. So is the number of times we watch the liquid stream impacting the sink. The circular flow pattern where the fast and shallow water film diverging from the impact point changes into a deeper, bubbling flow is too familiar to deserve attention. Very few people looking at the circular, bubbling pattern - referred to as a hydraulic jump by hydraulics specialists - are aware that they are contemplating a shock wave.

Closing the tap too quickly may result in a thud sound. This is the audible manifestation of the well-known water hammer phenomenon, a train of pressure waves propagating in the metal pipes as fast as hundreds or thousands of meters per second. The water hammer phenomenon is known to cause considerable damage to hydropower duct systems or water supply networks under the sudden operation of valves, pumps or turbines. The sound is heard because the vibrations of the duct system communicate with the ambient atmosphere, and from there with the operator's ears.

Everyone has once thrown stones into a pond, watching the concentric ripples propagate on the surface. Less visible and much slower than the ripples is the moving groundwater that displaces a pollutant front in a journey that may last for years.

As ubiquitous and familiar as wave propagation may be, the phenomenon is often poorly understood. The reason why intuition so often fails to grasp the mechanisms of wave propagation may lie in the commonly shared, instinctive

perception that waves are made of matter. This, however, is not true. In the example of the hydraulic jump in the sink, the water molecules move across an immobile wave. In the example of the ripples propagating on the free surface of a pond, the waves travel while the water remains immobile.

Waves appear when an object or a system (e.g. the molecules in a fluid, a rigid metallic structure) reacts to a perturbation and transmits it to its neighbors. In many cases, as in the example of the water ripples, the initially perturbed system returns to its initial equilibrium state, while the waves keep propagating. In this respect, waves may be seen as information. The ripples propagating in a pond are a sign that the water molecules “inform” their neighbors that the equilibrium state has been perturbed. A sound is nothing other than information about a perturbation occurring in the atmosphere.

Numerical techniques for wave propagation simulation have been the subject of intensive research over the last 50 years. The advent of fast computers has led to the development of efficient numerical techniques. Engineers and consultants now use simulation software packages for wave propagation on a daily basis. Whether for the purpose of acoustics, aerodynamics, flood wave propagation or contaminant transport studies, computer-based simulation tools have become indispensable to almost all domains of engineering. Such tools, however, remain instruments operated by human beings to execute tedious, repetitive operations previously carried out by hand. They cannot, and hopefully never will, replace the expert’s judgment and experience. Human presence remains necessary for the sound assessment of the relevance and accuracy of modeling results. Such an assessment, however, is possible only provided that the very specific type of reasoning required for the correct understanding of wave propagation phenomena has been acquired.

The main purpose of this book is to contribute to a better understanding of wave propagation phenomena and the most commonly used numerical techniques for its simulation. The first three chapters deal with the physics and mathematics of wave propagation. [Chapters 4, 5 and 10](#) provide insight into more theoretical notions, used in specific numerical techniques. [Chapters 6 to 9](#) are devoted to finite difference, finite volume and finite element techniques. [Chapter 11](#) is devoted to practical advice for the modeler. Basic notions of linear algebra and numerical methods are presented in [Appendices A to C](#). The various formulae used in the present book are summarized in [Appendix D](#).

## **What is the intended readership of this book?**

This book is intended for students of professional and research master's programs and those engaged in doctoral studies, the curriculum of which contains hydraulics and/or fluid mechanics-related subjects. Engineers and developers in the field of fluid mechanics and hydraulics are also a potential target group. This book was written with the following objectives:

(i) To introduce the physics of wave propagation, the governing assumptions and the derivation of the governing equations (in other words, the modeling process) in various domains of fluid mechanics. The application fields are as diverse as contaminant transport, open channel and free surface hydraulics, or aerodynamics.

(ii) To explain how the behavior of the physical systems can be analyzed using very simple mathematical techniques, thus allowing practical problems to be solved.

(iii) To introduce the main families of numerical techniques used in most simulation software packages. As today's

practicing engineers cannot afford not to master modeling packages, a basic knowledge of the existing numerical techniques appears as an indispensable engineering skill.

# How should this book be read?

Most of the chapters are made up of three parts:

- the first part of the chapter is devoted to the theoretical notions applied in the remainder of the chapter;
- the second part deals with the application of these theoretical notions to various hydraulics and fluid mechanics equations;
- the third part provides a summary of the key points developed in the chapter, as well as suggestions of application exercises.

The main purpose of the application exercises is to test the reader's ability to reuse the notions developed in the chapter and apply them to practical problems. The solutions to the exercises may be accessed at the following URL: <http://vincentguinot.free.fr/waves/exercises.htm>.

Try to resist the temptation to read the solution immediately. Solving the exercise by yourself should be the primary objective. The solution text is provided only as an aid, in case you cannot find a way to start and for you to check the validity of your reasoning after completing the exercise.



# Chapter 1

## Scalar Hyperbolic Conservation Laws in One Dimension of Space

### 1.1. Definitions

#### **1.1.1. *Hyperbolic scalar conservation laws***

A one-dimensional hyperbolic scalar conservation law is a Partial Differential Equation (PDE) that can be written in the form:

$$[1.1] \quad \frac{\partial U}{\partial t} + \frac{\partial F}{\partial x} = S$$

where  $t$  and  $x$  are respectively the time- and space-coordinates,  $U$  is the so-called conserved variable,  $F$  is the flux and  $S$  is the source term. Equation [1.1] is said to be the conservation, or divergent, form of the conservation law. The following definitions are used:

- the flux  $F$  is the amount of  $U$  that passes at the abscissa  $x$  per unit time due to the fact that  $U$  (also called the transported variable) is being displaced;
- the source term  $S$  is the amount of  $U$  that appears per unit time and per unit volume, irrespective of the amount transported via the flux  $F$ . If  $U$  represents the concentration

in a given chemical substance, the source term may express degradation phenomena, or radioactive decay.  $S$  is positive when the conserved variable appears in the domain, negative if  $U$  disappears from the domain;

- the conservation law is said to be scalar because it deals with only one dependent variable. When several equations in form [\[1.1\]](#) are satisfied simultaneously, the term “system of conservation laws” is used. Systems of conservation laws are dealt with in [Chapter 2](#).

Only hyperbolic conservation laws are dealt with in what follows. The conservation law is said to be hyperbolic if the flux  $F$  is a function of  $U$  (and none of its derivatives) and, possibly, of  $x$  and  $t$ . Such a dependence is expressed in the form:

$$\left. \begin{aligned} F &= F(U, x, t) \\ S &= S(U, x, t) \end{aligned} \right\} \text{[1.2]}$$

The function  $F(U, x, t)$  is called the “flux function”.

NOTE.- The expression  $F(U, x, t)$  in equation [\[1.2\]](#) indicates that  $F$  depends on  $U$  at the abscissa  $x$  at the time  $t$  and does not depend on such quantities as derivatives of  $U$  with respect to time or space. For instance, the following expression:

$$\text{[1.3]} \quad F = aU$$

is a permissible expression [\[1.2\]](#) for  $F$ , while the following, diffusion flux:

$$\text{[1.4]} \quad F = -D \frac{\partial U}{\partial x}$$

where  $D$  is the diffusion coefficient, does not yield a hyperbolic conservation law because the flux  $F$  is a function of the first-order derivative of  $U$  with respect to space.

In the case of a zero source term, equation [\[1.1\]](#) becomes

$$\text{[1.5]} \quad \frac{\partial U}{\partial t} + \frac{\partial F}{\partial x} = 0$$

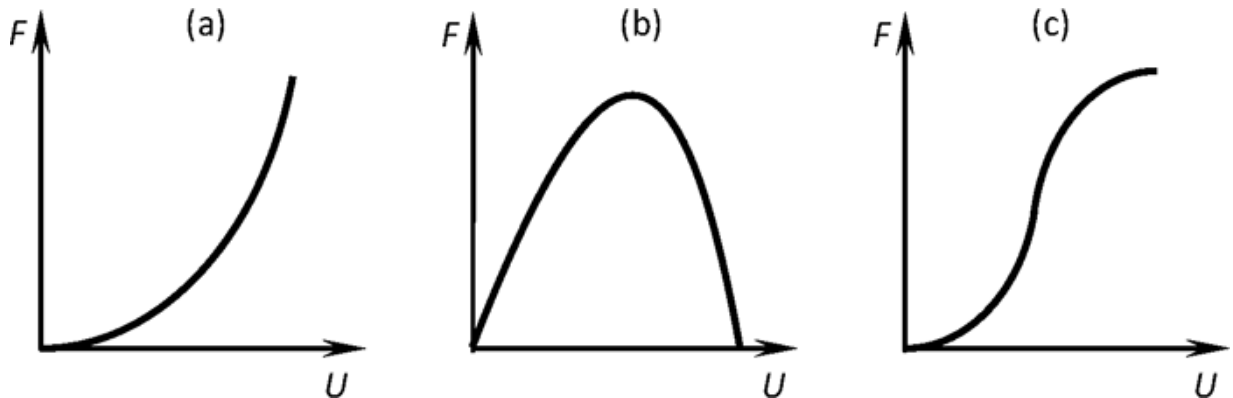
In such a case (see section [1.1.2](#)),  $U$  is neither created nor destroyed over the domain. The total amount of  $U$  over the

domain varies only due to the difference between the incoming and outgoing fluxes at the boundaries of the domain.

Depending on the expression of the flux function, the conservation law is said to be convex, concave or non-convex ([Figure 1.1](#)):

- the law is convex when the second-order derivative  $\partial^2 F / \partial U^2$  of the flux function with respect to  $U$  is positive for all  $U$ ;
- the law is concave when the second-order derivative  $\partial^2 F / \partial U^2$  of the flux function with respect to  $U$  is negative for all  $U$ ;
- the law is said to be non-convex when the sign of the second-order derivative  $\partial^2 F / \partial U^2$  of the flux function with respect to  $U$  changes with  $U$ .

**Figure 1.1.** Typical examples of flux functions: convex (a), concave (b), non-convex (c)

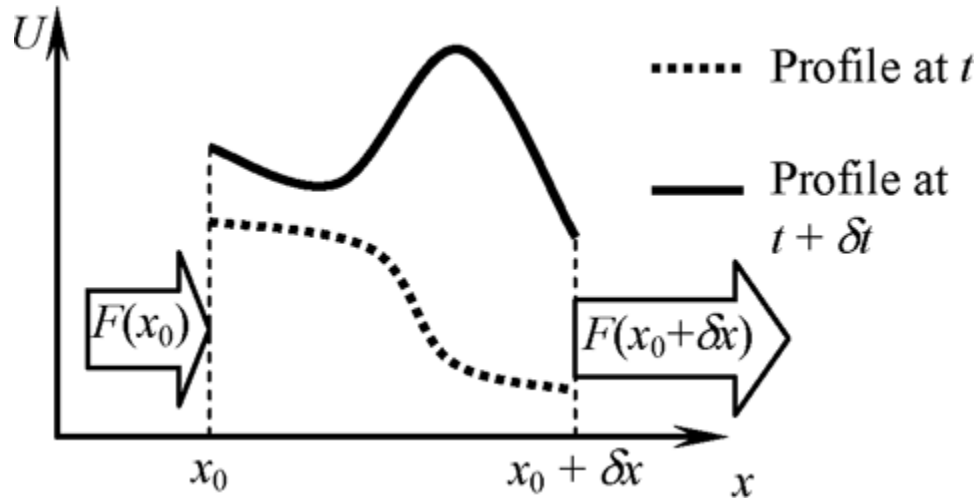


## 1.1.2. Derivation from general conservation principles

The conservation form [\[1.1\]](#) is derived from a balance over a control volume of unit section defined between  $x_0$  and  $x_0 + \delta x$  ([Figure 1.2](#)). The balance is carried out over the control volume between two times  $t_0$  and  $t_0 + \delta t$ . The variation in

the total amount of  $U$  contained in the control volume is then related to the derivatives  $\partial U / \partial t$  and  $\partial F / \partial x$  in the limit of vanishing  $\delta t$  and  $\delta x$ .

**Figure 1.2.** Definition sketch for the balance over a control volume



The total amount  $M(t_0)$  of  $U$  contained in the control volume at  $t = t_0$  is defined as:

$$[1.6] \quad M(t_0) = \int_{x_0}^{x_0 + \delta x} U(x, t_0) dx$$

At  $t = t_0 + \delta t$ , the total amount of  $U$  contained in the control volume is:

$$[1.7] \quad M(t_0 + \delta t) = \int_{x_0}^{x_0 + \delta x} U(x, t_0 + \delta t) dx$$

The variation  $\delta S$  in the amount of  $U$  induced by the source term  $S$  over the domain between  $t_0$  and  $t_0 + \delta t$  is given by:

$$[1.8] \quad \delta S = \int_{t_0}^{t_0 + \delta t} \int_{x_0}^{x_0 + \delta x} S(U, x, t) dx dt$$

The amount  $\delta F(x_0)$  of  $U$  brought by the flux  $F$  across the left-hand side boundary of the control volume between  $t_0$  and  $t_0 + \delta t$  is given by:

$$[1.9] \quad \delta F(x_0) = \int_{t_0}^{t_0 + \delta t} F(x_0, t) dt$$

A quantity  $\delta F(x_0 + \delta x)$  leaves the domain across the right-hand side boundary:

$$[1.10] \quad \delta F(x_0 + \delta x) = \int_t^{t+\delta t} F(x_0 + \delta x, t) dt$$

Stating the conservation of  $U$  over the control volume  $[x_0, x_0 + \delta x]$  between  $t_0$  and  $t_0 + \delta t$ , the following equality is obtained:

$$[1.11] \quad M(t_0 + \delta t) = M(t_0) + \delta F(x_0) - \delta F(x_0 + \delta x) + \delta S$$

Substituting equations [1.6] - [1.10] into equation [1.11] leads to:

$$[1.12] \quad \int_{x_0}^{x_0+\delta x} [U(x, t_0 + \delta t) - U(x, t_0)] dx = \int_{t_0}^{t_0+\delta t} [F(x_0, t) - F(x_0 + \delta x, t)] dt + \int_{t_0}^{t_0+\delta t} \int_{x_0}^{x_0+\delta x} S(x, t) dx dt$$

A first-order Taylor series expansion around  $(x_0, t_0)$  gives:

$$[1.13] \quad \left. \begin{aligned} U(x_0, t_0 + \delta t) - U(x_0, t_0) &= \delta t \frac{\partial U}{\partial t} + O(\delta t^2) \\ F(x_0, t_0) - F(x_0 + \delta x, t_0) &= -\delta x \frac{\partial F}{\partial x} + O(\delta x^2) \end{aligned} \right\}$$

where the quantities  $O(\delta t^2)$  and  $O(\delta x^2)$  are second- or higher-order polynomials with respect to  $\delta t$  and  $\delta x$  respectively. These polynomials contain the second- and higher-order derivatives of  $U$  and  $F$  with respect to  $t$  and  $x$ . When  $\delta t$  and  $\delta x$  tend to zero, the polynomial  $O(\delta t^2)$  becomes negligible compared to the quantity  $\delta t \partial U / \partial t$  because  $\delta t^2$  decreases faster than  $\delta t$ . The polynomial  $O(\delta x^2)$  becomes negligible compared to  $\delta x \partial F / \partial x$  for the same reason. Relationships [1.13] thus become:

$$[1.14] \quad \left. \begin{aligned} U(x_0, t_0 + \delta t) - U(x_0, t_0) &\approx \delta t \frac{\partial U}{\partial t} \\ F(x_0, t_0) - F(x_0 + \delta x, t_0) &\approx -\delta x \frac{\partial F}{\partial x} \end{aligned} \right\}$$

A similar reasoning leads to the following equivalence:

$$[1.15] \quad \int_{t_0}^{t_0+\delta t} \int_{x_0}^{x_0+\delta x} S(x,t) \, dx \, dt \underset{\substack{\delta t \rightarrow 0 \\ \delta x \rightarrow 0}}{=} \delta t \, \delta x \, S$$

Substituting equations [\[1.14\]](#) and [\[1.15\]](#) into equation [\[1.12\]](#) leads to

$$[1.16] \quad \delta t \frac{\partial U}{\partial t} \delta x = -\delta x \frac{\partial F}{\partial x} \delta t + \delta t \delta x \, S$$

Dividing equation [\[1.16\]](#) by  $\delta t \, \delta x$  yields the conservation form [\[1.1\]](#), recalled here:

$$\frac{\partial U}{\partial t} + \frac{\partial F}{\partial x} = S$$

The following remarks can be made:

- the Partial Differential Equation (PDE) [\[1.1\]](#) is a particular case of the more general, integral equation [\[1.12\]](#). Equation [\[1.1\]](#) is obtained from equation [\[1.12\]](#) using the assumption that  $\delta t$  and  $\delta x$  tend to zero. Equation [\[1.12\]](#) is the so-called weak form of equation [\[1.1\]](#) (see [Chapter 3](#) for more details);

- the conservation form [\[1.1\]](#) is based on the implicit assumption that  $F$  is differentiable with respect to  $x$  and  $U$  is differentiable with respect to  $t$ . Consequently, [\[1.1\]](#) is meaningful only when  $U$  is continuous in space and time. In contrast, equation [\[1.12\]](#) is meaningful even when  $U$  is discontinuous in space and/or time. This has consequences on the calculation of discontinuous solutions, as shown in [Chapter 3](#).

### 1.1.3. *Non-conservation form*

Equation [\[1.1\]](#) can be rewritten in the so-called non-conservation form that involves only derivatives of  $U$ . The non-conservation form of equation [\[1.1\]](#) is:

$$[1.17] \quad \frac{\partial U}{\partial t} + \lambda \frac{\partial U}{\partial x} = S'$$

where  $\lambda$  is called the wave speed, and  $S'$  is a source term that may be identical (but not necessarily) to the source

term  $S$  in equation [1.1]. Equation [1.17] is obtained from equation [1.1] by rewriting the derivative  $\partial F / \partial x$  as:

$$[1.18] \quad \frac{\partial F}{\partial x} = \frac{\partial F}{\partial U} \frac{\partial U}{\partial x} + F'$$

where the term  $F' = (\partial F / \partial x)_{U = \text{Const}}$  contains all the derivatives of  $F$  other than the derivative with respect to  $U$ . The expression of  $F$  being known,  $\partial F / \partial U$  and  $F'$  are easily determined. Substituting equation [1.18] into equation [1.1] yields:

$$[1.19] \quad \frac{\partial U}{\partial t} + \frac{\partial F}{\partial U} \frac{\partial U}{\partial x} + F' = S$$

that is:

$$[1.20] \quad \frac{\partial U}{\partial t} + \frac{\partial F}{\partial U} \frac{\partial U}{\partial x} = S - F'$$

Comparing equation [1.20] to equation [1.17] leads to the following definitions for  $\lambda$  and  $S'$

$$[1.21] \quad \left. \begin{aligned} \lambda &= \frac{\partial F}{\partial U} \\ S' &= S - \left( \frac{\partial F}{\partial x} \right)_{U=\text{Const}} \end{aligned} \right\}$$

The expressions of  $F$  and  $S$  being known, the knowledge of  $U$  at any point in time and space allows  $\lambda$  and  $S'$  to be calculated directly. From definition [1.21], in the case where the variations in  $F$  are due to variations in  $U$  only,  $F' = 0$  and  $S'$  is identical to  $S$ .

*Example:* assume that the flux function  $F$  is defined as in equation [1.3], recalled here:

$$F = aU$$

where  $a$  is a function of  $x$  and  $t$ . Equation [1.18] then becomes:

$$[1.22] \quad \frac{\partial F}{\partial x} = \frac{\partial}{\partial x} [a(x, t)U] = a \frac{\partial U}{\partial x} + U \frac{\partial a}{\partial x}$$

and  $\lambda$  and  $F'$  are given by:

$$[1.23] \quad \left. \begin{aligned} \lambda &= a \\ F' &= U \frac{\partial a}{\partial x} \end{aligned} \right\}$$

If  $a$  does not depend on  $x$ ,  $F = 0$  because  $\partial a / \partial x = 0$ .

## 1.1.4. Characteristic form - Riemann invariants

Writing a conservation law in non-conservation form leads to the notions of characteristic form and the Riemann invariant. Such notions are essential to the understanding of hyperbolic conservation laws. A very convenient way of determining the behavior of the solutions of hyperbolic conservation laws consists of identifying invariant quantities (that is, quantities that do not change) along certain trajectories, also called the “characteristic curves” (or more simply the “characteristics”). The solution is calculated by “following” the invariants along the characteristics, which allows the value of  $U$  to be determined at any point. To do so, the non-conservation form [\[1.17\]](#) is used:

$$\frac{\partial U}{\partial t} + \lambda \frac{\partial U}{\partial x} = S'$$

The purpose is to derive the expression of the variation  $\delta U$  in  $U$  observed by an observer travelling at a given speed  $v$ . A small time interval  $\delta t$  is considered, over which the traveler moves by a distance  $\delta x = v \delta t$ . The variation  $\delta U$  “seen” by the observer is given by:

$$\text{[1.24]} \quad \delta U = \frac{\partial U}{\partial t} \delta t + \frac{\partial U}{\partial x} \delta x = \left( \frac{\partial U}{\partial t} + v \frac{\partial U}{\partial x} \right) \delta t$$

Note that from the observer’s point of view,  $U$  is a function of time only, because the observer’s location  $x(t)$  is defined by  $dx/dt = v$ . When  $\delta t$  tends to zero, the ratio  $\delta U / \delta t$  tends to the so-called total derivative  $dU/dt$ . Therefore equation [\[1.24\]](#) becomes:

$$\text{[1.25]} \quad \frac{\delta U}{\delta t} \underset{\delta t \rightarrow 0}{\approx} \frac{dU}{dt} = \frac{\partial U}{\partial t} + v \frac{\partial U}{\partial x} \quad \text{for } \frac{dx}{dt} = v$$

In the particular case of an observer moving at a speed  $\lambda$ , equation [\[1.25\]](#) becomes:



$$[1.26] \quad \frac{dU}{dt} = \frac{\partial U}{\partial t} + \lambda \frac{\partial U}{\partial x} \quad \text{for } \frac{dx}{dt} = \lambda$$

Comparing equations [1.26] and [1.17] leads to:

$$[1.27] \quad \frac{dU}{dt} = S \quad \text{for } \frac{dx}{dt} = \lambda$$

Equation [1.27] is the so-called characteristic form of equation [1.1]. The trajectory, the equation of which is  $dx/dt = \lambda$ , is called a characteristic.  $\lambda$  is called the wave speed.

$S$  being a function of  $U$ ,  $x$  and  $t$ , its value may be calculated at any point  $(x, t)$  if the value of  $U$  is known. The first-order ordinary Differential Equation (ODE) [1.27] is applicable along the characteristic.

In the (important) particular case where the source term  $S$  is zero, equation [1.17] becomes:

$$[1.28] \quad \frac{\partial U}{\partial t} + \lambda \frac{\partial U}{\partial x} = 0$$

and equation [1.27] becomes:

$$[1.29] \quad \frac{dU}{dt} = 0 \quad \text{for } \frac{dx}{dt} = \lambda$$

Equation [1.29] can also be written as:

$$[1.30] \quad U = \text{Const} \quad \text{for } \frac{dx}{dt} = \lambda$$

Consequently, the quantity  $U$  is invariant to an observer moving at the speed  $\lambda$ .  $U$  is called a Riemann invariant.

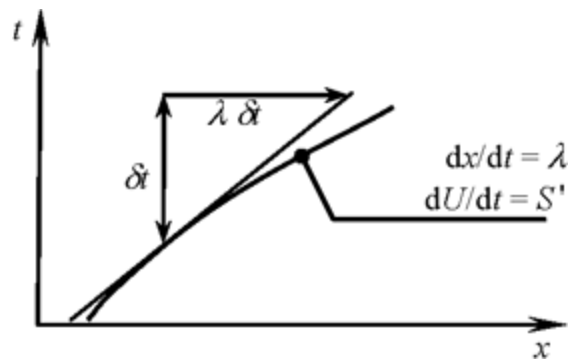
The physical meaning of the wave speed is the following. The wave speed is the speed at which the variations in  $U$  (and not  $U$  itself) propagate. A perturbation appearing in the profile of  $U$  at a given time propagates at the speed  $\lambda$ . The wave speed can be viewed as the speed at which "information", or "signals" created by variations in  $U$ , propagate in space.

## 1.2. Determination of the solution

## 1.2.1. Representation in the phase space

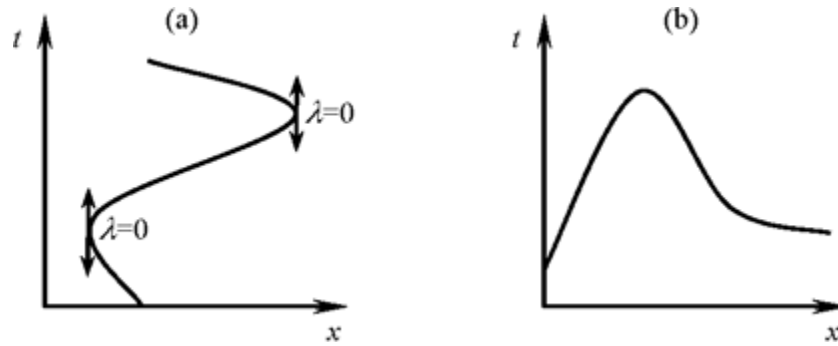
The phase space is a very useful tool in the determination of the behavior of the solutions of hyperbolic conservation laws. The term “phase space” indicates the  $(x, t)$  plane formed by the space coordinate  $x$  and the time coordinate  $t$  ([Figure 1.3](#)).

**Figure 1.3.** Representation of characteristic curves in the phase space



The trajectory  $dx/dt = \lambda$  is represented by a curve in the phase space. The distance  $\delta x$  covered by the characteristic over a time interval  $\delta t$  is given by  $\delta x = \lambda \delta t$ , therefore the slope of the line is  $\delta t/\delta x = 1/\lambda$ . Note that the sign of  $\lambda$  may change with time depending on the variations in  $U$  and the expressions of  $\lambda$  and  $S$ . When  $\lambda$  becomes zero the tangent to the characteristic curve is vertical in the phase space ([Figure 1.4a](#)). In contrast, an extremum with respect to time is not physically permissible ([Figure 1.4b](#)) because “travelling backwards in time” is not possible.

**Figure 1.4.** Physically permissible (a) and non-permissible (b) characteristics



The representation in the phase space may be used to determine the behavior of the solutions of conservation law [1.1] using the so-called “method of characteristics”. The following simple case is considered:

- the source term  $S$  in equation [1.1] is zero;
- the flux depends only on  $U$ , therefore  $F = 0$  in equations [1.18] - [1.20].

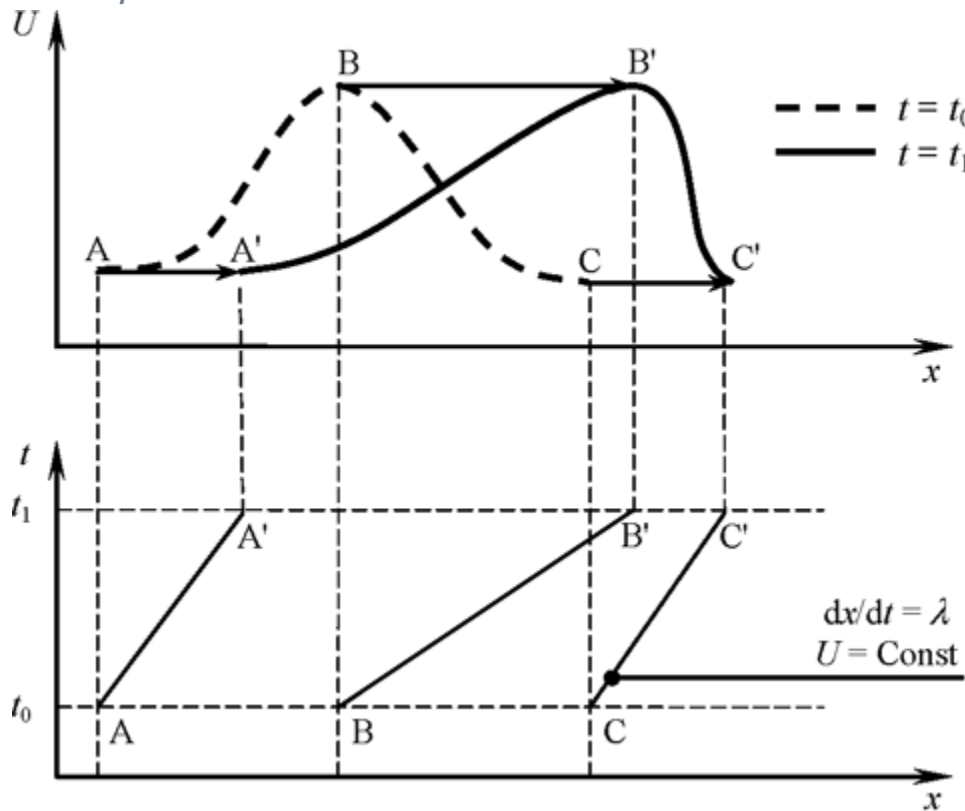
The characteristic form [1.27] then reduces to equation [1.30], recalled here:

$$U = \text{Const for } \frac{dx}{dt} = \lambda$$

$F$  being a function of  $U$  only,  $\lambda$  is also a function of  $U$  only. Consequently, if  $U$  is constant along a characteristic line,  $\lambda$  is also constant and the characteristic is a straight line in the phase space (Figure 1.5). Assume that the profile  $U(x, t_0)$  is known for all  $x$  at the time  $t_0$ . The purpose is to determine the profile  $U(x, t_1)$  for all  $x$  at the time  $t_1 > t_0$ . Consider the point A, the abscissa of which is denoted by  $x_A$ , at which the value of  $U$  at  $(x_A, t_0)$  is denoted by  $U_A$ . Since the wave speed  $\lambda$  depends on  $U$  only, the characteristic passing at A is a straight line. Its (constant) wave speed is  $\lambda_A = \partial F / \partial U (U_A)$ . At time  $t_1$ , the characteristic has moved to point A', the abscissa  $x_{A'}$  of which is given by:

$$[1.31] \quad x_{A'} = x_A + (t_1 - t_0)\lambda_A$$

**Figure 1.5.** Representation of the characteristics in the phase space (bottom) and behavior of the physical profile (top) in the particular case  $F' = S = 0$



From the property of invariance of  $U$  along the characteristic,  $U$  remains unchanged between  $A$  and  $A'$ :

$$[1.32] \quad U_{A'} = U_A$$

Extending the reasoning above to any value of  $x$ , the following relationship is obtained:

$$[1.33] \quad U(x + \lambda \Delta t, t + \Delta t) = U(x, t)$$

where  $\Delta t$  represents the quantity  $(t_1 - t_0)$  and  $\lambda$  is estimated at  $(x, t)$ .

[Figure 1.5](#) shows how the method of characteristics can be used to determine the evolution of a given profile  $[ABC]$ . The figure is drawn assuming that  $\lambda$  is an increasing function of  $U$ . Therefore, point  $B$  moves faster than points  $A$  and  $C$  because  $U_B$  is larger than  $U_A$  and  $U_C$ . Consequently, the region  $[AB]$  tends to spread in time, while the region  $[BC]$

becomes narrower. After a certain amount of time point B catches up with point C and the solution becomes discontinuous at point  $B' = C'$ . The derivatives  $\partial U / \partial t$  and  $\partial U / \partial x$  are no longer defined and a specific treatment must be applied to determine the solution at later times. Such a treatment is detailed in [Chapter 3](#).

In the general case,  $S$  and  $F$  are non-zero. Then relationship [\[1.33\]](#) cannot be used because:

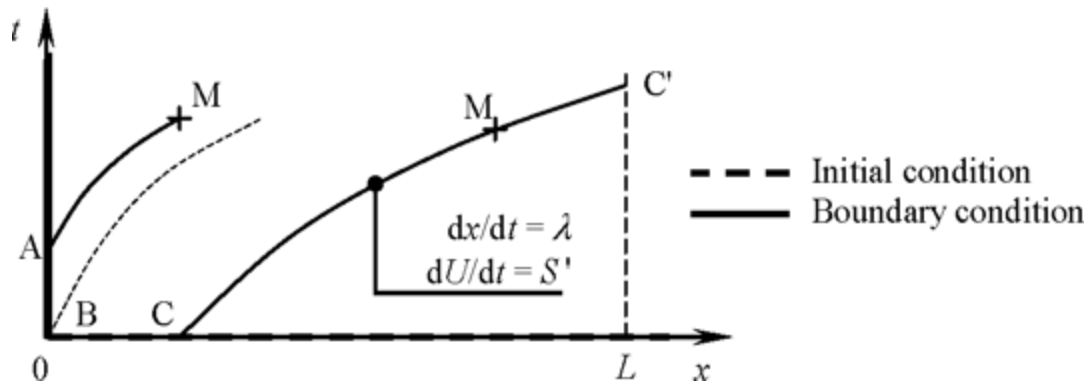
- $U$  is not invariant along a characteristic line;
- the characteristics are therefore curved lines, the slope of which depends on the local value of  $x$  and  $U$ .

Therefore, no simple relationship can be derived between the initial profile at  $t = t_0$  and the final profile at  $t = t_1$ . In most cases, the solution must be computed approximately using numerical methods. Such methods are dealt with in [Chapters 6](#) and [7](#).

## **1.2.2. *Initial conditions, boundary conditions***

In practical applications, the solution of equation [\[1.1\]](#) is sought over a domain of finite length. A key issue is the amount of information needed for the calculation of  $U$  at a point  $M(x, t)$  in the domain. This question is best answered using the phase space ([Figure 1.6](#)). The solution domain is assumed to extend from  $x = 0$  to  $x = L$ .

**[Figure 1.6.](#)** *Initial and boundary conditions in the phase space*



For the sake of clarity, the wave speed  $\lambda$  is assumed to be positive over the entire domain (the case where the sign of the wave speed changes is examined at the end of the section). Two possibilities arise:

- If point M is located on the right-hand side of the characteristic that passes at point B ( $x = 0, t = 0$ ), there exists a point C on the line ( $t = 0$ ) such that the characteristic passing at C passes at M. Point C is called the foot of the characteristic at  $t = 0$ . If the value of  $U$  is known at point C,  $U$  can be computed along the characteristic line by solving the characteristic form [1.27] using any analytical or numerical method. Therefore, the value of  $U$  can be computed at any point M located on the right-hand side of the characteristic that passes at B (0 0), provided that  $U(x, 0)$  is known for all  $x$  between 0 and  $L$ . The function that describes the profile  $U(x, 0)$  is called the initial condition. It is expressed as follows:

$$[1.34] \quad U_0(x) = U(x, t = 0), \quad x \in [0, L]$$

- If point M is located on the left-hand side of the characteristic passing at B, the value of  $U$  at M cannot be calculated from the initial condition and the knowledge of the value of  $U$  at all points A along the line ( $x = 0$ ) is necessary. The function that describes the profile  $U(0, t)$  is called a boundary condition. In the case of a positive  $\lambda$ , the characteristics enter the domain on the left-hand side and the left boundary condition must be used. It is expressed as follows: