

Lecture notes for Math 272, Winter 2021

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These notes will be continuously updated during the course.

The plan for this class is to cover the following topics:

I. Basic theory of Hamilton-Jacobi equations:

1. Existence and long time behavior for the viscous Hamilton-Jacobi equations.
2. Basic viscosity solutions theory for the first order Hamilton-Jacobi equations.
3. The Lions-Papanicolaou-Varadhan theorem and applications to periodic homogenization.
4. Long time behavior for the Lax-Oleinik semigroup, and very rudimentary aspects of the Fathi theory.

II. Hamilton-Jacobi equations with a constraint and applications to the biological modeling.

III. An introduction to mean-field games, based on the lecture notes by P. Cardaliaguet and A. Porretta.

Part I of these lecture notes is a draft of a chapter in a book in preparation with Sasha Kiselev and Jean-Michel Roquejoffre. The preliminary version of the draft of this chapter was written mostly by Jean-Michel. All mistakes are, obviously, mine.

The draft will be updated as we go, potentially with major re-writes back and forth. Because of that, I plan to update the lecture notes after each lecture, to reflect what was actually presented in class, and not upload the full draft of Chapter 2 of these notes from the start.

In addition, I include Chapter 1 (which is actual Chapter 2 of the book draft) in the lecture notes because some of the results of that chapter will be used in class, and it is easy to refer to them in this way. However, this content is included solely for your convenience, the class will not cover that chapter and will start with Chapter 2 of these notes (which is Chapter 3 of the book draft).

The texts of Chapter 1 and 2 have not been finalized so all comments are extremely welcome!

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Chapter 1

Diffusion equations

1.1 Introduction to the chapter

Parabolic equations of the form

$$\frac{\partial u}{\partial t} - \sum_{i,j=1}^n a_{ij}(x) \frac{\partial^2 u}{\partial x_i \partial x_j} + \sum_{j=1}^n b_j(x) \frac{\partial u}{\partial x_j} = f(x, u, \nabla u), \quad (1.1.1)$$

are ubiquitous in mathematics and various applications in physics, biology, economics and other fields. While there are many textbooks on the subject, ranging from the most elementary to extremely advanced, most of them concentrate on the highly non-trivial questions of the existence and regularity of the solutions. We have chosen instead to focus on some striking qualitative properties of the solutions that, nevertheless, can be proved with almost no background in analysis, using only the very basic regularity results. The unifying link in this chapter will be the parabolic maximum principle and the Harnack inequality. Together with the parabolic regularity, they will be responsible for the seemingly very special behavior that we will observe in the solutions of these equations.

The chapter starts with an informal probabilistic introduction. While we do not try to motivate the basic diffusion equations by models in the applied sciences here, an interested reader would have no difficulty finding the connections between such equations and models in physics, biology, chemistry and ecology in many basic textbooks. On the other hand, the parabolic equations have a deep connection with probability. Indeed, some of the most famous results in the parabolic regularity theory were proved by probabilistic tools. It is, therefore, quite natural to start the chapter by explaining how the basic linear models arise, in a very simple manner, from limits of a random walk. We reassure the reader that the motivation from the physical or life sciences will not be absent from this book, as some of the later chapters will precisely be motivated by problems in fluid mechanics or biology. We also keep the probabilistic considerations to an elementary level, without any use of stochastic analysis.

The probabilistic section is followed by a brief interlude on the maximum principle. There is nothing original in the exposition, and we do not even present the proofs, as they can be found in many textbooks on PDE. We simply recall the statements that we will need.

We then proceed to the section on the existence and regularity theory for the nonlinear heat equations: the reaction-diffusion equations and viscous Hamilton-Jacobi equations. They

arise in many models in physical and biological sciences, and our "true" interest is in the qualitative behavior of their solutions, as these reflect the corresponding natural phenomena. However, an unfortunate feature of the nonlinear partial differential equations is that, before talking knowledgeably about their solutions or their behavior, one first has to prove that they exist. This will, as a matter of fact, be a non-trivial problem in the last two chapters of this book, where we look at the fluid mechanics models, for which the existence of the solutions can be quite subtle. As the reaction-diffusion equations that we have in mind here and in Chapter ?? both belong to a very well studied class and are much simpler, it would not be inconceivable to brush their existence theory under the rug, invoking other books. This would not be completely right, for several reasons. The first is that we do not want to give the impression that the theory is inaccessible: it is quite simple and can be explained very easily. The second reason is that we wish to explain both the power and the limitation of the parabolic regularity theory, so that the difficulty of the existence issues for the fluid mechanics models in the latter chapters would be clearer to the reader. The third reason is more practical: even for the qualitative properties that we aim for, we still need to estimate derivatives. So, it is better to say how this is done.

The next section contains a rather informal guide to the regularity theory for the parabolic equations with inhomogeneous coefficients. We state the results we will need later, and outline the details of some of the main ideas needed for the proofs without presenting them in full – they can be found in the classical texts we mention below. We hope that by this point the reader will be able to study the proofs in these more advanced textbooks without losing sight of the main ideas. This section also contains the Harnack inequality. What is slightly different here is the statement of a (non-optimal) version of the Harnack inequality that will be of an immediate use to us in the first main application of this chapter, the convergence to the steady solutions in the one-dimensional Allen-Cahn equations on the line. The reason we have chosen this example is that it really depends on nothing else than the maximum principle and the Harnack inequality, illustrating how far reaching this property is. It is also a perfect example of how a technical information, such as bounds on the derivatives, has a qualitative implication – the long time behavior of the solutions.

The next section concerns the principal eigenvalue of the second order elliptic operators, a well-treated subject in its own right. We state the Krein-Rutman theorem and, just to show the reader that we are not using any machinery heavier than the results we want to prove, we provide a proof in the context of the second order elliptic and parabolic operators. It shares many features with the convergence proof of the next section, without its sometimes technically involved details. We hope the reader will realize the ubiquitous character of the ideas presented.

We end the chapter with the study of reaction-diffusion fronts. While it is, in its own right, a huge subject that is still advancing at the time of the writing of this chapter, we have decided that talking about them was a good way to follow the main pledge of this book: show the reader results that are interesting and representative of the theory, while not being the most advanced or up-to-date. With nothing else than the tools displayed in this chapter, we will see that we can say a lot about the large time organization of this class of models, a striking example being the convergence to pulsating waves: periodicity in space will generate a sort of time periodicity for the solutions.

This chapter is quite long so we ask the reader to be prepared to persevere through the

more technical places, but we feel that it is worth showing how far one may go with the sole aid of the maximum principle and a few estimates. We hope that in the end the reader will find the effort rewarding.

A note on notation. We will follow throughout the book the summation convention: the repeated indices are always summed over, unless specified otherwise. In particular, we will usually write equations such as (1.1.1) in the form

$$\frac{\partial u}{\partial t} - a_{ij}(x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(x) \frac{\partial u}{\partial x_j} = f(x, u, \nabla u), \quad (1.1.2)$$

or

$$\frac{\partial u}{\partial t} - a_{ij}(x) \partial_{x_i} \partial_{x_j} u + b_j(x) \partial_{x_j} u = f(x, u, \nabla u). \quad (1.1.3)$$

We hope the reader will get accustomed to this convention sufficiently fast so that it causes no confusion or inconvenience.

1.2 A probabilistic introduction to the evolution equations

Let us explain informally how the linear equations of the form (1.1.2), with $f \equiv 0$ arise from random walks, in a very simple way. One should emphasize that many of the qualitative properties of the solutions to the parabolic and integral equations, such as the maximum principle and regularity, on a very informal level, are an "obvious" consequence of the microscopic random walk model. For simplicity, we will mostly consider the one-dimensional case, the reader can, and should, generalize this approach to higher dimensions – this is quite straightforward.

Discrete equations and random walks

The starting point in our derivation of the evolution equations is a discrete time Markov jump process $X_{n\tau}$, with a time step $\tau > 0$, defined on a lattice with mesh size h :

$$h\mathbb{Z} = \{0, \pm h, \pm 2h, \dots\}.$$

The particle position evolves as follows: if the particle is located at a position $x \in h\mathbb{Z}$ at the time $t = n\tau$ then at the time $t = (n+1)\tau$ it jumps to a random position $y \in h\mathbb{Z}$, with the transition probability

$$P(X_{(n+1)\tau} = y | X_{n\tau} = x) = k(x-y), \quad x, y \in h\mathbb{Z}. \quad (1.2.1)$$

Here, $k(x)$ is a prescribed non-negative kernel such that

$$\sum_{y \in h\mathbb{Z}} k(y) = 1. \quad (1.2.2)$$

The classical symmetric random walk with a spatial step h and a time step τ corresponds to the choice $k(\pm h) = 1/2$, and $k(y) = 0$ otherwise – the particle may only jump to the nearest neighbor on the left and on the right, with equal probabilities.

In order to connect this process to an evolution equation, let us take a function $f : h\mathbb{Z} \rightarrow \mathbb{R}$, defined on our lattice, and introduce

$$u(t, x) = \mathbb{E}(f(X_t(x))). \quad (1.2.3)$$

Here, $X_t(x)$, $t \in \tau\mathbb{N}$, is the above Markov process starting at a position $X_0(x) = x \in h\mathbb{Z}$ at the time $t = 0$. If $f \geq 0$ then one may think of $u(t, x)$ as the expected value of a “prize” to be collected at the time t at a (random) location of $X_t(x)$ given that the process starts at the point x at the time $t = 0$. An important special case is when f is the characteristic function of a set A . Then, $u(t, x)$ is the probability that the jump process $X_t(x)$ that starts at the position $X_0 = x$ is inside the set A at the time t .

As the process $X_t(x)$ is Markov, the function $u(t, x)$ satisfies the following relation

$$u(t + \tau, x) = \mathbb{E}(f(X_{t+\tau}(x))) = \sum_{y \in h\mathbb{Z}} P(X_\tau = y | X_0 = x) \mathbb{E}(f(X_t(y))) = \sum_{y \in h\mathbb{Z}} k(x - y) u(t, y). \quad (1.2.4)$$

This is because after the initial step when the particle jumps at the time τ from the starting position x to a random position y , the process “starts anew”, and runs for time t between the times τ and $t + \tau$. Equation (1.2.4) can be re-written, using (1.2.2) as

$$u(t + \tau, x) - u(t, x) = \sum_{y \in h\mathbb{Z}} k(x - y) [u(t, y) - u(t, x)]. \quad (1.2.5)$$

The key point of this section is that the discrete equation (1.2.5) leads to various interesting continuum limits as $h \downarrow 0$ and $\tau \downarrow 0$, depending on the choice of the transition kernel $k(y)$, and on the relative size of the spatial mesh size h and the time step τ . In other words, depending on the microscopic model – the particular properties of the random walk – we will end up with different macroscopic continuous models.

The heat equation and random walks

Before showing how a general parabolic equation with non-constant coefficients can be obtained via a limiting procedure from a random walk on a lattice, let us show how this can be done for the heat equation

$$\frac{\partial u}{\partial t} = a \frac{\partial^2 u}{\partial x^2}, \quad (1.2.6)$$

with a constant diffusivity constant $a > 0$. We will assume that the transition probability kernel has the form

$$k(x) = \phi\left(\frac{x}{h}\right), \quad x \in h\mathbb{Z}, \quad (1.2.7)$$

with a non-negative function $\phi(m) \geq 0$ defined on \mathbb{Z} , such that

$$\sum_m \phi(m) = 1. \quad (1.2.8)$$

This form of $k(x)$ allows us to re-write (1.2.5) as

$$u(t + \tau, x) - u(t, x) = \sum_{y \in h\mathbb{Z}} \phi\left(\frac{x - y}{h}\right) [u(t, y) - u(t, x)], \quad (1.2.9)$$

or, equivalently,

$$u(t + \tau, x) - u(t, x) = \sum_{m \in \mathbb{Z}} \phi(m) [u(t, x - mh) - u(t, x)]. \quad (1.2.10)$$

In order to arrive to the heat equation in the limit, we will make the assumption that jumps are symmetric on average:

$$\sum_{m \in \mathbb{Z}} m \phi(m) = 0. \quad (1.2.11)$$

Then, expanding the right side of (1.2.10) in h and the left side in τ , we obtain

$$\tau \frac{\partial u(t, x)}{\partial t} = \frac{ah^2}{2} \frac{\partial^2 u}{\partial x^2}(t, x) + \text{lower order terms}, \quad (1.2.12)$$

with

$$a = \sum_m |m|^2 \phi(m). \quad (1.2.13)$$

To balance the left and the right sides of (1.2.12), we need to take the time step $\tau = h^2$ – note that the scaling $\tau = O(h^2)$ is essentially forced on us if we want to balance the two sides of this equation. Then, in the limit $\tau = h^2 \downarrow 0$, we obtain the heat equation

$$\frac{\partial u(t, x)}{\partial t} = \frac{a}{2} \frac{\partial^2 u(t, x)}{\partial x^2}. \quad (1.2.14)$$

The diffusion coefficient a given by (1.2.13) is the second moment of the jump size – in other words, it measures the “overall jumpiness” of the particles. This is a very simple example of how the microscopic information, the kernel $\phi(m)$, translates into a macroscopic quantity – the overall diffusion coefficient a in the macroscopic equation (1.2.14).

Exercise 1.2.1 Show that if (1.2.11) is violated and

$$b = \sum_{m \in \mathbb{Z}} m \phi(m) \neq 0, \quad (1.2.15)$$

then one needs to take $\tau = h$, and the (formal limit) is the advection equation

$$\frac{\partial u(t, x)}{\partial t} + b \frac{\partial u(t, x)}{\partial x} = 0, \quad (1.2.16)$$

without any diffusion.

Exercise 1.2.2 A reader familiar with the basic probability theory should relate the limit in (1.2.16) to the law of large numbers and explain the relation $\tau = h$ in these terms. How can (1.2.14) and the relation $\tau = h^2$ between the temporal and spatial steps be explained in terms of the central limit theorem?

Parabolic equations with variable coefficients and drifts and random walks

In order to connect a linear parabolic equation with inhomogeneous coefficients, such as (1.1.2) with the right side $f \equiv 0$:

$$\frac{\partial u}{\partial t} - a(x)\frac{\partial^2 u}{\partial x^2} + b(x)\frac{\partial u}{\partial x} = 0, \quad (1.2.17)$$

to a continuum limit of random walks, we consider a slight modification of the microscopic dynamics that led to the heat equation in the macroscopic limit. We go back to (1.2.4):

$$u(t + \tau, x) = \mathbb{E}(f(X_{t+\tau}(x))) = \sum_{y \in h\mathbb{Z}} P(X_\tau = y | X_0 = x) \mathbb{E}(f(X_t(y))) = \sum_{y \in h\mathbb{Z}} k(x, y) u(t, y). \quad (1.2.18)$$

Here, $k(x, y)$ is the probability to jump to the position y from a position x . Note that we no longer assume that the law of the jump process is spatially homogeneous: the transition probabilities depend not only on the difference $x - y$ but both on x and y . However, we will assume that $k(x, y)$ is "locally homogeneous". This condition translates into considering

$$k(x, y; h) = \phi\left(x, \frac{x - y}{h}; h\right). \quad (1.2.19)$$

The "slow" spatial dependence of the transition probability density is encoded in the dependence of the function $\phi(x, z, h)$ on the "macroscopic" variable x , while its "fast" spatial variations are described by the dependence of $\phi(x, z, h)$ on the variable z .

Exercise 1.2.3 Make sure you can interpret this point. Think of "freezing" the variable x and only varying the z -variable.

We will soon see why we introduce the additional dependence of the transition density on the mesh size h – this will lead to a non-trivial first order term in the parabolic equation we will obtain in the limit. We assume that the function $\phi(x, m; h)$, with $x \in \mathbb{R}$, $m \in \mathbb{Z}$ and $h \in (0, 1)$, satisfies

$$\sum_{m \in \mathbb{Z}} \phi(x, m; h) = 1 \text{ for all } x \in \mathbb{R} \text{ and } h \in (0, 1), \quad (1.2.20)$$

which leads to the analog of the normalization (1.2.2):

$$\sum_{y \in h\mathbb{Z}} k(x, y) = 1 \text{ for all } x \in h\mathbb{Z}. \quad (1.2.21)$$

This allows us to re-write (1.2.18) in the familiar form

$$u(t + \tau, x) - u(t, x) = \sum_{y \in h\mathbb{Z}} \phi\left(x, \frac{x - y}{h}; h\right) [u(t, y) - u(t, x)], \quad (1.2.22)$$

or, equivalently,

$$u(t + \tau, x) - u(t, x) = \sum_{m \in \mathbb{Z}} \phi(x, m; h) [u(t, x - mh) - u(t, x)], \quad (1.2.23)$$

We will make the assumption that the average asymmetry of the jumps is of the size h . In other words, we suppose that

$$\sum_{m \in \mathbb{Z}} m \phi(x, m; h) = b(x)h + O(h^2), \quad (1.2.24)$$

that is,

$$\sum_{m \in \mathbb{Z}} m \phi(x, m; 0) = 0 \text{ for all } x \in \mathbb{R},$$

and

$$b(x) = \sum_{m \in \mathbb{Z}} m \frac{\partial \phi(x, m; h=0)}{\partial h} \quad (1.2.25)$$

is a given smooth function. The last assumption we will make is that the time step is $\tau = h^2$, as before. Expanding the left and the right side of (1.2.23) in h now leads to the parabolic equation

$$\frac{\partial u}{\partial t} = -b(x) \frac{\partial u(t, x)}{\partial x} + a(x) \frac{\partial^2 u(t, x)}{\partial x^2}, \quad (1.2.26)$$

with

$$a(x) = \frac{1}{2} \sum_{m \in \mathbb{Z}} |m|^2 \phi(x, m; h=0). \quad (1.2.27)$$

This is a parabolic equation of the form (1.1.2) in one dimension. We automatically satisfy the condition $a(x) > 0$ (known as the ellipticity condition) unless $\phi(x, m; h=0) = 0$ for all $m \in \mathbb{Z} \setminus \{0\}$. That is, $a(x) = 0$ only at the positions where the particles are completely stuck and can not jump at all. Note that the asymmetry in (1.2.24), that is, the mismatch in the typical jump sizes to the left and right, leads to the first order term in the limit equation (1.2.26) – because of that the first-order coefficient $b(x)$ is known as the drift, while the second-order coefficient $a(x)$ (known as the diffusivity) measures "the overall jumpiness" of the particles, as seen from (1.2.27).

Exercise 1.2.4 Relate the above considerations to the method of characteristics for the first order linear equation

$$\frac{\partial u}{\partial t} + b(x) \frac{\partial u}{\partial x} = 0.$$

How does it arise from similar considerations?

Exercise 1.2.5 It is straightforward to generalize this construction to higher dimensions leading to general parabolic equations of the form (1.1.2). Verify that the diffusion matrices $a_{ij}(x)$ in (1.1.2) that arise in this fashion, will always be nonnegative, in the sense that for any $\xi \in \mathbb{R}^n$ and all x , we have (once again, as the repeated indices are summed over):

$$a_{ij}(x) \xi_i \xi_j \geq 0. \quad (1.2.28)$$

This is very close to the lower bound in the ellipticity condition on the matrix $a_{ij}(x)$ which says that there exists a constant $c > 0$ so that for any $\xi \in \mathbb{R}^n$ and $x \in \mathbb{R}^n$ we have

$$c|\xi|^2 \leq a_{ij}(x) \xi_i \xi_j \leq c^{-1}|\xi|^2. \quad (1.2.29)$$

We see that the ellipticity condition appears very naturally in the probabilistic setting.

Summarizing, we see that parabolic equations of the form (1.1.2) arise as limits of random walks that make jumps of the size $O(h)$, with a time step $\tau = O(h^2)$. Thus, the overall number of jumps by a time $t = O(1)$ is very large, and each individual jump is very small. The drift vector $b_j(x)$ appears from the local non-zero mean of the jump direction and size, and the diffusivity matrix $a_{ij}(x)$ measures the typical jump size. In addition, the diffusivity matrix is nonnegative-definite: condition (1.2.28) is satisfied.

Parabolic equations and branching random walks

Let us now explain how random walks can lead to parabolic equations with a zero-order term:

$$\frac{\partial u}{\partial t} - a_{ij}(x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(x) \frac{\partial u}{\partial x_j} + c(x)u = 0. \quad (1.2.30)$$

This will help us understand qualitatively the role of the coefficient $c(x)$. Once again, we will consider the one-dimensional case for simplicity, and will only give the details for the case

$$\frac{\partial u}{\partial t} - \frac{\partial^2 u}{\partial x^2} + c(x)u = 0, \quad (1.2.31)$$

as the non-constant diffusion matrix $a_{ij}(x)$ and drift $b_j(x)$ can be treated exactly as in the case $c = 0$.

In order to incorporate the zero order term we need to allow the particles not only jump but also branch – this is the reason why the zero-order term will appear in (1.2.30). As before, our particles make jumps on the lattice $h\mathbb{Z}$, at the discrete times $t \in \tau\mathbb{N}$. We start at $t = 0$ with one particle at a position $x \in h\mathbb{Z}$. Let us assume that at the time $t = n\tau$ we have a collection of N_t particles $X_1(t, x), \dots, X_{N_t}(t, x)$ (the number N_t is random, as will immediately see). At the time t , each particle $X_m(t, x)$ behaves independently from the other particles. With the probability

$$p_0 = 1 - |c(X_m(t))|\tau,$$

it simply jumps to a new location $y \in h\mathbb{Z}$, chosen with the transition probability $k(X_m(t) - y)$, as in the process with no branching. If the particle at $X_m(t, x)$ does not jump – this happens with the probability $p_1 = 1 - p_0$, there are two possibilities. If $c(X_m(t)) < 0$, then it is replaced by two particles at the same location $X_m(t, x)$ that remain at this position until the time $t + \tau$. If $c(X_m(t)) > 0$ and the particle does not jump, then it is removed. This process is repeated independently for all particles $X_1(t, x), \dots, X_{N_t}(t, x)$, giving a new collection of particles at the locations $X_1(t + \tau, x), \dots, X_{N_{t+\tau}}(t + \tau, x)$ at the time $t + \tau$. If $c(x) > 0$ at some positions, then the process can terminate when there are no particles left. If $c(x) \leq 0$ everywhere, then the process continues forever.

To connect this particle system to an evolution equation, given a function f , we define, for $t \in \tau\mathbb{N}$, and $x \in h\mathbb{Z}$,

$$u(t, x) = \mathbb{E}[f(X_1(t, x)) + f(X_2(t, x)) + \dots + f(X_{N_t}(t, x))].$$

The convention is that $f = 0$ inside the expectation if there are no particles left. This is similar to what we have done for particles with no branching. If f is the characteristic function of a set A , then $u(t, x)$ is the expected number of particles inside A at the time $t > 0$.

In order to get an evolution equation for $u(t, x)$, we look at the initial time when we have just one particle at the position x : if $c(x) \leq 0$, then this particle either jumps or branches, leading to the balance

$$u(t + \tau, x) = (1 + c(x)\tau) \sum_{y \in h\mathbb{Z}} k(x - y)u(t, y) - 2c(x)\tau u(t, x), \quad \text{if } c(x) \leq 0, \quad (1.2.32)$$

which is the analog of (1.2.4). If $c(x) > 0$ the particle either jumps or is removed, leading to

$$u(t + \tau, x) = (1 - |c(x)|\tau) \sum_{y \in h\mathbb{Z}} k(x - y)u(t, y). \quad (1.2.33)$$

In both cases, we can re-write the balances similarly to (1.2.5):

$$u(t + \tau, x) - u(t, x) = (1 - |c(x)|\tau) \sum_{y \in h\mathbb{Z}} k(x - y)(u(t, y) - u(t, x)) - c(x)\tau u(t, x). \quad (1.2.34)$$

We may now take the transition probability kernel of the familiar form

$$k(x) = \phi\left(\frac{x}{h}\right),$$

with a function $\phi(m)$ as in (1.2.7)-(1.2.8). Taking $\tau = h^2$ leads, as in (1.2.12), to the diffusion equation but now with a zero-order term:

$$\frac{\partial u}{\partial t} = \frac{a}{2} \frac{\partial^2 u}{\partial x^2} - c(x)u. \quad (1.2.35)$$

Thus, the zero-order coefficient $c(x)$ can be interpreted as the branching (or killing, depending on the sign of $c(x)$) rate of the random walk. The parabolic maximum principle for $c(x) \geq 0$ that we will discuss in the next section simply means, on this informal level, that if the particles never branch, and can only be removed, their expected number can not grow in time.

Exercise 1.2.6 Add branching to the random walk we have discussed in Section 1.2 of this chapter, and obtain a more general parabolic equation, in higher dimensions:

$$\frac{\partial u}{\partial t} - a_{ij}(x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(x) \frac{\partial u}{\partial x_j} + c(x)u = 0. \quad (1.2.36)$$

1.3 The maximum principle interlude: the basic statements

As the parabolic maximum principle underlies most of the parabolic existence and regularity theory, we first recall some basics on the maximum principle for parabolic equations. They are very similar in spirit to what we have described in the previous chapter for the Laplace and Poisson equations. This material can, once again, be found in many standard textbooks, such as [?], so we will not present most of the proofs but just recall the statements we will need.

We consider a (more general than the Laplacian) elliptic operator of the form

$$Lu(x) = -a_{ij}(t, x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(t, x) \frac{\partial u}{\partial x_j}, \quad (1.3.1)$$

in a bounded domain $x \in \Omega \subset \mathbb{R}^n$ and for $0 \leq t \leq T$. Note that the zero-order coefficient is set to be zero for the moment. The ellipticity of L means that the matrix $a_{ij}(t, x)$ is uniformly positive-definite and bounded. That is, there exist two positive constants $\lambda > 0$ and $\Lambda > 0$ so that, for any $\xi \in \mathbb{R}^n$, and $0 \leq t \leq T$, and any $x \in \Omega$, we have

$$\lambda |\xi|^2 \leq a_{ij}(t, x) \xi_i \xi_j \leq \Lambda |\xi|^2. \quad (1.3.2)$$

We also assume that all coefficients $a_{ij}(t, x)$ and $b_j(t, x)$ are continuous and uniformly bounded. Given a time $T > 0$, define the parabolic cylinder $\Omega_T = [0, T] \times \Omega$ and its parabolic boundary as

$$\Gamma_T = \{x \in \Omega, 0 \leq t \leq T : \text{either } x \in \partial\Omega \text{ or } t = 0\}.$$

In other words, Γ_T is the part of the boundary of Ω_T without “the top” $\{(t, x) : t = T, x \in \Omega\}$.

Theorem 1.3.1 (*The weak maximum principle*) *Let a function $u(t, x)$ satisfy*

$$\frac{\partial u}{\partial t} + Lu \leq 0, \quad x \in \Omega, \quad 0 \leq t \leq T, \quad (1.3.3)$$

and assume that Ω is a smooth bounded domain. Then $u(t, x)$ attains its maximum over Ω_T on the parabolic boundary Γ_T , that is,

$$\sup_{\Omega_T} u(t, x) = \sup_{\Gamma_T} u(t, x). \quad (1.3.4)$$

As in the elliptic case, we also have the strong maximum principle.

Theorem 1.3.2 (*The strong maximum principle*) *Let a smooth function $u(t, x)$ satisfy*

$$\frac{\partial u}{\partial t} + Lu \leq 0, \quad x \in \Omega, \quad 0 \leq t \leq T, \quad (1.3.5)$$

in a smooth bounded domain Ω . Then if $u(t, x)$ attains its maximum over $\bar{\Omega}_T$ at an interior point $(t_0, x_0) \notin \Gamma_T$ then $u(t, x)$ is equal to a constant in Ω_T .

We will not prove these results here, the reader may consult [?] or other standard textbooks on PDEs for a proof. One standard generalization of the maximum principle is to include the lower order term with a sign, as in the elliptic case – compare to Theorem ?? in Chapter ?. Namely, it is quite straightforward to show that if $c(x) \geq 0$ then the maximum principle still holds for parabolic equations (1.3.5) with an operator L of the form

$$Lu(x) = -a_{ij}(t, x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u. \quad (1.3.6)$$

The proof can, once again, be found in [?]. However, as we have seen in the elliptic case, in the maximum principles for narrow domains (Theorem ?? in Chapter ??) and domains of a small volume (Theorem ?? in the same chapter), the sign condition on the coefficient $c(t, x)$ is not necessary for the maximum principle to hold. Later in this chapter, we will discuss a more general condition that quantifies the necessary assumptions on the operator L for the maximum principle to hold in a unified way.

A consequence of the maximum principle is the comparison principle, a result that holds also for operators with zero order coefficients and in unbounded domains. In general, the comparison principle in unbounded domains holds under a proper restriction on the growth of the solutions at infinity. Here, for simplicity we assume that the solutions are uniformly bounded.

Theorem 1.3.3 *Let the smooth uniformly bounded functions $u(t, x)$ and $v(t, x)$ satisfy*

$$\frac{\partial u}{\partial t} + Lu + c(t, x)u \geq 0, \quad 0 \leq t \leq T, \quad x \in \Omega \quad (1.3.7)$$

and

$$\frac{\partial v}{\partial t} + Lv + c(t, x)v \leq 0, \quad 0 \leq t \leq T, \quad x \in \Omega, \quad (1.3.8)$$

in a smooth (and possibly unbounded) domain Ω , with a bounded function $c(t, x)$. Assume that $u(0, x) \geq v(0, x)$ and

$$u(t, x) \geq v(t, x) \text{ for all } 0 \leq t \leq T \text{ and } x \in \partial\Omega.$$

Then, we have

$$u(t, x) \geq v(t, x) \text{ for all } 0 \leq t \leq T \text{ and all } x \in \Omega.$$

Moreover, if in addition, $u(0, x) > v(0, x)$ on an open subset of Ω then $u(t, x) > v(t, x)$ for all $0 < t < T$ and all $x \in \Omega$.

The assumption that both $u(t, x)$ and $v(t, x)$ are uniformly bounded is important if the domain Ω is unbounded – without this condition even the Cauchy problem for the standard heat equation in \mathbb{R}^n may have more than one solution, while the comparison principle implies uniqueness trivially. An example of non-uniqueness is discussed in detail in [?] – such solutions grow very fast as $|x| \rightarrow +\infty$ for any $t > 0$, while satisfying the initial condition $u(0, x) \equiv 0$. The extra assumption that $u(t, x)$ is bounded allows to rule out this non-uniqueness issue. Note that the special case $\Omega = \mathbb{R}^n$ is included in Theorem 1.3.3, and in that case only the comparison at the initial time $t = 0$ is needed for the conclusion to hold for bounded solutions. Once again, a reader not interested in treating the proof as an exercise should consult [?], or another of his favorite basic PDE textbooks. We should stress that in the rest of this book we will only consider solutions for which the uniqueness holds.

A standard corollary of the parabolic maximum principle is the following estimate.

Exercise 1.3.4 Let Ω be a (possibly unbounded) smooth domain, and $u(t, x)$ be the solution of the initial boundary value problem

$$\begin{aligned} u_t + Lu + c(t, x)u &= 0, \quad \text{in } \Omega, \\ u(t, x) &= 0 \text{ for } x \in \partial\Omega, \\ u(0, x) &= u_0(x). \end{aligned} \quad (1.3.9)$$

Assume (to ensure the uniqueness of the solution) that u is locally in time bounded: for all $T > 0$ there exists $C_T > 0$ such that $|u(t, x)| \leq C_T$ for all $t \in [0, T]$ and $x \in \Omega$. Assume that the function $c(t, x)$ is bounded, with $c(t, x) \geq -M$ for all $x \in \Omega$, and show that then $u(t, x)$ satisfies

$$|u(t, x)| \leq \|u_0\|_{L^\infty} e^{Mt}, \quad \text{for all } t > 0 \text{ and } x \in \Omega. \quad (1.3.10)$$

The estimate (1.3.10) on the possible growth (or decay) of the solution of (1.3.9) is by no means optimal, and we will soon see how it can be improved.

We also have the parabolic Hopf Lemma, of which we will only need the following version.

Lemma 1.3.5 (*The parabolic Hopf Lemma*) *Let $u(t, x) \geq 0$ be a solution of*

$$u_t + Lu + c(t, x)u = 0, \quad 0 \leq t \leq T,$$

in a ball $B(z, R)$. Assume that there exists $t_0 > 0$ and $x_0 \in \partial B(z, R)$ such that $u(t_0, x_0) = 0$, then we have

$$\frac{\partial u(t_0, x_0)}{\partial \nu} < 0. \quad (1.3.11)$$

The proof is very similar to that of the elliptic Hopf Lemma, and can be found, for instance, in [?].

1.4 The forced linear heat equation

The regularity theory for the parabolic equations is an extremely rich and fascinating subject that is often misunderstood as "technical". To keep things relatively simple, we are not going to delve into it head first. Rather, we focus in this section on the regularity results for the forced linear heat equation in the whole space:

$$u_t - \Delta u = g(t, x), \quad t > 0, \quad x \in \mathbb{R}^n, \quad (1.4.1)$$

with an initial condition $u(0, x) = u_0(x)$. As we will see almost immediately, in Proposition 1.4.3, the contribution of the initial condition can be treated in a very simple way, and the main question is what can we say about the regularity of $u(t, x)$ in terms of the prescribed regularity of $g(t, x)$. In Section 1.5, the answers to these seemingly technical and "boring" issues will allow us to address the question of existence and regularity of solutions to "much more interesting" nonlinear equations, in a very large class. The completely explicit results for the heat equation we describe in this section also explain quite well how one can approach general inhomogeneous parabolic equations – we explain this at the qualitative level in Section 1.6.

This section is both longer and more technical than what the reader has encountered so far in the book. This techniques are mostly elementary but still require us to get our hands dirty and the computations reveal some of the very important cancellations that underline the regularity theory in the general case discussed in Section 1.6. We proceed in several steps. First we show that if $g(t, x)$ is bounded, without any other assumptions on its regularity,

then the function $u(t, x)$ is Hölder continuous both in t and in x , and the corresponding Hölder norms of u are bounded by the L^∞ -norm of g . This is done in Section 1.4.2, and the main result there is Proposition 1.4.7. Next, in Section 1.4.3 we assume that $g(t, x)$ is Hölder continuous and show that then $u(t, x)$ is once differentiable in t and twice in x , with the corresponding bounds on the derivatives in terms of the Hölder norm of g . This is made precise in Proposition 1.4.13, and a generalization to higher order derivatives is explained in Proposition 1.4.20. Finally, in Section 1.4.4 we show that if $g(t, x)$ is Hölder continuous then the first derivative in time and second derivatives in space are not just bounded but actually themselves Hölder continuous – this is stated in Proposition 1.4.18. Of course, in the first place, the reader may wonder what we mean by a solution to the heat equation that is not necessarily differentiable. This is explained in Section 1.4.1 in terms of the Duhamel formula.

The proofs of all these results are painfully computational but they open the gates to beautiful results in the theory of nonlinear diffusion equations, so the payoff for the hard work in this section is quite high. The hope is that the reader will emerge at the end of this section with the understanding that the parabolic regularity theory does require some calculations but is by no means mysterious or inaccessible. As we will see later, the results it provides are not light but worth their weight in gold.

Recommendation. This section contains many exercises that are computational in nature and may at the first look appear somewhat unappealing to the reader. We strongly encourage you to do them as they show the machinery and details behind the beautiful theory.

1.4.1 The Duhamel formula

We consider the forced linear heat equation

$$u_t = \Delta u + g(t, x), \tag{1.4.2}$$

posed in the whole space $x \in \mathbb{R}^n$, and with an initial condition

$$u(0, x) = u_0(x). \tag{1.4.3}$$

The basic question for us in this section is how regular the solution of (1.4.2)-(1.4.3) is, in terms of the regularity of the initial condition $u_0(x)$ and the forcing term $g(t, x)$. The function $u(t, x)$ is given explicitly by the Duhamel formula

$$u(t, x) = v(t, x) + \int_0^t w(t, x; s) ds. \tag{1.4.4}$$

Here, $v(t, x)$ is the solution to the homogeneous heat equation

$$v_t = \Delta v, \quad x \in \mathbb{R}^n, \quad t > 0, \tag{1.4.5}$$

with the initial condition $v(0, x) = u_0(x)$, and $w(t, x; s)$ is the solution of the Cauchy problem

$$w_t(t, x; s) = \Delta w(t, x; s), \quad x \in \mathbb{R}^n, \quad t > s, \tag{1.4.6}$$

that runs starting at the time s , and is supplemented by the initial condition at $t = s$:

$$w(t = s, x; s) = g(s, x). \tag{1.4.7}$$

Exercise 1.4.1 If the reader has not previously encountered the Duhamel formula, you should consider it in a more general setting of a forced problem on a Banach space X :

$$\frac{du}{dt} = Lu + g, \quad t > 0, \quad (1.4.8)$$

with an initial condition $u(0) = u_0 \in X$, for some general linear operator $L : X \rightarrow X$ and forcing $g \in C([0, T]; X)$. The basic assumption on the operator L is that for any $v_0 \in X$ the initial value problem

$$\frac{dv}{dt} = Lv, \quad t > 0, \quad (1.4.9)$$

with the initial condition $v(0) = v_0 \in X$ has a unique bounded solution $u(t) \in X$ for all $t \geq 0$. Show that for all $0 \leq t \leq T$ the function $u(t)$ can be written as

$$u(t) = v(t) + \int_0^t w(t; s) ds.$$

Here, $v(t)$ is the solution to the initial value problem (1.4.9) with $v(0) = u_0$, and $w(t; s)$ solves the initial value problem starting at a time $s < t$:

$$\frac{dw}{dt} = Lw, \quad t > s, \quad (1.4.10)$$

with the initial condition $w(s) = g(s)$.

Let us denote the solution of the Cauchy problem (1.4.5) as

$$v(t, x) = e^{t\Delta} u_0. \quad (1.4.11)$$

This defines the operator $e^{t\Delta}$. It maps the initial condition of the heat equation to its solution at the time t , and is given explicitly as

$$e^{t\Delta} f(x) = \frac{1}{(4\pi t)^{n/2}} \int_{\mathbb{R}^n} e^{-(x-y)^2/(4t)} f(y) dy. \quad (1.4.12)$$

With this notation, another way to write the Duhamel formula (1.4.2) is

$$u(t, x) = e^{t\Delta} u_0(x) + \int_0^t e^{(t-s)\Delta} g(s, x) ds, \quad (1.4.13)$$

or, more explicitly:

$$u(t, x) = \frac{1}{(4\pi t)^{n/2}} \int e^{-(x-y)^2/(4t)} u_0(y) dy + \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} g(s, y) dy ds. \quad (1.4.14)$$

Of course, we can make these expressions much shorter and more elegant if we introduce the heat kernel

$$G(t, x) = \frac{1}{(4\pi t)^{n/2}} e^{-|x|^2/(4t)}, \quad (1.4.15)$$

and rewrite them in terms of convolutions with $G(t, x)$. We keep the formulas for $u(t, x)$ as explicit as feasible on purpose, to keep the potential singularities as visible as possible, so that the reader would be alert of the potential dangers in the estimates.

Here is an exercise on the Duhamel formula for a different partial differential equation.

Exercise 1.4.2 Consider the one-dimensional wave equation

$$u_{tt} - u_{xx} = g(t, x), \quad t > 0, \quad x \in \mathbb{R}, \quad (1.4.16)$$

with zero initial condition $u(0, x) = u_t(0, x) = 0$. Show that its solution is given by the Duhamel formula

$$J_{wave}(t, x) = \frac{1}{2} \int_0^t \int_{x-(t-s)}^{x+(t-s)} g(s, y) dy ds. \quad (1.4.17)$$

The first term in (1.4.14) is rather benign as far as regularity is concerned. We use the notation

$$|k| = k_1 + \dots + k_n,$$

for a multi-index $k = (k_1, \dots, k_n)$, and

$$D_x^k u = \frac{\partial^{|k|} u}{\partial x_1^{k_1} \dots \partial x_n^{k_n}}.$$

Proposition 1.4.3 Let u_0 be a bounded and continuous function and set

$$v(t, x) = e^{t\Delta} u_0 = \frac{1}{(4\pi t)^{n/2}} \int e^{-(x-y)^2/(4t)} u_0(y) dy. \quad (1.4.18)$$

Show that for any $t > 0$ and for any multi-index k with $|k| = m$ there exists $C_m > 0$ that depends only on m so that

$$|D_x^k v(t, x)| \leq \frac{C_m}{t^{m/2}} \|u_0\|_{L^\infty}, \quad |\partial_t^m v(t, x)| \leq \frac{C_m}{t^m} \|u_0\|_{L^\infty}, \quad \text{for all } t > 0 \text{ and } x \in \mathbb{R}^n. \quad (1.4.19)$$

The reader should note the following simple observations. First, the estimates on the derivatives of $v(t, x)$ in (1.4.19) blow-up as $t \downarrow 0$. This is expected – we only assume that $u_0(x)$ is continuous. More importantly, the estimates on the derivatives at a positive time $t > 0$ depend only on the L^∞ -norm of u_0 – this is the instant regularization effect of the heat equation.

Exercise 1.4.4 Prove Proposition 1.4.3. The proof involves nothing but calculus and remembering when an integral with an integrand that depends on a parameter can be differentiated in this parameter.

For the rest of this section we focus on the second term in (1.4.14),

$$J(t, x) = \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} g(s, y) dy ds. \quad (1.4.20)$$

Note that the function $J(t, x)$, if it is sufficiently regular, is the solution to the forced linear heat equation

$$J_t = \Delta J + g(t, x), \quad (1.4.21)$$

posed in the whole space $x \in \mathbb{R}^n$, with the initial condition $J(0, x) \equiv 0$. However, as a priori we do not know that $J(t, x)$ is differentiable, for now, we can not be quite sure that (1.4.21) makes classical sense as stated. It is potentially problematic because of the term $(t-s)^{-n/2}$ in (1.4.20) that blows up as $s \uparrow t$. In particular, a naive attempt to differentiate the integrand in t or x would lead to expressions that are too singular to be absolutely integrable without some cancellations.

Exercise 1.4.5 Differentiate the integrand in $J(t, x)$ in t blindly, observe the singularity as $s \rightarrow t$ and get stuck.

However, to see that the singularity is not as dangerous as it may naively seem, observe that a simple change of variables shows that if $g(t, x)$ is bounded then so is $J(t, x)$:

$$|J(t, x)| \leq \|g\|_{L^\infty} \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} dy ds = \frac{t\|g\|_{L^\infty}}{\pi^{n/2}} \int e^{-z^2} dz = t\|g\|_{L^\infty}. \quad (1.4.22)$$

We used above the simple change of variables

$$z = \frac{x-y}{2\sqrt{t-s}}, \quad (1.4.23)$$

that will be a recurring refrain throughout this section and beyond. In other words, the integral (1.4.20) defines an L^∞ function $J(t, x)$ if $g(t, x)$ itself is an L^∞ function. The reader should informally think of $J(t, x)$ as the solution to (1.4.21) even if it does not have sufficient regularity to be a classical solution. In the remainder of this section we focus exactly to the question of regularity of $J(t, x)$.

Exercise 1.4.6 Deduce the upper bound (1.4.22) for $J(t, x)$ directly from the parabolic maximum principle, without any explicit computations.

1.4.2 Regularity gain: from a bounded $g(t, x)$ to a Hölder $J(t, x)$

The estimate (1.4.22) can be restated as an $L^\infty - L^\infty$ bound:

$$\|J\|_{L^\infty} \leq t\|g\|_{L^\infty}. \quad (1.4.24)$$

Such bounds are useful but they do not give any better regularity for the function $J(t, x)$ than for $g(t, x)$: it says that if g is bounded then so is J . On the other hand, the following proposition gives a quantifiable way to say that if $g(t, x)$ is bounded, and without any assumptions on the continuity of g , then the function J is Hölder continuous in t and differentiable in x . Hence, it is more regular than the assumed regularity of g . This is a very simple example of the general phenomenon of parabolic regularity: solution is better than the input data, such as the initial condition or forcing.

Recall the notion of the Hölder norm of a function $g(t, x)$ defined for $(t, x) \in [0, T] \times \mathbb{R}^n$:

$$\|g\|_{C_t^\alpha C_x^\beta} = \|g\|_{L^\infty} + \sup \frac{|g(t, x) - g(t', x')|}{|t - t'|^\alpha + |x - x'|^\beta}, \quad (1.4.25)$$

with the supremum taken over all $0 \leq t, t' \leq T$ and $x, x' \in \mathbb{R}^n$ such that $(t, x) \neq (t', x')$. We will use the notation $C_t^\alpha C_x^\beta([0, T] \times \mathbb{R}^n)$, or $C_t^\alpha C_x^\beta$ for short, for the space of Hölder continuous functions on $[0, T] \times \mathbb{R}^n$ with a finite Hölder norm. We apologize to the reader for the use of this cumbersome notation but it allows us to distinguish between the regularity in time and space and avoid various other notational pitfalls.

Proposition 1.4.7 Let $g(t, x)$ be a measurable bounded function, so that $g \in L^\infty([0, T] \times \mathbb{R}^n)$, and $J(t, x)$ be given by (1.4.20). Then,

(i) the function $J(t, x)$ is once differentiable in x for all $t > 0$ and $x \in \mathbb{R}^n$,

(ii) for any $\alpha \in (0, 1)$, the function $J(t, x)$ is C^α -Hölder continuous in t for all $t > 0$ and all $x \in \mathbb{R}^n$, and

(iii) for all $1 \leq k \leq n$, the derivatives $\partial_{x_k} J(t, x)$ are C^α -Hölder continuous in x , for all $t > 0$ and $x \in \mathbb{R}^n$, for all $\alpha \in (0, 1)$.

Moreover, there exist $C > 0$ and $C_\alpha > 0$ so that for all $t, t' \in [0, T]$ and $x, x' \in \mathbb{R}^n$ we have

$$\begin{aligned} |\partial_{x_k} J(t, x)| &\leq C\sqrt{t}\|g\|_{L^\infty([0, T] \times \mathbb{R}^n)}, \\ |\partial_{x_k} J(t, x) - \partial_{x_k} J(t, x')| &\leq C_\alpha\|g\|_{L^\infty([0, T] \times \mathbb{R}^n)}|x - x'|^\alpha, \\ |J(t, x) - J(t', x)| &\leq C_\alpha\|g\|_{L^\infty([0, T] \times \mathbb{R}^n)}|t - t'|^\alpha. \end{aligned} \tag{1.4.26}$$

The difference in regularity of $J(t, x)$ in t and x is not an artifact of the proof. It is easy to see that $J(t, x)$ need not be differentiable in t if all we know about $g(t, x)$ is that it is bounded. Indeed, the reader can simply think of $g(t, x) = \text{sgn}(t - 1)$, in which case $J(t, x) = |t - 1| - 1$ and is not differentiable at $t = 1$ but is Hölder continuous for all $t \geq 0$. The next exercise shows that neither can one expect the function $J(t, x)$ to be twice continuously differentiable in x under the assumption that $g(t, x)$ is bounded and not necessarily continuous. Hence, the claimed regularity of $J(t, x)$ in Proposition 1.4.7 is "reasonably optimal".

Exercise 1.4.8 Give an example of a bounded function $g(t, x)$, $t \geq 0$, $x \in \mathbb{R}$, such that $J(t, x)$ is not twice continuously differentiable in x even though the derivative $\partial_x J(t, x)$ is α -Hölder continuous in x for any $\alpha \in (0, 1)$.

The next exercise asks you to compare the gain of regularity for the heat equation and for the wave equation.

Exercise 1.4.9 Does the result of Proposition 1.4.7 apply to the solution to the wave equation given in Exercise 1.4.2?

Proof of Proposition 1.4.7. Let us freeze $t > 0$, fix some $1 \leq i \leq n$, and prove that $J(t, x)$ is differentiable in x_i . The first inclination may be to simply differentiate the integrand in (1.4.20), as suggested, albeit with a warning, in Exercise 1.4.4. This can not be done in the t -variable, simply because we have seen that $J(t, x)$ need not be differentiable in t . Such differentiation in the x -variable can be justified, but it is also instructive to work from scratch with the finite differences, as we will need to do that with the time increments anyway. Let e_i be the unit vector in the x_i -direction and write

$$\frac{J(t, x + he_i) - J(t, x)}{h} = \frac{1}{h} \int_0^t \int_{\mathbb{R}^n} \left[e^{-(x+he_i-y)^2/4(t-s)} - e^{-(x-y)^2/4(t-s)} \right] g(s, y) \frac{dy ds}{(4\pi(t-s))^{n/2}}. \tag{1.4.27}$$

The familiar change of variables (1.4.23)

$$z = \frac{x - y}{2\sqrt{t - s}}, \tag{1.4.28}$$

leads to

$$\begin{aligned} \frac{J(t, x + he_i) - J(t, x)}{h} &= \frac{1}{h} \int_0^t \int_{\mathbb{R}^n} \left[e^{-(z+he_i/(2\sqrt{t-s}))^2} - e^{-z^2} \right] g(s, x - 2z\sqrt{t-s}) \frac{dz ds}{(4\pi)^{n/2}} \\ &= \int_0^t Q_{h,i}(t, s) \frac{ds}{\sqrt{t-s}}. \end{aligned} \quad (1.4.29)$$

with

$$Q_{h,i}(t, s) = \frac{\sqrt{t-s}}{h} \int_{\mathbb{R}^n} \left[e^{-(z+he_i/(2\sqrt{t-s}))^2} - e^{-z^2} \right] g(s, x - 2z\sqrt{t-s}) \frac{dz}{(4\pi)^{n/2}}. \quad (1.4.30)$$

Exercise 1.4.10 Show that if $g \in L^\infty([0, T] \times \mathbb{R}^n)$, then for almost every $0 < s \leq t$ fixed we have

$$\lim_{h \rightarrow 0} Q_{h,i}(t, s) = \bar{Q}_i(t, s) := - \int_{\mathbb{R}^n} z_i e^{-z^2} g(s, x - 2z\sqrt{t-s}) \frac{dz}{(4\pi)^{n/2}}, \quad (1.4.31)$$

and that there exists $C > 0$ so that for almost every $0 < s \leq t$ and all $h \in (0, 1)$ we have

$$|Q_{h,i}(t, s)| \leq C \|g\|_{L^\infty}. \quad (1.4.32)$$

The result of Exercise 1.4.10 allows us to use the Lebesgue dominated convergence theorem and pass to the limit $h \rightarrow 0$ in (1.4.29) and conclude that

$$\frac{\partial J(t, x)}{\partial x_i} = - \int_0^t \int_{\mathbb{R}^n} z_i e^{-z^2} g(s, x - 2z\sqrt{t-s}) \frac{dz}{(4\pi)^{n/2}} \frac{ds}{\sqrt{t-s}}. \quad (1.4.33)$$

This shows both that $J(t, x)$ is differentiable in x and that the first bound in (1.4.26) holds:

$$\left| \frac{\partial J(t, x)}{\partial x_i} \right| \leq C \sqrt{t} \|g\|_{L^\infty}. \quad (1.4.34)$$

The above argument can not be repeated for the time derivative: if we differentiate the integrand in time, and make the same change of variable to z as in (1.4.28), we would get

$$\frac{\partial}{\partial t} \left(\frac{e^{-|x-y|^2/4(t-s)}}{(4\pi(t-s))^{n/2}} \right) = \frac{|x-y|^2}{4(t-s)^2} \frac{e^{-|x-y|^2/4(t-s)}}{(4\pi(t-s))^{n/2}} = \frac{1}{4(t-s)} \frac{|z^2| e^{-|z|^2}}{(4\pi(t-s))^{n/2}}, \quad z = \frac{x-y}{\sqrt{t-s}}.$$

Thus, the change of variables to z would bring a non-integrable $(t-s)^{-1}$ singularity instead of a $(t-s)^{-1/2}$ term that appears in (1.4.33).

Exercise 1.4.11 Verify that differentiating the integrand twice in x leads to the same kind of (seemingly non-integrable) singularity in $(t-s)$ as differentiating once in t .

Nevertheless, the Hölder continuity of $J(t, x)$ in time is proved by a very similar, except slightly longer, argument. We again compute a partial difference. Assume, for convenience, that $t' \geq t$, and write

$$\begin{aligned} J(t, x) - J(t', x) &= \int_0^t \int_{\mathbb{R}^n} \left(\frac{e^{-|x-y|^2/4(t-s)}}{(4\pi(t-s))^{n/2}} - \frac{e^{-|x-y|^2/4(t'-s)}}{(4\pi(t'-s))^{n/2}} \right) g(s, y) dy ds \\ &\quad - \int_t^{t'} \int_{\mathbb{R}^n} \frac{e^{-|x-y|^2/4(t'-s)}}{(4\pi(t'-s))^{n/2}} g(s, y) dy ds = I_1(t, t', x) + I_2(t, t', x). \end{aligned} \quad (1.4.35)$$

The second term above satisfies the simple estimate

$$|I_2(t, t', x)| \leq \|g\|_{L^\infty} |t' - t|, \quad (1.4.36)$$

obtained via the by now automatic change of variables as in (1.4.28). As for I_1 , we write, using the Newton-Leibniz formula in the t variable, for a fixed $s \in [0, t]$

$$\frac{e^{-|x-y|^2/4(t'-s)}}{(4\pi(t'-s))^{n/2}} - \frac{e^{-|x-y|^2/4(t-s)}}{(4\pi(t-s))^{n/2}} = \int_t^{t'} \frac{h(z)}{(4\pi(\tau-s))^{n/2+1}} d\tau, \quad z = \frac{x-y}{\sqrt{\tau-s}}, \quad (1.4.37)$$

with an integrable function

$$h(z) = \left(-\frac{n}{2} + \frac{|z|^2}{4} \right) e^{-|z|^2}.$$

Thus, we have, changing the variables $y \rightarrow z$ in the integral over \mathbb{R}^n , and integrating z out, using integrability of $h(z)$:

$$\begin{aligned} |I_1(t, t', x)| &\leq C \|g\|_{L^\infty} \int_0^t \int_t^{t'} \frac{d\tau}{\tau-s} ds = C \|g\|_{L^\infty} \int_0^t \log\left(\frac{t'-s}{t-s}\right) ds \\ &= C \|g\|_{L^\infty} (t' \log t' - t \log t - (t' - t) \log(t' - t)). \end{aligned} \quad (1.4.38)$$

This proves that

$$|I_1(t, t', x)| \leq C \|g\|_{L^\infty} |t' - t|^\alpha, \quad (1.4.39)$$

for all $\alpha \in (0, 1)$.

Exercise 1.4.12 Consider the partial differences

$$\frac{\partial J(t, x + he_j)}{\partial x_i} - \frac{\partial J(t, x)}{\partial x_i}$$

using expression (1.4.33) for $\partial_{x_i} J(t, x)$ and use a trick similar to (1.4.37) to show that $\partial_{x_i} J(t, x)$ is Hölder continuous and the last estimate in (1.4.26) holds.

This exercise finishes the proof of Proposition 1.4.7. \square

Let us stress again that the logarithmic term $\log(t - t')$ that appears in (1.4.38) is not a fluke of the proof: it represents a genuine obstacle to differentiability of $J(t, x)$ in time if $g(t, x)$ is just bounded and not Hölder continuous in t . This fact is absolutely crucial in the parabolic regularity theory, and not just in the present reasonably simple context.

1.4.3 Regularity gain: from Hölder $g(t, x)$ to differentiable $J(t, x)$

Proposition 1.4.7 shows that if we assume that $g(t, x)$ is bounded then $J(t, x)$ is differentiable in space and Hölder continuous in time, and, as we have seen, one can not expect a better regularity for $J(t, x)$ without further assumptions on the function $g(t, x)$. We now assume that $g(t, x)$ itself is Hölder continuous in time and space, and show that then $J(t, x)$ is differentiable in t and twice differentiable in x .

Proposition 1.4.13 Assume that $g(t, x) \in C_t^{\alpha/2} C_x^\alpha([0, T] \times \mathbb{R}^n)$, so that there exists $K > 0$ such that for all $0 \leq t, t' \leq T$ and $x, x' \in \mathbb{R}^n$ we have

$$|g(t, x) - g(t', x')| \leq K(|t - t'|^{\alpha/2} + |x - x'|^\alpha) \quad (1.4.40)$$

for some $\alpha \in (0, 1)$. Then $J(t, x)$ given by (1.4.20) is twice continuously differentiable in x , and once continuously differentiable in t over $(0, T) \times \mathbb{R}^n$. Moreover, there exists $C > 0$ so that

$$\|D_x^2 J\|_{L^\infty([0, T] \times \mathbb{R}^n)} + \|\partial_t J\|_{L^\infty([0, T] \times \mathbb{R}^n)} \leq C \|g\|_{C^{\alpha/2, \alpha}([0, T] \times \mathbb{R}^n)}. \quad (1.4.41)$$

Note the difference in the Hölder exponents in t and x in the assumption (1.4.40) on the function $g(t, x)$. It is of course not necessary as any Hölder assumptions in x and t would lead to result. But it is very natural, as will be clear from the argument below. In a similar fashion, the gain of regularity in time and space for J is different: one derivative in time and two derivatives in space. Both instances are related to the different scaling of the heat equation and other parabolic problems in time and space.

Continuing the theme started in Exercise 1.4.9, we ask the reader to consider the following question.

Exercise 1.4.14 Does the result of Proposition 1.4.13 apply to the solution $J_{wave}(t, x)$ to the wave equation given in Exercise 1.4.2?

Proof. One could look again at the partial differences, as in the proof of Proposition 1.4.7. However, we will use a different strategy, to illustrate another method. We will take $\delta \in (0, t)$ small, and consider an approximation

$$J_\delta(t, x) = \int_0^{t-\delta} e^{(t-s)\Delta} g(s, \cdot)(x) ds = \int_0^{t-\delta} \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} g(s, y) dy ds. \quad (1.4.42)$$

Note that $J_\delta(s, x)$ is the solution to the Cauchy problem (in the variable s , with t fixed)

$$\frac{\partial J_\delta}{\partial s} = \Delta J_\delta + H(t - s - \delta)g(s, x), \quad (1.4.43)$$

with the initial condition $J_\delta(0, x) = 0$. Here, we have introduced the cut-off $H(s) = 1$ for $s < 0$ and $H(s) = 0$ for $s > 0$.

The function $J_\delta(t, x)$ is smooth both in t and x for all $\delta > 0$ – this is easy to check simply by differentiating the integrand in (1.4.42) in t and x , since that does not produce any singularity due to $t - s > \delta$. Moreover, $J_\delta(t, x)$ converges uniformly to $J(t, x)$ as $\delta \downarrow 0$ – this follows from the estimate

$$|J(t, x) - J_\delta(t, x)| \leq \delta \|g\|_{L^\infty}, \quad (1.4.44)$$

that can be checked as in (1.4.22).

Exercise 1.4.15 Check that (1.4.44) holds.

As a consequence of (1.4.44), the derivatives of $J_\delta(t, x)$ converge weakly, in the sense of distributions, to the corresponding weak derivatives of $J(t, x)$. Thus, to show that, say, the second derivatives (understood in the sense of distributions) $\partial_{x_i x_j} J(t, x)$ are actually continuous functions, it suffices to prove that the partial derivatives $\partial_{x_i x_j} J_\delta(t, x)$ converge uniformly to a continuous function, and that is what we will do. In other words, we are relying on the following real analysis exercise.

Exercise 1.4.16 Assume that $f_n(x)$, $x \in \mathbb{R}$, is a sequence of infinitely differentiable functions that converges uniformly on \mathbb{R} to a limit $g \in C(\mathbb{R})$. Show that then $f'_n \rightarrow g'$ in the sense of distributions. Suppose, in addition, that there is a function $p \in C(\mathbb{R})$ such that $f'_n \rightarrow p$, also uniformly on \mathbb{R} . Show that then $g(x)$ is continuously differentiable and $p(x) = g'(x)$ for all $x \in \mathbb{R}$.

We will look in detail at $\partial_{x_i x_j} J_\delta$, with $i \neq j$. As the integrand for J_δ has no singularity at $s = t$, we may simply differentiate under the integral sign

$$\frac{\partial^2 J_\delta(t, x)}{\partial x_i \partial x_j} = \int_0^{t-\delta} \int_{\mathbb{R}^n} \frac{(x_i - y_i)(x_j - y_j)}{4(t-s)^2 (4\pi(t-s))^{n/2}} e^{-|x-y|^2/4(t-s)} g(s, y) ds dy.$$

The extra factor $(t-s)^2$ in the denominator can not be removed simply by the change of variable (1.4.28) – as the reader can immediately check, this would still leave a non-integrable extra factor of $(t-s)^{-1}$ that would cause an obvious problem in passing to the limit $\delta \downarrow 0$.

A very simple but absolutely crucial observation that will come to our rescue here is that, as $i \neq j$, we have

$$\int_{\mathbb{R}^n} (x_i - y_i)(x_j - y_j) e^{-|x-y|^2/4(t-s)} dy = 0. \quad (1.4.45)$$

This allows us to write

$$\frac{\partial^2 J_\delta(t, x)}{\partial x_i \partial x_j} = \int_0^{t-\delta} \int_{\mathbb{R}^n} \frac{(x_i - y_i)(x_j - y_j)}{4(t-s)^2 (4\pi(t-s))^{n/2}} e^{-|x-y|^2/4(t-s)} (g(s, y) - g(t, x)) ds dy. \quad (1.4.46)$$

Note that we use here crucially the fact that $\delta > 0$ and all integrals are finite because of that. Now, we can use the regularity of $g(s, y)$ to help us. In particular, the Hölder continuity assumption (1.4.40) gives

$$\begin{aligned} & \left| \frac{(x_i - y_i)(x_j - y_j)}{4(t-s)^2 (4\pi(t-s))^{n/2}} e^{-|x-y|^2/4(t-s)} (g(s, y) - g(t, x)) \right| \\ & \leq \frac{C|z|^2 e^{-|z|^2} (|t-s|^{\alpha/2} + |x-y|^\alpha)}{(t-s)(t-s)^{n/2}} \|g\|_{C_t^{\alpha/2} C_x^\alpha} \leq \frac{C}{(t-s)^{1-\alpha/2}} \frac{k(z)}{(4\pi(t-s))^{n/2}} \|g\|_{C_t^{\alpha/2} C_x^\alpha}, \end{aligned} \quad (1.4.47)$$

still with $z = (x-y)/\sqrt{t-s}$, as in (1.4.28), and

$$k(z) = |z|^2 e^{-|z|^2/4} (1 + |z|^\alpha).$$

As before, the factor of $(t-s)^{n/2}$ in the right side of (1.4.47) goes into the volume element

$$dz = \frac{dy}{(t-s)^{n/2}},$$

and we only have the factor $(t-s)^{1-\alpha/2}$ left in the denominator in (1.4.47), which is integrable in s , unlike the factor $(t-s)^{-1}$ one would get without using the cancellation in (1.4.45) and the Hölder regularity of $g(t, x)$. Thus, after accounting for the Jacobian factor, the integrand

in the expression for $\partial_{x_i x_j} J_\delta$ is dominated by an integrable function in z . This has two consequences. First, the Lebesgue dominated convergence theorem implies that

$$\frac{\partial^2 J_\delta(t, x)}{\partial x_i \partial x_j} \rightarrow Z_{ij}(t, x) := \int_0^t \int_{\mathbb{R}^n} \frac{(x_i - y_i)(x_j - y_j)}{4(t-s)^2 (4\pi(t-s))^{n/2}} e^{-|x-y|^2/4(t-s)} (g(s, y) - g(t, x)) ds dy, \quad (1.4.48)$$

as $\delta \rightarrow 0$, pointwise in t and x . In addition, the bound on the integrand in (1.4.47) implies that the convergence in (1.4.48) is uniform in $x \in \mathbb{R}^n$. In particular, the continuity of the limit $Z_{ij}(t, x)$ follows as well. Invoking the claim of Exercise 1.4.16, we now deduce that

$$\frac{\partial^2 J(t, x)}{\partial x_i \partial x_j} = Z_{ij}(t, x) \text{ for } i \neq j \text{ and all } t > 0 \text{ and } x \in \mathbb{R}^n, \quad (1.4.49)$$

and that these mixed second derivatives are continuous. In addition, we also see from (1.4.47) that

$$\left| \frac{\partial^2 J_\delta(t, x)}{\partial x_i \partial x_j} \right| \leq C \int_0^t \frac{ds}{(t-s)^{1-\alpha/2}} \|g\|_{C_t^{\alpha/2} C_x^\alpha} \leq C t^{\alpha/2} \|g\|_{C_t^{\alpha/2} C_x^\alpha}, \quad (1.4.50)$$

and thus the derivatives of $J(t, x)$ obey the same bound:

$$\left| \frac{\partial^2 J(t, x)}{\partial x_i \partial x_j} \right| \leq C t^{\alpha/2} \|g\|_{C_t^{\alpha/2} C_x^\alpha}. \quad (1.4.51)$$

Exercise 1.4.17 Complete the argument by looking at the remaining derivatives $\partial_t J(t, x)$ and $\partial_{x_i x_i} J(t, x)$. In both cases, one would start with J_δ , find a cancellation such as in (1.4.45), leading to a version of (1.4.46), and then pass to the limit $\delta \downarrow 0$ using the Hölder regularity of $g(t, x)$.

1.4.4 Regularity gain: from a Hölder g to Hölder derivatives of J

Proposition 1.4.13 is slightly sub-optimal: it says that if g is Hölder continuous then J is twice differentiable in x and once in t but says nothing about the continuity or regularity of these derivatives. We now show that they are actually themselves Hölder continuous then the Hölder continuity passes on to the corresponding derivatives of the solution.

Proposition 1.4.18 *Assume that $g(t, x) \in C_t^{\alpha/2} C_x^\alpha([0, T] \times \mathbb{R}^n)$, so that there exists $K > 0$ such that for all $0 \leq t, t' \leq T$ and $x, x' \in \mathbb{R}^n$ we have*

$$|g(t, x) - g(t', x')| \leq K(|t - t'|^{\alpha/2} + |x - x'|^\alpha) \quad (1.4.52)$$

for some $\alpha \in (0, 1)$. Then, there exists $C > 0$ such that

$$\|D_x^2 J\|_{C^{\alpha/2, \alpha}([0, T] \times \mathbb{R}^n)} + \|\partial_t J\|_{C_t^{\alpha/2} C_x^\alpha([0, T] \times \mathbb{R}^n)} \leq C \|g\|_{C_t^{\alpha/2} C_x^\alpha([0, T] \times \mathbb{R}^n)}. \quad (1.4.53)$$

There is a subtle but important point here. Our first result, Proposition 1.4.7 said that if g is in L^∞ then you nearly gain one derivative in time and two derivatives in space, but only nearly: the true result is that J and $D_x J$ are α -Hölder continuous in t and x for any $\alpha \in (0, 1)$ but it is not true that $\partial_t J$ and $D_x^2 J$ necessarily exist. Proposition 1.4.18 says, on the other

hand, that if g is Hölder continuous and not just bounded, so that $g \in C_t^{\alpha/2}C_x^\alpha$, then you fully gain one derivative in t and two in x : $\partial_t J$ and $D_x^2 J$ are bounded in the same space $C_t^{\alpha/2}C_x^\alpha$ as g . This result is optimal, one can not expect anything better, as can be seen simply from the form of the heat equation

$$J_t - \Delta J = g. \quad (1.4.54)$$

Warning. In the proof below we denote by the constants C, C' etc. various universal constants that do not depend on anything but elementary calculus, and, in particular, not on the function g, t or x . We make no attempt to optimize them. We also set, for some brevity

$$M_g = \|g\|_{C_t^{\alpha/2}C_x^\alpha}. \quad (1.4.55)$$

Proof. Our analysis follows what we did in Section 1.4.3 except we have to look at the Hölder differences for the second derivatives. The function $J(t, x)$ is given by the Duhamel formula

$$J(t, x) = \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} g(s, y) dy ds. \quad (1.4.56)$$

As in the proof of Proposition 1.4.13, we are going to examine only $\partial_{x_i x_j} J$, with $i \neq j$, leaving the other derivatives to the reader as a lengthy but straightforward exercise. Let us set

$$h_{ij}(z) = \frac{z_i z_j}{(4\pi)^{n/2}} e^{-|z|^2}, \quad D(s, t, x, y) = h_{ij}\left(\frac{x-y}{2\sqrt{t-s}}\right) \frac{g(s, y) - g(t, x)}{(t-s)^{n/2+1}},$$

so that we may write (1.4.48)-(1.4.49) as

$$\frac{\partial^2 J(t, x)}{\partial x_i \partial x_j} = \int_0^t \int_{\mathbb{R}^n} D(s, t, x, y) ds dy. \quad (1.4.57)$$

Now, for $0 < t \leq t' \leq T$ and x, x' in \mathbb{R}^n , we have

$$\begin{aligned} \frac{\partial^2 J(t', x')}{\partial x_i \partial x_j} - \frac{\partial^2 J(t, x)}{\partial x_i \partial x_j} &= \int_t^{t'} \int_{\mathbb{R}^n} D(s, t', x', y) ds dy \\ &+ \int_0^t \int_{\mathbb{R}^n} (D(s, t', x', y) - D(s, t, x, y)) ds dy = J_1(t, t', x') + J_2(t, x, t', x'). \end{aligned} \quad (1.4.58)$$

Exercise 1.4.19 Verify that no additional ideas other than what has already been developed in the proof of Propositions 1.4.7 and 1.4.13 are required to prove that the integral J_1 satisfies an inequality of the form

$$|J_1(t, t', x')| \leq CM_g |t - t'|^{\alpha/2}. \quad (1.4.59)$$

As for the integral J_2 , we need to look at it a little deeper. The change of variables

$$z = \frac{x-y}{2\sqrt{t-s}}$$

transforms (1.4.57) into

$$\frac{\partial^2 J(t, x)}{\partial x_i \partial x_j} = \int_0^t \int_{\mathbb{R}^n} h_{ij}(z) \frac{(g(s, x - 2\sqrt{t-s}z) - g(t, x)) ds dz}{t-s} \frac{1}{\pi^{n/2}},$$

and J_2 becomes

$$\begin{aligned}
& J_2(t, t', x, x') \\
&= \int_0^t \int_{\mathbb{R}^n} h_{ij}(z) \left[\frac{g(s, x' - 2\sqrt{t' - sz}) - g(t', x')}{t' - s} - \frac{g(s, x - 2\sqrt{t - sz}) - g(t, x)}{t - s} \right] \frac{dsdz}{\pi^{n/2}} \\
&= \int_0^t \int_{\mathbb{R}^n} h_{ij}(z) \left[\frac{g(s, x' - 2\sqrt{t' - sz}) - g(t', x')}{t' - s} - \frac{g(s, x - 2\sqrt{t' - sz}) - g(t', x)}{t' - s} \right] \frac{dsdz}{\pi^{n/2}} \\
&+ \int_0^t \int_{\mathbb{R}^n} h_{ij}(z) \left[\frac{g(s, x - 2\sqrt{t' - sz}) - g(t', x)}{t' - s} - \frac{g(s, x - 2\sqrt{t - sz}) - g(t, x)}{t - s} \right] \frac{dsdz}{\pi^{n/2}} \\
&= J_{21}(t, t', x, x') + J_{22}(t, t', x, x'). \tag{1.4.60}
\end{aligned}$$

We estimate each term separately.

The estimate of $J_{22}(t, t', x, x)$. We split the time integration domain $0 \leq s \leq t$ into the intervals

$$A = \{s : t - (t' - t) \leq s \leq t\}, \quad B = \{0 \leq s \leq t - (t' - t)\}.$$

Note that if $t' - t \geq t$, then $A = [0, t]$ and B is empty. The Hölder regularity of $g(t, x)$ in (1.4.52) implies that

$$|g(s, x - 2\sqrt{t' - sz}) - g(t', x)| \leq 2M_g(t' - s)^{\alpha/2}(1 + |z|^\alpha), \tag{1.4.61}$$

and

$$|g(s, x - 2\sqrt{t - sz}) - g(t, x)| \leq 2M_g(t - s)^{\alpha/2}(1 + |z|^\alpha). \tag{1.4.62}$$

Note that for $s \in A$ we have

$$t' - s \leq 2(t' - t), \quad t - s \leq (t' - t).$$

Hence, the contribution to J_{22} by the integral over the interval A can be bounded as

$$\begin{aligned}
J_{22}^A(t, t', x, x') &\leq 2M_g \int_{t-(t'-t)}^t \int_{\mathbb{R}^n} |h_{ij}(z)|(1 + |z|^\alpha) \left[\frac{1}{(t' - s)^{1-\alpha/2}} + \frac{1}{(t - s)^{1-\alpha/2}} \right] \frac{dsdz}{\pi^{n/2}} \\
&\leq C_\alpha (t' - t)^{\alpha/2} M_g, \tag{1.4.63}
\end{aligned}$$

with a constant C_α that depends only on $\alpha \in (0, 1)$. We used (1.4.61) and (1.4.62) above.

To estimate the contribution to J_{22} by the integral over the interval B , note that for $s \in B$ both increments $t - s$ and $t' - s$ are strictly positive, so that the integrand is not singular. Let us also recall that h_{ij} has zero integral. Thus, we may remove both $g(t, x)$ and $g(t', x')$ from the integral. This allows us to rewrite J_{22}^B as

$$\begin{aligned}
J_{22}^B(t, t', x, x') &= \int_0^{t-(t'-t)} \int_{\mathbb{R}^n} \left(\frac{g(s, x - 2\sqrt{t' - sz})}{t' - s} - \frac{g(s, x - 2\sqrt{t - sz})}{t - s} \right) h_{ij}(z) \frac{dsdz}{\pi^{n/2}} \\
&= \int_0^{t-(t'-t)} \int_{\mathbb{R}^n} \left(\frac{g(s, x - 2\sqrt{t' - sz}) - g(s, x - 2\sqrt{t - sz})}{t' - s} \right) h_{ij}(z) \frac{dsdz}{\pi^{n/2}} \\
&+ \int_0^{t-(t'-t)} \int_{\mathbb{R}^n} \left(\frac{g(s, x - 2\sqrt{t - sz})}{t - s} - \frac{g(s, x - 2\sqrt{t' - sz})}{t' - s} \right) h_{ij}(z) \frac{dsdz}{\pi^{n/2}} \\
&= J_{221}^B + J_{222}^B. \tag{1.4.64}
\end{aligned}$$

Note that the integrand in the term J_{221}^B can be bounded from above by

$$CM_g |z|^2 e^{-|z|^2} \frac{|\sqrt{t'-s} - \sqrt{t-s}| z|^\alpha}{t'-s}, \quad (1.4.65)$$

with a constant $C > 0$ that only depends on $\alpha \in (0, 1)$. Integrating out the z -variable then gives

$$\begin{aligned} J_{221}^B(t, t', x, x') &\leq CM_g \int_0^{t-(t'-t)} \frac{(\sqrt{t'-s} - \sqrt{t-s})^\alpha ds}{t'-s} \leq CM_g \int_0^{t-(t'-t)} \frac{1}{t-s} \frac{(t-t)^\alpha}{(t-s)^{\alpha/2}} ds \\ &\leq CM_g (t'-t)^{\alpha/2}. \end{aligned} \quad (1.4.66)$$

To estimate J_{222}^B we again use the zero integral property of $h_{ij}(z)$ to write this term as

$$J_{222}^B(t, t', x, x') = \int_0^{t-(t'-t)} \int_{\mathbb{R}^n} (g(s, x - 2\sqrt{t-s}z) - g(s, x)) \left(\frac{1}{t'-s} - \frac{1}{t-s} \right) h_{ij}(z) \frac{ds dz}{\pi^{n/2}}. \quad (1.4.67)$$

The integrand in (1.4.69) can be bounded by

$$CM_g |t-s|^{\alpha/2} |z|^\alpha |z|^2 e^{-|z|^2} \frac{t-t}{(t-s)^2}. \quad (1.4.68)$$

Integrating out the z -variable and then the s variable, we obtain

$$|J_{222}^B(t, t', x, x')| \leq CM_g |t-t'|^{\alpha/2}. \quad (1.4.69)$$

We conclude that

$$J_{22}(t, t', x, x') \leq CM_g (t'-t)^{\alpha/2}, \quad 0 < t \leq t'. \quad (1.4.70)$$

The estimate of $J_{21}(t, t', x, x')$. Now, we estimate

$$\begin{aligned} &J_{21}(t, t', x, x') \\ &= \int_0^t \int_{\mathbb{R}^n} h_{ij}(z) \left[\frac{g(s, x' - 2\sqrt{t'-s}z) - g(t', x')}{t'-s} - \frac{g(s, x - 2\sqrt{t'-s}z) - g(t', x)}{t'-s} \right] \frac{ds dz}{\pi^{n/2}} \\ &= J_{21}^A + J_{21}^B. \end{aligned} \quad (1.4.71)$$

The two terms above refer to the integration over the time interval $A = \{t - |x - x'|^2 \leq s \leq t\}$ and its complement B . As before, if $t \leq |x - x'|^2$, then we only have $A = \{0 \leq s \leq t\}$. In the first domain, we just use the bounds

$$|g(s, x' - 2\sqrt{t'-s}z) - g(t', x')| \leq CM_g (t'-s)^{\alpha/2} (1 + |z|^\alpha) \quad (1.4.72)$$

and

$$|g(s, x - 2\sqrt{t'-s}z) - g(t', x)| \leq CM_g (t'-s)^{\alpha/2} (1 + |z|^\alpha). \quad (1.4.73)$$

After integrating out the z -variable, this leads to

$$|J_{21}^A(t, t', x, x')| \leq CM_g \int_{t-|x-x'|^2}^t (t'-s)^{-1+\frac{\alpha}{2}} ds \leq CM_g (|t'-t|^{\alpha/2} + |x-x'|^\alpha). \quad (1.4.74)$$

Next, as h_{ij} has zero mass and $t' - s$ is strictly positive when $s \in B$, we can drop the terms involving $g(t', x')$ and $g(t', x)$ leading to

$$\begin{aligned} J_{21}^B(t, t', x, x') &= \int_0^{t-|x-x'|^2} \int_{\mathbb{R}^n} h_{ij}(z) \frac{g(s, x' - 2\sqrt{t'-s}z) - g(s, x - 2\sqrt{t'-s}z)}{t' - s} \frac{dsdz}{\pi^{n/2}} \\ &= \int_0^{t-|x-x'|^2} \int_{\mathbb{R}^n} \left(h_{ij}\left(\frac{x' - y}{2\sqrt{t'-s}}\right) - h_{ij}\left(\frac{x - y}{2\sqrt{t'-s}}\right) \right) \frac{g(s, y)}{t' - s} \frac{dsdy}{(4\pi(t'-s))^{n/2}}. \end{aligned} \quad (1.4.75)$$

Once again, because h_{ij} has zero mass we have

$$J_{21}^B(t, t', x, x') = \int_0^{t-|x-x'|^2} \int_{\mathbb{R}^n} \left(h_{ij}\left(\frac{x' - y}{2\sqrt{t'-s}}\right) - h_{ij}\left(\frac{x - y}{2\sqrt{t'-s}}\right) \right) \frac{g(s, y) - g(t', x')}{t' - s} \frac{dsdy}{(4\pi(t'-s))^{n/2}}.$$

The integrand above can be re-written as

$$\begin{aligned} &\left(h_{ij}\left(\frac{x' - y}{2\sqrt{t'-s}}\right) - h_{ij}\left(\frac{x - y}{2\sqrt{t'-s}}\right) \right) \frac{g(s, y) - g(t', x')}{t' - s} \\ &= \frac{1}{2} \int_0^1 \frac{g(s, y) - g(s, x_\sigma) + g(s, x_\sigma) - g(t', x')}{(t' - s)^{3/2}} (x' - x) \cdot \nabla h_{ij}\left(\frac{x_\sigma - y}{2\sqrt{t'-s}}\right) d\sigma, \end{aligned} \quad (1.4.76)$$

with $x_\sigma = \sigma x + (1 - \sigma)x'$. It follows that

$$\begin{aligned} |J_{21}^B(t, t', x, x')| &\leq CM_g |x - x'| \int_0^{t-|x-x'|^2} \int_0^1 \int_{\mathbb{R}^n} \left| \nabla h_{ij}\left(\frac{x_\sigma - y}{2\sqrt{t'-s}}\right) \right| \\ &\quad \times \frac{|y - x_\sigma|^\alpha + |x' - x_\sigma|^\alpha + |t' - s|^{\alpha/2}}{(t' - s)^{3/2}} \frac{dsdyd\sigma}{(t' - s)^{n/2}}. \end{aligned} \quad (1.4.77)$$

Using the estimates

$$|\nabla h(z)| \leq C|z|^3 e^{-|z|^2},$$

and $|x' - x_\sigma| \leq |x - x'|$, and making the usual change of variable

$$z = \frac{x_\sigma - y}{2\sqrt{t'-s}},$$

and integrating out the z and σ variables, we arrive at

$$|J_{21}^B(t, t', x, x')| \leq CM_g |x - x'| \int_0^{t-|x-x'|^2} \left(\frac{1}{(t-s)^{(3-\alpha)/2}} + \frac{|x-x'|^\alpha}{(t-s)^{3/2}} \right) ds. \quad (1.4.78)$$

Integrating out the s variable, we obtain

$$|J_{21}^B(t, t', x, x')| \leq CM_g |x - x'| (|x - x'|^{\alpha-1} + |x - x'|^\alpha |x - x'|^{-1}) \leq CM_g |x - x'|^\alpha, \quad (1.4.79)$$

thus J_{21} is also Hölder continuous, finishing the proof. \square

Higher order derivatives

The previous results can be generalized to the higher order derivatives of $J(t, x)$ assuming that the corresponding derivatives of $g(t, x)$ exist and are Hölder continuous. Recall that we use the notation

$$|k| = k_1 + \dots + k_n,$$

for a multi-index $k = (k_1, \dots, k_n)$, and

$$D_x^k u = \frac{\partial^{|k|} u}{\partial x_1^{k_1} \dots \partial x_n^{k_n}}.$$

The following result will be indispensable in the analysis of nonlinear equations, despite its seemingly technical nature.

Proposition 1.4.20 *Assume that the function $g(t, x)$ is M times continuously differentiable in t and K times continuously differentiable in x , and $\partial_t^M D_x^k g(t, x) \in C_t^{\alpha/2} C_x^\alpha([0, T] \times \mathbb{R}^n)$ for all multi-indices k with $|k| = K$, for some $\alpha \in (0, 1)$. Then $J(t, x)$ given by (1.4.20) is $M + 1$ times continuously differentiable in t , and $K + 2$ times continuously differentiable in x for all $0 \leq t \leq T$ and $x \in \mathbb{R}^n$. Moreover, there exists $C > 0$ that depends on M and K so that for any multi-indices k and k' with $|k| = K$ and $|k'| = K + 2$ we have*

$$\|\partial_t^{M+1} D_x^k J\|_{C_t^{\alpha/2} C_x^\alpha([0, T] \times \mathbb{R}^n)} + \|\partial_t^M D_x^{k'} J\|_{C_t^{\alpha/2} C_x^\alpha([0, T] \times \mathbb{R}^n)} \leq C \sup_{0 \leq |r| \leq K, 0 \leq m \leq M} \|D_t^m D_x^r g\|_{C_t^{\alpha/2} C_x^\alpha([0, T] \times \mathbb{R}^n)}. \quad (1.4.80)$$

In particular, if $g(t, x)$ is infinitely differentiable with each derivative uniformly bounded in t and x then so is $J(t, x)$.

Exercise 1.4.21 Provide the proof of Proposition 1.4.20. One way to run the argument is to solve $\partial_t v - \Delta v = D_t^m D_x^k g$, $v(0) = 0$, apply the results we proved above to v and then show that $D_t^m D_x^k u = v$.

A remark on the constant coefficients case

To finish this section, consider solutions to general constant coefficients equations of the form

$$u_t - a_{ij} \partial_{x_i x_j} u + b_j \partial_{x_j} u + cu = f(t, x). \quad (1.4.81)$$

We assume that a_{ij} , b_i and c are constants, and the matrix $A := (a_{ij})$ is positive definite: there exists a constant $\lambda > 0$ so that for any vector $\xi \in \mathbb{R}^n$ we have

$$a_{ij} \xi_i \xi_j \geq \lambda |\xi|^2. \quad (1.4.82)$$

Assume also that f is an α -Hölder function over $[0, T] \times \mathbb{R}^n$, and take the initial condition $v(0, x) \equiv 0$. The function $v(t, x) = u(t, x + Bt) \exp(ct)$, with $B = (b_1, \dots, b_n)$, solves

$$v_t - a_{ij} \partial_{x_i x_j} v = f(t, x + Bt). \quad (1.4.83)$$

The change of variable $w(t, x) = v(t, \sqrt{A}x)$ brings us back to the forced heat equation:

$$w_t - \Delta w = f(t, \sqrt{A}(x + Bt)). \quad (1.4.84)$$

We see that the conclusion of Proposition 1.4.18 also applies to other parabolic equations with constant coefficients, as long as the ellipticity condition (1.4.82) holds.

Exercise 1.4.22 Consider the solutions of the equation

$$u_t - u_{xx} + u_y = f(t, x, y), \quad (1.4.85)$$

in \mathbb{R}^2 and use this example to convince yourself that the ellipticity condition is necessary for the Hölder regularity as in Proposition 1.4.18 to hold.

Congratulations. We congratulate the reader who managed to follow the lengthy computations in this section!

1.5 Regularity for the nonlinear heat equations

In this section, we reap the fruit of our labour in the previous section and prove global in time existence of solutions to some nonlinear parabolic equations. We will not strive to achieve the sharpest results. Rather, we have in mind two particular classes of nonlinear parabolic equations, for which eventually we would like to understand the large time behavior: the semi-linear and quasi-linear equations of the simplest form. The truth is that the two examples we consider here contain some of the main features under which the more general global existence and regularity results hold: the Lipschitz behavior of the nonlinearity, and the smooth spatial dependence of the coefficients in the equation. Thus, after reading this section the reader should be well prepared to digest the more general results described in other, more specialized books.

1.5.1 Existence and regularity for a semi-linear diffusion equation

First, we consider semi-linear parabolic equations of the form

$$u_t = \Delta u + f(t, x, u). \quad (1.5.1)$$

Such equations are generally known as the reaction-diffusion equations, and are very common in biological and physical sciences. We will discuss the origins of such equations, and the behavior of the solutions to a class of such equations in great detail in Chapter ??.

We will consider (1.5.1) posed in \mathbb{R}^n , and equipped with a bounded and continuous initial condition

$$u(0, x) = u_0(x). \quad (1.5.2)$$

As in the theory of nonlinear ordinary differential equations, we need to assume some Lipschitz property of the function $f(t, x, u)$ in the u -variable. Otherwise we may run into blow-up issues, familiar from the solutions to the ordinary differential equation

$$\frac{du}{dt} = u^2, \quad u(0) = u_0. \quad (1.5.3)$$

Recall that if $u_0 > 0$ then solution to (1.5.3) exists only until the time $T_0 = 1/u_0$ and

$$\lim_{t \uparrow T_0} u(t) = +\infty. \quad (1.5.4)$$

This is something we would like to avoid in this expository section.

There are two possible assumptions that will ensure that solutions to (1.5.1) exist and do not blow up in a finite time. First, we may simply assume that the function f is smooth in all its variables and globally Lipschitz in u : there exists a constant $C_f > 0$ so that

$$|f(t, x, u) - f(t, x, u')| \leq C_f |u - u'|, \text{ for all } t \geq 0, x \in \mathbb{R}^n \text{ and } u, u' \in \mathbb{R}. \quad (1.5.5)$$

Alternatively, we may assume that $f(t, x, u)$ is smooth in all its variables, and locally Lipschitz in u : for every $K > 0$ there exists $C_K > 0$ such that

$$|f(t, x, u) - f(t, x, u')| \leq C_K |u - u'|, \text{ for all } t \geq 0, x \in \mathbb{R}^n \text{ and } |u|, |u'| \leq K, \quad (1.5.6)$$

and, in addition, there exist $M_1 < M_2$ so that

$$f(t, x, M_1) = f(t, x, M_2) = 0 \text{ for all } t \geq 0 \text{ and } x \in \mathbb{R}^n. \quad (1.5.7)$$

Under this assumption we will show that solutions corresponding to initial data u_0 such that $M_1 \leq u_0 \leq M_2$ will be globally regular. One reason why (1.5.6)-(1.5.7) is a useful alternative to the global Lipschitz assumption in (1.5.5) is the Fisher-KPP equation

$$u_t = \Delta u + u(1 - u), \quad (1.5.8)$$

with the predator-prey nonlinearity $f(u) = u(1 - u)$ that does not satisfy (1.5.5) but which does obey (1.5.6)-(1.5.7). We refer the reader to Chapter ?? for the discussion of how this equation arises in the biological modeling and other applications, as well as to the explanation of its name.

Another important example is the time-dependent version of the Allen-Cahn equation we have encountered in Chapter ??:

$$u_t = \Delta u + u - u^3. \quad (1.5.9)$$

Here, once again, the nonlinearity $f(u) = u - u^3$ satisfies (1.5.6)-(1.5.7) but not (1.5.5). We will prove the following existence result under assumptions (1.5.6)-(1.5.7).

Theorem 1.5.1 *Assume that assumptions (1.5.6)-(1.5.7) hold with some $M_1 < M_2$, and the initial condition $u_0(x)$ is bounded and smooth, and*

$$M_1 \leq u_0(x) \leq M_2 \text{ for all } x \in \mathbb{R}^n. \quad (1.5.10)$$

Then, there exists a unique bounded smooth solution $u(t, x)$ to (1.5.1)-(1.5.2), which, in addition, satisfies

$$M_1 < u(t, x) < M_2 \text{ for all } t > 0 \text{ and } x \in \mathbb{R}^n. \quad (1.5.11)$$

Moreover, for all $T > 0$ each derivative of u is uniformly bounded over $[T, +\infty) \times \mathbb{R}^n$.

Assumption (1.5.7) may seem too stringent to the reader. Its role is to ensure that $u(t, x)$ satisfies the uniform bounds in (1.5.11). We explain in Exercise 1.5.4 how this assumption can be relaxed, while still ensuring that (1.5.11) holds. We also do not need to assume that the initial condition $u_0(x)$ is smooth – it suffices to assume that it is bounded and continuous. This is the subject of Exercise 1.5.5.

Let us first explain the two simpler claims in Theorem 1.5.1: the bounds in (1.5.11) and uniqueness. Let $u(t, x)$ be a bounded smooth solution to (1.5.1)-(1.5.2) with an initial condition $u_0(x)$ that satisfies (1.5.10). Consider the function $v(t, x) = u(t, x) - M_1$. This function satisfies

$$v_t = \Delta v + c(t, x)v, \quad (1.5.12)$$

with

$$c(t, x) = \frac{f(t, x, u(t, x)) - f(t, x, M_1)}{u(t, x) - M_1}. \quad (1.5.13)$$

As the function $u(t, x)$ is bounded, assumptions (1.5.7) and (1.5.6) imply that $c(t, x)$ is also bounded. Hence, the comparison principle in Theorem 1.3.3 can be applied to (1.5.12). As $v(0, x) \geq 0$ for all $x \in \mathbb{R}^n$, it follows that $v(t, x) > 0$ for all $t > 0$, so that $u(t, x) > M_1$ for all $x \in \mathbb{R}^n$. The other inequality in (1.5.11) can be proved similarly.

Uniqueness of bounded solutions is proved in an analogous fashion. Assume that $u_1(x)$ and $u_2(x)$ are two smooth bounded solutions to the Cauchy problem (1.5.1)-(1.5.2). Then $w = u_1 - u_2$ satisfies

$$w_t = \Delta w + c(t, x)w, \quad (1.5.14)$$

with the initial condition $w(0, x) = 0$ and a bounded function

$$c(t, x) = \frac{f(t, x, u_1(t, x)) - f(t, x, u_2(t, x))}{u_1(t, x) - u_2(t, x)}.$$

The comparison principle then implies that both $w(t, x) \leq 0$ and $w(t, x) \geq 0$, thus $w(t, x) \equiv 0$, proving the uniqueness.

Thus, the main issue in the proof of Theorem 1.5.1 is to prove the existence of a bounded solution to (1.5.1)-(1.5.2). As the function $f(t, x, u)$ is not necessarily globally Lipschitz, we are going to use the following trick based on the fact that f satisfies (1.5.7). Consider a function $\tilde{f}(t, x, u)$ such that

$$f(t, x, u) = \tilde{f}(t, x, u) \text{ for all } x \in \mathbb{R}^n \text{ and } M_1 \leq u \leq M_2, \quad (1.5.15)$$

and there exists $K > 0$ so that

$$|\tilde{f}(t, x, u)| \leq K \text{ for all } x \in \mathbb{R}^n \text{ and } u \in \mathbb{R}. \quad (1.5.16)$$

We may also ensure that $\tilde{f}(t, x, u)$ is globally Lipschitz: there exists $C_f > 0$ so that

$$|\tilde{f}(t, x, u_1) - \tilde{f}(t, x, u_2)| \leq C_f |u_1 - u_2|, \text{ for all } t \geq 0, x \in \mathbb{R}^n \text{ and all } u_1, u_2 \in \mathbb{R}, \quad (1.5.17)$$

as compared to $f(t, x, u)$ that is only locally Lipschitz in u . Note that (1.5.7) holds automatically for $\tilde{f}(t, x, u)$ because of (1.5.15). Hence, by what we have just shown, any smooth bounded solution to the initial value problem

$$\begin{aligned} u_t &= \Delta u + \tilde{f}(t, x, u), \\ u(0, x) &= u_0(x), \end{aligned} \quad (1.5.18)$$

with an initial condition $u_0(x)$ that satisfies assumption (1.5.10) will obey (1.5.11):

$$M_1 < u(t, x) < M_2 \text{ for all } t > 0 \text{ and } x \in \mathbb{R}^n. \quad (1.5.19)$$

It follows that $\tilde{f}(t, x, u(x, t)) \equiv f(t, x, u(x, t))$, and thus any bounded solution to (1.5.18) is a solution to (1.5.1) with the same initial condition. Therefore, it suffices to construct a bounded solution to (1.5.18).

A typical approach to the existence proofs in nonlinear problems is to use a fixed point argument. To this end, it is useful, and standard, to rephrase the parabolic initial value problem (1.5.1)-(1.5.2) as an integral equation, using the Duhamel formula. This is done as follows. Given a fixed $T > 0$ and initial condition $u_0(x)$, we define an operator \mathcal{T} as a mapping of the space $C([0, T] \times \mathbb{R}^n)$ to itself via

$$\begin{aligned} [\mathcal{T}u](t, x) &= e^{t\Delta}u_0(x) + \int_0^t e^{(t-s)\Delta}\tilde{f}(s, \cdot, u(s, \cdot))(x)ds \\ &= e^{t\Delta}u_0(x) + \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} \tilde{f}(s, y, u(s, y)) dy ds, \end{aligned} \quad (1.5.20)$$

with the operator $e^{t\Delta}$ defined in (1.4.12):

$$e^{t\Delta}\eta(x) = \frac{1}{(4\pi t)^{n/2}} \int_{\mathbb{R}^n} e^{-(x-y)^2/(4t)} \eta(y) dy. \quad (1.5.21)$$

The Duhamel formula for the solution to the Cauchy problem (1.5.18) can be now succinctly restated as

$$u(t, x) = [\mathcal{T}u](t, x). \quad (1.5.22)$$

In other words, any smooth bounded solution to the initial value problem is a fixed point of the operator \mathcal{T} .

Exercise 1.5.2 Show that if a function $u(t, x)$ that is continuously differentiable in t and twice continuously differentiable in x satisfies (1.5.22), then $u(t, x)$ is a solution to the initial value problem (1.5.18).

Thus, to prove the existence part of Theorem 1.5.1 we need to show that a fixed point of the operator \mathcal{T} exists and is sufficiently regular to differentiate it once in t and twice in x .

Existence of a fixed point: the Picard iteration argument on a short time interval

The first step is to prove the existence of a fixed point of \mathcal{T} in $C([0, T] \times \mathbb{R}^n)$ for $T > 0$ sufficiently small. In the second step, we will extend the existence result to all $T > 0$.

We will use the standard Picard iteration approach: set $u^{(0)} = 0$ and define

$$u^{(n+1)}(t, x) = \mathcal{T}u^{(n)}(t, x). \quad (1.5.23)$$

In particular, we have

$$u^{(1)}(t, x) = e^{t\Delta}u_0. \quad (1.5.24)$$

As the initial condition $u_0(x)$ is continuous and bounded, the function $u^{(1)}(t, x)$ is infinitely differentiable in t and x . Proposition 1.4.20, combined with a simple induction argument, shows that $u^{(n)}(t, x)$ are smooth for $t > 0$ and $x \in \mathbb{R}^n$, for all $n \geq 1$.

The global Lipschitz property (1.5.17) of $\tilde{f}(t, x, u)$ allows us to write

$$\begin{aligned} |\mathcal{T}u(t, x) - \mathcal{T}v(t, x)| &\leq \int_0^t \int_{\mathbb{R}^n} \frac{e^{-(x-y)^2/(4(t-s))}}{(4\pi(t-s))^{n/2}} |\tilde{f}(s, y, u(s, y)) - \tilde{f}(s, y, v(s, y))| dy ds \\ &\leq C_f \int_0^t \int_{\mathbb{R}^n} \frac{e^{-(x-y)^2/(4(t-s))}}{(4\pi(t-s))^{n/2}} |u(s, y) - v(s, y)| dy ds \\ &\leq C_f T \sup_{0 \leq s \leq T, y \in \mathbb{R}^n} |u(s, y) - v(s, y)|. \end{aligned} \tag{1.5.25}$$

This shows that if $T < C_f^{-1}$, then the mapping \mathcal{T} is a contraction on $C([0, T] \times \mathbb{R}^n)$ and thus has a unique fixed point in $C([0, T] \times \mathbb{R}^n)$. Before we extend this result to all $T > 0$ we first show that the fixed point is a smooth function hence a classical solution to (1.5.18) on $[0, T]$.

The bootstrap argument

Smoothness of the fixed point $u(t, x)$ is proved using what is commonly called a boot-strap argument. The key observation is the following.

Lemma 1.5.3 *Let $u(t, x)$ be a fixed point of the operator \mathcal{T} in $C([0, T]; \mathbb{R}^n)$, so that it satisfies (1.5.22), and is bounded and continuous on $[0, T] \times \mathbb{R}^n$. Then $u(t, x)$ is infinitely differentiable in t and x for all $t > 0$ and $x \in \mathbb{R}^n$.*

Proof. Let us write (1.5.22) as

$$u(t, x) = u^{(1)}(t, x) + D[u](t, x), \tag{1.5.26}$$

with

$$u^{(1)}(t, x) = e^{t\Delta} u_0, \tag{1.5.27}$$

and

$$D[u](t, x) = \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} g(s, y) dy ds, \tag{1.5.28}$$

where

$$g(t, x) = \tilde{f}(t, x, u(t, x)). \tag{1.5.29}$$

As we have noted, the function $u^{(1)}(t, x)$ is infinitely differentiable for any $t > 0$ and $x \in \mathbb{R}^n$ simply because it is a solution to the heat equation with a bounded and continuous initial condition u_0 . Thus, we only need to deal with the Duhamel term $D[u](t, x)$. To treat this term, we will use Propositions 1.4.7 and 1.4.18. The function $g(t, x)$ defined in (1.5.29) is bounded on $[0, T] \times \mathbb{R}^n$ because of (1.5.16). Hence, we may apply Proposition 1.4.7 and deduce that $D[u](t, x)$ is actually Hölder continuous on $[0, T] \times \mathbb{R}^n$, and we have a priori bounds

$$\begin{aligned} \|D[u]\|_{L^\infty} &\leq C \|g\|_{L^\infty}, \\ \|\partial_{x_k} D[u]\|_{L^\infty} &\leq C \|g\|_{L^\infty} \end{aligned} \tag{1.5.30}$$

and for any $\alpha \in (0, 1)$ there is C_α so that

$$|D[u](t, x) - D[u](t', x)| \leq C_\alpha \|g\|_{L^\infty} |t - t'|^\alpha. \quad (1.5.31)$$

From this and (1.5.26), we conclude that u itself satisfies the same bounds:

$$\begin{aligned} \|u\|_{L^\infty} &\leq \|u_0\|_{L^\infty} + C\|g\|_{L^\infty}, \\ \|\partial_{x_k} u\|_{L^\infty} &\leq \|D_x u_0\|_{L^\infty} + C\|g\|_{L^\infty} \\ |u(t, x) - u(t', x)| &\leq \|D_x^2 u_0\|_{L^\infty} |t - t'| + C_\alpha \|g\|_{L^\infty} |t - t'|^\alpha. \end{aligned} \quad (1.5.32)$$

Therefore, $u(t, x)$ is not just continuous and bounded, but also Hölder continuous in t and x , with explicit bounds above. Then, so is $g(t, x) = \tilde{f}(t, x, u(t, x))$, and then Proposition 1.4.18 tells us that $D[u](t, x)$ is differentiable once in t and twice in x and the derivatives $\partial_t u$ and $D_x^2 u$ are themselves Hölder continuous. Then, (1.5.26), in turn, implies that $u(t, x)$ is differentiable in t and twice differentiable in x , with Hölder continuous derivatives, and thus $g(t, x)$ possesses the same regularity. We may iterate this argument, using Proposition 1.4.18, each time gaining derivatives in t and x , and conclude that, actually, $u(t, x)$ is infinitely differentiable in t and x . This is known as a boot-strap argument.

The global in time existence

To show that existence of a solution can be extended to all $T > 0$, note that, as we have shown that the fixed point $u(t, x)$ of the mapping \mathcal{T} is smooth, we know that $u(t, x)$ is a classical solution to the initial value problem (1.5.18) on the time interval $0 \leq t \leq T$, hence it satisfies

$$M_1 \leq u(T, x) \leq M_2, \quad \text{for all } x \in \mathbb{R}^n. \quad (1.5.33)$$

Moreover, the existence time T does not depend on u_0 . Therefore, we can repeat the Picard iteration argument on the time intervals $[T, 2T]$, $[2T, 3T]$, and so on, eventually constructing a global in time solution to the Cauchy problem. This finishes the proof of Theorem 1.5.1.

Exercise 1.5.4 Assumption (1.5.7) is more stringent than necessary. Show that the claim of Theorem 1.5.1 holds also if instead of (1.5.7) we assume that there exist M_1 and M_2 such that

$$f(t, x, M_1) \operatorname{sgn}(M_1) \leq 0 \text{ and } f(t, x, M_2) \operatorname{sgn}(M_2) \leq 0 \quad \text{for all } t \geq 0 \text{ and } x \in \mathbb{R}^n, \quad (1.5.34)$$

and the initial condition $u_0(x)$ satisfies (1.5.10).

Exercise 1.5.5 As the reader may have noticed, the only place where we have used the assumption that the initial condition $u_0(x)$ is smooth is in estimates (1.5.32). It was used there to bound the derivatives of the contribution $u^{(1)}(t, x) = e^{t\Delta} u_0$ to $u(t, x)$ in (1.5.26). However, this term is smooth even if $u_0(x)$ is just continuous and not necessarily smooth. Use this to show that the conclusion of Theorem 1.5.1 holds if we only assume that u_0 is bounded and continuous.

1.5.2 The regularity of the solutions to a quasi-linear heat equation

One may wonder if the treatment that we have given to the semi-linear heat equation (1.5.1) is too specialized. To dispel this concern, we show how the above approach can be extended to equations with a drift and quasi-linear heat equations of the form

$$u_t - \Delta u = f(t, x, \nabla u), \quad (1.5.35)$$

posed for $t > 0$ and $x \in \mathbb{R}^n$. The nonlinearity is now stronger: it depends not on u itself but on its gradient ∇u . We ask that the nonlinear term $f(t, x, p)$ satisfies the following two hypotheses: first, there exists $C_1 > 0$ so that

$$|f(t, x, 0)| \leq C_1 \text{ for all } t \geq 0 \text{ and } x \in \mathbb{R}^n, \quad (1.5.36)$$

and, second, f is uniformly Lipschitz in the p -variable: there exists $C_2 > 0$ so that

$$|f(t, x, p_1) - f(t, x, p_2)| \leq C_2 |p_1 - p_2|, \text{ for all } t \geq 0 \text{ and } x, p_1, p_2 \in \mathbb{R}^n. \quad (1.5.37)$$

One consequence of (1.5.36) and (1.5.37) is a uniform bound

$$|f(t, x, p)| \leq C_3(1 + |p|), \text{ for all } t \geq 0, x, p \in \mathbb{R}^n, \quad (1.5.38)$$

showing that $f(t, x, p)$ grows at most linearly in p . We also require that $f(t, x, p)$ is smooth in t, x and p , and obeys the estimates

$$|\partial_t^m f(t, x, p)| + |D_x^k f(t, x, p)| \leq C_{m,k}(1 + |p|), \text{ for all } t \geq 0, x, p \in \mathbb{R}^n, \quad (1.5.39)$$

for any $m \geq 1$ and multi-index $k \in \mathbb{Z}^n$. This smoothness assumption can be greatly relaxed but we are not concerned with the optimal results here.

Two standard examples of equations of the form (1.5.35) are parabolic equations with constant diffusion and nonuniform drifts, such as

$$u_t = \Delta u + b_j(t, x) \frac{\partial u}{\partial x_j}, \quad (1.5.40)$$

with a prescribed drift $b(t, x) = (b_1(t, x), \dots, b_n(t, x))$, and viscous regularizations of the Hamilton-Jacobi equations, such as

$$u_t = \Delta u + f(|\nabla u|). \quad (1.5.41)$$

We will encounter both of them in the sequel. Our goal is to prove the following.

Theorem 1.5.6 *Under the above assumptions, equation (1.5.35), equipped with a bounded continuous initial condition u_0 , has a unique smooth solution $u(t, x)$ over $(0, +\infty) \times \mathbb{R}^n$, which is bounded with all its derivatives over every set of the form $(\varepsilon, T) \times \mathbb{R}^n$, with $0 < \varepsilon < T$.*

We will use the ideas displayed in the proof of Theorem 1.5.1. However, a serious additional difficulty for a quasi-linear equation (1.5.35) compared to a semi-linear equation such as (1.5.1) is that it involves a nonlinear function of the gradient of the function u , which, a priori, may not be smooth at all. That is, if u is not smooth, and its gradient is only a distribution,

giving the meaning to a nonlinear function $f(x, \nabla u)$ becomes problematic. Note that there is no problem of that sort with the Laplacian Δu in (1.5.35), as we may interpret it in the sense of distributions. In addition, if we try to write down the Duhamel formula for (1.5.35), an analog to (1.5.20)-(1.5.18), it would take the form

$$u(t, x) = [\mathcal{T}u](t, x), \quad (1.5.42)$$

with the operator \mathcal{T} given now by

$$[\mathcal{T}u](t, x) = e^{t\Delta}u_0(x) + \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} f(s, y, \nabla u(s, y)) dy ds. \quad (1.5.43)$$

This operator can not be considered as a mapping on $C([0, T] \times \mathbb{R}^n)$ because of the term involving ∇u . Hence, the strategy in the proof of Theorem 1.5.1 needs to be modified.

A natural and standard idea is to regularize the nonlinear term, and then pass to the limit, removing the regularization. We will consider the following nonlocal approximation to (1.5.35):

$$u_t^\varepsilon - \Delta u^\varepsilon = f(t, x, \nabla v^\varepsilon), \quad v^\varepsilon = e^{\varepsilon\Delta}u^\varepsilon. \quad (1.5.44)$$

When $\varepsilon > 0$ is small, one expects the solutions to (1.5.35) and (1.5.44) to be close as

$$e^{\varepsilon\Delta}\psi \rightarrow \psi, \text{ as } \varepsilon \rightarrow 0. \quad (1.5.45)$$

Exercise 1.5.7 For ψ in which function spaces does the convergence in (1.5.45) hold? For instance, does it hold in L^2 or L^∞ ? How about $C^1(\mathbb{R})$?

A damper on our expectations is that the convergence in (1.5.45) does not automatically translate into the convergence of the corresponding gradients, unless we already know that ψ is differentiable. In other words, there is no reason to expect that

$$\nabla(e^{\varepsilon\Delta}\psi) \rightarrow \nabla\psi,$$

simply because the right side may not exist. Unfortunately, a result of this kind is exactly what we need in order to understand the convergence of the term $f(x, \nabla v^\varepsilon)$ in (1.5.44).

Nevertheless, a huge advantage of (1.5.44) over (1.5.35) is that the function v^ε that appears inside the nonlinearity is smooth if u^ε is merely continuous, as long as $\varepsilon > 0$. This can be used to show that the Cauchy problem for (1.5.44) has a unique smooth solution.

Exercise 1.5.8 Show that, for every $\varepsilon > 0$ and every bounded function $u(x)$, we have

$$\|\nabla(e^{\varepsilon\Delta}u)\|_{L^\infty} \leq \frac{C}{\sqrt{\varepsilon}}\|u\|_{L^\infty}. \quad (1.5.46)$$

Use this fact, and the strategy in the proof of Theorem 1.5.1, to prove that (1.5.44), equipped with a bounded continuous initial condition u_0 , has a unique smooth solution u^ε over a set of the form $(0, T_\varepsilon] \times \mathbb{R}^n$, with a time $T_\varepsilon > 0$ that depends on $\varepsilon > 0$ but not on the initial condition u_0 .

Recommendation. The reader should take this exercise very seriously. You do not need any tools beyond what has been already done in this chapter, and it presents a good opportunity to check your understanding so far.

Having constructed solutions to (1.5.44) on a finite time interval $[0, T_\varepsilon]$, in order to obtain a global in time solution to the original equation (1.5.35), we need to do two things: (1) extend the existence of the solutions to the approximate equation (1.5.44) to all $t > 0$, and (2) pass to the limit $\varepsilon \rightarrow 0$ and show that the limit of u^ε exists (possibly along a sub-sequence) and satisfies “the true equation” (1.5.35). The latter step will require uniform bounds on ∇u^ε that do not depend on ε – something much better than what is required in Exercise 1.5.8. The last step will be to prove uniqueness of such global in time smooth solution to (1.5.35) but that is much simpler.

Global in time existence of the approximate solution

To show that the solution to (1.5.44) exists for all $t > 0$, and not just on the interval $[0, T_\varepsilon]$, we use the Duhamel formula

$$u^\varepsilon(t, x) = e^{t\Delta}u_0(x) + \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} f(s, y, \nabla v^\varepsilon(s, y)) dy ds. \quad (1.5.47)$$

Assumption (1.5.38), together with the gradient bound (1.5.46), implies an estimate

$$|f(t, x, \nabla v^\varepsilon(t, x))| \leq C(1 + |\nabla v^\varepsilon(t, x)|) \leq C\left(1 + \frac{\|u^\varepsilon(t, \cdot)\|_{L^\infty}}{\sqrt{\varepsilon}}\right), \quad (1.5.48)$$

that can be used in (1.5.47) to yield a Grownwall inequality

$$\|u^\varepsilon(t, \cdot)\|_{L^\infty} \leq \|u_0\|_{L^\infty} + Ct + \frac{C}{\sqrt{\varepsilon}} \int_0^t \|u^\varepsilon(s, \cdot)\|_{L^\infty} ds. \quad (1.5.49)$$

We used the maximum principle to bound the first term in the right side of (1.5.47), and (1.5.48) together with the standard change of variables (1.4.28):

$$z = \frac{x - y}{2\sqrt{t - s}}, \quad (1.5.50)$$

to estimate the integral in the right side of (1.5.47).

We set

$$Z_\varepsilon(t) = \int_0^t \|u^\varepsilon(s, \cdot)\|_{L^\infty} ds,$$

and write (1.5.49) as

$$\frac{dZ_\varepsilon}{dt} \leq \|u_0\|_{L^\infty} + Ct + \frac{C}{\sqrt{\varepsilon}} Z_\varepsilon. \quad (1.5.51)$$

Multiplying (1.5.51) by $\exp(-Ct/\sqrt{\varepsilon})$ and integrating, keeping in mind that $Z_\varepsilon(0) = 0$, gives

$$Z_\varepsilon(t) \leq \frac{\sqrt{\varepsilon}}{C} e^{Ct/\sqrt{\varepsilon}} (\|u_0\|_{L^\infty} + Ct). \quad (1.5.52)$$

Using this bound in (1.5.49) gives the estimate

$$\|u^\varepsilon(t, \cdot)\|_{L^\infty} \leq (\|u_0\|_{L^\infty} + Ct)(1 + e^{Ct/\sqrt{\varepsilon}}). \quad (1.5.53)$$

Therefore, the L^∞ -norm of the solution can grow by at most a fixed factor over the time interval $[0, T_\varepsilon]$. This estimate, together with the result of Exercise 1.5.8 allows us to restart the Cauchy problem on the time interval $[T_\varepsilon, 2T_\varepsilon]$, and then on $[2T_\varepsilon, 3T_\varepsilon]$, and so on, showing that the regularized problem (1.5.44) admits a global in time solution.

Passing to the limit $\varepsilon \downarrow 0$

A much more serious challenge than proving the global in time existence of u^ε is to send $\varepsilon \downarrow 0$, and recover a smooth solution of the original equation (1.5.35) in the limit. Note that the upper bound (1.5.53) deteriorates very badly as $\varepsilon \downarrow 0$. Hence, we need to come up with much better bounds than that in order to pass to the limit $\varepsilon \downarrow 0$. To do this, we will obtain the Hölder estimates for u^ε and its derivatives up to the second order in space and the first order in time, that will be independent of ε . The Ascoli-Arzelà theorem will then provide us with the compactness of the family u^ε , and allow us to pass to the limit along a subsequence and obtain a solution to (1.5.35).

Exercise 1.5.9 Assume that there exists $\alpha \in (0, 1)$ such that, for all $\delta > 0$ and $T > \delta$, there is $C_\delta(T) > 0$, that is independent of $\varepsilon \in (0, 1)$, for which we have the following Hölder regularity estimates:

$$\left| \frac{\partial}{\partial t} (u^\varepsilon(t, x) - u^\varepsilon(t', x')) \right| + \left| D_x^2 (u^\varepsilon(t, x) - u^\varepsilon(t', x')) \right| \leq C_\delta(T) \left(|t - t'|^{\alpha/2} + |x - x'|^\alpha \right), \quad (1.5.54)$$

for all $t, t' \in [\delta, T]$ and $x, x' \in \mathbb{R}^n$, together with a uniform bound

$$|u^\varepsilon(t, x)| \leq C(T), \text{ for all } 0 \leq t \leq T \text{ and all } x \in \mathbb{R}^n. \quad (1.5.55)$$

Write down a complete proof that then there exists a subsequence $u^{\varepsilon_k}(t, x)$ that converges to a limit $u(t, x)$ as $k \rightarrow +\infty$, and, moreover, that $\nabla v^\varepsilon \rightarrow \nabla u$. In which space does the convergence take place? Show that the limit $u(t, x)$ is twice continuously differentiable in space, and once continuously differentiable in time, and is a solution to (1.5.35). For now, we leave open the question of why the limit satisfies the initial conditions as well.

This exercise gives us the road map to the construction of a solution to (1.5.35): we “only” need to establish the Hölder estimates (1.5.54) for the solutions to the approximate equation (1.5.44). We will use the following lemma, that is a slight generalization of the Gronwall lemma, and which is very useful in estimating the derivatives for the solutions of the parabolic equations.

Lemma 1.5.10 *Let $\varphi(t)$ be a nonnegative bounded function that satisfies, for all $0 \leq t \leq T$:*

$$\varphi(t) \leq \frac{a}{\sqrt{t}} + b \int_0^t \frac{\varphi(s)}{\sqrt{t-s}} ds. \quad (1.5.56)$$

Then, for all $T > 0$, there is $C(T, b) > 0$ that depends on T and b but not on $\|\varphi\|_{L^\infty}$, such that

$$\varphi(t) \leq \frac{C(T, b)a}{\sqrt{t}}, \text{ for all } 0 < t \leq T. \quad (1.5.57)$$

Proof. First, note that we can write $\varphi(t) = a\psi(t)$, leading to

$$\psi(t) \leq \frac{1}{\sqrt{t}} + b \int_0^t \frac{\psi(s)}{\sqrt{t-s}} ds. \quad (1.5.58)$$

Then, iterating (1.5.58) we obtain

$$\psi(t) \leq \sum_{k=0}^n I_k(t) + R_{n+1}(t), \quad (1.5.59)$$

for any $n \geq 0$, with

$$I_{n+1}(t) = b \int_0^t \frac{I_n(s)}{\sqrt{t-s}} ds, \quad I_0(t) = \frac{1}{\sqrt{t}}, \quad (1.5.60)$$

and

$$R_{n+1}(t) = b \int_0^t \frac{R_n(s)}{\sqrt{t-s}} ds, \quad R_0(t) = \psi(t). \quad (1.5.61)$$

We claim that there exist a constant $c > 0$, and $p > 1$ so that

$$I_n(t) \leq \frac{1}{\sqrt{t}} \frac{(ct)^{n/2}}{(n!)^{1/p}}. \quad (1.5.62)$$

Indeed, this bound holds for $n = 0$, and if it holds for $I_n(t)$, then we have

$$\begin{aligned} I_{n+1}(t) &= b \int_0^t \frac{I_n(s)}{\sqrt{t-s}} ds \leq \frac{bc^{n/2}}{(n!)^{1/p}} \int_0^t \frac{s^{(n-1)/2}}{\sqrt{t-s}} ds = \frac{bc^{n/2}t^{(n+1)/2}}{\sqrt{t}(n!)^{1/p}} \int_0^1 \frac{\tau^{(n-1)/2}}{\sqrt{1-\tau}} d\tau \\ &\leq \frac{bc^{n/2}t^{(n+1)/2}}{(n!)^{1/p}\sqrt{t}} \left(\int_0^1 \tau^{3(n-1)/2} d\tau \right)^{1/3} \left(\int_0^1 \frac{d\tau}{(1-\tau)^{3/4}} \right)^{2/3} \\ &= \frac{bc^{n/2}t^{(n+1)/2}}{\sqrt{t}(n!)^{1/p}} \frac{4^{2/3}}{(3n/2 - 1/2)^{1/3}} \leq \frac{bc^{n/2}t^{(n+1)/2}}{\sqrt{t}(n!)^{1/p}} \frac{4}{(n+1)^{1/3}}. \end{aligned} \quad (1.5.63)$$

We used above the Hölder inequality with exponents 3 and 3/2. Thus, the bound (1.5.62) holds with $p = 3$ and $c = 16b^2$.

As we assume that $\varphi(t)$ is bounded, so is $R_0(t) = \psi(t)$. This leads to a better bound for $R_n(t)$ than for $I_n(t)$: we claim that there exist a constant $c > 0$, and $p > 1$ so that

$$R_n(t) \leq \frac{(ct)^{n/2}}{(n!)^{1/p}} \|\psi\|_{L^\infty}. \quad (1.5.64)$$

The computation is very similar to (1.5.63): we know that (1.5.64) holds for $n = 0$, and if it holds for some n , then we have

$$\begin{aligned} R_{n+1}(t) &= b \int_0^t \frac{R_n(s)}{\sqrt{t-s}} ds \leq \frac{bc^{n/2}}{(n!)^{1/p}} \int_0^t \frac{s^{n/2}}{\sqrt{t-s}} ds = \frac{bc^{n/2}t^{(n+1)/2}}{(n!)^{1/p}} \int_0^1 \frac{\tau^{(n-1)/2}}{\sqrt{1-\tau}} d\tau \\ &\leq \frac{bc^{n/2}t^{(n+1)/2}}{(n!)^{1/p}} \left(\int_0^1 \tau^{3(n-1)/2} d\tau \right)^{1/3} \left(\int_0^1 \frac{d\tau}{(1-\tau)^{3/4}} \right)^{2/3} \\ &= \frac{bc^{n/2}t^{(n+1)/2}}{(n!)^{1/p}} \frac{4^{2/3}}{(3n/2 - 1/2)^{1/3}} \leq \frac{bc^{n/2}t^{(n+1)/2}}{(n!)^{1/p}} \frac{4}{(n+1)^{1/3}}. \end{aligned} \quad (1.5.65)$$

Once again, we can take $p = 3$ and $c = 16b^2$. We conclude that

$$R_n(t) \rightarrow 0 \text{ as } n \rightarrow +\infty, \text{ uniformly on } [0, T]. \quad (1.5.66)$$

Going back to (1.5.59), we see that

$$\varphi(t) \leq \frac{a}{\sqrt{t}} + a \sum_{n=1}^{\infty} I_n(t). \quad (1.5.67)$$

Now, the desired estimate (1.5.57) follows from (1.5.67) and (1.5.62). \square

With the claim of Lemma 1.5.10 in hand, let us go back to the Duhamel formula (1.5.47)

$$u^\varepsilon(t, x) = e^{t\Delta} u_0(x) + \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} f(s, y, \nabla v^\varepsilon(s, y)) dy ds. \quad (1.5.68)$$

We first get a Hölder bound on ∇u^ε . The maximum principle implies that

$$\|e^{t\Delta} u_0\|_{L^\infty} \leq \|u_0\|_{L^\infty}, \quad (1.5.69)$$

and also that the gradient

$$\nabla v^\varepsilon = e^{\varepsilon\Delta} \nabla u^\varepsilon,$$

satisfies the bound

$$\|\nabla v^\varepsilon(t, \cdot)\|_{L^\infty} \leq \|\nabla u^\varepsilon(t, \cdot)\|_{L^\infty}. \quad (1.5.70)$$

We use these estimates, together with assumption (1.5.38) on the function $f(t, x, p)$, and the change of variables (1.5.50), in the Duhamel formula (1.5.68), leading to

$$\|u^\varepsilon(t, \cdot)\|_{L^\infty} \leq \|u_0\|_{L^\infty} + Ct + C \int_0^t \|\nabla u^\varepsilon(s, \cdot)\|_{L^\infty} ds. \quad (1.5.71)$$

The next step is to take the gradient of the Duhamel formula:

$$\begin{aligned} \nabla u^\varepsilon(t, x) &= \nabla(e^{t\Delta} u_0(x)) + \nabla \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} f(s, y, \nabla v^\varepsilon(s, y)) dy ds \\ &= \nabla(e^{t\Delta} u_0(x)) + \int_0^t \nabla \left[e^{(t-s)\Delta} f(s, \cdot, \nabla v^\varepsilon(s, \cdot)) \right] (x) ds. \end{aligned} \quad (1.5.72)$$

The first term in the right side is estimated as in (1.5.46):

$$\|\nabla(e^{t\Delta} u_0)\|_{L^\infty} \leq \frac{C}{\sqrt{t}} \|u_0\|_{L^\infty}. \quad (1.5.73)$$

To bound the gradient of the integral term in (1.5.72), we note that (1.5.73), together with assumption (1.5.38) give

$$\|\nabla e^{(t-s)\Delta} f(s, \cdot, \nabla v^\varepsilon(s, \cdot))\|_{L^\infty} \leq \frac{C}{\sqrt{t-s}} \|f(s, \cdot, \nabla v^\varepsilon(s, \cdot))\|_{L^\infty} \leq \frac{C}{\sqrt{t-s}} (1 + \|\nabla v^\varepsilon(s, \cdot)\|_{L^\infty}). \quad (1.5.74)$$

Using (1.5.70) once again and putting together (1.5.72), (1.5.73) and (1.5.74), leads to

$$\|\nabla u^\varepsilon(t, \cdot)\|_{L^\infty} \leq \frac{C}{\sqrt{t}} \|u_0\|_{L^\infty} + C\sqrt{t} + C \int_0^t \frac{\|\nabla u^\varepsilon(s, \cdot)\|_{L^\infty}}{\sqrt{t-s}} ds. \quad (1.5.75)$$

Writing

$$\frac{C}{\sqrt{t}} \|u_0\|_{L^\infty} + C\sqrt{t} \leq \frac{C\|u_0\|_{L^\infty} + CT}{\sqrt{t}}, \quad 0 \leq t \leq T,$$

we can put (1.5.75) into the form of (1.5.56). Lemma 1.5.10 implies then that there exists a constant $C(T) > 0$, independent of ε , such that

$$\|\nabla u^\varepsilon(t, \cdot)\|_{L^\infty} \leq \frac{C(T)}{\sqrt{t}}, \quad 0 < t \leq T. \quad (1.5.76)$$

This bound, which is uniform in $\varepsilon \in (0, 1)$, is absolutely crucial and allows us to proceed relatively effortlessly. Note that even though the right side of (1.5.76) blows up as $t \downarrow 0$, we can not expect any better bound than (1.5.76) as we only assume that the initial condition $u_0(x)$ is continuous and not necessarily differentiable.

The first simple observation is that using the estimate (1.5.76) in (1.5.71) gives a uniform bound on u^ε itself:

$$\|u^\varepsilon(t, \cdot)\|_{L^\infty} \leq C(T), \quad 0 < t \leq T. \quad (1.5.77)$$

This is the uniform bound (1.5.55) in Exercise 1.5.9. In other words, for $t \in (\delta, \infty)$ for any $\delta > 0$, the family $u^\varepsilon(t, \cdot)$ is uniformly bounded in the Sobolev space $W^{1,\infty}(\mathbb{R}^n)$ – the space of L^∞ functions with gradients (in the sense of distributions) that are also L^∞ functions:

$$\|u^\varepsilon(t, \cdot)\|_{W^{1,\infty}} \leq \frac{C(T)}{\sqrt{t}}, \quad 0 < t \leq T. \quad (1.5.78)$$

The constant $C(T)$ depends only on T , the constant C_3 in (1.5.38) and $\|u_0\|_{L^\infty}$.

The uniform bound on the gradient in (1.5.76) seems a far cry from what we need in Exercise 1.5.9 – there, we require a Hölder estimate on the second derivatives in x , and so far we only have a uniform bound on the first derivative. We do not even know yet that the first derivatives are Hölder continuous. Surprisingly, the end of the proof is actually not far off. Take some $1 \leq i \leq n$, and set

$$z_i^\varepsilon(t, x) = \frac{\partial u^\varepsilon(t, x)}{\partial x_i}.$$

Note that such differentiation is perfectly legal since the functions u^ε are smooth. The equation for z_i^ε is (using, as usual, the summation convention for repeated indices)

$$\partial_t z_i^\varepsilon - \Delta z_i^\varepsilon = \partial_{x_i} f(t, x, \nabla v^\varepsilon) + \partial_{p_j} f(t, x, \nabla v^\varepsilon) \partial_{x_j} q_i^\varepsilon, \quad q_i^\varepsilon = e^{\varepsilon \Delta} z_i^\varepsilon. \quad (1.5.79)$$

We look at (1.5.79) as an equation for z_i^ε , with a given function $\nabla v^\varepsilon(t, x)$ that satisfies the already proved uniform bound

$$\|\nabla v^\varepsilon(t, \cdot)\|_{L^\infty} \leq \frac{C(T)}{\sqrt{t}}, \quad 0 < t \leq T, \quad (1.5.80)$$

that follows immediately from (1.5.76). Thus, (1.5.79) is of the form

$$\partial_t z_i^\varepsilon - \Delta z_i^\varepsilon = G(t, x, \nabla q_i^\varepsilon), \quad q_i^\varepsilon = e^{\varepsilon \Delta} z_i^\varepsilon, \quad (1.5.81)$$

with

$$G(t, x, p) = \partial_{x_i} f(t, x, \nabla v^\varepsilon(t, x)) + \partial_{p_j} f(t, x, \nabla v^\varepsilon(t, x)) p_j. \quad (1.5.82)$$

The function $G(t, x, p)$ satisfies the assumptions on the nonlinearity $f(t, x, p)$ stated at the beginning of this section – it is simply a linear function in the variable p , and the gradient bound (1.5.80), together with the smoothness assumptions on $f(t, x, p)$ in (1.5.39), and the Lipschitz estimate (1.5.37), implies that

$$|G(t, x, p)| \leq \frac{C(T)}{\sqrt{t}}(1 + |p|). \quad (1.5.83)$$

Hence, on any time interval $[\delta, T]$ with $\delta > 0$, the function z_i^ε satisfies an equation of the type we have just analyzed for u^ε , and our previous analysis shows that

$$\|\nabla z_i^\varepsilon(t, \cdot)\|_{L^\infty} \leq \frac{C(T, \delta)}{\sqrt{t - \delta}}, \quad \delta < t \leq T. \quad (1.5.84)$$

The constant $C(T, \delta)$ depends on $\delta > 0$ because it depends on $\|z_i^\varepsilon(\delta, \cdot)\|_{L^\infty}$ and because (1.5.83) produces an upper bound on $G(t, x, p)$ for $t > \delta$ that depends on $\delta > 0$. Rephrasing (1.5.84), we have the bound

$$\|D^2 u^\varepsilon(t, \cdot)\|_{L^\infty} \leq \frac{C(T, \delta)}{\sqrt{t - \delta}}, \quad \delta < t \leq T, \quad (1.5.85)$$

with a constant $C(T, \delta)$ that depends on $\delta > 0$, $T > 0$ and $\|u_0\|_{L^\infty}$.

This is almost what we need in (1.5.54) – we also need to show that $D^2 u^\varepsilon$ are Hölder continuous, and deal with the time derivative and Hölder continuity in t . With the information we have already obtained, we know that the right side of (1.5.81) is a uniformly bounded function, on any time interval $[\delta, T]$, with $\delta > 0$. Proposition 1.4.7 implies then immediately that $\nabla z_i(t, x)$ is Hölder continuous in x and $z_i(t, x)$ itself is Hölder continuous in t on the time interval $[2\delta, T]$, with bounds that do not depend on $\varepsilon > 0$. In addition, the uniform bound on $\|D^2 u^\varepsilon\|_{L^\infty}$ and equation (1.5.44) itself imply a uniform bound on the time derivative:

$$\|\partial_t u^\varepsilon(t, \cdot)\|_{L^\infty} \leq \frac{C(T, \delta)}{\sqrt{t - \delta}}, \quad \delta < t \leq T. \quad (1.5.86)$$

To get a Hölder bound on the time derivative

$$\zeta^\varepsilon(t, x) = \frac{\partial u^\varepsilon(t, x)}{\partial t}, \quad (1.5.87)$$

we differentiate (1.5.44) in time to get the following equation

$$\partial_t \zeta^\varepsilon - \Delta \zeta^\varepsilon = \partial_t f(t, x, \nabla v^\varepsilon) + \partial_{p_j} f(t, x, \nabla v^\varepsilon) \partial_{x_j} \eta^\varepsilon, \quad \eta^\varepsilon = e^{\varepsilon \Delta} \zeta^\varepsilon. \quad (1.5.88)$$

This equation has the same form as equation (1.5.81) for $z_i^\varepsilon(t, x)$. In addition, (1.5.86) gives an a priori bound for $\zeta^\varepsilon(\delta, \cdot)$. Hence, arguing as above gives an analog of (1.5.84):

$$\|\nabla \zeta^\varepsilon(t, \cdot)\|_{L^\infty} \leq \frac{C(T, \delta)}{\sqrt{t - \delta}}, \quad \delta < t \leq T. \quad (1.5.89)$$

This allows us to bound $\nabla\eta^\varepsilon$, so that we can view (1.5.88) as an equation of the form

$$\partial_t \zeta^\varepsilon - \Delta \zeta^\varepsilon = F^\varepsilon(t, x), \quad (1.5.90)$$

with a function $F(t, x)$ that satisfies a uniform bound

$$\|F^\varepsilon(t, \cdot)\|_{L^\infty} \leq \frac{C(T, \delta)}{\sqrt{t - \delta}}, \quad \delta < t \leq T. \quad (1.5.91)$$

As we know that $\zeta^\varepsilon(t, x)$ is uniformly bounded by (1.5.86), Proposition 1.4.7 can be used again, this time to deduce that $\zeta^\varepsilon(t, x)$ is Hölder continuous in t on any time interval $[2\delta, T]$, with a bound that does not depend on $\varepsilon > 0$. Thus, (1.5.54) is finally proved.

Now, given any $\delta > 0$, Exercise 1.5.9 allows us to find a sequence $\varepsilon_k \rightarrow 0$ so that u^{ε_k} converges to a limit $u(t, x)$ locally uniformly, and ∇u^{ε_k} converge to $\nabla u(t, x)$ on any time interval $[\delta, T]$. A standard diagonal argument allows us to pass to the limit on $(0, T)$. The limit is also twice continuously differentiable in x and once continuously differentiable in t , and these derivatives themselves are Hölder continuous. Passing to the limit in (1.5.44)

$$u_t^\varepsilon - \Delta u^\varepsilon = f(t, x, \nabla v^\varepsilon), \quad v^\varepsilon = e^{\varepsilon \Delta} u^\varepsilon, \quad (1.5.92)$$

leads to

$$u_t - \Delta u = f(t, x, \nabla u), \quad (1.5.93)$$

as desired. In order to prove that $u(t, x)$ satisfies the initial condition, we go back to the Duhamel formula (1.5.47) to obtain, with the help of (1.5.80):

$$\begin{aligned} |u^\varepsilon(t, x) - e^{t\Delta} u_0(x)| &\leq \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} |f(s, y, \nabla v^\varepsilon(s, y))| dy ds \\ &\leq C(T) \int_0^t \int_{\mathbb{R}^n} \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))} \frac{dy ds}{\sqrt{s}} \leq C(T) \sqrt{t}. \end{aligned} \quad (1.5.94)$$

Passing to the limit $\varepsilon_k \rightarrow 0$ implies that $u(0, x) = u_0(x)$, so that the initial condition is satisfied.

Exercise 1.5.11 Differentiate the equation for u and iterate the above argument, showing that the solution is actually infinitely differentiable.

All that is left in the proof of Theorem 1.5.6 is to prove the uniqueness of a smooth solution. We will invoke the maximum principle again. Recall that we are looking for smooth solutions, so the difference $w = u_1 - u_2$ between any two solutions u_1 and u_2 simply satisfies an equation with a drift:

$$w_t - \Delta w = b(t, x) \cdot \nabla w, \quad (1.5.95)$$

with a smooth drift $b(t, x)$ such that

$$f(x, \nabla u_1(t, x)) - f(x, \nabla u_2(t, x)) = b(t, x) \cdot [\nabla u_1(t, x) - \nabla u_2(t, x)].$$

As $w(0, x) \equiv 0$, the comparison principle of Theorem 1.3.3 implies that $w(t, x) \equiv 0$ and $u_1 \equiv u_2$. This completes the proof of Theorem 1.5.6. \square

Exercise 1.5.12 Prove that, if u_0 is smooth, then smoothness holds up to $t = 0$. Prove that equation (1.5.35) holds up to $t = 0$, that is:

$$u_t(0, x) = \Delta u_0(x) + f(x, \nabla u_0(x)).$$

1.5.3 Applications to linear equations with a drift

Let us now discuss how the above results can be made more quantitative for linear equations of the form

$$u_t = \Delta u + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u, \quad t > 0, \quad x \in \mathbb{R}^n, \quad (1.5.96)$$

with smooth coefficients $b_j(t, x)$ and $c(t, x)$. We recall, once again, that the repeated indices are summed. When $c(t, x) = 0$, this equation has the form of the non-linear equation (1.5.35) considered in the previous section. In particular, Theorem 1.5.6 implies immediately that given any initial condition $u_0(x)$ that is a bounded continuous function, the equation (1.5.96) has a unique solution $u(t, x)$ that is infinitely differentiable for all $t > 0$ and $x \in \mathbb{R}^n$ such that $u(0, x) = u_0(x)$. The same result holds, with essentially an identical proof when $c(t, x)$ is smooth.

Exercise 1.5.13 Extend the result of Theorem 1.5.6 to equations of the form (1.5.96) with smooth coefficients $b_j(t, x)$ and $c(t, x)$.

A more important claim is that the quantitative regularity results formulated in Proposition 1.4.3 for the linear heat equation in the whole space also hold essentially verbatim for (1.5.96).

1.6 A survival kit in the jungle of regularity

In our noble endeavor to carry out as explicit computations as possible, we have not touched the question of regularity of solutions to inhomogeneous equations where the diffusivity can be not constant. An inhomogeneous drift has been treated in Section 1.5.2. We address here the following question: given a linear inhomogeneous equation of the form

$$u_t - a_{ij}(t, x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u = f(t, x), \quad (1.6.1)$$

the coefficients a_{ij} , b_j , c and the right side f having a certain given degree of smoothness, what is the best regularity that one may expect from u ? The question is a little different from what we did for the nonlinear equations, where one would first prove a certain, possibly small, amount of regularity, in the hope that this would be sufficient for a bootstrap argument leading to a much better regularity than in one iteration step. The answer to the question of maximal regularity is, in a nutshell: if the coefficients have a little bit of continuity, such as the Hölder continuity, then the derivatives u_t and D^2u have the same regularity as f . This, however, is true up to some painful exceptions: continuity for f will not entail, in general, the continuity of u_t and D^2u . This is exactly what we have seen in Section 1.4 for the forced heat equation, so this should not come as a surprise to the reader.

The question of the maximal regularity for linear parabolic equations has a certain degree of maturity, an interested reader should consult [?] to admire the breadth, beauty and technical complexity of the available results. Our goal here is much more modest: we will explain why the Hölder continuity of f will entail the Hölder continuity of u_t and D^2u – the result we have already seen for the heat equation using the explicit computations with the Duhamel term.

When $a_{ij}(t, x) = \delta_{ij}$ (the Kronecker symbol), the second order term in (1.6.1) is the Laplacian, and our work was already almost done in Theorem 1.5.6 even though we have not formulated the precise Hölder estimates on the solution in the case when equation (1.5.35) happens to be linear. Nevertheless, the reader should be able to extract them from the proof of that theorem and discover a version of Proposition 1.4.20 for an equation of the form

$$u_t - \Delta u + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u = f(t, x), \quad (1.6.2)$$

with smooth coefficients $b_j(t, x)$, $j = 1, \dots, n$ and $c(t, x)$, and $f(t, x) \in C_t^{\alpha/2} C_x^\alpha$. We will try to convince the reader, without giving the full details of all the proofs, that this carries over to variable diffusion coefficients, and, importantly, to problems with boundary conditions. Our main message here is that all the ideas necessary for the various proofs have already been displayed, and that "only" technical complexity and dexterity are involved. Our discussion follows Chapter 4 of [?], which presents various results with much more details. Let us emphasize again that in this section, we will only give a sketch of the proofs, and sometimes we will not state the results in a formal way.

When the diffusion coefficients are not continuous, but merely bounded, the methods described in this chapter break down. Chapter ??, based on the Nash inequality, explains to some extent how to deal with such problems by a very different approach.

The Cauchy problem for the inhomogeneous coefficients

We have all the ideas to understand the first big piece of this section, the Cauchy problem for the parabolic equations with variable coefficients in the whole space, without any forcing:

$$\begin{aligned} u_t - a_{ij}(t, x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u &= 0, \quad t > 0, \quad x \in \mathbb{R}^n, \\ u(0, x) &= u_0(x), \quad x \in \mathbb{R}^n. \end{aligned} \quad (1.6.3)$$

We make the following assumptions on the coefficients: first, they are sufficiently regular – the functions $(a_{ij}(t, x))_{1 \leq i, j \leq N}$, $(b_j(t, x))_{1 \leq j \leq N}$ and $c(t, x)$, all α -Hölder continuous over $[0, T] \times \mathbb{R}^n$. Second, we make the ellipticity assumption, generalizing (1.4.82): there exist $\lambda > 0$ and $\Lambda > 0$ so that for any vector $\xi \in \mathbb{R}^n$ and any $x \in \mathbb{R}^n$ we have

$$\lambda |\xi|^2 \leq a_{ij}(t, x) \xi_i \xi_j \leq \Lambda |\xi|^2. \quad (1.6.4)$$

We assume that the initial condition $u_0(x)$ is a continuous function – this assumption can be very much weakened but we do not focus on it right now.

Theorem 1.6.1 *The Cauchy problem (1.6.3) has a unique solution $u(t, x)$, whose Hölder norm on the sets of the form $[\varepsilon, T] \times \mathbb{R}^n$ is controlled by the L^∞ norm of u_0 .*

The statement of this theorem is deliberately vague – the correct statement should become clear to the reader after we outline the ideas of the proof.

Exercise 1.6.2 Show that the uniqueness of the solution is an immediate consequence of the maximum principle.

Thus, the main issue is the construction of a solution with the desired regularity. The idea is to construct the fundamental solution of (1.6.3), that is, the solution $E(t, s, x, y)$ of (1.6.3) on the time interval $s \leq t \leq T$, instead of $0 \leq t \leq T$:

$$\partial_t E - a_{ij}(t, x) \frac{\partial E}{\partial x_i} x_j + b_j(t, x) \frac{\partial E}{\partial x_j} + c(t, x)E = 0, \quad t > s, \quad x \in \mathbb{R}^n, \quad (1.6.5)$$

with the initial condition

$$E(t = s, s, x, y) = \delta(x - y), \quad (1.6.6)$$

the Dirac mass at $x = y$. The solution of (1.6.3) can then be written as

$$u(t, x) = \int_{\mathbb{R}^n} E(t, 0, x, y) u_0(y) dy. \quad (1.6.7)$$

If can show that $E(t, s, x, y)$ is smooth enough (at least away from $t = s$), $u(t, x)$ will satisfy the desired estimates as well – they can be obtained by differentiating or taking partial differences in (1.6.7). Note that regularity of E for $t > s$ is a very strong property: the initial condition in (1.6.6) at $t = s$ is a measure – and we need to show that for all $t > s$ the solution is actually a smooth function. On the other hand, this is exactly what happens for the heat equation

$$u_t = \Delta u,$$

where the fundamental solution is

$$E(t, s, x, y) = \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))},$$

and is smooth for all $t > s$.

Exercise 1.6.3 Go back to the equation

$$u_t - u_{xx} + u_y = 0.$$

considered in Exercise 1.4.22. Show that its fundamental solution is not a smooth function in the y -variable. Thus, the ellipticity condition is important for this property.

The understanding of the regularity of the solutions of the Cauchy problem is also a key to the inhomogeneous problem because of the Duhamel principle.

Exercise 1.6.4 Let $f(t, x)$ be a Hölder-continuous function over $[0, T] \times \mathbb{R}^n$. Use the Duhamel principle to write down the solution of

$$\begin{aligned} u_t - a_{ij}(t, x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u &= f(t, x), \quad t > 0, \quad x \in \mathbb{R}^n, \\ u(0, x) &= u_0(x), \quad x \in \mathbb{R}^n, \end{aligned} \quad (1.6.8)$$

in terms of $E(t, s, x, y)$.

Thus, everything boils down to constructing the fundamental solution $E(t, s, x, y)$, and a way to do it is via the parametrix method. Let us set $b_j = c = 0$ – this does not affect the essence of the arguments but simplifies the notation. The philosophy is that the possible singularities of $E(t, s, x, y)$ are localized at $t = s$ and $x = y$ (as for the heat equation). Therefore, in order to capture the singularities of $E(t, s, x, y)$ we may try to simply freeze the coefficients in the equation at $t = s$ and $x = y$, and compare $E(t, s, x, y)$ to the fundamental solution $E_0(s, t, x, y)$ of the resulting equation:

$$\begin{aligned}\partial_t E_0 - a_{ij}(s, y) \frac{\partial^2 E_0}{\partial x_i \partial x_j} &= 0, \quad t > s, \quad x \in \mathbb{R}^n, \\ E_0(t = s, x) &= \delta(x - y), \quad x \in \mathbb{R}^n.\end{aligned}\tag{1.6.9}$$

There is no reason to expect the two fundamental solutions to be close – they satisfy different equations. Rather, the expectation is that E will be a smooth perturbation of E_0 – and, since E_0 solves an equation with constant coefficients (remember that s and y are fixed here), we may compute it exactly.

To this end, let us write the equation for $E(t, s, x, y)$ as

$$\begin{aligned}\partial_t E - a_{ij}(s, y) \frac{\partial^2 E}{\partial x_i \partial x_j} &= F(t, x), \quad t > s, \quad x \in \mathbb{R}^n, \\ E(t = s, x) &= \delta(x - y), \quad x \in \mathbb{R}^n,\end{aligned}\tag{1.6.10}$$

with the right side

$$F(t, x, y) = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 E}{\partial x_i \partial x_j}.\tag{1.6.11}$$

The difference

$$R_0 = E - E_0$$

satisfies

$$\partial_t R_0 - a_{ij}(s, y) \frac{\partial^2 R_0}{\partial x_i \partial x_j} = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 E_0}{\partial x_i \partial x_j} + F_0(t, x), \quad t > s,\tag{1.6.12}$$

with the initial condition $R_0(t = s, x) = 0$, and

$$F_0(t, x) = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 R_0}{\partial x_i \partial x_j}.\tag{1.6.13}$$

Let us further decompose

$$R_0 = E_1 + R_1.$$

Here, E_1 is the solution of

$$\partial_t E_1 - a_{ij}(s, y) \frac{\partial^2 E_1}{\partial x_i \partial x_j} = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 E_0}{\partial x_i \partial x_j}, \quad t > s,\tag{1.6.14}$$

with the initial condition $E_1(t = s, x) = 0$. The remainder R_1 solves

$$\partial_t R_1 - a_{ij}(s, y) \frac{\partial^2 R_1}{\partial x_i \partial x_j} = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 E_1}{\partial x_i \partial x_j} + F_1(t, x), \quad t > s,\tag{1.6.15}$$

with $R_1(t = s, x) = 0$, and

$$F_1(t, x) = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 R_1}{\partial x_i \partial x_j}. \quad (1.6.16)$$

Equation (1.6.14) for E_1 is a forced parabolic equation with constant coefficients – as we have seen, its solutions behave exactly like those of the standard heat equation with a forcing, except for rotations and dilations. We may assume without loss of generality that $a_{ij}(s, y) = \delta_{ij}$, so that the reference fundamental solution is

$$E_0(t, s, x, y) = \frac{1}{(4\pi(t-s))^{n/2}} e^{-(x-y)^2/(4(t-s))}, \quad (1.6.17)$$

and (1.6.14) is simply a forced heat equation:

$$\partial_t E_1 - \Delta E_1 = [a_{ij}(t, x) - \delta_{ij}] \frac{\partial^2 E_0(t, s, x, y)}{\partial x_i \partial x_j}, \quad t > s, \quad x \in \mathbb{R}^n. \quad (1.6.18)$$

The functions $a_{ij}(t, x)$ Hölder continuous, with $a_{ij}(s, y) = \delta_{ij}$. The regularity of E_1 can be approached by the tools of the previous sections – after all, (1.6.14) is just another forced heat equation! The next exercise may be useful for understanding what is going on.

Exercise 1.6.5 Consider, instead of (1.6.14) the solution of

$$\partial_t z - \Delta z = \frac{\partial^2 E_0(t, s, x, y)}{\partial x_i \partial x_j}, \quad t > s, \quad x \in \mathbb{R}^n, \quad (1.6.19)$$

with $z(t = s, x) = 0$. How does its regularity compare to that of E_0 ? Now, what can you say about the regularity of the solution to (1.6.18), how does the factor $[a_{ij}(t, x) - \delta_{ij}]$ help to make E_1 more regular than z ? In which sense is E_1 more regular than E_0 ?

With this understanding in hand, one may consider the iterative process: write

$$R_1 = E_2 + R_2,$$

with E_2 the solution of

$$\partial_t E_2 - a_{ij}(s, y) \frac{\partial^2 E_2}{\partial x_i \partial x_j} = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 E_1}{\partial x_i \partial x_j}, \quad t > s, \quad (1.6.20)$$

with $E_2(t = s, x) = 0$, and R_2 the solution of

$$\partial_t R_2 - a_{ij}(s, y) \frac{\partial^2 R_2}{\partial x_i \partial x_j} = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 E_2}{\partial x_i \partial x_j} + F_2(t, x), \quad t > s, \quad (1.6.21)$$

with $R_2(t = s, x) = 0$, and

$$F_2(t, x) = (a_{ij}(t, x) - a_{ij}(s, y)) \frac{\partial^2 R_2}{\partial x_i \partial x_j}. \quad (1.6.22)$$

Continuing this process, we have a representation for $E(t, s, x, y)$ as

$$E = E_0 + E_1 + \cdots + E_n + R_n, \tag{1.6.23}$$

with each next term E_j more regular than E_0, \dots, E_{j-1} . Regularity of all E_j can be inferred as in Exercise 1.6.5. One needs, of course, also to estimate the remainder R_n to obtain a "true theorem" but we leave this out of this chapter, to keep the presentation short. An interested reader should consult the aforementioned references for full details. We do, however, offer the reader another (non trivial) exercise.

Exercise 1.6.6 Prove that $E(s, t, x, y)$ obeys Gaussian estimates of the form:

$$m \frac{e^{-|x-y|^2/Dt}}{(t-s)^{n/2}} \leq E(s, t, x, y) \leq M \frac{e^{-|x-y|^2/dt}}{(t-s)^{n/2}},$$

for all $0 < s < t, T$ and $x, y \in \mathbb{R}^n$. The constants m and M , unfortunately, depend very much on T ; however the constants d and D do not.

Interior regularity

So far, we have considered parabolic problems in the whole space \mathbb{R}^n , without any boundaries. One of the miracles of the second order diffusion equations is that the regularity properties are *local*. That is, the regularity of the solutions in a given region only depends on how regular the coefficients are in a slightly larger region. Consider, again, the inhomogeneous parabolic equation

$$u_t - a_{ij}(t, x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u = f(t, x), \quad t > 0, \tag{1.6.24}$$

and assume that the coefficients $a_{ij}(t, x)$, $b_j(t, x)$ and $c(t, x)$, and forcing $f(t, x)$, are α -Hölder in $S = [0, T] \times B_R(x_0)$. It turns out that the derivatives $D^2 u(t, x)$ and $\partial_t u(t, x)$ are α -Hölder in a smaller set of the form $S = [\varepsilon, T] \times B_{(1-\varepsilon)R}(x_0)$, for any $\varepsilon > 0$. The most important point is that the Hölder norm of u in S is controlled only by ε , R , and the Hölder norms of the coefficients and the L^∞ bound of u , both inside the larger set S . Note that the Hölder estimates on u in terms of the L^∞ -norm of u over S do not hold in the original set S , we need a small margin, going down to the smaller set S_ε . This is very similar to what happens for the heat equation: the bounded solution to

$$u_t = \Delta u, \quad t > 0, \quad x \in \mathbb{R}^n, \tag{1.6.25}$$

with an initial condition $u(0, x) = \nabla u$ satisfies a bound

$$\|\nabla u(t, \cdot)\|_{L^\infty} \leq \frac{C}{\sqrt{t}} \|u_0\|_{L^\infty}, \tag{1.6.26}$$

that gives information only for $t > 0$ – this is exactly the margin we have discussed above.

Exercise 1.6.7 Prove this fact. One standard way to do it is to pick a nonnegative and smooth function $\gamma(x)$, equal to 1 in $B_{R/2}(x_0)$ and 0 outside of $B_R(x)$, and to write down an equation for $v(t, x) = \gamma(x)u(t, x)$. Note that this equation is now posed on $(0, T] \times \mathbb{R}^n$, and that the whole space theory can be applied. The computations should be, at times cumbersome. If in doubt, consult [?]. Looking ahead, we will use this strategy in the proof of Proposition 1.7.10 in Section 1.7 below, so the reader may find it helpful to read this proof now.

Regularity up to the boundary

Specifying the Dirichlet boundary conditions allows to get rid of this small margin, and this is the last issue that we are going to discuss in this section. Let us consider equation (1.6.24), posed this time in $(0, T] \times \Omega$, where Ω is bounded smooth open subset of \mathbb{R}^n . As a side remark, it is not crucial that Ω be bounded. However, if Ω is unbounded, we should ask its boundary to oscillate not too much at infinity. Let us supplement (1.6.24) by an initial condition $u(0, x) = u_0(x)$ in Ω , with a continuous function u_0 , and the Dirichlet boundary condition

$$u(t, x) = 0 \text{ for } 0 \leq t \leq T \text{ and } x \in \partial\Omega. \quad (1.6.27)$$

Theorem 1.6.8 *Assume $a_{ij}(t, x)$, $b_j(t, x)$, $c(t, x)$, and $f(t, x)$ are α -Hölder in $(0, T] \times \bar{\Omega}$ – note that, here, we do need the closure of Ω . The above initial-boundary value problem has a unique solution $u(t, x)$ such that $D^2u(t, x)$ and $\partial_t u(t, x)$ are α -Hölder in $[\varepsilon, T] \times \bar{\Omega}$, with their Hölder norms controlled by the L^∞ bound of u_0 , and the Hölder norms of the coefficients and f .*

The way to prove this result parallels the way we followed to establish Theorem 1.6.1. First, we write down an explicit solution on a model situation. Then, we prove the regularity in the presence of a Hölder forcing in the model problem. Once this is done, we turn to general constant coefficients. Then, we do the parametrix method on the model situation. Finally, we localize the problem and reduce it to the model situation.

Let us be more explicit. The model situation is the heat equation in a half space

$$\Omega_n = \mathbb{R}_+^n := \{x = (x_1, \dots, x_n) \in \mathbb{R}^n : x_n > 0\}.$$

Setting $x' = (x_1, \dots, x_{n-1})$, we easily obtain the solution of the initial boundary value problem

$$\begin{aligned} u_t - \Delta u &= 0, & t > 0, & x \in \Omega_n, \\ u(t, x', 0) &= 0, \\ u(0, x) &= u_0(x), \end{aligned} \quad (1.6.28)$$

as

$$u(t, x) = \int_{\mathbb{R}^n} E_0(t, x, y) u_0(y) dy, \quad (1.6.29)$$

with the fundamental solution

$$E_0(t, x, y) = \frac{e^{-(x'-y')^2/4t}}{(4\pi t)^{n/2}} \left(e^{-(x_n-y_n)^2/4t} - e^{-(x_n+y_n)^2/4t} \right). \quad (1.6.30)$$

Let us now generalize step by step: for an equation with a constant drift

$$u_t - \Delta u + b_j \partial_{x_j} u = 0, \quad t > 0, \quad x \in \Omega_n, \quad (1.6.31)$$

the change of unknowns $u(t, x) = e^{x_n b_n / 2} v(t, x)$ transforms the equation into

$$v_t - \Delta v + b_j \partial_{x_j} v - \frac{b_n^2}{4} v = 0, \quad t > 0, \quad x \in \Omega_n. \quad (1.6.32)$$

Thus, the fundamental solution, for (1.6.31) is

$$E(t, x, y) = e^{t b_n^2 / 4 - x b_n / 2} E_0(t, x - t B', y), \quad B' = (b_1, \dots, b_{n-1}, 0). \quad (1.6.33)$$

For an equation of the form

$$u_t - a_{ij} \partial_{x_i} \partial_{x_j} u = 0, \quad t > 0, \quad x \in \Omega_n, \quad (1.6.34)$$

with a constant positive-definite diffusivity matrix a_{ij} , we use the fact that the function

$$u(t, x) = v(t, \sqrt{A^{-1}} x),$$

with $v(t, x)$ a solution of the heat equation

$$v_t = \Delta v,$$

solves (1.6.34). For an equation mixing the two sets of coefficients, one only has to compose the transformations. At that point, one can, with a nontrivial amount of computations, prove the desired regularity for the solutions of

$$u_t - a_{ij} \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j \frac{\partial u}{\partial x_j} + cu = f(t, x) \quad (1.6.35)$$

with constant coefficients, and the Dirichlet boundary conditions on $\partial\Omega_n$. Then, one can use the parametrix method to obtain the result for general inhomogeneous coefficients. This is how one proves Theorem 1.6.8 for $\Omega_n = \mathbb{R}_+^n$.

How does one pass to a general Ω ? Unfortunately, the work is not at all finished yet. One still has to prove a local version of the already proved theorem in Ω_n , in the spirit of the local regularity in \mathbb{R}^n , up to the fact that we must not avoid the boundary. Once this is done, consider a general Ω . Cover its boundary $\partial\Omega$ with balls such that, in each of them, $\partial\Omega$ is a graph in a suitable coordinate system. By using this new coordinate system, one retrieves an equation of the form (1.6.8), and one has to prove that the diffusion coefficients satisfy a coercivity inequality. At this point, maximal regularity for the Dirichlet problem is proved.

Of course, all kinds of local versions (that is, versions of Theorem 1.6.8 where the coefficients are α -Hölder only in a part of $\bar{\Omega}$) are available. Also, most of the above material is valid for the Neumann boundary conditions

$$\partial_\nu u = 0 \text{ on } \partial\Omega,$$

or Robin boundary conditions

$$\partial_\nu u + \gamma(t, x)u = 0 \text{ on } \partial\Omega.$$

We encourage the reader who might still be interested in the subject to try to produce the full proofs, with an occasional help from the books we have mentioned.

The Harnack inequalities

We will only touch here on the Harnack inequalities, a very deep and involved topic of parabolic equations. In a nutshell, the Harnack inequalities allow to control the infimum of a positive solution of a parabolic equation by a universal fraction of its maximum, modulo a time shift. They provide one possible, and very beautiful, path to prove regularity, but we will ignore this aspect here. They are also mostly responsible for the behaviors that are very specific to the diffusion equations, as will be seen in the next section.

We are going to prove what is, in a sense, a "poor man's" version. It is not as scale invariant as one would wish, and uses the regularity theory instead of proving it. It is, however, suited to what we wish to do, and already gives a good account of what is going on. Consider our favorite equation

$$u_t - \Delta u + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u = 0, \quad (1.6.36)$$

with smooth coefficients b_j and c , posed for $t \in (0, T)$, and $x \in B_{R+1}(0)$. We stress that the variable smooth diffusion coefficients could be put in the picture.

Theorem 1.6.9 *Let $u(t, x) \geq 0$ be a non-negative bounded solution of (1.6.36) for $0 \leq t \leq T$ and $x \in B_{R+1}(0)$, and assume that for all $t \in [0, T]$:*

$$\sup_{|x| \leq R+1} u(t, x) \leq k_2, \quad \sup_{|x| \leq R} u(t, x) \geq k_1. \quad (1.6.37)$$

There is a constant $h_R > 0$, that does not depend on T , but that depends on k_1 and k_2 , such that, for all $t \in [1, T]$:

$$h_R \leq \inf_{|x| \leq R} u(t, x). \quad (1.6.38)$$

Proof. The proof is by contradiction. Assume that there exists a sequence u_n of the solutions of (1.6.36) satisfying (1.6.37), and $t_n \in [1, T]$, and $x_n \in B_R(0)$, such that

$$\lim_{n \rightarrow +\infty} u_n(t_n, x_n) = 0. \quad (1.6.39)$$

Up to a possible extraction of a subsequence, we may assume that

$$t_n \rightarrow t_\infty \in [1, T] \text{ and } x_n \rightarrow x_\infty \in B_R(0).$$

The Hölder estimates on u_n and its derivatives in Theorem 1.6.8 together with the Ascoli-Arzelà theorem, imply that the sequence u_n is relatively compact in $C^2([t_\infty - 1/2, T] \times B_{R+1/2}(0))$. Hence, again, after a possible extraction of a subsequence, we may assume that u_n converges to $u_\infty \in C^2([t_\infty - 1/2, T] \times B_{R+1/2}(0))$, together with its first two derivatives in x and the first derivatives in t . Thus, the limit $u_\infty(t, x)$ satisfies (1.6.36) for $t_\infty - 1/2 \leq t \leq T$, and $x \in B_{R+1/2}(0)$, and is non-negative. It also satisfies the bounds in (1.6.37), hence it is not identically equal to zero. Moreover it satisfies $u_\infty(t_\infty, x_\infty) = 0$. This contradicts the strong maximum principle. \square

1.7 The principal eigenvalue for elliptic operators and the Krein-Rutman theorem

One consequence of the strong maximum principle is the existence of a positive eigenfunction for an elliptic operator in a bounded domain with the Dirichlet or Neumann boundary conditions. Such eigenfunction necessarily corresponds to the eigenvalue with the smallest real part. A slightly different way to put it is that the strong maximum principle makes the Krein-Rutman Theorem applicable, which in turn, implies the existence of such eigenfunction. In this section, we will prove this theorem in the context of parabolic operators with time periodic coefficients. We then deduce, in an easy way, some standard properties of the principal elliptic eigenvalue.

1.7.1 The periodic principal eigenvalue

The maximum principle for elliptic and parabolic problems has a beautiful connection to the eigenvalue problems, which also allows to extend it to operators with a zero-order term. We will first consider the periodic eigenvalue problems, that is, elliptic equations where the coefficients are 1-periodic in every direction in \mathbb{R}^n , and the sought for solutions are all 1-periodic in \mathbb{R}^n . It would, of course, be easy to deduce, by dilating the coordinates, the same results for coefficients with general periods T_1, \dots, T_n in the directions e_1, \dots, e_n . We will consider operators of the form

$$Lu(x) = -\Delta u + b_j(x) \frac{\partial u}{\partial x_j} + c(x)u, \quad (1.7.1)$$

with bounded, smooth and 1-periodic coefficients $b_j(x)$ and $c(x)$. We could also consider more general operators of the form

$$Lu(x) = -a_{ij}(x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(x) \frac{\partial u}{\partial x_j} + c(x)u,$$

with uniformly elliptic (and 1-periodic) , and regular coefficients a_{ij} , with the help of the elliptic regularity theory. This will not, however, be needed for our purposes. In order to avoid repeating that the coefficients and the solutions are 1-periodic, we will just say that $x \in \mathbb{T}^n$, the n -dimensional unit torus.

The key spectral property of the operator L comes from the comparison principle. To this end, let us recall the Krein-Rutman theorem. It says that if M is a compact operator in a strongly ordered Banach space X (that is, there is a solid cone K which serves for defining an order relation: $u \leq v$ iff $v - u \in K$), that preserves K : $Mu \in K$ for all $u \in K$, and maps the boundary of K into its interior, then M has an eigenfunction ϕ that lies in this cone:

$$M\phi = \lambda\phi. \quad (1.7.2)$$

Moreover, the corresponding eigenvalue λ has the largest real part of all eigenvalues of the operator M . The classical reference [?] has a nice and clear presentation of this theorem but one can find it in other textbooks, as well.

How can we apply this theorem to the elliptic operators? The operator L given by (1.7.1) is not compact, nor does it preserve any interesting cone. However, let us assume momentarily that $c(x)$ is continuous and $c(x) > 0$ for all $x \in \mathbb{T}^n$. Then the problem

$$Lu = f, \quad x \in \mathbb{T}^n \tag{1.7.3}$$

has a unique solution, and, in addition, if $f(x) \geq 0$ and $f \not\equiv 0$, then $u(x) > 0$ for all $x \in \mathbb{T}^n$. Indeed, let $v(t, x)$ be the solution of the initial value problem

$$v_t + Lv = 0, \quad t > 0, \quad x \in \mathbb{T}^n, \tag{1.7.4}$$

with $v(0, x) = f(x)$. The comparison principle implies a uniform upper bound

$$|v(t, x)| \leq e^{-\bar{c}t} \|f\|_{L^\infty}, \tag{1.7.5}$$

with

$$\bar{c} = \inf_{x \in \mathbb{T}^n} c(x) > 0. \tag{1.7.6}$$

This allows us to define

$$u(x) = \int_0^\infty v(t, x) dt. \tag{1.7.7}$$

Exercise 1.7.1 Verify that if $c(x) > 0$ for all $x \in \mathbb{T}^n$, then $u(x)$ given by (1.7.7) is a solution to (1.7.3). Use the maximum principle to show that (1.7.3) has a unique solution.

This means that we may define the inverse operator $M = L^{-1}$. This operator preserves the cone of the positive functions, and maps its boundary (non-negative functions that vanish somewhere in Ω) into its interior – this is a consequence of the strong maximum principle that holds if $c(x) > 0$. In addition, M is a compact operator from $C(\mathbb{T}^n)$ to itself. Hence, the inverse operator satisfies the assumptions of the Krein-Rutman theorem.

Exercise 1.7.2 Compactness of the inverse M follows from the elliptic regularity estimates. One way to convince yourself of this fact is to consult Evans [?]. Another way is to go back to Theorem 1.5.6, use it to obtain the Hölder regularity estimates on $v(t, x)$, and translate them in terms of $u(x)$ to show that, if f is continuous, then ∇u is α -Hölder continuous, for all $\alpha \in (0, 1)$. The Arzela-Ascoli theorem implies then compactness of M . Hint: be careful about the regularity of $v(t, x)$ as $t \downarrow 0$.

Thus, there exists a positive function f and $\mu \in \mathbb{R}$ so that the function $u = \mu f$ satisfies (1.7.3). Positivity of f implies that the solution of (1.7.3) is also positive, hence $\mu > 0$. As μ is the eigenvalue of L^{-1} with the largest real part, $\lambda = \mu^{-1}$ is the eigenvalue of L with the smallest real part. In particular, it follows that all eigenvalues λ_k of the operator L have a positive real part.

If the assumption $c(x) \geq 0$ does not hold, we may take $K > \|c\|_{L^\infty}$, and consider the operator

$$L'u = Lu + Ku.$$

The zero-order coefficient of L' is

$$c'(x) = c(x) + K \geq 0.$$

Hence, we may apply the previous argument to the operator L' and conclude that L' has an eigenvalue μ_1 that corresponds to a positive eigenfunction, and has the smallest real part among all eigenvalues of L' . The same is true for the operator L , with the eigenvalue

$$\lambda_1 = \mu_1 - K.$$

We say that λ_1 is the principal periodic eigenvalue of the operator L .

1.7.2 The Krein-Rutman theorem: the periodic parabolic case

As promised, we will prove the Krein-Rutman Theorem in the context of the periodic eigenvalue problems. Our starting point will be a slightly more general problem with time-periodic coefficients:

$$u_t - \Delta u + b_j(t, x) \frac{\partial u}{\partial x_j} + c(t, x)u = 0, \quad x \in \mathbb{T}^n. \quad (1.7.8)$$

Here, the coefficients $b_j(t, x)$ and $c(t, x)$ are smooth, 1-periodic in x and T -periodic in t . Let $u(t, x)$ be the solution of the Cauchy problem for (1.7.8), with a 1-periodic, continuous initial condition

$$u(t, x) = u_0(x). \quad (1.7.9)$$

We define the "time T " operator S_T as

$$[S_T u_0](x) = u(T, x). \quad (1.7.10)$$

Exercise 1.7.3 Use the results of Section 1.5 to show that S_T is compact operator on $C(\mathbb{T}^n)$ that preserves the cone of positive functions.

We are going to prove the Krein-Rutman Theorem for S_T first.

Theorem 1.7.4 *The operator S_T has an eigenvalue $\bar{\mu} > 0$ that corresponds to a positive eigenfunction $\phi_1(x) > 0$. The eigenvalue $\bar{\mu}$ is simple: the only solutions of*

$$(S_T - \bar{\mu})u = 0, \quad x \in \mathbb{T}^n$$

are multiples of ϕ_1 . If μ is another (possibly non-real) eigenvalue of S_T , then $|\mu| < \bar{\mu}$.

Proof. Let us pick any positive function $\phi_0 \in C(\mathbb{T}^n)$, set $\psi_0 = \phi_0 / \|\phi_0\|_{L^\infty}$, and consider the iterative sequence (ϕ_n, ψ_n) :

$$\phi_{n+1} = S_T \psi_n, \quad \psi_{n+1} = \frac{\phi_{n+1}}{\|\phi_{n+1}\|_{L^\infty}}.$$

Note that, because ϕ_0 is positive, both ϕ_n and ψ_n are positive for all n , by the strong maximum principle. For every n , let μ_n be the smallest μ such that

$$\phi_{n+1}(x) \leq \mu\psi_n(x), \quad \text{for all } x \in \mathbb{T}^n. \quad (1.7.11)$$

Note that (1.7.11) holds for large μ , because each of the ϕ_n is positive, hence the smallest such μ exists. It is also clear that $\mu_n \geq 0$. We claim that the sequence μ_n is non-increasing. To see that, we apply the operator S_T to both sides of the inequality (1.7.11) with $\mu = \mu_n$, written as

$$S_T\psi_n(x) \leq \mu_n\psi_n(x), \quad \text{for all } x \in \mathbb{T}^n. \quad (1.7.12)$$

and use the fact that S_T preserves positivity, to get

$$(S_T \circ S_T)\psi_n(x) \leq \mu_n S_T\psi_n(x), \quad \text{for all } x \in \mathbb{T}^n, \quad (1.7.13)$$

which is

$$S_T\phi_{n+1}(x) \leq \mu_n\phi_{n+1}(x), \quad \text{for all } x \in \mathbb{T}^n. \quad (1.7.14)$$

Dividing both sides by $\|\phi_{n+1}\|_{L^\infty}$. we see that

$$S_T\psi_{n+1}(x) \leq \mu_n\psi_{n+1}(x), \quad \text{for all } x \in \mathbb{T}^n, \quad (1.7.15)$$

hence

$$\phi_{n+2}(x) \leq \mu_n\psi_{n+1}(x), \quad \text{for all } x \in \mathbb{T}^n. \quad (1.7.16)$$

It follows that $\mu_{n+1} \leq \mu_n$.

Thus, μ_n converges to a limit $\bar{\mu}$.

Exercise 1.7.5 Show that, up to an extraction of a subsequence, the sequence ψ_n converges to a limit ψ_∞ , with $\|\psi_\infty\|_{L^\infty} = 1$.

The corresponding subsequence ϕ_{n_k} converges to $\phi_\infty = S_T\psi_\infty$, by the continuity of S_T . And we have, by (1.7.11):

$$S_T\psi_\infty \leq \bar{\mu}\psi_\infty. \quad (1.7.17)$$

If we have the equality in (1.7.17):

$$S_T\psi_\infty(x) = \bar{\mu}\psi_\infty(x) \text{ for all } x \in \mathbb{T}^n, \quad (1.7.18)$$

then ψ_∞ is a positive eigenfunction of S_T corresponding to the eigenvalue $\bar{\mu}$. If, on the other hand, we have

$$S_T\psi_\infty(x) < \bar{\mu}\psi_\infty(x), \text{ on an open set } U \subset \mathbb{T}^n, \quad (1.7.19)$$

they we may use the fact that S_T maps the boundary of the cone of non-negative functions into its interior. In other words, we use the strong maximum principle here. Applying S_T to both sides of (1.7.17) gives then:

$$S_T\phi_\infty < \bar{\mu}\phi_\infty \text{ for all } x \in \mathbb{T}^n. \quad (1.7.20)$$

This contradicts, for large n , the minimality of μ_n . Thus, (1.7.19) is impossible, and $\bar{\mu}$ is the sought for eigenvalue. We set, from now on, $\phi_1 = \psi_\infty$:

$$S_T\phi_1 = \bar{\mu}\phi_1, \quad \phi_1(x) > 0 \text{ for all } x \in \mathbb{T}^n. \quad (1.7.21)$$

Exercise 1.7.6 So far, we have shown that $\bar{\mu} \geq 0$. Why do we know that, actually, $\bar{\mu} > 0$?

Let ϕ be an eigenfunction of S_T that is not a multiple of ϕ_1 , corresponding to an eigenvalue μ :

$$S_T\phi = \mu\phi.$$

Let us first assume that μ is real, and so is the eigenfunction ϕ . If $\mu \geq 0$, after multiplying ϕ by an appropriate factor, we may assume without loss of generality that $\phi_1(x) \geq \phi(x)$ for all $x \in \mathbb{T}^n$, $\phi_1 \not\equiv \phi$, and there exists $x_0 \in \mathbb{T}^n$ such that $\phi_1(x_0) = \phi(x_0)$. The strong comparison principle implies that then

$$S_T\phi_1(x) > S_T\phi(x) \quad \text{for all } x \in \mathbb{T}^n.$$

It follows, in particular, that

$$\bar{\mu}\phi_1(x_0) > \mu\phi(x_0),$$

hence $\bar{\mu} > \mu \geq 0$, as $\phi_1(x_0) = \phi(x_0) > 0$. This argument also shows that $\bar{\mu}$ is a simple eigenvalue.

If $\mu < 0$, then we can choose ϕ (after multiplying it by a, possibly negative, constant) so that, first,

$$\phi_1(x) \geq \phi(x), \quad \phi(x) \geq -\phi_1(x), \quad \text{for all } x \in \mathbb{T}^n, \quad (1.7.22)$$

and there exists $x_0 \in \mathbb{T}^n$ such that

$$\phi(x_0) = \phi_1(x_0).$$

Applying S_T to the second inequality in (1.7.22) gives, in particular,

$$\mu\phi(x_0) > -\bar{\mu}\phi_1(x_0), \quad (1.7.23)$$

thus $\bar{\mu} > -\mu$. In both cases, we see that $|\mu| < \bar{\mu}$.

Exercise 1.7.7 Use a similar consideration for the case when μ is complex. In that case, it helps to write the corresponding eigenfunction as

$$\phi = u + iv,$$

and consider the action of S_T on the span of u and v , using the same comparison idea. Show that $|\mu| < \bar{\mu}$. If in doubt, consult [?].

This completes the proof of Theorem 1.7.4. \square

1.7.3 Back to the principal periodic elliptic eigenvalue

Consider now again the operator L given by (1.7.1):

$$Lu(x) = -\Delta u + b_j(x) \frac{\partial u}{\partial x_j} + c(x)u, \quad (1.7.24)$$

with bounded, smooth and 1-periodic coefficients $b_j(x)$ and $c(x)$. One consequence of Theorem 1.7.4 is the analogous result for the principal periodic eigenvalue for L . We will also refer to the following as the Krein-Rutman theorem.

Theorem 1.7.8 *The operator L has a unique eigenvalue λ_1 associated to a positive function ϕ_1 . Moreover, each eigenvalue of L has a real part strictly larger than λ_1 .*

Proof. The operator L falls, of course, in the realm of Theorem 1.7.4, since its time-independent coefficients are T -periodic for all $T > 0$. We are also going to use the formula

$$L\phi = -\lim_{t \downarrow 0} \frac{S_t\phi - \phi}{t}, \quad (1.7.25)$$

for smooth $\phi(x)$, with the limit in the sense of uniform convergence. This is nothing but an expression of the fact that the function $u(t, x) = [S_t\phi](x)$ is the solution of

$$u_t + Lu = 0, \quad (1.7.26)$$

with the initial condition $u(0, x) = \phi(x)$, and if ϕ is smooth, then (1.7.26) holds also at $t = 0$.

Given $n \in \mathbb{N}$, let $\bar{\mu}_n$ be the principal eigenvalue of the operator $S_{1/n}$, with the principal eigenfunction $\phi_n > 0$:

$$S_{1/n}\phi_n = \bar{\mu}_n\phi_n,$$

normalized so that $\|\phi_n\|_\infty = 1$.

Exercise 1.7.9 Show that

$$\lim_{n \rightarrow \infty} \bar{\mu}_n = 1$$

directly, without using (1.7.27) below.

As $(S_{1/n})^n = S_1$ for all n , we conclude that ϕ_n is a positive eigenfunction of S_1 with the eigenvalue $(\bar{\mu}_n)^n$. By the uniqueness of the positive eigenfunction, we have

$$\bar{\mu}_n = (\bar{\mu}_1)^{1/n}, \quad \phi_n = \phi_1. \quad (1.7.27)$$

Note that, by the parabolic regularity, ϕ_1 is infinitely smooth, simply because it is a multiple of $S_1\phi_1$, which is infinitely smooth. Hence, (1.7.25) applies to ϕ_1 , and

$$L\phi_1 = -\lim_{n \rightarrow +\infty} n(S_{1/n} - I)\phi_1 = -\lim_{n \rightarrow +\infty} n(\bar{\mu}_1^{1/n} - 1)\phi_1 = -(\log \bar{\mu}_1)\phi_1.$$

We have thus proved the existence of an eigenvalue $\lambda_1 = -\log \bar{\mu}_1$ of L that corresponds to a positive eigenfunction. It is easy to see that if

$$L\phi = \lambda\phi,$$

then

$$S_1\phi = e^{-\lambda}\phi.$$

It follows that L can have only one eigenvalue corresponding to a positive eigenfunction. As we know that all eigenvalues μ of S_1 satisfy $|\mu| < \bar{\mu}_1$, we conclude that λ_1 is the eigenvalue of L with the smallest real part. \square

If L is symmetric – that is, it has the form

$$Lu = -\frac{\partial}{\partial x_i} \left(a_{ij}(x) \frac{\partial u}{\partial x_j} \right) + c(x)u, \quad (1.7.28)$$

with $a_{ij} = a_{ji}$, then the first eigenvalue is given by the minimization over $H^1(\mathbb{T}^n)$ of the Rayleigh quotient

$$\lambda_1 = \inf_{u \in H^1(\mathbb{T}^n)} \frac{\int_{\mathbb{T}^n} (a_{ij}(x)(\partial_i u)(\partial_j u) + c(x)u^2(x))dx}{\int_{\mathbb{T}^n} u^2(x)dx}. \quad (1.7.29)$$

The existence and uniqueness (up to a factor) of the minimizer is a classical exercise that we do not reproduce here. As for the positivity of the minimizer, we notice that, if ϕ is a minimizer of the Rayleigh quotient, then $|\phi_1|$ is also a minimizer, thus the unique minimizer is a positive function.

We end this section with a proposition that may look slightly academic, because it has to do with lowering the smoothness of the coefficients - something that we have not been so much interested in so far. A first reason to state it here is that it involves a nice juggling of estimates. Another reason is that it will have a true application in the next chapter. For an \mathbb{R}^n -valued function $v(x)$ we denote the divergence of $v(x)$ by

$$\nabla \cdot v = \sum_{j=1}^n \frac{\partial v_j}{\partial x_j}.$$

Proposition 1.7.10 *Let $b(x)$ be a smooth vector field over \mathbb{T}^n . The linear equation*

$$-\Delta e + \nabla \cdot (eb) = 0, \quad x \in \mathbb{T}^n, \quad (1.7.30)$$

has a unique solution $e_1^(x)$ normalized so that*

$$\|e_1^*\|_{L^\infty} = 1, \quad (1.7.31)$$

and such that $e_1^ > 0$ on \mathbb{T}^n . Moreover, for all $\alpha \in (0, 1)$, the function e_1^* is α -Hölder continuous, with the α -Hölder norm bounded by a universal constant depending only on $\|b\|_{L^\infty(\mathbb{T}^n)}$.*

A key point here is that the Hölder regularity of the solution only depends on the L^∞ -norm of $b(x)$ but not on its smoothness or any of its derivatives – this is typical for equations in the divergence form, and we will see much more of this in Chapter ???. This is very different from what we have seen so far in this chapter: we have always relied on the assumption that the coefficients are smooth, and the Hölder bounds for the solutions depended on the regularity of the coefficients. A very remarkable fact is that for equations in the divergence form, such as (1.7.30), one may often obtain bounds on the regularity of the solutions that depend only on the L^∞ -norm of the coefficients but not on their smoothness. Such bounds are much harder to get for equations in the non-divergence form.

Proof of Proposition 1.7.10

Let us denote

$$L\phi = -\Delta\phi - b_j(x)\frac{\partial\phi}{\partial x_j}. \quad (1.7.32)$$

The constant functions are the principal periodic eigenfunctions of L and zero is the principal eigenvalue:

$$L1 = 0. \quad (1.7.33)$$

Thus, by Theorem 1.7.8, the operator L has no other eigenvalue with a non-positive real part, which entails the same result for the operator

$$L^*\phi = -\Delta\phi + \nabla \cdot (b(x)\phi).$$

In particular, zero is the principal eigenvalue of L^* , associated to a positive eigenfunction $e_1^*(x) > 0$:

$$L^*e_1^* = 0, \quad \text{for all } x \in \mathbb{T}^n,$$

and we can normalize e_1^* so that that (1.7.31) holds. Thus, existence of $e_1^*(x)$ is the easy part of the proof.

The challenge is, of course, to bound the Hölder norms of e_1^* in terms of $\|b\|_{L^\infty(\mathbb{T}^n)}$ only. We would like to use a representation formula, as we already did many times in this chapter. In other words, we would like to treat the term $\nabla \cdot (e_1^*b)$ as a force, and convolve it with the fundamental solution of the Laplace equation in \mathbb{R}^n . For that, we need various quantities to be sufficiently integrable, so we first localize the equation, and then write a representation formula. This is very similar to the proof of the interior regularity estimates that we have mentioned very briefly in Section 1.6 – see Exercise 1.6.7. We recommend the reader to go back to this Section after finishing the current proof, and attempt this exercise again, setting $a_{ij}(t, x) = \delta_{ij}$ in (1.6.24) for simplicity.

Let $\Gamma(x)$ be a nonnegative smooth cut-off function such that $\Gamma(x) \equiv 1$ for $x \in [-2, 2]^n$ and $\Gamma(x) \equiv 0$ outside $(-3, 3)^n$. The function $v(x) = \Gamma(x)e_1^*(x)$ satisfies

$$-\Delta v = -2\nabla\Gamma \cdot \nabla e_1^* - e_1^*\Delta\Gamma - \Gamma\nabla \cdot (e_1^*b), \quad x \in \mathbb{R}^n. \quad (1.7.34)$$

Remember that e_1^* is bounded in L^∞ , thus so is v . As we will see, nothing should be feared from the cumbersome quantities like $\Delta\Gamma$ or $\nabla\Gamma$. We concentrate on the space dimensions $n \geq 2$, leaving $n = 1$ as an exercise. Let $E(x)$ be the fundamental solution of the Laplace equation in \mathbb{R}^n : the solution of

$$-\Delta u = f, \quad x \in \mathbb{R}^n, \quad (1.7.35)$$

is given by

$$u(x) = \int_{\mathbb{R}^n} E(x-y)u(y)dy. \quad (1.7.36)$$

Then we have

$$v(x) = \int_{\mathbb{R}^n} E(x-y) \left[-2\nabla\Gamma(y) \cdot \nabla e_1^*(y) - e_1^*(y)\Delta\Gamma(y) - \Gamma(y)\nabla \cdot (e_1^*(y)b(y)) \right] dy. \quad (1.7.37)$$

After an integration by parts, we obtain

$$v(x) = \int_{\mathbb{R}^n} \left((\nabla E(x-y) \cdot \nabla\Gamma(y))e_1^*(y) + E(x-y)e_1^*(y)\Delta\Gamma(y) + \nabla(E(x-y)\Gamma(y)) \cdot b(y)e_1^*(y) \right) dy. \quad (1.7.38)$$

The key point is that no derivatives of $b(x)$ or $e_1^*(x)$ appear in the right side of (1.7.38) – this is important as the only a priori information that we have on these functions is that they are bounded in L^∞ . Thus, the main point is to prove that integrals of the form

$$P(x) = \int_{\mathbb{R}^n} E(x-y)G(y)dy, \quad (1.7.39)$$

with a bounded and compactly supported function $G(x)$, and

$$I(x) = \int_{\mathbb{R}^n} \nabla E(x-y) \cdot F(y)dy, \quad (1.7.40)$$

with a bounded and compactly supported vector-valued function $F : \mathbb{R}^n \mapsto \mathbb{R}^n$, are α -Hölder continuous for all $\alpha \in (0, 1)$, with the Hölder constants depending only on α and the L^∞ -norms of F and G . Both F and G are supported inside the cube $[-3, 3]^n$. We will only consider the integral $I(x)$, as this would also show that $\nabla P(x)$ is α -Hölder. Using the expression

$$\nabla E(z) = c_n \frac{z}{|z|^n},$$

we see that

$$|I(x) - I(x')| \leq c_n \|F\|_{L^\infty} K(x, x'), \quad (1.7.41)$$

with

$$K(x, x') = \int_{(-3,3)^n} \left| \frac{x-y}{|x-y|^n} - \frac{x'-y}{|x'-y|^n} \right| dy. \quad (1.7.42)$$

Pick now $\alpha \in (0, 1)$. We estimate K by splitting the integration domain into two disjoint pieces:

$$A_x = \{y \in (-3, 3)^n : |x-y| \leq |x-x'|^\alpha\}, \quad B_x = \{y \in (-3, 3)^n : |x-y| > |x-x'|^\alpha\},$$

and denote by $K_A(x, x')$ and $K_B(x, x')$ the contribution to $K(x, x')$ by the integration over each of these two regions. To avoid some unnecessary trouble, we assume that $|x-x'| \leq l_\alpha$, with l_α such that

$$3l \leq l^\alpha \quad \text{for all } l \in [0, l_\alpha]. \quad (1.7.43)$$

With this choice, we have

$$|x'-y| \leq |x'-x| + |x-y| \leq 2|x-x'|^\alpha \quad \text{if } y \in A_x, \quad (1.7.44)$$

and

$$|x'-y| \geq |x-y| - |x'-x| \geq 2|x-x'| \quad \text{if } y \in B_x. \quad (1.7.45)$$

It follows that

$$K_A(x, x') \leq \int_{|x-y| \leq |x-x'|^\alpha} \frac{dy}{|x-y|^{n-1}} + \int_{|x'-y| \leq 2|x-x'|^\alpha} \frac{dy}{|x'-y|^{n-1}} \leq C|x-x'|^\alpha. \quad (1.7.46)$$

To estimate K_B , we write

$$\left| \frac{x-y}{|x-y|^n} - \frac{x'-y}{|x'-y|^n} \right| \leq C|x-x'| \int_0^1 \frac{d\sigma}{|x_\sigma - y|^n}, \quad x_\sigma = \sigma x + (1-\sigma)x'. \quad (1.7.47)$$

Note that for all $y \in B_x$ we have

$$|x_\sigma - y| \geq |x - y| - |x - x_\sigma| \geq |x - x'|^\alpha - |x - x'| \geq 2|x' - x|,$$

and $|y| \leq 3\sqrt{n}$, hence

$$\begin{aligned} K_B(x, x') &\leq |x - x'| \int_0^1 d\sigma \int_{B_x} \frac{dy}{|x_\sigma - y|^n} \leq |x - x'| \int_0^1 d\sigma \int_{|x_\sigma - y| \geq |x - x'|} \frac{\chi(|y| \leq 3\sqrt{n}) dy}{|x_\sigma - y|^n} \\ &\leq C|x - x'| \log |x - x'|, \end{aligned} \tag{1.7.48}$$

which implies the uniform α -Hölder bound for $I(x)$, for all $\alpha \in (0, 1)$. \square

The Dirichlet principal eigenvalue, related issues

We have so far talked about the principal eigenvalue for spatially periodic elliptic problems. This discussion applies equally well to problems in bounded domains, with the Dirichlet or Neumann boundary conditions. In the rest of this book, we will often encounter the Dirichlet problems, so let us explain this situation. Let Ω be a smooth bounded open subset of \mathbb{R}^n , and consider our favorite elliptic operator

$$Lu = -\Delta u + b_j(x) \frac{\partial u}{\partial x_j} + c(x)u, \tag{1.7.49}$$

with smooth coefficients $b_j(x)$ and $c(x)$. One could easily look at the more general problem

$$Lu = -a_{ij}(x) \frac{\partial^2 u}{\partial x_i \partial x_j} + b_j(x) \frac{\partial u}{\partial x_j} + c(x)u, \tag{1.7.50}$$

with essentially identical results, as long as the matrix $a_{ij}(x)$ is uniformly elliptic – we will avoid this just to keep the notation simpler. We are interested in the eigenvalue problem

$$\begin{aligned} Lu &= \lambda u && \text{in } \Omega, \\ u &= 0 && \text{on } \partial\Omega, \end{aligned} \tag{1.7.51}$$

and, in particular, in the existence of a positive eigenfunction $\phi > 0$ in Ω . The strategy will be as in the periodic case, to look at the initial value problem

$$\begin{aligned} u_t - \Delta u + b_j(x) \frac{\partial u}{\partial x_j} + c(x)u &= 0, && t > 0, x \in \Omega, \\ u &= 0, && t > 0, x \in \partial\Omega, \\ u(0, x) &= u_0(x). \end{aligned} \tag{1.7.52}$$

The coefficients b_j and c are smooth in (t, x) and T -periodic in t . Again, we set

$$(S_T u_0)(x) = u(T, x).$$

The main difference with the periodic case is that, here, the cone of continuous functions which are positive in Ω and vanish on $\partial\Omega$ has an empty interior, so we can not repeat verbatim the proof of the Krein-Rutman theorem for the operators on \mathbb{T}^n .

Exercise 1.7.11 Revisit the proof of the Krein-Rutman theorem in that case and identify the place where the proof would fail for the Dirichlet boundary conditions.

What will save the day is the strong maximum principle and the Hopf Lemma. We are not going to fully repeat the proof of Theorems 1.7.4 and 1.7.8, but we are going to prove a key proposition that an interested reader can use to prove the Krein-Rutman theorem for the Dirichlet problem.

Proposition 1.7.12 *Assume $u_0 \in C^1(\bar{\Omega})$ – that is, u_0 has derivatives that are continuous up to $\partial\Omega$, and that $u_0 > 0$ in Ω , and both $u_0 = 0$ and $\partial_\nu u_0 < 0$ on $\partial\Omega$. Then there is $\mu_1 > 0$ defined by the formula*

$$\mu_1 = \inf\{\mu > 0 : S_T u_0 \leq \mu u_0\}. \quad (1.7.53)$$

Moreover, if $\mu_2 > 0$ is defined as

$$\mu_2 = \inf\{\mu > 0 : (S_T \circ S_T)u_0 \leq \mu S_T u_0\}, \quad (1.7.54)$$

then either $\mu_1 > \mu_2$, or $\mu_1 = \mu_2$, and in the latter case $(S_T \circ S_T)u_0 \equiv \mu_2 S_T u_0$.

Proof. For the first claim, the existence of the infimum in (1.7.53), we simply note that

$$\mu u_0 \geq S_T u_0,$$

as soon as $\mu > 0$ is large enough, because $\partial_\nu u_0 < 0$ on $\partial\Omega$, $u_0 > 0$ in Ω , and $S_T u_0$ is a smooth function up to the boundary. As for the second item, let us first observe that

$$u(t+T, x) \leq \mu_1 u(t, x), \quad (1.7.55)$$

for any $t > 0$, by the maximum principle. Let us assume that

$$u(2T, x) \not\equiv \mu_1 u(T, x). \quad (1.7.56)$$

Then the maximum principle implies that

$$u(2T, x) < \mu_1 u(T, x) \text{ for all } x \in \Omega. \quad (1.7.57)$$

As

$$\max_{x \in \bar{\Omega}} [u(2T, x) - \mu_1 u(T, x)] = 0,$$

the parabolic Hopf lemma, together with (1.7.55) and (1.7.56), implies the existence of $\delta > 0$ such that

$$\partial_\nu (u(2T, x) - \mu_1 u(T, x)) \geq \delta > 0, \quad \text{for all } x \in \partial\Omega. \quad (1.7.58)$$

It follows that for $\varepsilon > 0$ sufficiently small, we have

$$u(2T, x) - \mu_1 u(T, x) \leq -\frac{\delta}{2} d(x, \partial\Omega) \quad \text{for } x \in \Omega \text{ such that } d(x, \partial\Omega) < \varepsilon.$$

On the other hand, once again, the strong maximum principle precludes a touching point between $u(2T, x)$ and $\mu_1 u(T, x)$ inside

$$\bar{\Omega}_\varepsilon = \{x \in \Omega : d(x, \partial\Omega) \geq \varepsilon\}.$$

Therefore, there exists δ_1 such that

$$u(2T, x) - \mu_1 u(T, x) \leq -\delta_1, \quad \text{for all } x \in \overline{\Omega}_\varepsilon.$$

We deduce that there is a – possibly very small – constant $c > 0$ such that

$$u(2T, x) - \mu_1 u(T, x) \leq -cd(x, \partial\Omega) \quad \text{in } \Omega.$$

However, $u(T, x)$ is controlled from above by $Cd(x, \partial\Omega)$, for a possibly large constant $C > 0$. All in all, we have

$$u(2T, x) \leq \left(\mu_1 - \frac{c}{C}\right)u(T, x),$$

hence (1.7.56) implies that $\mu_2 < \mu_1$, which proves the second claim of the proposition. \square

Exercise 1.7.13 Deduce from Proposition 1.7.12 the versions of Theorems 1.7.4 and 1.7.8 for operators S_T and L , this time with the Dirichlet boundary conditions.

Thus, the eigenvalue problem (1.7.51), has a principal eigenvalue that enjoys all the properties we have proved in the periodic one: it has the least real part among all eigenvalues, and is the only eigenvalue associated to a positive eigenfunction.

Exercise 1.7.14 Assume that L is symmetric; it has the form

$$Lu = -\frac{\partial}{\partial x_j} \left(a_{ij}(x) \frac{\partial u}{\partial x_i} \right) + c(x)u \quad (1.7.59)$$

Then, the principal eigenvalue is given by the minimization of the Rayleigh quotient over the Sobolev space $H_0^1(\Omega)$:

$$\lambda_1 = \inf_{u \in H_0^1(\Omega), \|u\|_{L^2} = 1} \int_{\Omega} \left(a_{ij}(x) \frac{\partial u}{\partial x_i} \frac{\partial u}{\partial x_j} + c(x)u^2(x) \right) dx. \quad (1.7.60)$$

Exercise 1.7.15 Adapt the preceding discussion to prove the existence of a principal eigenvalue to the Neumann eigenvalue problem

$$\begin{aligned} Lu &= \lambda u, & x &\in \Omega, \\ \partial_\nu u &= 0, & x &\in \partial\Omega. \end{aligned} \quad (1.7.61)$$

Exercise 1.7.16 (*The principal eigenvalue in an unbounded domain*) Consider the differential operator

$$L = -\Delta + c(x), \quad x \in \mathbb{R}^n.$$

Assume that c is bounded and uniformly continuous, Assume the existence of $c_\infty > 0$ such that

$$\lim_{|x| \rightarrow +\infty} c(x) = c_\infty. \quad (1.7.62)$$

Also assume that $c(x) < c_\infty$ for all $x \in \mathbb{R}^n$. The goal of the exercise is to prove a Krein-Rutman type theorem for L . For $n \geq 1$, let μ_n be the principal eigenvalue of L in $B_n(0)$, with Dirichlet conditions.

1. Show that $(\mu_n)_n$ is a decreasing sequence, bounded by $-\|c\|_\infty$. Let μ_∞ be its limit.
2. Let $\psi_{\varepsilon,n}(x)$ solve

$$\begin{aligned} -\Delta\psi + \varepsilon\psi &= 0 & (B_n(0)\setminus B_R(0)) \\ \psi &= 1 & (\partial B_R(0)) \\ \psi &= 0 & (\partial B_n(0)). \end{aligned}$$

Show that $|\partial_r\psi_{\varepsilon,n}| = O(\sqrt{\varepsilon})$ on $\partial B_R(0)$, as soon as n is large enough.

3. Let ϕ_{2R} be the first Dirichlet eigenfunction in $B_{2R}(0)$, that is equal to 1 on $\partial B_R(0)$ (why is ϕ_{2R} radial?). Let $\underline{\phi}_n$ be equal to ϕ_{2R} in $B_R(0)$, and $\psi_{\varepsilon,n}$ in $B_n(0)\setminus B_R(0)$. Show that, if $\varepsilon > 0$ is small enough, R large and n very large, then we have

$$L\underline{\phi}_n \leq (c_\infty - \varepsilon)\underline{\phi}_n.$$

4. Deduce that $\mu_\infty \leq c_\infty - \varepsilon$.
5. Conclude that L has the Krein-Rutman property.
6. Show that the first eigenfunction decays exponentially fast at infinity.

Exercise 1.7.17 Set

$$L = -\Delta + c(x), \quad x \in \mathbb{R}^n,$$

the function c satisfying (1.7.62) for some positive c_∞ . We do not, however, assume $c(x) < c_\infty$ anymore. Find as many properties of the preceding exercise as possible that would fail without this assumption.

Exercise 1.7.18 Redo the existence part of Exercise 1.7.16 with the aid of the Rayleigh quotients, without any approximation on a finite domain.

These three exercises give just a glimpse at what happens to the principal eigenvalue in unbounded domains – an interested reader should investigate further, starting with the variational formulations of [?] and [?], and continuing with the more recent papers [?, ?].

1.7.4 The principal eigenvalue and the comparison principle

Let us now connect the principal eigenvalue and the comparison principle. Since we are at the moment dealing with the Dirichlet problems, let us remain in this context. There would be nothing significantly different about the periodic problems.

The principal eigenfunction $\phi_1 > 0$, solution of

$$L\phi_1 = \lambda_1\phi_1, \text{ in } \Omega, \tag{1.7.63}$$

$$\phi_1 = 0 \text{ on } \partial\Omega, \tag{1.7.64}$$

with

$$Lu = -\Delta u + b_j(x)\frac{\partial u}{\partial x_j} + c(x)u, \tag{1.7.65}$$

in particular, provides a special solution

$$\psi(t, x) = e^{-\lambda_1 t} \phi_1(x) \quad (1.7.66)$$

for the linear parabolic problem

$$\begin{aligned} \psi_t + L\psi &= 0, \quad t > 0, x \in \Omega \\ \psi &= 0 \text{ on } \partial\Omega. \end{aligned} \quad (1.7.67)$$

Consider then the Cauchy problem

$$\begin{aligned} v_t + Lv &= 0, \quad t > 0, x \in \Omega \\ v &= 0 \text{ on } \partial\Omega, \\ v(0, x) &= g(x), \quad x \in \Omega, \end{aligned} \quad (1.7.68)$$

with a smooth bounded function $g(x)$ that vanishes at the boundary $\partial\Omega$. We can find a constant $M > 0$ so that

$$-M\phi_1(x) \leq g(x) \leq M\phi_1(x), \quad \text{for all } x \in \Omega.$$

The comparison principle then implies that for all $t > 0$ we have a bound

$$-M\phi_1(x)e^{-\lambda_1 t} \leq v(t, x) \leq M\phi_1(x)e^{-\lambda_1 t}, \quad \text{for all } x \in \Omega, \quad (1.7.69)$$

which is very useful, especially if $\lambda_1 > 0$. The assumption that the initial condition g vanishes at the boundary $\partial\Omega$ is not necessary but removes the technical step of having to show that even if $g(x)$ does not vanish on the boundary, then for any positive time $t_0 > 0$ we can find a constant C_0 so that $|v(t_0, x)| \leq C_0\phi_1(x)$. This leads to the bound (1.7.69) for all $t > t_0$.

Let us now apply the above considerations to the solutions of the elliptic problem

$$\begin{aligned} Lu &= g(x), \text{ in } \Omega, \\ u &= 0 \text{ on } \partial\Omega, \end{aligned} \quad (1.7.70)$$

with a non-negative function $g(x)$. When can we conclude that the solution $u(x)$ is also non-negative? The solution of (1.7.70) can be formally written as

$$u(x) = \int_0^\infty v(t, x) dt. \quad (1.7.71)$$

Here, the function $v(t, x)$ satisfies the Cauchy problem (1.7.68). If the principal eigenvalue λ_1 of the operator L is positive, then the integral (1.7.71) converges for all $x \in \Omega$ because of the estimates (1.7.69), and the solution of (1.7.70) is, indeed, given by (1.7.71). On the other hand, if $g(x) \geq 0$ and $g(x) \not\equiv 0$, then the parabolic comparison principle implies that $v(t, x) > 0$ for all $t > 0$ and all $x \in \Omega$. It follows that $u(x) > 0$ in Ω .

Therefore, we have proved the following theorem that succinctly relates the notions of the principal eigenvalue and the comparison principle.

Theorem 1.7.19 *If the principal eigenvalue of the operator L is positive then solutions of the elliptic equation (1.7.70) satisfy the comparison principle: $u(x) > 0$ in Ω if $g(x) \geq 0$ in Ω and $g(x) \not\equiv 0$.*

This theorem allows to look at the maximum principle in narrow domains introduced in the previous chapter from a slightly different point of view: the narrowness of the domain implies that the principal eigenvalue of L is positive no matter what the sign of the free coefficient $c(x)$ is. This is because the “size” of the second order term in L increases as the domain narrows, while the “size” of the zero-order term does not change. Therefore, in a sufficiently narrow domain the principal eigenvalue of L will be positive (recall that the required narrowness does depend on the size of $c(x)$). A similar philosophy applies to the maximum principle for the domains of a small volume.

We conclude this topic with another characterization of the principal eigenvalue of an elliptic operator in a bounded domain, which we leave as an (important) exercise for the reader. Let us define

$$\mu_1(\Omega) = \sup\{\lambda : \exists \phi \in C^2(\Omega) \cap C^1(\bar{\Omega}), \phi > 0 \text{ and } (L - \lambda)\phi \geq 0 \text{ in } \Omega\}, \quad (1.7.72)$$

and

$$\mu'_1(\Omega) = \inf\{\lambda : \exists \phi \in C^2(\Omega) \cap C^1(\bar{\Omega}), \phi = 0 \text{ on } \partial\Omega, \phi > 0 \text{ in } \Omega, \text{ and } (L - \lambda)\phi \leq 0 \text{ in } \Omega\}. \quad (1.7.73)$$

Exercise 1.7.20 *Let L be an elliptic operator in a smooth bounded domain Ω , and let λ_1 be the principal eigenvalue of the operator L , and $\mu_1(\Omega)$ and $\mu'_1(\Omega)$ be as above. Show that*

$$\lambda_1 = \mu_1(\Omega) = \mu'_1(\Omega). \quad (1.7.74)$$

As a hint, say, for the equality $\lambda_1 = \mu_1(\Omega)$, we suggest, assuming existence of some $\lambda > \lambda_1$ and $\phi > 0$ such that

$$(L - \lambda)\phi \geq 0,$$

to consider the Cauchy problem

$$u_t + (L - \lambda)u = 0, \quad \text{in } \Omega$$

with the initial data $u(0, x) = \phi(x)$, and with the Dirichlet boundary condition $u(t, x) = 0$ for $t > 0$ and $x \in \partial\Omega$. One should prove two things: first, that $u_t(t, x) \leq 0$ for all $t > 0$, and, second, that there exists some constant $C > 0$ so that

$$u(t, x) \geq C\phi_1(x)e^{(\lambda - \lambda_1)t},$$

where ϕ_1 is the principal Dirichlet eigenfunction of L . This will lead to a contradiction. The second equality in (1.7.74) is proved in a similar way.

1.8 Reaction-diffusion waves

As a conclusion to this chapter, we will be interested here in one-dimensional models of the form

$$u_t - u_{xx} = f(x, u), \quad t > 0, \quad x \in \mathbb{R}, \quad (1.8.1)$$

the assumptions on f being made precise as the study develops. We will see in this section how the possibility of comparing two solutions of the same problem will imply their convergence in the long time limit, putting to work the two main characters we have seen so far in this chapter: the comparison principle and the Harnack inequality. We will also put to good use the ideas developed for the existence of the principal eigenvalues for elliptic operators, they are the same - although they will sometimes be imbedded in more or less technical considerations. We will start by the simplest possible model: the Allen-Cahn equation, with $f(x, u) = f(u) = u(1 - u)^2$. There is an explicit steady solution, and we will show in detail how every solution of the problem, that vaguely looks like the steady solution at both ends at time $t = 0$, will converge to a translate of it for large times. The rest of the chapter will be devoted to showing that the idea is universal, and helps the understanding of seemingly more complicated, or unrelated models. We will first treat nonlinearities that are less symmetric than the Allen-Cahn one, giving rise to travelling waves, that will attract the whole dynamics of the solutions. We will finally investigate the large-time behavior of (1.8.1) with an $f(x, u)$ periodic in x . This is going to give raise to pulsating waves, i.e. waves that look time-periodic in some Galilean reference frame. These waves will be shown to be globally attracting, thus giving some substance to the advertisemnt (that we made in the introduction) about how space periodicity generates time-periodicity.

1.8.1 The long time behavior for the Allen-Cahn equation

We consider the one-dimensional Allen-Cahn equation

$$u_t - u_{xx} = f(u), \quad (1.8.2)$$

with

$$f(u) = u - u^3. \quad (1.8.3)$$

Recall that we have already considered the steady solutions of this equation in Section ?? of Chapter ??, and, in particular, the role of its explicit time-independent solutions

$$\phi(x) = \tanh\left(\frac{x}{\sqrt{2}}\right), \quad (1.8.4)$$

and its translates $\phi_{x_0}(x) := \phi(x + x_0)$, $x_0 \in \mathbb{R}$.

Exercise 1.8.1 We have proved in Chapter ?? that, if $\psi(x)$ is a steady solution to (1.8.2) that satisfies

$$\lim_{x \rightarrow -\infty} \psi(x) = -1, \quad \lim_{x \rightarrow +\infty} \psi(x) = 1,$$

then ψ is a translate of ϕ . For an alternative proof, draw the phase portrait of the equation

$$-\psi'' = f(\psi) \quad (1.8.5)$$

in the (ψ, ψ') plane. For an orbit (ψ, ψ') connecting $(-1, 0)$ to $(1, 0)$, show that the solution tends to $(-1, 0)$ exponentially fast. Multiply then (1.8.5) by ψ' , integrate from $-\infty$ to x and conclude.

Recall that the Allen-Cahn equation is a simple model for a physical situation when two phases are stable, corresponding to $u = \pm 1$. The time dynamics of the initial value problem for (1.8.2) corresponds to a competition between these two states. The fact that

$$\int_{-1}^1 f(u) du = 0 \tag{1.8.6}$$

means that the two states are "equally stable" – this is a necessary condition for (1.8.2) to have a time-independent solution $\phi(x)$ such that

$$\phi(x) \rightarrow \pm 1, \quad \text{as } x \rightarrow \pm\infty. \tag{1.8.7}$$

In other words, such connection between $+1$ and -1 exists only if (1.8.6) holds.

Since the two phases $u = \pm 1$ are equally stable, one expects that if the initial condition $u_0(x)$ for (1.8.2) satisfies

$$\lim_{x \rightarrow -\infty} u_0(x) = -1, \quad \lim_{x \rightarrow +\infty} u_0(x) = 1, \tag{1.8.8}$$

then, as $t \rightarrow +\infty$, the solution $u(t, x)$ will converge to a steady equilibrium, that has to be a translate of ϕ . This is the subject of the next theorem, that shows, in addition, that the convergence rate is exponential.

Theorem 1.8.2 *There exists $\omega > 0$ such that for any uniformly continuous and bounded initial condition u_0 for (1.8.2) that satisfies (1.8.8), we can find $x_0 \in \mathbb{R}$ and $C_0 > 0$ such that*

$$|u(t, x) - \phi(x + x_0)| \leq C_0 e^{-\omega t}, \quad \text{for all } x \in \mathbb{R} \text{ and } t > 0. \tag{1.8.9}$$

Since there is a one parameter family of steady solutions, naturally, one may ask how the solution of the initial value problem chooses a particular translation of ϕ in the long time limit. In other words, one would like to know how the shift x_0 depends on the initial condition u_0 . However, this dependence is quite implicit and there is no simple expression for x_0 .

There are at least two ways to prove Theorem 1.8.2, both of them need the forthcoming Lemma 1.8.4, that bounds the level sets of the solution. Once this is at hand, a first option is to solve the following

Exercise 1.8.3 Assume Lemma 1.8.4 to be true.

1. Verify that the energy functional

$$J(u) = \int_{\mathbb{R}} \left(\frac{1}{2} |u_x|^2 - F(u) \right) dx, \quad F(u) = \int_{-1}^u f(v) dv,$$

decreases in time for any solution $u(t, x)$ of (1.8.2).

2. With the aid of Lemma 1.8.4 and the preceding question, show that the solution eventually comes very close to a translate $\phi_{x_0}(x)$, uniformly on \mathbb{R} .

3. Prove a Krein-Rutman type property for the operator

$$\mathcal{M}u = -u_{xx} - f'(\phi_{x_0})u.$$

What is the principal eigenvalue, and what is an associated eigenfunction?

4. If $v(x)$ is close to ϕ_{x_0} , show that one may decompose it uniquely as

$$v(x) = \phi_{x_0}(x + X) + w,$$

X small and w small, orthogonal to the null space of M .

5. Assume that u_0 is close to ϕ_{x_0} .

- Show the existence of $T > 0$ such that the decomposition

$$u(t, x) = \phi_{x_0}(x + X(t)) + w(t, x)$$

X small and w small, orthogonal to the null space of M , holds at least up to time T .

- Write a system of equations for $(X(t), w(t, x))$.
- Deduce that T can be chosen infinite and that $X(t)$ converges, exponentially in time, to some x_1 close to x_0 .

6. Round up everything and conclude.

This is, more or less, the method devised in the beautiful paper of Fife and McLeod [?]. It has been generalized to gradient systems in a remarkable paper of Risler [?], which proves very precise spreading estimates of the leading edge of the solutions, only based on a one-dimensional set of energy functionals. Risler's ideas were put to work on the simpler example (1.8.2) in a paper by Gally and Risler [?].

We chose to present an alternative method, entirely based on sub and super-solutions that come closer and closer to each other. It avoids the spectral arguments, and is more flexible as there are many reaction-diffusion problems where the comparison principle and the Harnack inequality are available but the energy functionals do not exist. The reader should also be aware that there are many problems, such as many reaction-diffusion systems, where the situation is the opposite: the energy functional exists but the comparison principle is not applicable.

Before we begin, we note that the function f satisfies

$$f'(u) \leq -1 \quad \text{for } |u| \geq 5/6, \quad f'(u) \leq -3/2 \quad \text{for } |u| \geq 11/12. \quad (1.8.10)$$

We will also take $R_0 > 0$ such that

$$|\phi(x)| \geq 11/12 \quad \text{for } |x| \geq R_0. \quad (1.8.11)$$

A bound on the level sets

The first ingredient is to prove that the level sets of $u(t, x)$ do not, indeed, go to infinity, so that the region of activity, where $u(t, x)$ is not too close to ± 1 , happens, essentially, in a compact set. This crucial step had already been identified by Fife and McLeod, and we reproduce here their argument. The idea is to squish $u(t, x)$ between two different translates of ϕ , with a correction that goes to zero exponentially in fast time.

Lemma 1.8.4 *Let u_0 satisfy the assumptions of the theorem. There exist $\xi_\infty^\pm \in \mathbb{R}$, and $q_0 > 0$, such that*

$$\phi(x + \xi_\infty^-) - q_0 e^{-t} \leq u(t, x) \leq \phi(x + \xi_\infty^+) + q_0 e^{-t}, \quad (1.8.12)$$

for all $t \geq 0$ and $x \in \mathbb{R}$.

Proof. For the upper bound, we are going to devise two functions $\xi^+(t)$ and $q(t)$ such that

$$\bar{u}(t, x) = \phi(x + \xi^+(t)) + q(t) \quad (1.8.13)$$

is a super-solution to (1.8.2), with an increasing but bounded function $\xi^+(t)$, and an exponentially decreasing function $q(t) = q_0 \exp(-t)$. One would also construct, in a similar way, a sub-solution of the form

$$\underline{u}(t, x) = \phi(x + \xi^-(t)) - q(t), \quad (1.8.14)$$

possibly increasing q a little, with a decreasing but bounded function $\xi^-(t)$.

Let us denote

$$N[u] = \partial_t u - u_{xx} - f(u). \quad (1.8.15)$$

Now, with $\bar{u}(t, x)$ as in (1.8.13), we have

$$N[\bar{u}] = \dot{q} + \dot{\xi}^+ \phi'(\zeta) - f(\phi(\zeta) + q) + f(\phi(\zeta)), \quad (1.8.16)$$

with $\zeta = x + \xi^+(t)$. Our goal is to choose $\xi^+(t)$ and $q(t)$ so that

$$N[\bar{u}] \geq 0, \quad \text{for all } t \geq 0 \text{ and } x \in \mathbb{R}, \quad (1.8.17)$$

so that $\bar{u}(t, x)$ is a super-solution to (1.8.2). We will consider separately the regions $|\zeta| \leq R_0$ and $|\zeta| \geq R_0$.

Step 1. The region $|\zeta| \geq R_0$. First, we have

$$\phi(\zeta) + q(t) \geq 11/12 \text{ for } \zeta \geq R_0,$$

as $q(t) \geq 0$. If we assume that $q(0) \leq 1/12$ and make sure that $q(t)$ is decreasing in time, then we also have

$$\phi(\zeta) + q \leq -5/6 \text{ for } \zeta \leq -R_0.$$

We have, therefore, as long as $\xi^+(t)$ is increasing, using (1.8.10):

$$N[\bar{u}] \geq \dot{q} - f(\phi(\zeta) + q) + f(\phi) \geq \dot{q} + q, \text{ for } |\zeta| \geq R_0. \quad (1.8.18)$$

It suffices, therefore, to choose

$$q(t) = q(0)e^{-t}, \quad (1.8.19)$$

with $q(0) \leq 1/12$, and an increasing $\xi^+(t)$, to ensure that

$$N[\bar{u}] \geq 0, \quad \text{for all } t \geq 0 \text{ and } |\zeta| \geq R_0. \quad (1.8.20)$$

Step 2. The region $|\zeta| \leq R_0$. This time, we have to choose $\xi^+(t)$ properly. We write

$$N[\bar{u}] \geq \dot{q} + \dot{\xi}^+ \phi'(\zeta) - M_f q, \quad M_f = \|f'\|_{L^\infty}, \quad (1.8.21)$$

and choose

$$\dot{\xi}^+ = \frac{1}{k_0} (-\dot{q} + M_f q), \quad k_0 = \inf_{|\zeta| \leq R_0} \phi'(\zeta), \quad (1.8.22)$$

to ensure that the right side of (1.8.21) is non-negative. Using expression (1.8.19) for $q(t)$, we obtain

$$\xi^+(t) = \xi^+(0) + \frac{q(0)}{k_0} (1 + M_f)(1 - e^{-t}). \quad (1.8.23)$$

To summarize, with the above choices of $q(t)$ and $\xi^+(t)$, we know that \bar{u} satisfies (1.8.17).

It remains to choose $q(0)$ and $\xi^+(0)$ so that $\bar{u}(t, x)$ is actually above $u(t, x)$ – as we have already established (1.8.17), the comparison principle tells us that we only need to verify that

$$\bar{u}(0, x) \geq u_0(x), \quad \text{for all } x \in \mathbb{R}. \quad (1.8.24)$$

Because u_0 tends to ± 1 at $\pm\infty$, there exists ξ_0^+ (possibly quite large), and $q_0 \in (0, 1/12)$ such that

$$u_0(x) \leq \phi(x + \xi_0^+) + q_0. \quad (1.8.25)$$

Thus, it is enough to choose $q(0) = q_0$, $\xi^+(0) = \xi_0^+$. \square

Exercise 1.8.5 Follow the same strategy to construct a sub-solution $\underline{u}(t, x)$ as in (1.8.14).

Lemma 1.8.4 traps nicely the level sets of u . But will this imply convergence to a steady solution, or will the level sets of $u(t, x)$ oscillate inside a bounded set? First, let us restate our findings in a more precise way. We have shown the following

Corollary 1.8.6 *Assume that we have*

$$\phi(x + \xi_0^-) - q_0 \leq u_0(x) \leq \phi(x + \xi_0^+) + q_0, \quad (1.8.26)$$

with $0 \leq q_0 \leq 1/12$. Then, we have

$$\phi(x + \xi^-(t)) - q(t) \leq u_0(x) \leq \phi(x + \xi^+(t)) + q(t). \quad (1.8.27)$$

with $q(t) = q_0 e^{-t}$, and

$$\xi^+(t) = \xi_0^+ + \frac{q_0}{k_0} (1 + M_f)(1 - e^{-t}), \quad \xi^-(t) = \xi_0^- - \frac{q_0}{k_0} (1 + M_f)(1 - e^{-t}). \quad (1.8.28)$$

One issue here is that the gap between $\xi^+(t)$ and $\xi^-(t)$ is not decreasing in time but rather increasing – the opposite of what we want! Our goal is to show that we can actually choose $\xi^+(t)$ and $\xi^-(t)$ in (1.8.27) so that the "sub-solution/super-solution gap" $\xi^+(t) - \xi^-(t)$ would decrease to zero as $t \rightarrow +\infty$ – this will prove convergence of the solution to a translate of ϕ .

The mechanism to decrease this difference will be kindly provided by the strong maximum principle. The idea is to iteratively trap the solutions, at an increasing sequence of times, between translates of ϕ_0 , that will come closer and closer to each other, thus implying the convergence. However, as there will be some computations, it is worth explaining beforehand what the main idea is, and which difficulties we will see.

Let us consider for the moment a slightly better situation than in Lemma 1.8.4 – assume that $u_0(x)$ is actually trapped between $\phi(x + \xi_0^-)$ and $\phi(x + \xi_0^+)$, without the need for an additional term $q(t)$:

$$\phi(x + \xi_0^-) \leq u_0(x) \leq \phi(x + \xi_0^+). \quad (1.8.29)$$

Then, $u(t, x)$ is at a positive distance from one of the two translates, on compact sets, at least for $0 \leq t \leq 1$, say, $\phi(x + \xi_0^+)$. This is where the strong maximum principle strikes: at $t = 1$, it will make the infimum of $\phi(x + \xi_0^+) - u(t, x)$ strictly positive, at least on a large compact set. We would like to think that then we may translate $\phi(x + \xi_0^+)$ to the right a little, decreasing ξ_0^+ , while keeping it above $u(1, x)$. The catch is that, potentially, the tail of $u(1, x)$ – that we do not control very well at the moment – might go over $\phi(x + \xi)$, as soon as ξ is just a little smaller than ξ_0^+ . Let us ignore this, and assume that magically we have

$$\phi(x + \xi_0^-) \leq u(1, x) \leq \phi(x + \xi_1^+), \quad (1.8.30)$$

with

$$\xi_1^+ = \xi_0^+ - \delta(\xi_0^+ - \xi_0^-), \quad (1.8.31)$$

with some $\delta > 0$. If we believe in this scenario, we might just as well hope that the situation may be iterated: at the time $t = n$, we have

$$\phi(x + \xi_n^-) \leq u(n, x) \leq \phi(x + \xi_n^+), \quad (1.8.32)$$

with

$$\xi_{n+1}^+ - \xi_{n+1}^- \leq (1 - \delta)(\xi_n^+ - \xi_n^-). \quad (1.8.33)$$

This would imply a geometric decay of $\xi_n^+ - \xi_n^-$ to zero, which, in turn, would imply the exponential convergence of $u(t, x)$ to a translate of ϕ .

The gap in the previous argument is, of course, in our lack of control of the tail of $u(t, x)$ that prevents us from being sure that (1.8.30), with ξ_1^+ as in (1.8.31), holds everywhere on \mathbb{R} rather than on a compact set. There is no way we can simply ignore it: we will see in Chapter ?? that the dynamics of many respectable equations is controlled exactly by the tail of its solutions. Such will not be the case here, but we will have to go through the pain of controlling the tail of u at every step. This leads to the somewhat heavy proof that follows, which is itself a simplified version of [?], where global exponential stability of transition waves is shown. However, there is essentially no other idea than what we have just explained, the rest are just technical embellishments. The reader should also recall that we have already encountered a tool for the tail-control in the Allen-Cahn equation: Corollary ?? in Chapter ?? served exactly that purpose in the proof of Theorem ?. We are going to use something very similar here.

The proof of Theorem 1.8.2

As promised, the strategy is a refinement of the proof of Lemma 1.8.4. We will construct a sequence of sub-solutions \underline{u}_n and super-solutions \bar{u}_n defined for $t \geq T_n$, such that

$$\underline{u}_n(t, x) \leq u(t, x) \leq \bar{u}_n(t, x) \text{ for } t \geq T_n. \quad (1.8.34)$$

Here, $T_n \rightarrow +\infty$ is a sequence of times with

$$T_n + T \leq T_{n+1} \leq T_n + 2T, \quad (1.8.35)$$

and the time step $T > 0$ to be specified later on. The sub- and super-solutions will be of the familiar form (1.8.27)-(1.8.28):

$$\underline{u}_n(t, x) = \phi(x + \xi_n^-(t)) - q_n e^{-(t-T_n)}, \quad \bar{u}_n(t, x) = \phi(x + \xi_n^+(t)) + q_n e^{-(t-T_n)}, \quad t \geq T_n, \quad (1.8.36)$$

with $\xi_n^\pm(t)$ as in (1.8.28):

$$\xi_n^+(t) = \xi_n^+ + \frac{q_n}{k_0}(1 + M_f)(1 - e^{-(t-T_n)}), \quad \xi_n^-(t) = \xi_n^- - \frac{q_n}{k_0}(1 + M_f)(1 - e^{-(t-T_n)}). \quad (1.8.37)$$

The reader has surely noticed a slight abuse of notation: we denote by ξ_n^\pm the values of $\xi_n^\pm(t)$ at the time $t = T_n$. This allows us to avoid introducing further notation, and we hope it does not cause too much confusion.

Our plan is to switch from one pair of sub- and super-solutions to another at the times T_n , and improve the difference in the two shifts at the "switching" times, to ensure that

$$\xi_{n+1}^+ - \xi_{n+1}^- \leq (1 - \delta)(\xi_n^+ - \xi_n^-), \quad (1.8.38)$$

with some small but fixed constant $\delta > 0$ such that

$$e^{-T} \leq c_T \delta \leq \frac{1}{4}. \quad (1.8.39)$$

The constant c_T will also be chosen very small in the end – one should think of (1.8.39) as the requirement that the time step T is very large. This is natural: we can only hope to improve on the difference $\xi_n^+ - \xi_n^-$, as in (1.8.38), after a very large time step T . The shifts can be chosen so that they are uniformly bounded:

$$|\xi_n^\pm| \leq M, \quad (1.8.40)$$

with a sufficiently large M – this follows from the bounds on the level sets of $u(t, x)$ that we have already obtained. As far as q_n are concerned, we will ask that

$$0 \leq q_n \leq c_q \delta (\xi_n^+ - \xi_n^-), \quad (1.8.41)$$

with another small constant c_q to be determined. Note that at $t = 0$ we may ensure that q_0 satisfies (1.8.41) simply by taking ξ_0^+ sufficiently positive and ξ_0^- sufficiently negative.

As we have uniform bounds on the location of the level sets of $u(t, x)$, and the shifts ξ_n^\pm will be chosen uniformly bounded, as in (1.8.40), after possibly increasing R_0 in (1.8.11), we can ensure that

$$\phi(x + \xi_n^\pm(t)) \geq 11/12, \quad u(t, x) \geq 11/12, \quad \text{for } x \geq R_0 \text{ and } t \geq T_n, \quad (1.8.42)$$

and

$$-1 < \phi(x + \xi_n^\pm(t)) \leq 11/12, \quad -1 < u(t, x) \leq -11/12, \quad \text{for } x \leq -R_0 \text{ and } t \geq T_n, \quad (1.8.43)$$

which implies

$$f'(\phi(x + \xi_n^\pm(t))) \leq -1, \quad f'(u(t, x)) \leq -1, \quad \text{for } |x| \geq R_0 \text{ and } t \geq T_n. \quad (1.8.44)$$

Let us now assume that at the time $t = T_n$ we have the inequality

$$\phi(x + \xi_n^-) - q_n \leq u(T_n, x) \leq \phi(x + \xi_n^+) + q_n, \quad (1.8.45)$$

with the shift q_n that satisfies (1.8.41). Our goal is to find a time $T_{n+1} \in [T_n + T, T_n + 2T]$, and the new shifts ξ_{n+1}^\pm and q_{n+1} , so that (1.8.45) holds with n replaced by $n+1$ and the new gap $\xi_{n+1}^+ - \xi_{n+1}^-$ satisfies (1.8.38). We will consider two different cases.

Case 1: the solution gets close to the super-solution. Let us first assume that there is a time $\tau_n \in [T_n + T, T_n + 2T]$ such that the solution $u(\tau_n, x)$ is "very close" to the super-solution $\bar{u}_n(\tau_n, x)$ on the interval $\{|x| \leq R_0 + 1\}$. More precisely, we assume that

$$\sup_{|x| \leq R_0 + 1} \left(\bar{u}_n(\tau_n, x) - u(\tau_n, x) \right) \leq \delta(\xi_n^+ - \xi_n^-). \quad (1.8.46)$$

We will show that in this case we may take $T_{n+1} = \tau_n$, and set

$$\xi_{n+1}^+ = \xi_n^+(\tau_n), \quad \xi_{n+1}^- = \xi_n^- + (\xi_n^+(\tau_n) - \xi_n^+) + \delta(\xi_n^+ - \xi_n^-), \quad (1.8.47)$$

as long as δ is sufficiently small, making sure that

$$\xi_{n+1}^+ - \xi_{n+1}^- = (1 - \delta)(\xi_n^+ - \xi_n^-), \quad (1.8.48)$$

and also choose q_{n+1} so that

$$q_{n+1} = c_q \delta (\xi_{n+1}^+ - \xi_{n+1}^-). \quad (1.8.49)$$

As far as the super-solution is concerned, we note that

$$\begin{aligned} u(\tau_n, x) &\leq \phi(x + \xi_n^+(\tau_n)) + q_n e^{-(t-T_n)} \leq \phi(x + \xi_n^+(\tau_n)) + c_q \delta (\xi_n^+ - \xi_n^-) e^{-T} \\ &\leq \phi(x + \xi_n^+(\tau_n)) + q_{n+1}, \end{aligned} \quad (1.8.50)$$

for all $x \in \mathbb{R}$, provided that T is sufficiently large, independent of n .

For the sub-solution, we first look at what happens for $|x| \leq R_0 + 1$ and use (1.8.46):

$$u(\tau_n, x) \geq \phi(x + \xi_n^+(\tau_n)) + q_n e^{-(\tau_n - T_n)} - \delta(\xi_n^+ - \xi_n^-), \quad \text{for all } |x| \leq R_0 + 1. \quad (1.8.51)$$

Thus, for $|x| \leq R_0 + 1$ we have

$$u(\tau_n, x) \geq \phi(x + \xi_n^+(\tau_n)) - \delta(\xi_n^+ - \xi_n^-) \geq \phi(x + \xi_n^+ - C_R \delta (\xi_n^+ - \xi_n^-)) \geq \phi(x + \xi_{n+1}^-), \quad (1.8.52)$$

with the constant C_R that depends on R_0 , as long as $\delta > 0$ is sufficiently small.

It remains to look at $|x| \geq R_0 + 1$. To this end, recall that

$$u(\tau_n, x) \geq \phi(x + \xi_n^-(\tau_n)) - q_n e^{-(\tau_n - T_n)}, \text{ for all } x \in \mathbb{R}, \quad (1.8.53)$$

so that, as follows from the definition of $\xi_n^-(t)$, we have

$$u(\tau_n, x) \geq \phi(x + \xi_n^- - Cq_n) - q_n e^{-2T}, \text{ for all } x \in \mathbb{R}. \quad (1.8.54)$$

Observe that, as $\phi(x)$ is approaching ± 1 as $x \rightarrow \pm\infty$ exponentially fast, there exist $\omega > 0$ and $C > 0$ such that, taking into account (1.8.41) we can write for $|x| \geq R_0 + 1$:

$$\begin{aligned} \phi(x + \xi_n^- - Cq_n) &\geq \phi(x + \xi_n^- + (\xi_n^+(\tau_n) - \xi_n^+) + \delta(\xi_n^+ - \xi_n^-)) - C\delta e^{-\omega R_0}(\xi_n^+ - \xi_n^-) \\ &\geq \phi(x + \xi_{n+1}^-) - q_{n+1}, \end{aligned} \quad (1.8.55)$$

as long as R_0 is large enough. Here, we have used ξ_{n+1}^- and q_{n+1}^- as in (1.8.47) and (1.8.49). We conclude that

$$u(\tau_n, x) \geq \phi(x + \xi_{n+1}^-) - q_{n+1}, \text{ for } |x| \geq R_0 + 1. \quad (1.8.56)$$

Summarizing, if (1.8.46) holds, we set $T_{n+1} = \tau_n$, define the new shifts ξ_{n+1}^\pm as in (1.8.47) and (1.8.49), which ensures that the "shift gap" is decreased by a fixed factor, so that (1.8.48) holds, and we can restart the argument at $t = T_{n+1}$, because

$$\phi(x + \xi_{n+1}^-) - q_{n+1} \leq u(T_{n+1}, x) \leq \phi(x + \xi_{n+1}^+) + q_{n+1}, \text{ for all } x \in \mathbb{R}. \quad (1.8.57)$$

Of course, if at some time $\tau_n \in [T_n + T, T_n + 2T]$ we have, instead of (1.8.46) that

$$\sup_{|x| \leq R_0 + 1} (u(\tau_n, x) - \underline{u}(\tau_n, x)) \leq \delta(\xi_n^+ - \xi_n^-), \quad (1.8.58)$$

then we could repeat the above argument essentially verbatim, using the fact that now the solution is very close to the sub-solution on a very large interval.

Case 2: the solution and the super-solution are never too close. Next, let us assume that for all $t \in [T_n + T, T_n + 2T]$, we have

$$\sup_{|x| \leq R_0 + 1} (\bar{u}_n(t, x) - u(t, x)) \geq \delta(\xi_n^+ - \xi_n^-). \quad (1.8.59)$$

Because $\xi_n^+(t)$ is increasing, we have, for all $|x| \leq R_0 + 1$ and $t \in [T_n + T, T_n + 2T]$:

$$\bar{u}_n(t, x) \leq \phi(x + \xi_n^+(T_n + 2T)) + q_n e^{-T} \leq \phi(x + \xi_n^+(T_n + 2T)) + q_n e^{-T} \rho_0, \quad (1.8.60)$$

with

$$\rho_0 = \left(\inf_{|x| \leq R_0 + M + 10} \phi'(x) \right)^{-1}. \quad (1.8.61)$$

Here, M is the constant in the upper bound (1.8.40) for ξ_n^\pm . Note that by choosing T sufficiently large we can make sure that the argument in ϕ in the right side of (1.8.60) is within the range of the infimum in (1.8.61). The function

$$w_n(t, x) = \phi(x + \xi_n^+(T_n + 2T)) + q_n e^{-T} \rho_0 - u(t, x).$$

that appears in the right side of (1.8.60) solves a linear parabolic equation

$$\partial_t w_n - \partial_{xx} w_n + a_n(t, x) w_n = 0, \quad (1.8.62)$$

with the coefficient a_n that is bounded in n , t and x :

$$a_n(t, x) = -\frac{f(\phi(x + \xi_n^+(T_n + 2T) + q_n e^{-T} \rho_0)) - f(u(t, x))}{\phi(x + \xi_n^+(T_n + 2T) + q_n e^{-T} \rho_0) - u(t, x)}. \quad (1.8.63)$$

It follows from assumption (1.8.59) and (1.8.60) that

$$\sup_{|x| \leq R_0+1} w_n(t, x) \geq \delta(\xi_n^+ - \xi_n^-), \quad \text{for all } t \in [T_n + T, T_n + 2T], \quad (1.8.64)$$

but in order to improve the shift, we would like to have not the supremum but the infimum in the above inequality. And here the Harnack inequality comes to the rescue: we will use Theorem 1.6.9 for the intervals $|x| \leq R_0 + 1$ and $|x| \leq R_0$. For that, we need to make sure that at least a fraction of the supremum in (1.8.64) is attained on $[-R_0, R_0]$: there exists k_1 so that

$$\sup_{|x| \leq R_0} w_n(t, x) \geq k_1 \delta(\xi_n^+ - \xi_n^-), \quad \text{for all } T_n + T \leq t \leq T_n + 2T. \quad (1.8.65)$$

However, if there is a time $T_n + T \leq s_n \leq T_n + 2T$ such that

$$\sup_{|x| \leq R_0} w_n(s_n, x) \leq \frac{\delta}{2}(\xi_n^+ - \xi_n^-), \quad (1.8.66)$$

then we have

$$\bar{u}(s_n, x) - u(s_n, x) \leq \frac{\delta}{2}(\xi_n^+ - \xi_n^-) \quad \text{for all } |x| \leq R_0. \quad (1.8.67)$$

This is the situation we faced in Case 1, and we can proceed as in that case. Thus, we may assume that

$$\sup_{|x| \leq R_0} w_n(t, x) \geq \frac{\delta}{2}(\xi_n^+ - \xi_n^-) \quad \text{for all } T_n + T \leq t \leq T_n + 2T. \quad (1.8.68)$$

In that case, we may apply the Harnack inequality of Theorem 1.6.9 to (1.8.62) on the intervals $|x| \leq R_0 + 1$ and $|x| \leq R_0$: there exists a Harnack constant h_{R_0} that is independent of T , such that

$$w_n(t, x) \geq h_{R_0} \delta(\xi_n^+ - \xi_n^-), \quad \text{for all } t \in [T_n + T + 1, T_n + 2T] \text{ and } |x| \leq R_0. \quad (1.8.69)$$

Exercise 1.8.7 Show that, as a consequence, we can find $\rho_1 > 0$ that depends on R_0 but not on n such that for $|x| \leq R_0$ and $T_n + T + 1 \leq t \leq T_n + 2T$, we have

$$\tilde{w}_n(t, x) = \phi(x + \xi_n^+(T_n + 2T) + \rho_0 e^{-T} q_n - \rho_1 h_{R_0} \delta(\xi_n^+ - \xi_n^-)) - u(t, x) \geq 0. \quad (1.8.70)$$

Let us now worry about what \tilde{w}_n does for $|x| \geq R_0$. In this range, the function \tilde{w}_n solves another linear equation of the form

$$\partial_t \tilde{w}_n - \partial_{xx} \tilde{w}_n + \tilde{a}_n(t, x) \tilde{w}_n = 0, \quad (1.8.71)$$

with $\tilde{a}_n(t, x) \geq 1$ that is an appropriate modification of the expression for $a_n(t, x)$ in (1.8.63). In addition, at the boundary $|x| = R_0$, we have $\tilde{w}_n(t, x) \geq 0$, and at the time $t = T_n + T$, we have an estimate of the form

$$\tilde{w}_n(T_n + T, x) \geq -K(\xi_n^+ - \xi_n^-), \quad |x| \geq R_0. \quad (1.8.72)$$

Exercise 1.8.8 What did we use to get (1.8.72)?

Therefore, the maximum principle applied to (1.8.71) implies that

$$\tilde{w}_n(T_n + 2T, x) \geq -K e^{-T}(\xi_n^+ - \xi_n^-), \quad |x| \geq R_0. \quad (1.8.73)$$

We now set $T_{n+1} = T_n + 2T$. The previous argument shows that we have

$$u(T_{n+1}, x) \leq \phi(x + \xi_n^+(T_{n+1}) + \rho_0 e^{-T} q_n - \rho_1 h_{R_0} \delta(\xi_n^+ - \xi_n^-)) + q_{n+1}, \quad (1.8.74)$$

with

$$0 \leq q_{n+1} \leq K e^{-T}(\xi_n^+ - \xi_n^-). \quad (1.8.75)$$

In addition, we still have the lower bound:

$$u(T_n + 2T) \geq \phi(x + \xi_n^-(T_{n+1})) - e^{-T} q_n. \quad (1.8.76)$$

It only remains to define ξ_{n+1}^\pm and q_{n+1} properly, to convert (1.8.74) and (1.8.76) into the form required to restart the iteration process. We take

$$q_{n+1} = \max(e^{-T} q_n, K e^{-T}(\xi_n^+ - \xi_n^-)), \quad \xi_{n+1}^- = \xi_n^-(T_{n+1}), \quad (1.8.77)$$

and

$$\xi_{n+1}^+ = \xi_n^+(T_{n+1}) + \rho_0 e^{-T} q_n - h_{R_0} \rho_1 \delta(\xi_n^+ - \xi_n^-). \quad (1.8.78)$$

It is easy to see that assumption (1.8.41) holds for q_{n+1} provided we take T sufficiently large, so that

$$e^{-T} \ll c_q. \quad (1.8.79)$$

The main point to verify is that the contraction in (1.8.38) does happen with the above choice. We recall (1.8.37):

$$\xi_n^+(T_{n+1}) = \xi_n^+ + \frac{q_n}{k_0}(1 + M_f)(1 - e^{-2T}), \quad \xi_n^-(T_{n+1}) = \xi_n^- - \frac{q_n}{k_0}(1 + M_f)(1 - e^{-2T}). \quad (1.8.80)$$

Hence, in order to ensure that

$$\xi_{n+1}^+ - \xi_{n+1}^- \leq (1 - \frac{h_{R_0} \rho_1 \delta}{2})(\xi_n^+ - \xi_n^-), \quad (1.8.81)$$

it suffices to make sure that the term $h_{R_0} \rho_1 \delta(\xi_n^+ - \xi_n^-)$ dominates all the other multiples of $\delta(\xi_n^+ - \xi_n^-)$ in the expression for the difference $\xi_{n+1}^+ - \xi_{n+1}^-$ that come with the opposite sign. However, all such terms are multiples of q_n , thus it suffices to make sure that the constant c_q is small, which, in turn, can be accomplished by taking T sufficiently large. This completes the proof. \square

1.8.2 Spreading in unbalanced Allen-Cahn equations, and related models

Let us now discuss, informally, what one would expect, from the physical considerations, to happen to the solution of the initial value problem if the balance condition (1.8.6) fails, that is,

$$\int_{-1}^1 f(u) du \neq 0. \quad (1.8.82)$$

To be concrete, let us consider the nonlinearity $f(u)$ of the form

$$f(u) = (u + 1)(u + a)(1 - u), \quad (1.8.83)$$

with $a \in (0, 1)$. so that $u = \pm 1$ are still the two stable solutions of the ODE

$$\dot{u} = f(u),$$

but instead of (1.8.6) we have

$$\int_{-1}^1 f(u) du > 0.$$

As an indication of what happens we give the reader the following exercises. They are by no means short but they can all be done with the tools of this section, and we strongly recommend them to a reader interested in understanding this material well.

Exercise 1.8.9 To start, show that for $f(u)$ given by (1.8.83), we can find a special solution $u(t, x)$ of the Allen-Cahn equation (1.8.2):

$$u_t = u_{xx} + f(u), \quad (1.8.84)$$

of the form

$$u(t, x) = \psi(x + ct), \quad (1.8.85)$$

with $c > 0$ and a function $\psi(x)$ that satisfies

$$c\psi' = \psi'' + f(\psi), \quad (1.8.86)$$

together with the boundary condition

$$\psi(x) \rightarrow \pm 1, \quad \text{as } x \rightarrow \pm\infty. \quad (1.8.87)$$

Solutions of the form (1.8.85) are known as traveling waves. Show that such c is unique, and ψ is unique up to a translation: if $\psi_1(x)$ is another solution of (1.8.86)-(1.8.87) with c_1 replaced c , then $c = c_1$ and there exists $x_1 \in \mathbb{R}$ such that $\psi_1(x) = \psi(x + x_1)$.

Exercise 1.8.10 Try to modify the proof of Lemma 1.8.4 to show that if $u(t, x)$ is the solution of the Allen-Cahn equation (1.8.84) with an initial condition $u_0(x)$ that satisfies (1.8.8):

$$u_0(x) \rightarrow \pm 1, \quad \text{as } x \rightarrow \pm\infty, \quad (1.8.88)$$

then we have

$$u(t, x) \rightarrow 1 \text{ as } t \rightarrow +\infty, \text{ for each } x \in \mathbb{R} \text{ fixed.} \quad (1.8.89)$$

It should be helpful to use the traveling wave solution to construct a sub-solution that will force (1.8.89). Thus, in the "unbalanced" case, the "more stable" of the two states $u = -1$ and $u = +1$ wins in the long time limit. Show that the convergence in (1.8.89) is not uniform in $x \in \mathbb{R}$.

Exercise 1.8.11 Let $u(t, x)$ be a solution of (1.8.84) with an initial condition $u_0(x)$ that satisfies (1.8.88). Show that for any $c' < c$ and $x \in \mathbb{R}$ fixed, we have

$$\lim_{t \rightarrow +\infty} u(t, x - c't) = 1, \quad (1.8.90)$$

and for any $c' > c$ and $x \in \mathbb{R}$ fixed, we have

$$\lim_{t \rightarrow +\infty} u(t, x - c't) = -1. \quad (1.8.91)$$

Exercise 1.8.12 Let $u(t, x)$ be a solution of (1.8.84) with an initial condition $u_0(x)$ that satisfies (1.8.88). Show that there exists $x_0 \in \mathbb{R}$ (which depends on u_0) so that for all $x \in \mathbb{R}$ fixed we have

$$\lim_{t \rightarrow +\infty} u(t, x - ct) = \psi(x + x_0). \quad (1.8.92)$$

1.8.3 When the medium is inhomogeneous: pulsating waves

We will be interested, in this final section, in equations of the form

$$u_t - u_{xx} = f(x, u), \quad t > 0, \quad x \in \mathbb{R}, \quad (1.8.93)$$

with f 1-periodic in the variable x . We will assume the following form f : there is $\theta \in (0, 1)$ such that $f(x, u) \equiv 0$ if $u < \theta$, and $f(x, u) > 0$ if $u > \theta$. In the vicinity of θ there holds

$$f_u(x, u) \geq \frac{f(x, u)}{u}. \quad (1.8.94)$$

Finally we assume that $f(x, 1) \equiv 0$, and that there is $\alpha > 0$ such that $f_u(x, 1) < -\alpha$ for all $x \in \mathbb{R}$. Of course, the set of such functions is by no means empty. For instance, (1.8.94) is true when

$$f(x, u) = a(x)(u - \theta)_+^p,$$

with $a(x) > 0$, 1-periodic, and p a sufficiently large integer.

We wish to understand the large time behavior of a solution $u(t, x)$ which, at time $t = 0$, tends to 1 as $x \rightarrow -\infty$ and to 0 as $x \rightarrow +\infty$. Clearly, the large time asymptotics cannot be given by a traveling wave, as the function f depends explicitly on x . It turns out that, for such a nonlinearity, there is a special class of solutions generalizing traveling waves, these are called pulsating waves. The reason for that is simple: these solutions will be shown to be periodic in a well chosen Galilian reference frame. We will also show that, in fact, they attract all the solutions that, initially, has the afore-mentioned behavior.

The notion of pulsating waves was introduced by Xin at the beginning of the 90's, see for instance the review paper [?]. It was much extended and generalized by Berestycki and Hamel, especially for models posed in the whole space, where things are more subtle than in the one-dimensional model that we present. See [?] for a much detailed account of the theory - that is, by the way, still evolving. We note that, even in one space dimension, relaxing the assumptions that we have (for instance, asking $f(x, 1) \equiv 0$ instead of a more natural looking assumption allowing uniform boundedness of the solutions) made may not necessarily modify the nature of the results that we are about to present, but it would certainly involve a good deal of additional work, which would not be in the spirit of this chapter. Once again, we refer to [?] for an account of what happens in the most general situations.

Exercise 1.8.13 Under the stated assumptions on f , the only solutions $\phi(x) \in [0, 1]$ of

$$-\phi'' = f(x, \phi), \quad x \in \mathbb{R}, \quad (1.8.95)$$

are all the constants between 0 and θ , and 1.

Exercise 1.8.14 Consider a nonlinearity $f(x, u)$ that satisfies all the above assumptions, except $f(x, 1) \equiv 0$. Assume, though, the existence of $f_0(x, u)$ that satisfies them all, and that is additionally close to f in the C^1 norm.

1. Show that the (1.8.95) has a unique minimal nonconstant solution $\phi_+(x)$, that is C^2 -close to 1.
2. Show that the first periodic eigenvalue of $-\partial_{xx} - f_u(x, \phi_+)$ is negative.

Why have we suddenly shifted from nonlinearities of the Allen-Cahn type to nonnegative nonlinearities of the afore-mentioned type? As the reader may have guessed in view of the previous exercise, we will construct waves that connect various solutions of (1.8.95), but we will not be necessarily very easy to count - a very interesting exercise of course, but once again outside the scope of this section. We may still propose the following exercise to the interested reader.

Exercise 1.8.15 Let $f(x, u)$ be C^2 -close to an unbalanced Allen-Cahn nonlinearity. Find all the solutions of (1.8.95), as well as, for each of them, the sign of $-\partial_{xx} - f_u(x, \phi_+)$. Hint: there is a catch, the limiting Allen-Cahn equation has nonconstant solutions!

These informal preliminaries being dealt with, let us now define precisely the object that will be in our preoccupations.

Definition 1.8.16 A function $\phi(t, x)$ is a pulsating wave of (1.8.93), with speed $c > 0$, and connecting 1 to 0 if it satisfies the following properties:

1. ϕ solves (1.8.93) of $\mathbb{R} \times \mathbb{R}$,
2. it is $1/c$ -periodic in a Galilean reference frame with speed c ,
3. we have

$$\lim_{x \rightarrow -\infty} \phi(t, x) = 1, \quad \lim_{t \rightarrow +\infty} \phi(t, x) = 0,$$

pointwise in t .

Let us make the following simple remarks.

Remark 1.8.17 • *The parabolic regularity and the time periodicity imply that the limits are in fact uniform in t , as soon as one looks at the phenomenon in the reference frame with speed c .*

- *For a pulsating wave with speed c , the time-periodicity $1/c$ is the only possible one. This forced by the 1-periodicity in x : in the reference frame with speed ct , the function $f(x, \phi)$ becomes $f(x+ct, \phi)$. If ϕ is T -periodic, then $f(x+ct, \cdot)$ should also be T -periodic. Another way to view it is the following: the speed being c , the wave takes the time $1/c$ to cover the cell of length 1. When this is achieved, it retrieves its original shape.*
- *If $\phi(t, x)$ is a pulsating wave with speed c , then any translate in time $\phi(t + t_0, x)$ is also a pulsating wave with speed c .*

And we may state the main achievement of the section, namely the existence and uniqueness of pulsating waves.

Theorem 1.8.18 *Problem (1.8.93) has a one-dimensional family of pulsating waves $\phi(t, x)$ (one can be deduced from another by a translate in t) with speed c , connecting 1 to the left to 0 to the right. The speed c is unique. moreover we have $\partial_t \phi > 0$.*

The last statement is a striking parallel with the main property of the traveling waves, namely that they are decreasing in x . Of course, here, monotonicity in x is not true, what replaces it is monotonicity on t .

Monotonicity and uniqueness are very easily proved by the sliding ideas that we have exposed at length in the first chapter, it is therefore a good time to propose a last refresh to the reader. The idea is the same in both cases: take two different waves, and prove that a sufficiently large translate of one is below the other. Then, translate back until it is not possible anymore, and derive a contradiction. Let us give some more details and prove monotonicity first. let ϕ be a wave with speed $c > 0$, we infer, simply by the fact that they have limits as $x \rightarrow \pm\infty$, the existence of a large $T > 0$ such that

$$\phi(t + T, x) \leq \phi(t, x) \text{ for all } (t, x) \in [0, \frac{1}{c}] \times [-M, M],$$

and we may choose $M > 0$ such that

$$\phi(t, x) \leq \frac{\theta}{2} \text{ for } t \in [0, \frac{1}{c}], x \geq M, \quad \phi(t, x) \geq 1 - \delta \text{ for } t \in [0, \frac{1}{c}], x \leq -M,$$

with $f_u(x, u) \leq \frac{\alpha}{2}$ for $u \geq 1 - \delta$. In the reference frame moving with velocity c , both the wave (still denoted by ϕ) and its translate solve:

$$\partial_t \phi - c \partial_x \phi - \partial_{xx} \phi = f(x - ct, \phi), \tag{1.8.96}$$

and $\psi(t, x) : \phi(t + T, x) - \phi(t, x)$ solve

$$\partial_t \psi - c \partial_x \psi - \partial_{xx} \psi - a(t, x) \psi = 0, \tag{1.8.97}$$

where $a(t, x)$ is, as usual, some convex combination of $f_u(x - ct, \phi)$. Note that now, both ϕ and its translate are $1/c$ -periodic in t . Therefore, $\psi(t, x) \geq 0$ for $t \geq 0$, $-M \leq x \leq M$. For $x \leq -M$ we have $a(t, x) \geq \alpha/2$, while for $x \geq M$ we have an inequation of the form

$$\partial_t \psi - c \partial_x \psi - \partial_{xx} \psi \geq 0. \quad (1.8.98)$$

Sending $t \rightarrow +\infty$ allows us to infer, from (1.8.97) and (1.8.98):

$$\liminf_{t \rightarrow +\infty} \psi(t, x) \geq 0.$$

Exercise 1.8.19 Work out the details by placing a subsolution below ψ on $\mathbb{R}_+ \times (-\infty, -M)$ and $(M, +\infty)$. Nothing difficult here, the only intermediate step is to prove that the solution of the Dirichlet advection-diffusion equation

$$\begin{aligned} v_t - v_{xx} - cv_x &= 0, & t > 0, & x \geq M \\ v(t, 0) &= 0 \\ \lim_{t \rightarrow +\infty} v(0, x) &= 0 \end{aligned} \quad (1.8.99)$$

tends to 0 as t increases to infinity, uniformly in x . One may proceed as follows.

1. Show that it is enough to assume $v(0, x) \geq 0$ and to prove that the lim sup is zero.
2. Construct a super-solution of the form $e^{-\beta t - \gamma x}$, give explicit values to β and γ .
3. Prove the result if $v(0, x)$ is compactly supported.
4. Show that (1.8.99) generates a weakly contracting semigroup, that is: $\|v_1(t, \cdot) - v_2(t, \cdot)\|_\infty \leq \|v_1(0, \cdot) - v_2(0, \cdot)\|_\infty$.
5. Conclude.

In the area $\{t > 0, x \geq -M\}$, use the fact that $a(t, x) \leq -\alpha/2$.

Sending t to infinity and remembering the $1/c$ -periodicity of ψ allows us to infer that, actually, we have $\phi(t + T, x) \geq \phi(t, x)$ everywhere. And so, there is $t_{min} \geq 0$ such that $\phi(t + t_{min}, x) \leq \phi(t, x)$ everywhere. Assume $t_{min} > 0$, the inequality has to be strict at at least one point, otherwise we would have $\phi(t + t_{min}, x) = \phi(t, x)$, something that the positive speed of propagation as well as the limits as $x \rightarrow \pm\infty$ oppose vehemently. But then, the \leq sign should be replaced by a $<$ sign, just to appease the strong maximum principle. So, on a large compact set that we still call $[-M, M]$, we may translate a little more and obtain, for small $\delta > 0$:

$$\phi(t + t_{min} - \delta, x) \leq \phi(t, x), \quad t \in \mathbb{R}, \quad -M \leq x \leq M.$$

The previous argument can be repeated, so that the inequality holds in fact on $\mathbb{R} \times \mathbb{R}$. This contradicts the minimality of t_{min} , ensuring $\phi(t + \delta, x) \geq \phi(t, x)$ for all $\delta > 0$. This entails the monotonicity of the wave, and the same argument may be used to show uniqueness.

The uniqueness of the speed is hardly any more difficult. Let c_1 and c_2 be two potential wave speeds, write down (1.8.93) in the reference frame moving with speed c_1 . Translate ϕ_1 enough so that it is below ϕ_2 in a large compact set, then use the following

Exercise 1.8.20 Construct a sub-solution to (1.8.96) under the form

$$\underline{\phi}_1(t, x) = \phi_1(t + \xi(t), x) - qe^{-cx/2 - c^2t/4}\Gamma(x) - qe^{-\alpha t/2}(1 - \Gamma(x - x_1)),$$

where Γ is a smooth function that is equal to 1 on \mathbb{R}_- and 0 on $[1, +\infty)$. The constant x_1 should be large, and the constant q should be small. The inspiration should be taken from Lemma 1.8.4. The monotonicity in x of the traveling wave should now be replaced by the monotonicity in t of the pulsating wave.

Clearly, we may choose the constants to have $\underline{\phi}_1 \leq \phi_2$. Sending $t \rightarrow +\infty$ and using the $1/c$ periodicity yields $\phi_1(t + t_0, x) \leq \phi_2(t, x)$, which implies $c_1 \geq c_2$. And so, $c_1 = c_2$ by symmetry.

It is now a good time to explain why pulsating waves exist. The construction that we are going to provide relies, once again, on rather simple comparison arguments. The idea is to solve the Cauchy Problem for (1.8.93), with a sub-solution as an initial datum. Thus we will obtain, for the resulting solution, monotonicity for free. We will then examine its behavior as $|x|$ becomes very large, and see that it has the correct limits, and that they are taken uniformly with respect to t . This can be viewed as a finite thickness property for the front, and it will provide a sufficient amount of compactness for us to prove that the limiting solution is nontrivial. A last effort, where we will again use comparison and sliding, will tell us that we have put our hand on the sought for pulsating wave.

The starting point is therefore the construction of a sub-solution. Pick $\theta_1 \in (\theta, 1)$ close enough to 1, $\beta > 0$ small to be chosen later, and $\alpha > 0$ small so that we have

$$f(x, u) \geq \alpha^2(1 - u), \quad u \in [\theta_1, 1],$$

and define

$$\underline{f}(u) = \begin{cases} \alpha^2(1 - u), & u \in [\theta_1, 1] \\ -\beta^2 u, & u \in [0, \theta_1]. \end{cases}$$

Hence $\underline{f}(u) \leq f(x, u)$. A solution $v(x)$ of

$$-v'' = \underline{f}(v), \quad \lim_{x \rightarrow -\infty} v(x) = 1$$

is easily computed: we have

$$v(x) = \frac{1}{2} \left(\theta_1 + \frac{1 - \theta_1}{\beta} \right) e^{-\beta x} + \frac{1}{2} \left(\theta_1 - \frac{1 - \theta_1}{\beta} \right) e^{\beta x}.$$

The sought for sub-solution is simply chosen as $\underline{u}(x) = v^+(x)$. Solve (1.8.93) with $u(0, x) = \underline{u}(x)$: we obtain a solution $u(t, x)$ that satisfies $\partial_t u \geq 0$, moreover $u(t, x)$ assumes the correct limits as $x \rightarrow \pm\infty$; the trouble is, however, that there is no uniformity at this stage. Our main tool will be the following

Proposition 1.8.21 *There is a function $\lambda \mapsto q(\lambda)$, bounded and bounded away from 0 on every compact of $[0, 1)$, such that we have*

$$u_t(t, x) \geq q(\lambda)u(t, x) \text{ if } u(t, x) = \lambda. \quad (1.8.100)$$

Proof. To see that it is true on every compact set of $(\theta, 1)$ is not so difficult: pick $\lambda \in (\theta, 1)$ and assume that there is a sequence (t_n, x_n) such that $u(t_n, x_n) = \lambda$ and $\lim_{n \rightarrow +\infty} u_t(t_n, x_n) = 0$. Because $u_t \geq 0$, the sequence of functions

$$u_n(t, x) = u(t + t_n, x + x_n)$$

converges, uniformly on every compact, and up to the extraction of a subsequence, to a function $u_\infty(t, x)$ that satisfies (1.8.93). Because of the Harnack inequality we have $(\partial_t u_n)_n$ converges to 0, so that u_∞ does not depend on x , moreover $u_\infty(0, 0) = \lambda$. The positivity of f yields a contradiction. Thus, the proposition is true for this range of λ .

Now, pick $\delta > 0$ small, and let us study what happens for the remaining values of λ . Define

$$\Omega = \{(t, x) \in \mathbb{R} \times \mathbb{R} : u(t, x) < \theta + \delta\}.$$

Obviously, (1.8.100) holds true on

$$\{(t, x) \in \mathbb{R} \times \mathbb{R} : u(t, x) = \theta + \delta\},$$

provided $q(\theta + \delta)$ is chosen appropriately. The inequality also holds at $t = 0$, by the definition of \underline{f} , with $q(\lambda) = \beta^2$ for all $\lambda \in [0, \theta]$. For short, we set

$$q = \inf_{\lambda \in [0, \theta + \delta]} q(\lambda).$$

The function $v(t, x) = u_t(t, x) - qu(t, x)$ solves, in Ω :

$$\begin{aligned} v_t - v_{xx} &= f_u(x, u)v + (f_u(x, u) - \frac{f(x, u)}{u})u \\ &\geq f_u(x, u)v \text{ because of (1.8.94).} \end{aligned}$$

Thus, $v(t, x) \geq 0$ in Ω , what we wanted to prove. \square

Proposition 1.8.21 is a very important property that will allow us to conclude almost effortlessly. Let us now denote, this time, Ω the set $\{u < \theta\}$, and Γ its boundary within $\mathbb{R}_+^* \times \mathbb{R}$; because $u_t \geq q\theta$ on Γ it is a smooth curve $\{(\tau(x), x), x \in \mathbb{R}\}$. We choose, once and for all:

$$q = \inf\{q(\lambda), 0 \leq \lambda \leq \frac{1 + \theta}{2}\}. \quad (1.8.101)$$

We may always assume that τ is defined on $[0, +\infty)$, with $\tau(0) = 0$. The main consequence of Proposition 1.8.21 is the following statement, which says that the front has finite thickness.

Corollary 1.8.22 *The level set Γ may be described as a smooth curve of the form $(t, X(t))$, with $\dot{X}(t) \geq \mu > 0$.*

$$\lim_{x-X(t) \rightarrow +\infty} u(t, x) = 0, \quad \lim_{x-X(t) \rightarrow -\infty} u(t, x) = 1, \quad (1.8.102)$$

uniformly with respect to $t \in \mathbb{R}_+$.

Proof. For small $t > 0$, say, $t \in [0, t_0]$, continuity implies that $u_x(t, x) > 0$ when $(t, x) \in \Gamma$. This implies, for every $t \in [0, t_0]$, the existence and uniqueness of the function $X(t)$ defined

as in the proposition. And, for any t in this interval, the equation for u to the left of $X(t)$ becomes

$$-u_{xx} + qu \leq u_t - u_{xx} = 0,$$

so that $u(t, x) \leq \theta e^{-\sqrt{q}(x-X(t))}$, and $u_x(t, X(t)) \leq -\theta\sqrt{q}$. What would make this beautiful estimate break down would be a time $t_{max} \geq t_0$ and a point x_{max} such that $t_{max} = \tau(x_{max})$ and $\tau'(x_{max}) = 0$, in other words $u_x(t_{max}, x_{max}) = 0$. However, at time t_{max} , x_{max} is still the rightmost point (or, at least, can be made the rightmost one if τ' vanishes on an interval) x such that $u(t_{max}, x) = \theta$; this makes the previous argument work, and, in particular, $u_x(t_{max}, x_{max}) \leq -\theta\sqrt{q}$. This is an impossibility. Consequently, the function $X(t)$ is defined for all time, moreover we have

$$\dot{X}(t) = -\frac{u_t(t, X(t))}{u_x(t, X(t))} \in \left[\frac{q}{\|u_x\|_\infty}, \frac{\|u_t\|_\infty}{\inf_\Gamma u_x} \right], \quad (1.8.103)$$

and each of the above quantities is a positive constant. Notice that we have also ruled out a scenario where a hole of points where u would drop below θ would appear to the left of Γ , simply because $u_t > 0$.

Our argument also shows that $u(t, x)$ goes to 0 exponentially fast as $x - X(t)$ becomes negative, so the only thing that should worry us now is the limit $x - X(t) \rightarrow +\infty$. For this, it is convenient to write the equation for u in the reference frame following $X(t)$:

$$u_t - u_{xx} - \dot{X}(t)u_x = f(x + X(t), u), \quad (1.8.104)$$

by the Implicit function Theorem there is $\delta > 0$ and $\theta_1 \in (\theta, 1)$ such that

$$u(t, x) \geq \theta_1 \text{ for } x \leq -\delta.$$

The assumptions implies the existence of a small $\alpha > 0$ such that

$$f(x, u) \geq \alpha^2(1 - u) \text{ for } u \geq \theta_1. \quad (1.8.105)$$

The maximum principle implies that $1 - u(t, x) \leq v(t, x)$ on $\mathbb{R}_+ \times (-\infty, \delta]$, where

$$\begin{aligned} v_t - v_{xx} - \dot{X}(t)v_x + \alpha^2v &= 0, & t > 0, x < -\delta \\ v(t, -\delta) &= 1 - \theta_1 \\ v(0, x) &= 1 - \underline{u}(x). \end{aligned} \quad (1.8.106)$$

The function $\bar{v}(x) = (1 - \theta_1)e^{\varepsilon x}$ is a super-solution to (1.8.106) as soon as $0 < \varepsilon \leq 2\sqrt{\alpha}$, which implies the uniform exponential decay to 0 of $1 - u$. This proves the corollary. \square

Exercise 1.8.23 At the beginning of the proof, we merrily said "continuity implies that $u_x(t, x) > 0$ when $(t, x) \in \Gamma$ ". Make this continuity argument a little more explicit.

To construct the pulsating wave, we send t to $+\infty$; more precisely, we consider a sequence $(t_n)_n$ going to $+\infty$ such that the sequence

$$u_n(t, x) = u(t_n + t, [X(t_n)] + x)$$

converges, locally on every compact, to a function that we denote $\phi(t, x)$, as well as all its derivatives. We still denote by $X(t)$ the θ -level set of ϕ , we note that $\phi(t, x)$ and $X(t, x)$ enjoy all the properties listed in Proposition 1.8.21 and its corollary. Moreover, ϕ is now defined on the whole plane \mathbb{R}^2 . It remains to see that it is the sought for pulsating wave, and this is where sliding will make a last appearance. In fact, our argument will be quite similar to the one we used to prove the Krein-Rutman Theorem. We claim the existence of a large $T > 0$ such that

$$\phi(t + T, x) \geq \phi(t, x - 1), \text{ for all } (t, x) \in \mathbb{R}^2. \quad (1.8.107)$$

For all $M > 0$, we set

$$\Omega_M = \{(t, x) \in \mathbb{R}^2 : -M \leq x - X(t) \leq M\}.$$

From Corollary 1.8.22, we may choose $M > 0$ large enough so that $u(t, x) = 0$ if (t, x) is at the left of Ω_M , whereas $u(t, x) \geq \theta_1$ if (t, x) is at the right of Ω_M . This being done, we notice that indeed, there is a large $T > 0$ such that 1.8.107 holds on Ω_M . To see that the inequality holds everywhere, notice that an inequation for $v(t, x) = \phi(t + T, x) - \phi(t, x - 1)$ in the reference frame of $X(t)$ is

$$v_t - v_{xx} - \dot{X}(t)v_x + a(t, x) \geq 0,$$

with $a(t, x) = 0$ if $x \geq M$, and $a(t, x) \geq \alpha^2$ if $x \leq -M$. Notice finally the existence of $\mu > 0$ such that $\dot{X}(t) \geq \mu$. This said, pick any time $t \in \mathbb{R}$, we may assume for simplicity that $t = 0$, just by translation. At the time $-n$ (n is any integer) there is $\varepsilon > 0$ independent of n such that

$$u(-n, x) \leq \theta e^{-\varepsilon(x-M)} \text{ for } x \geq M,$$

and

$$1 - u(-n, x) \leq \theta_1 e^{\varepsilon(x+M)} \text{ for } x \leq -M.$$

Exercise 1.8.24 Show the existence of a constant $K > 0$ and $\delta > 0$ such that

$$v(t, x) \geq -Ke^{-\delta(t+n)} \text{ for } t \geq -n \text{ and } x \leq -M,$$

whereas

$$v(t, x) \geq -Ke^{-\varepsilon x/2 - \delta(t+n)} \text{ for } t \geq -n \text{ and } x \geq M.$$

The second estimate is the one where $\dot{X}(t) \geq \mu$ is needed.

Sending n to $+\infty$ reveals that $v(0, x) \geq 0$, so that the T we have defined satisfies (1.8.107). The reader who has accepted to follow us up to this point will have no problem with this penultimate hurdle:

Exercise 1.8.25 Call $1/c$ the smallest T such that (1.8.107) holds. Then, for $T = 1/c$, the inequality becomes an equality.

Simply by uniqueness for the Cauchy Problem, we have

$$\phi(t + 1/c, x) = \phi(t, x - 1).$$

Thus, ϕ is $1/c$ -periodic in a Galilean reference frame with speed c , and we have found our pulsating wave.

Exercise 1.8.26 Prove that any solution of the Cauchy Problem (1.8.93) converges, exponentially fast in time, to a pulsating wave, provided that (i) $u(0, x)$ converges to 1 as $x \rightarrow \infty$, and (ii) $u(0, x)$ converges exponentially fast in x to 0 as $x \rightarrow +\infty$. It will be useful to write the equation in the reference frame moving like $X(t)$, and to recycle the proof of Theorem 1.8.2 with the family of sub and super-solutions defined in Exercise 1.8.20. If you are challenged by the task, look at [?].

With this last result, we are leaving the fascinating world of invasion in reaction-diffusion models. We are not leaving the topic forever, we will meet it again in a few chapters, with equations that look very similar, but where both the mechanisms and the tools used for their investigation will be very different from those displayed here. As always, we do not claim that we have exhausted the subject, as multi-dimensional models - and also some questions in one dimensions - require arguments that are sometimes much more sophisticated than those presented here. However, the material of this chapter will be a good basis to a reader who will want to attack the numerous open questions of the field.

Chapter 2

Inviscid Hamilton-Jacobi equations

2.1 Introduction to the chapter

We will consider in this chapter the Hamilton-Jacobi equations

$$u_t + H(x, \nabla u) = 0 \tag{2.1.1}$$

on the unit torus $\mathbb{T}^n \subset \mathbb{R}^n$, or, sometimes, in all of \mathbb{R}^n . Note that here, unlike in the viscous Hamilton-Jacobi equations we have considered in Chapter 1, the diffusion coefficient vanishes. There are two reasons to do that in this book, where diffusion is remarkably present everywhere. The first is to emphasize some of the difficulties and phenomena that one encounters in the absence of diffusion. Another is that, as we will see, a physically reasonable class of solutions to (2.1.1) behave very much like the solutions to a regularized problem

$$u_t^\varepsilon + H(x, \nabla u^\varepsilon) = \varepsilon \Delta u^\varepsilon, \tag{2.1.2}$$

with a small diffusivity $\varepsilon > 0$. Most of the techniques we have introduced so far rely on the positivity of the diffusion coefficient and will deteriorate badly when the diffusion coefficient is small. However, we will see that some of the bounds may survive even as the diffusion term vanishes, because they are helped by the nonlinear Hamiltonian $H(x, \nabla u)$. Obviously, not every nonlinearity is beneficial: for example, solutions to the linear advection equation

$$u_t + b(x) \cdot \nabla u(x) = 0, \tag{2.1.3}$$

are typically no better than the initial condition $u_0(x) = u(0, x)$, no matter how smooth the drift $b(x)$ is. Therefore, we will have to restrict ourselves to some class of Hamiltonians $H(x, p)$ that do help to regularize the problem. This nonlinear regularization effect is one of the main points of this chapter.

The organization of the chapter

We begin with an informal derivation of the Hamilton-Jacobi equations in Section 2.2. Next, we discuss in Section 2.3 a class of viscous Hamilton-Jacobi equations, that is, equations of the form (2.1.1) with a Laplacian added:

$$u_t + H(x, \nabla u) = \Delta u. \tag{2.1.4}$$

Armed with the knowledge gathered in Chapter 1, we will (not so easily, but also not at an immense cost) construct a particular type of solutions, that we will call wave solutions. Quite similarly to what happens for traveling waves for the parabolic equations, these waves will be unique up to a translation in time, and solutions to the Cauchy problem will converge exponentially fast to them. In the remainder of the chapter, we will keep these features in mind to investigate how far the behavior of solutions to the inviscid Hamilton-Jacobi equations deviates from this simple picture.

We will then go through the most naive approach, looking for the smooth solutions to (2.1.1) in Section 2.4. However, a reader familiar with the theory of conservation laws, would see immediately the connection between them and the Hamilton-Jacobi equations: in one dimension, $n = 1$, differentiating (2.1.1) in x , we get a conservation law for $v = u_x$:

$$v_t + (H(x, v))_x = 0. \quad (2.1.5)$$

The basic one-dimensional conservation laws theory tells us that it is reasonable to expect that $v(t, x)$ becomes discontinuous in x at a finite time t , forming a shock, which means that we can not hope that solutions to the inviscid Hamilton-Jacobi equations are better than Lipschitz continuous generally. We will say a few words about the classical theory and explain why it breaks down very quickly. This is well-known, see for instance [?], where it is done very nicely. For the reader's convenience, we recall the basics here. This somehow pessimistic message should, however, be softened: there are (perhaps, less well-known) instances where a nice classical theory can be developed, and we are going to discuss one such example here. This material will, hopefully, be a good introduction to the more abstract theory to come next.

We proceed with a discussion of the viscosity solutions of the first order Hamilton-Jacobi equations in Section 2.5. Similarly to the parabolic regularity theory in Chapter 1, it is impossible to give a reasonable overview of the state of the art in this field. Rather, we will focus on its two most elementary aspects: first, that a viscosity solution is a solution obtained by sending a diffusion coefficient (viscosity) to zero, and, second, that the viscosity solutions exhibit the power of the comparison principle, something we have very much seen in the previous chapters. Our main goal will be to convince the reader that, although the viscous terms will have disappeared from the equations, some nontrivial features remain, such as the large time convergence to a steady state. One may call this the Cheshire cat smile effect [?]. This is explored, once again, in stages, where we first give a relatively accurate account of the Cauchy problem without dwelling too much on technicalities.

With the solution theory of Sections 2.4 and 2.5 in hand, one may start looking for the long time behavior of the solutions we have constructed, and their convergence to plane waves. The first step in this direction, as in the viscous case, is to construct the wave solutions, and consider the stationary version of (2.1.1):

$$H(x, \nabla u) = c, \quad x \in \mathbb{T}^n. \quad (2.1.6)$$

This we will be done in Section 2.6. After what we will have done in Section 2.3, it should be clear to the reader why (2.1.6) has a constant c in the right side – solutions to (2.1.6) lead to the wave solutions for the time-dependent problem (2.1.1). As in the viscous case, we will prove that under reasonable assumptions, solutions to (2.1.6) exist only for a unique value

of c which has no reason to be equal to zero. Thus, the “standard” steady equation

$$H(x, \nabla u) = 0$$

typically would have no solutions.

The case of a strictly convex Hamiltonian is quite interesting, and has strong connections with dynamical systems. We are going to dwell on it in Section 2.7, and show surprising regularizing properties that are not due to diffusion anymore but to the convexity of the Hamiltonian. After that, we will come back to the large time behavior of the solutions to the Cauchy problem, in Sections 2.8 and 2.9. We will first settle on a particular, and important class of equations in Section 2.8, for which we will prove, just as in the viscous case, the convergence to a wave. Alas, even though the speed c is unique, we will lose the uniqueness of the profile of the steady solutions – unlike in the diffusive case, (2.1.6) may have non-unique solutions, even up to a translation. We are going to investigate this phenomenon, that may be considered as the second main point of this chapter, in some detail. This major difference with the diffusive Hamilton-Jacobi equations will not be enough to prevent the large time convergence, but will force us to find a selection mechanism that will make up for the loss of diffusion.

In the last section, we use these new ideas to explain that, in fact, convergence to a wave holds for general equations of the form (2.1.1), as long as the Hamiltonian H is strictly convex in its second variable. In order to achieve this objective, we will (although we do not pretend to give a comprehensive treatment of this vast subject, that is still evolving at the time of the writing) give a reasonably comprehensive view of the issues posed by these deceptively simple models.

2.2 An informal derivation of the Hamilton-Jacobi equations

We begin by providing an informal derivation of the Hamilton-Jacobi equations, in the spirit of what we have done in Section 1.2 for the linear and semi-linear parabolic equations. The material of this section will reappear in Section 2.7 in the form of the Lax-Oleinik formula for the solutions to the Hamilton-Jacobi equations.

As in Section 1.2, we start with a random walk on a lattice of size h in \mathbb{R}^n , and a time step τ . The walker evolves as follows. If the walker is located at a position $X(t) \in h\mathbb{Z}^n$ at a time $t = m\tau$, $m \in \mathbb{N}$, then at a time $t + \tau$ it finds itself at a position

$$X(t + \tau) = X(t) + v(t)\tau + h\xi(t). \tag{2.2.1}$$

Here, $\xi(t) \in \mathbb{R}^n$ is an \mathbb{R}^n -valued random variable such that each of the coordinates $\xi_k(t)$, with $k = 1, \dots, n$, are independent and take the values ± 1 with probabilities equal to $1/2$, so that

$$\mathbb{E}(\xi_k(t)) = 0, \quad \mathbb{E}(\xi_k(t)\xi_m(t')) = \delta_{km}\delta_{t,t'}, \tag{2.2.2}$$

for all $1 \leq k, m \leq n$ and all t, t' . The velocity $v(t)$ is known as a control, that the walker can choose from a set \mathcal{A} of admissible velocities. The choice of the velocity v on the time

interval $[t, t + \tau]$ comes with a cost $L(v)\tau$, where $L(v)$ is a prescribed cost function. At the terminal time $T = N\tau$ the walker finds itself at a position $X(T)$ and pays the terminal cost $f(X(T))$, where $f(x)$ is also a given function. The total cost of the trajectory that starts at a time $t = m\tau$ at a position x and continues until the time $T = N\tau$ is

$$w(t, x; V) = \sum_{k=m}^N L(v(k\tau))\tau + f(X(N\tau)). \quad (2.2.3)$$

Note that the total cost involves both the running cost and the terminal cost. We have denoted here by $V = (v(t), v(t + \tau), \dots, v((N - 1)\tau))$ the whole sequence of the controls (velocities) chosen by the walker between the times t and T . The quantity of interest is the least possible average cost, optimized over all choices of the velocities:

$$u(t, x) = \inf_{V \in \mathcal{A}_{t,T}} \mathbb{E} w(t, x; V) = \inf_{V \in \mathcal{A}_{t,T}} \mathbb{E} \sum_{k=m}^N L(v(k\tau))\tau + f(X(N\tau)). \quad (2.2.4)$$

Here, the expectation \mathbb{E} is taken with respect to the random variables $\xi(s)$, for all $s = k\tau$ with $m \leq k < N$ that describe the random contribution at each of the time steps between t and T . The set $\mathcal{A}_{t,T}$ is the set of all possible controls chosen between the times $t = m\tau$ and $T = N\tau$. The velocities v are viewed as not random, as they can be chosen by the walker. The function $u(t, x)$ is known as the value function and is the basic object of study in the control theory.

As the velocities $v(s)$ are chosen separately by the walker at each time s between t and T , and the random variables $\xi(s)$ and $\xi(s')$ are independent for $s \neq s'$, the function $u(t, x)$ satisfies the following relation:

$$u(t, x) = \inf_{v \in \mathcal{A}} \mathbb{E} [L(v)\tau + u(t + \tau, x + v\tau + h\xi(t))]. \quad (2.2.5)$$

This is the simplest version of a dynamic programming principle, a fundamental notion of the control theory. Here, v is the velocity chosen at the initial time t and the expectation is taken solely with respect to the random variable $\xi(t)$.

A version of the dynamic programming principle, such as (2.2.5), is a very common starting point for the derivation of the Hamilton-Jacobi and other related types of equations. To illustrate this idea, let us assume that $u(t, x)$ is a sufficiently smooth function. Expanding the right side of (2.2.5) in $h \ll 1$ and $\tau \ll 1$ gives

$$\begin{aligned} u(t, x) &= \inf_{v \in \mathcal{A}} \mathbb{E} [L(v)\tau + u(t + \tau, x + v\tau + h\xi(t))] = u(t, x) + \tau u_t + \frac{\tau^2}{2} u_{tt}(t, x) \\ &+ \inf_{v \in \mathcal{A}} \mathbb{E} \left[L(v)\tau + (v\tau + h\xi(t)) \cdot \nabla u(t, x) + \tau(v\tau + h\xi(t)) \cdot \nabla u_t(t, x) \right. \\ &\left. + \frac{1}{2} \sum_{i,j=1}^n (v_i\tau + h\xi_i(t))(v_j\tau + h\xi_j(t)) \frac{\partial^2 u(t, x)}{\partial x_i \partial x_j} \right] + l.o.t. \end{aligned} \quad (2.2.6)$$

Note that the terms of the order $O(1)$ in the left and the right sides of (2.2.6) cancel automatically. In addition, the terms that are linear in $\xi(t)$ vanish after taking the expectation.

It is easy to see then that, as in the random walk approximation of the diffusion equations we have encountered in Section 1.2, the interesting choice of the temporal and spatial steps τ and h is

$$h^2 = 2D\tau, \tag{2.2.7}$$

with a diffusion coefficient D . Then, after taking into account the aforementioned cancellations, the leading order terms in (2.2.6) are of the order $O(\tau) = O(h^2)$. Keeping in mind (2.2.2), we see that they combine to give the following equation for $u(t, x)$:

$$u_t(t, x) + \inf_{v \in \mathcal{A}} [L(v) + v \cdot \nabla u(t, x)] + D\Delta u(t, x) = 0. \tag{2.2.8}$$

Let us introduce the function

$$H(p) = \inf_{v \in \mathcal{A}} [L(v) + v \cdot p], \tag{2.2.9}$$

defined for $p \in \mathbb{R}^n$. Then (2.2.8) can be written as

$$u_t + H(\nabla u) + D\Delta u = 0. \tag{2.2.10}$$

This equation should be supplemented by the terminal condition $u(T, x) = f(x)$ that comes from the definition of the value function. Recall that $f(x)$ is the terminal cost function.

Equation (2.2.10) is backward in time. It is convenient to define the function

$$\bar{u}(t, x) = u(T - t, x),$$

which satisfies the forward in time Cauchy problem:

$$\begin{aligned} \bar{u}_t &= H(\nabla \bar{u}) + D\Delta \bar{u}, \quad t > 0, \\ \bar{u}(0, x) &= f(x), \end{aligned} \tag{2.2.11}$$

and for the sake of convenience we will focus on this forward in time Cauchy problem.

This is how the viscous Hamilton-Jacobi equations can be derived informally. Their rigorous derivation starting with a continuous in space and time stochastic control problem is not very different but requires the use of the stochastic calculus and the Ito formula. The inviscid equations of the form

$$u_t = H(\nabla u), \tag{2.2.12}$$

are derived in a very similar way but the random walk is taken to be purely deterministic, driven solely by the control v , with $\xi(t) = 0$.

Exercise 2.2.1 Generalize the above derivation to obtain a spatially inhomogeneous Hamilton-Jacobi equation of the form

$$u_t = H(x, \nabla u) + D\Delta u. \tag{2.2.13}$$

Exercise 2.2.2 Show that the function $H(p)$ defined in (2.2.9) is concave.

This exercise explains why we will often consider below the Hamilton-Jacobi equations of the form

$$u_t + H(x, \nabla u) = D\Delta u, \tag{2.2.14}$$

with a convex Hamiltonian $H(p)$, either with $D > 0$ or $D = 0$.

2.3 The simple world of viscous Hamilton-Jacobi equations

As a warm-up to the chapter, we are going to use the knowledge gathered in Chapter 1 for the study of the long time behavior of the solutions to the viscous Hamilton-Jacobi equations. This problem falls in the same class as what we did in Section 1.8.1, where we proved, essentially with the sole aid of the strong maximum principle and the Harnack inequality, the convergence of the solutions to the Cauchy problem for the Allen-Cahn equation to a translate of a stationary solution. The main difference is that now we will have to fight a little to show the existence of a steady state, while the long time convergence will be relatively effortless.

We are interested in the large time behavior of the solutions to the Cauchy problem

$$u_t - \Delta u = H(x, \nabla u), \quad t > 0, \quad x \in \mathbb{T}^n, \quad (2.3.1)$$

with a given initial condition $u(0, x) = u_0(x)$. This is an equation of the form (1.5.35) that we have considered in Section 1.5.2, and we make the same assumptions on the nonlinearity, that we now denote by H , the standard notation in the theory of the Hamilton-Jacobi equations, as in that section. First, we assume that H is smooth and 1-periodic in x . We also make the Lipschitz assumption on the function $H(x, p)$: there exists $C_L > 0$ so that

$$|H(x, p_1) - H(x, p_2)| \leq C_L |p_1 - p_2|, \quad \text{for all } x, p_1, p_2 \in \mathbb{R}^n. \quad (2.3.2)$$

In addition, we assume that H is growing linearly in p at infinity: there exist $\alpha > 0$ and $\beta > 0$ so that

$$0 < \alpha \leq \liminf_{|p| \rightarrow +\infty} \frac{H(x, p)}{|p|} \leq \limsup_{|p| \rightarrow +\infty} \frac{H(x, p)}{|p|} \leq \beta < +\infty, \quad \text{uniformly in } x \in \mathbb{T}^n. \quad (2.3.3)$$

One consequence of (2.3.3) is that $H(x, p)$ is uniformly bounded from below. Note also that if $u(t, x)$ solves (2.3.1) then $u(t, x) + Kt$ solves (2.3.1) with the Hamiltonian $H(x, p)$ replaced by $H(x, p) + K$. Therefore, we may assume without loss of generality that there exist $C_{1,2} > 0$ so that

$$C_1(1 + |p|) \leq H(x, p) \leq C_2(1 + |p|), \quad \text{for all } x \in \mathbb{T}^n \text{ and } p \in \mathbb{R}^n, \quad (2.3.4)$$

so that, in particular,

$$H(x, p) > C_1 \text{ for all } x \in \mathbb{T}^n \text{ and } p \in \mathbb{R}^n. \quad (2.3.5)$$

As we have seen in Section 1.5.2, these assumptions ensure the existence of a unique smooth 1-periodic solution $u(t, x)$ to (2.3.1) supplemented by a continuous, 1-periodic initial condition $u_0(x)$. In order to discuss its long time behavior, we need to introduce a special class of solutions of (2.3.1).

Theorem 2.3.1 *Under the above assumptions, there exists a unique $c \in \mathbb{R}$ so that (2.3.1) has solutions (that we will call the wave solutions) of the form*

$$w(t, x) = ct + \phi(x), \quad (2.3.6)$$

with a 1-periodic function $\phi(x)$. The profile $\phi(x)$ is unique up to an additive constant: if $w_1(t, x)$ and $w_2(t, x)$ are two such solutions then there exists $k \in \mathbb{R}$ so that $\phi_1(x) - \phi_2(x) \equiv k$ for all $x \in \mathbb{T}^n$.

The constant c is often referred to as the speed of the plane wave. The reason is that the solutions to the Hamilton-Jacobi equations, apart from the optimal control theory context that we have discussed above, also often describe the height of an interface, so that c may be thought of as the speed at which the height of the interface is moving up, and $\phi(x)$ as the fixed profile of that interface as it moves up at a constant speed.

Exercise 2.3.2 Consider the following discrete growing interface model, defined on the lattice $h\mathbb{Z}$, with a time step τ , with $u(t, x)$ being the interface height at the time t at the position x :

$$\begin{aligned} u(t + \tau, x) = & \frac{1}{2} [u(t, x - h) + u(t, x + h)] \\ & + \frac{1}{2} [F(u(t, x + h) - u(t, x)) + F(u(t, x) - u(t, x - h))] + \delta V(t, x), \end{aligned} \quad (2.3.7)$$

with a given function $F(p)$, and a prescribed source $V(t, x)$. The terms in the right side of (2.3.7) can be interpreted as follows: (1) the first term has an equilibrating effect, leveling the interface out, (2) the second term says that the rate of the interface growth depends on its slope – things falling from above can stick to the interface, and (3) the last term is an outside source of the interface growth (things falling from above). Find a scaling limit that relates τ , h and δ so that in the limit you get a Hamilton-Jacobi equation of the form

$$u_t = \Delta u + H(x, \nabla u) + V(t, x). \quad (2.3.8)$$

The large time behavior of $u(t, x)$ is summarized in the next theorem.

Theorem 2.3.3 *Let $u(t, x)$ be the solution of the Cauchy problem for (2.3.1) with a continuous 1-periodic initial condition u_0 . There is a wave solution $w(t, x)$ of the form (2.3.6), a constant $\omega > 0$ that does not depend on u_0 and $C_0 > 0$ that depends on u_0 such that*

$$|u(t, x) - w(t, x)| \leq C_0 e^{-\omega t}, \quad (2.3.9)$$

for all $t \geq 0$ and $x \in \mathbb{T}^n$.

We will first prove the existence part of Theorem 2.3.1, and that will occupy most of the rest of this section, while its uniqueness part and the convergence claim of Theorem 2.3.3 will be proved together rather quickly in the end.

2.3.1 The wave solutions

Outline of the existence proof

We first present an outline of the existence proof, before going into the details of the argument. Plugging the ansatz (2.3.6) into (2.3.1) and integrating over \mathbb{T}^n gives a relation

$$c = \int_{\mathbb{T}^n} H(x, \nabla \phi(x)) dx. \quad (2.3.10)$$

The equation for ϕ can, therefore, be written as

$$-\Delta \phi = H(x, \nabla \phi(x)) - \int_{\mathbb{T}^n} H(z, \nabla \phi(z)) dz, \quad (2.3.11)$$

and this will be the starting point of our analysis.

We are going to use a continuation method. As this strategy is typical for the existence proofs for many nonlinear PDEs, it is worth sketching out the general plan. Instead of just looking at (2.3.11) with a given Hamiltonian $H(x, p)$, we consider a family of equations

$$-\Delta\phi_\sigma = H_\sigma(x, \nabla\phi_\sigma) - \int_{\mathbb{T}^n} H_\sigma(z, \nabla\phi_\sigma) dz, \quad (2.3.12)$$

with the Hamiltonians

$$H_\sigma(x, p) = (1 - \sigma)H_0(x, p) + \sigma H(x, p), \quad (2.3.13)$$

parametrized by $\sigma \in [0, 1]$. At $\sigma = 0$, we start with a particular choice of $H_0(x, p)$ for which we know that (2.3.12) has a solution: we take

$$H_0(x, p) = \sqrt{1 + |p|^2},$$

so that $\phi_0(x) \equiv 0$ is an explicit solution to (2.3.12) with $\sigma = 0$. At $\sigma = 1$, we end with

$$H_1(x, p) = H(x, p). \quad (2.3.14)$$

The goal is show that (2.3.12) has a solution for all $\sigma \in [0, 1]$ and not just for $\sigma = 0$ by showing that the set Σ of σ for which (2.3.12) has a solution is both open and closed in $[0, 1]$.

Showing that Σ is closed requires a priori bounds on the solution ϕ_σ of (2.3.12) that would both be uniform in $\sigma \in [0, 1]$ and ensure the compactness of the sequence ϕ_{σ_n} of solutions to (2.3.12) as $\sigma_n \rightarrow \sigma$ in a suitable function space. The estimates should be strong enough to ensure that the limit ϕ_σ is a solution to (2.3.12).

In order to show that Σ is open, one relies on the implicit function theorem. Let us assume that (2.3.12) has a solution $\phi_\sigma(x)$ for some $\sigma \in [0, 1]$. In order to show that (2.3.12) has a solution for $\sigma + \varepsilon$, with a sufficiently small ε , we are led to consider the linearized problem

$$-\Delta h - \frac{\partial H_\sigma(x, \nabla\phi_\sigma)}{\partial p_j} \frac{\partial h}{\partial x_j} + \int_{\mathbb{T}^n} \frac{\partial H_\sigma(z, \nabla\phi_\sigma)}{\partial p_j} \frac{\partial h(z)}{\partial z_j} dz = f, \quad (2.3.15)$$

with

$$f(x) = H(x, \nabla\phi_\sigma) - H_0(x, \nabla\phi_\sigma) - \int_{\mathbb{T}^n} H(z, \nabla\phi_\sigma(z)) dz + \int_{\mathbb{T}^n} H_0(z, \nabla\phi_\sigma(z)) dz. \quad (2.3.16)$$

The implicit function theorem guarantees existence of the solution $\phi_{\sigma+\varepsilon}$, provided that the linearized operator in the left side of (2.3.15) is invertible, with the norm of the inverse a priori bounded in σ . This will show that the set Σ of $\sigma \in [0, 1]$ for which the solution to (2.3.12) exists is open.

The bounds on the operator that maps $f \rightarrow h$ in (2.3.15) also require the a priori bounds on ϕ_σ . Thus, both proving that Σ is open and that it is closed require us to prove the a priori uniform bounds on ϕ_σ . Therefore, our first step will be to assume that a solution $\phi_\sigma(x)$ to (2.3.12) exists and obtain a priori bounds on ϕ_σ . Note that if ϕ_σ is a solution to (2.3.12), then $k + \phi_\sigma$ is also a solution for any $k \in \mathbb{R}$. Thus, it is more natural to obtain a priori bounds

on $\nabla\phi_\sigma$ than on ϕ_σ itself, and then normalize the solution so that $\phi_\sigma(0) = 0$ to ensure that ϕ_σ is bounded.

It is important to observe that the Hamiltonians $H_\sigma(x, p)$ obey the Lipschitz bound (2.3.2), with a Lipschitz constant C_L that does not depend on $\sigma \in [0, 1]$, and estimate (2.3.4) also holds for H_σ with the same $C_{1,2} > 0$ for all $\sigma \in [0, 1]$. The key bound to prove will be to show that there exists a constant $K > 0$ that depends only on the Lipschitz constant of H in (2.3.2) and the two constants in the linear growth estimate (2.3.4) such that any solution to (2.3.12) satisfies

$$\|\nabla\phi_\sigma\|_{L^\infty(\mathbb{T}^n)} \leq K. \quad (2.3.17)$$

We stress that this bound will be obtained not just for one Hamiltonian but for all Hamiltonians with the same Lipschitz constant C_L in (2.3.2) that satisfy (2.3.4) with the same $C_{1,2} > 0$. The estimate (2.3.17) will turn out to be sufficient to apply the argument we have outlined above.

An a priori L^1 -bound on the gradient

Before establishing the L^∞ -bound (2.3.17), let us first prove that there exists a constant $C > 0$ that only depends on C_L in (2.3.2) and $C_{1,2}$ in (2.3.4) such that any solution $\phi_\sigma(x)$ of (2.3.12) satisfies

$$\int_{\mathbb{T}^n} H_\sigma(x, \nabla\phi_\sigma) dx \leq C. \quad (2.3.18)$$

Because of the lower bound in (2.3.3), this is equivalent to an a priori L^1 -bound on $|\nabla\phi_\sigma|$:

$$\int_{\mathbb{T}^n} |\nabla\phi_\sigma(x)| dx \leq C, \quad (2.3.19)$$

with a possibly different $C > 0$ that still depends only on C_L and $C_{1,2}$. To prove (2.3.18), we will rely on Proposition 1.7.10 in Chapter 1 that we recall here for the convenience of the reader.

Proposition 2.3.4 *Let $b(x)$ be a smooth vector field over \mathbb{T}^n . The linear equation*

$$-\Delta e + \nabla \cdot (eb) = 0, \quad x \in \mathbb{T}^n, \quad (2.3.20)$$

has a unique solution $e_1^(x)$ normalized so that*

$$\|e_1^*\|_{L^\infty} = 1, \quad (2.3.21)$$

and such that $e_1^ > 0$ on \mathbb{T}^n . Moreover, for all $\alpha \in (0, 1)$, the function e_1^* is α -Hölder continuous, with the α -Hölder norm bounded by a universal constant depending only on $\|b\|_{L^\infty(\mathbb{T}^n)}$.*

Let us first see why it implies (2.3.18). An immediate consequence of the normalization (2.3.21) and the claim about the Hölder norm of e_1^* , together with the positivity of e_1^* is that

$$\int_{\mathbb{T}^n} e_1^*(x) dx \geq K_1 > 0, \quad (2.3.22)$$

with a constant $K_1 > 0$ that depends only on $\|b\|_{L^\infty}$. Now, given a solution $\phi_\sigma(x)$ of (2.3.12), set

$$b_j(x) = \int_0^1 \partial_{p_j} H_\sigma(x, r \nabla \phi_\sigma(x)) dr, \quad (2.3.23)$$

so that

$$b(x) \cdot \nabla \phi_\sigma(x) = \sum_{j=1}^n b_j(x) \frac{\partial \phi_\sigma}{\partial x_j} = H_\sigma(x, \nabla \phi_\sigma) - H_\sigma(x, 0), \quad (2.3.24)$$

and (2.3.12) can be re-stated as

$$-\Delta \phi_\sigma - b_j(x) \frac{\partial \phi_\sigma}{\partial x_j} = H_\sigma(x, 0) - \int_{\mathbb{T}^n} H_\sigma(z, \nabla \phi_\sigma) dz. \quad (2.3.25)$$

Note that while $b(x)$ does depend on $\nabla \phi_\sigma$, the L^∞ -norm of $b(x)$ depends only on the Lipschitz constant C_L of the function $H_\sigma(x, p)$ in the p -variable. Let now e_1^* be the solution to (2.3.20) given by Proposition 2.3.4, with the above $b(x)$. Multiplying (2.3.25) by e_1^* and integrating over \mathbb{T}^n yields

$$0 = \int_{\mathbb{T}^n} e_1^*(x) H_\sigma(x, 0) dx - \left(\int_{\mathbb{T}^n} e_1^*(x) dx \right) \left(\int_{\mathbb{T}^n} H_\sigma(z, \nabla \phi_\sigma) dz \right), \quad (2.3.26)$$

hence

$$\int_{\mathbb{T}^n} H_\sigma(x, \nabla \phi_\sigma) dx \leq \left(\int_{\mathbb{T}^n} e_1^*(x) dx \right)^{-1} \int_{\mathbb{T}^n} e_1^*(x) H_\sigma(x, 0) dx, \quad (2.3.27)$$

and (2.3.19) follows from (2.3.22) and (2.3.4). As the constant K_1 in (2.3.22) depends only on the L^∞ -norm of $b(x)$ that, in turn, depends only on C_L , the constant C in the right side of (2.3.18), indeed, depends only on C_L and $C_{1,2}$.

An a priori L^∞ bound on the gradient

So far, we have obtained an a priori L^1 -bound (2.3.19) for the gradient of any solution ϕ_σ to (2.3.12). Now, we improve this estimate to an L^∞ bound.

Proposition 2.3.5 *There is a constant $C > 0$ that depends only on the constants C_L and $C_{1,2}$, such that any solution ϕ_σ to*

$$-\Delta \phi_\sigma = H_\sigma(x, \nabla \phi_\sigma) - \int_{\mathbb{T}^n} H_\sigma(z, \nabla \phi_\sigma) dz, \quad (2.3.28)$$

satisfies

$$\|\nabla \phi_\sigma\|_{L^\infty(\mathbb{T}^n)} \leq C. \quad (2.3.29)$$

As a consequence, if ϕ_σ is normalized such that $\phi_\sigma(0) = 0$, then we also have $\|\phi_\sigma\|_{L^\infty(\mathbb{T}^n)} \leq C$.

Proof. We borrow the strategy in the proof of Proposition 1.7.10. Let ϕ_σ be a solution to (2.3.28) such that $\phi_\sigma(0) = 0$. The only estimate we have so far is the L^1 -bound (2.3.19) for $\nabla \phi_\sigma$ – the idea is to estimate the L^∞ -norm $\|\nabla \phi_\sigma\|_{L^\infty(\mathbb{T})}$ solely from the L^1 -norm of ϕ_σ and the equation.

Let $\Gamma(x)$ be as in the proof of Proposition 1.7.10: a nonnegative smooth function equal to 1 in the cube $[-2, 2]^n$ and to zero outside of the cube $(-3, 3)^n$, and set $\psi(x) = \Gamma(x)\phi_\sigma(x)$. The function $\psi(x)$ satisfies an equation similar to what we have seen in (1.7.34):

$$-\Delta\psi = -2\nabla\Gamma \cdot \nabla\phi_\sigma - \phi_\sigma\Delta\Gamma + F(x), \quad x \in \mathbb{R}^n, \quad (2.3.30)$$

with

$$F(x) = \Gamma(x) \left[H_\sigma(x, \nabla\phi_\sigma(x)) - \int_{\mathbb{T}^n} H_\sigma(z, \nabla\phi_\sigma(z)) dz \right]. \quad (2.3.31)$$

The only a priori information we have about $F(x)$ and the term $\nabla\Gamma \cdot \nabla\phi_\sigma(x)$ so far is that they are supported inside $[-3, 3]^n$ and are uniformly bounded in $L^1(\mathbb{R}^n)$ via (2.3.18) and (2.3.19). Here, we use the assumption (2.3.4) that the Hamiltonian $H(x, p)$ is uniformly positive. It helps to combine these two terms:

$$G(x) = F(x) - 2\nabla\Gamma(x) \cdot \nabla\phi_\sigma(x), \quad (2.3.32)$$

with $G(x)$ supported inside $[-3, 3]^n$, and

$$\int_{\mathbb{R}^n} |G(x)| dx \leq C, \quad (2.3.33)$$

with a constant $C > 0$ that depends only on C_L and $C_{1,2}$, due to (2.3.18) and (2.3.19). We also know that

$$|G(x)| \leq C(1 + |\nabla\phi_\sigma(x)|), \quad (2.3.34)$$

because of (2.3.4).

Next, we use the fundamental solution $E(x)$ to the Laplace equation in \mathbb{R}^n to write

$$\psi(x) = \int_{\mathbb{R}^n} E(x-y)[G(y) - \phi_\sigma(y)\Delta\Gamma(y)] dy. \quad (2.3.35)$$

Differentiating (2.3.35) in x gives

$$\nabla\psi(x) = \int_{\mathbb{R}^n} \nabla E(x-y)[G(y) - \phi_\sigma(y)\Delta\Gamma(y)] dy. \quad (2.3.36)$$

Exercise 2.3.6 Note that the function $E(x-y)$ has a singularity at $y = x$. Show that nevertheless one can differentiate in (2.3.35) under the integral sign to obtain (2.3.36).

The function $\nabla E(x-y)$ has an integrable singularity at $y = x$, of the order $|x-y|^{-n+1}$, and is bounded everywhere else. Thus, for all $\varepsilon > 0$ we have, with the help of (2.3.33) and (2.3.34):

$$\begin{aligned} \left| \int_{\mathbb{R}^n} G(y)\nabla E(x-y) dy \right| &\leq \left| \int_{|x-y|\leq\varepsilon} G(y)\nabla E(x-y) dy \right| + \left| \int_{|x-y|\geq\varepsilon} G(y)\nabla E(x-y) dy \right| \\ &\leq C(1 + \|\nabla\phi_\sigma\|_{L^\infty}) \int_{|x-y|\leq\varepsilon} \frac{dy}{|x-y|^{n-1}} + C\varepsilon^{-n+1} \int_{|x-y|\geq\varepsilon} |G(y)| dy \\ &\leq C\varepsilon(1 + \|\nabla\phi_\sigma\|_{L^\infty}) + C\varepsilon^{1-n}. \end{aligned} \quad (2.3.37)$$

The integral in (2.3.36) also contains a factor of ϕ_σ , whereas our bounds so far deal with $\nabla\phi_\sigma$. However, we have assumed without loss of generality that $\phi_\sigma(0) = 0$, hence for any $\delta > 0$ we may write

$$\phi_\sigma(y) = \int_0^1 y \cdot \nabla\phi_\sigma(sy)ds = \int_0^\delta y \cdot \nabla\phi_\sigma(sy)ds + \int_\delta^1 y \cdot \nabla\phi_\sigma(sy)ds,$$

so that both, as $|y| \leq 1$, we have

$$|\phi_\sigma(y)| \leq \|\nabla\phi_\sigma\|_{L^\infty}, \quad (2.3.38)$$

and

$$\begin{aligned} \int_{\mathbb{T}^n} |\phi_\sigma(y)|dy &\leq C\delta\|\nabla\phi_\sigma\|_{L^\infty} + \int_\delta^1 \int_{\mathbb{T}^n} |y|\nabla\phi_\sigma(sy)dyds \\ &\leq C\delta\|\nabla\phi_\sigma\|_{L^\infty} + C \int_\delta^1 \int_{s\mathbb{T}^n} |y|\nabla\phi_\sigma(y)dy \frac{ds}{s^{n+1}} \leq C\delta\|\nabla\phi_\sigma\|_{L^\infty} + C \int_\delta^1 \frac{ds}{s^{1+n}} \\ &\leq C\delta\|\nabla\phi_\sigma\|_{L^\infty} + C\delta^{-n}. \end{aligned} \quad (2.3.39)$$

We used above the a priori bound (2.3.19) on $\|\nabla\phi\|_{L^1(\mathbb{T}^n)}$. Combining (2.3.38) and (2.3.39), we obtain, as in (2.3.37):

$$\begin{aligned} \left| \int_{\mathbb{R}^n} \phi_\sigma(y)\Delta\Gamma(y)\nabla E(x-y)dy \right| &\leq \int_{|x-y|\leq\varepsilon} |\phi_\sigma(y)||\Delta\Gamma(y)|\nabla E(x-y)dy \\ &+ \int_{|x-y|\geq\varepsilon} |\phi_\sigma(y)||\Delta\Gamma(y)|\nabla E(x-y)dy \leq C\varepsilon\|\phi_\sigma\|_{L^\infty} + C\varepsilon^{1-n} \int_{\mathbb{T}^n} |\phi_\sigma(y)|dy \\ &\leq C\varepsilon\|\nabla\phi_\sigma\|_{L^\infty} + C\varepsilon^{1-n}\delta\|\nabla\phi_\sigma\|_{L^\infty} + C\varepsilon^{1-n}\delta^{-n}. \end{aligned} \quad (2.3.40)$$

Together, (2.3.37) and (2.3.40) tell us that

$$\|\nabla\psi\|_{L^\infty} \leq C\varepsilon(1+\|\nabla\phi_\sigma\|_{L^\infty}) + C\varepsilon^{1-n} + C\varepsilon\|\nabla\phi_\sigma\|_{L^\infty} + C\varepsilon^{1-n}\delta\|\nabla\phi_\sigma\|_{L^\infty} + C\varepsilon^{1-n}\delta^{-n}. \quad (2.3.41)$$

Next, observe that, because $\Gamma \equiv 1$ in $[-2, 2]^n$ and ϕ_σ is 1-periodic, we have

$$\|\nabla\phi_\sigma\|_{L^\infty(\mathbb{T}^n)} = \|\nabla(\Gamma\phi_\sigma)\|_{L^\infty([-1,1]^n)} \leq \|\nabla(\Gamma\phi_\sigma)\|_{L^\infty([-3,3]^n)} = \|\nabla\psi\|_{L^\infty}. \quad (2.3.42)$$

Thus, if we take $\delta = \varepsilon^n$ in (2.3.41), we would obtain

$$\|\nabla\phi_\sigma\|_{L^\infty} \leq C\varepsilon\|\nabla\phi_\sigma\|_{L^\infty} + C_\varepsilon, \quad (2.3.43)$$

with a universal constant $C > 0$ and C_ε that does depend on ε . Now, the proof of (2.3.29) is concluded by taking $\varepsilon > 0$ small enough. \square

Going back to equation (2.3.11) for ϕ :

$$-\Delta\phi = H(x, \nabla\phi) - \int_{\mathbb{T}^n} H(x, \nabla\phi)dx, \quad (2.3.44)$$

the reader should do the following exercise.

Exercise 2.3.7 Use the L^∞ -bound on $\nabla\phi$ in Proposition 2.3.5 to deduce from (2.3.44) that, under the assumption that $H(x, p)$ is smooth (infinitely differentiable) in both variables x and p , the function $\phi(x)$ is, actually, infinitely differentiable, with all its derivatives of order n bounded by a priori constants C_n that do not depend on ϕ .

The linearized problem

We need one more ingredient to finish the proof of the existence part of Theorem 2.3.1: to set-up an application of the implicit function theorem. Let ϕ_σ be a solution to (2.3.12) and let us consider the linearized problem, with an unknown h :

$$-\Delta h - \partial_{p_j} H_\sigma(x, \nabla \phi_\sigma) \partial_{x_j} h + \int_{\mathbb{T}^n} \partial_{p_j} H_\sigma(y, \nabla \phi_\sigma) \partial_{x_j} h(y) dy = f \quad x \in \mathbb{T}^n. \quad (2.3.45)$$

We assume that $f \in C^{1,\alpha}(\mathbb{T}^n)$ for some $\alpha \in (0, 1)$, and f has zero mean over \mathbb{T}^n :

$$\int_{\mathbb{T}^n} f(x) dx = 0,$$

and require that the solution h to (2.3.45) also has zero mean:

$$\int_{\mathbb{T}^n} h(x) dx = 0. \quad (2.3.46)$$

Proposition 2.3.8 *Equation (2.3.45) has a unique solution $h \in C^{3,\alpha}(\mathbb{T}^n)$ with zero mean. The mapping $f \mapsto h$ is continuous from the set of $C^{1,\alpha}$ functions with zero mean to the set of $C^{3,\alpha}(\mathbb{T}^n)$ functions with zero mean.*

Proof. The Laplacian is a one-to-one map between the set of $C^{m+2,\alpha}$ functions with zero mean and the set of $C^{m,\alpha}(\mathbb{T}^n)$ functions with zero mean, for any $m \in \mathbb{N}$. Thus, we may talk about its inverse that we denote by $(-\Delta)^{-1}$. Equation (2.3.45) is thus equivalent to

$$(I + K)h = (-\Delta)^{-1}f, \quad (2.3.47)$$

with the operator

$$Kh = (-\Delta)^{-1} \left(-\partial_{p_j} H_\sigma(x, \nabla \phi_\sigma) \partial_{x_j} h + \int_{\mathbb{T}^n} \partial_{p_j} H_\sigma(y, \nabla \phi_\sigma) \partial_{x_j} h(y) dy \right). \quad (2.3.48)$$

Exercise 2.3.9 Show that K is a compact operator on the set of functions in $C^{1,\alpha}(\mathbb{T}^n)$ with zero mean.

The problem has been now reduced to showing that the only solution of

$$(I + K)h = 0 \quad (2.3.49)$$

with $h \in C^{1,\alpha}(\mathbb{T}^n)$ with zero mean is $h \equiv 0$. Note that (2.3.49) simply says that h is a solution of (2.3.45) with $f \equiv 0$. Let $e_1^* > 0$ be given by Proposition 2.3.4, with

$$b_j(x) = -\partial_{p_j} H_\sigma(x, \nabla \phi_\sigma). \quad (2.3.50)$$

That is, e_1^* is the positive solution of the equation

$$-\Delta e_1^* + \nabla \cdot (e_1^* b) = 0, \quad (2.3.51)$$

normalized so that $\|e_1^*\|_{L^\infty(\mathbb{T}^n)} = 1$. The uniform Lipschitz bound on $H_\sigma(x, p)$ in the p -variable implies that $b(x)$ is in $L^\infty(\mathbb{T}^n)$, and thus Proposition 2.3.4 can be applied. Multiplying (2.3.45) with $f = 0$ by e_1^* and integrating gives, as $e_1^* > 0$:

$$\int_{\mathbb{T}^n} \partial_{p_j} H_\sigma(y, \nabla \phi_\sigma) \partial_{x_j} h(y) dy = 0.$$

But then, the equation for h becomes simply

$$-\Delta h + b_j(x) \partial_{x_j} h = 0, \quad x \in \mathbb{T}^n,$$

which entails that h is constant, by the Krein-Rutman theorem. Because h has zero mean, we get $h \equiv 0$. \square

Exercise 2.3.10 Let $H_0(x, p)$ satisfy the assumptions of Theorem 2.3.3, and assume that equation (2.3.11), with $H = H_0$,

$$-\Delta \phi_0 = H_0(x, \nabla \phi_0) - \int_{\mathbb{T}^n} H_0(z, \nabla \phi_0) dz, \quad (2.3.52)$$

has a solution $\phi_0 \in C(\mathbb{T}^n)$. Consider $H_1(x, p) \in C^\infty(\mathbb{T} \times \mathbb{R}^n)$. Prove, with the aid of Propositions 2.3.5 and 2.3.8, and the implicit function theorem that there exist $R_0 > 0$ and $\varepsilon_0 > 0$ such that if

$$|H_1(x, p)| \leq \varepsilon_0, \quad \text{for } x \in \mathbb{T}^n \text{ and } |p| \leq R_0, \quad (2.3.53)$$

then equation (2.3.11) with $H = H_0 + H_1$:

$$-\Delta \phi = H(x, \nabla \phi) - \int_{\mathbb{T}^n} H(z, \nabla \phi) dz, \quad (2.3.54)$$

has a solution ϕ .

Existence of the solution

We finally prove the existence part of Theorem 2.3.1. Consider $H(x, p)$ satisfying the assumptions of the theorem. As before, we set

$$H_0(x, p) = \sqrt{1 + |p|^2} - 1,$$

and

$$H_\sigma(x, p) = (1 - \sigma)H_0(x, p) + \sigma H(x, p),$$

so that $H_1(x, p) = H(x, p)$, and consider existence of a solution to (2.3.12):

$$-\Delta \phi_\sigma = H_\sigma(x, \nabla \phi_\sigma) - \int_{\mathbb{T}^n} H_\sigma(z, \nabla \phi_\sigma) dz, \quad (2.3.55)$$

Consider the set

$$\Sigma = \{\sigma \in [0, 1] : \text{equation (2.3.55) has a solution}\}.$$

Our goal is to show that $\Sigma = [0, 1]$. We know that Σ is non empty, because $0 \in \Sigma$: indeed, $\phi_0(x) \equiv 0$ is a solution to (2.3.55) at $\sigma = 0$. Thus, if we show that Σ is both open and closed in $[0, 1]$, this will imply that $\Sigma = [0, 1]$, and, in particular, establish the existence of a solution to (2.3.55) for $H_1(x, p) = H(x, p)$.

Now that we know that the linearized problem is invertible, the openness of Σ is a direct consequence of the inverse function theorem, as explained in Exercise 2.3.10. Closedness of Σ is not too difficult to see either: consider a sequence $\sigma_n \in [0, 1]$ converging to $\bar{\sigma} \in [0, 1]$, and let ϕ_n be a solution to (2.3.55) with $H(x, p) = H_{\sigma_n}(x, p)$, normalized so that

$$\phi_n(0) = 0. \quad (2.3.56)$$

Proposition 2.3.5 implies that

$$\|\nabla\phi_n\|_{L^\infty(\mathbb{T}^n)} \leq C,$$

and thus

$$\|H(x, \nabla\phi_n)\|_{L^\infty} \leq C.$$

However, this means that ϕ_n solve an equation of the form

$$-\Delta\phi_n = F_n(x), \quad x \in \mathbb{T}^n, \quad (2.3.57)$$

with a uniformly bounded function

$$F_n(x) = H_{\sigma_n}(x, \nabla\phi_n) - \int_{\mathbb{T}^n} H_{\sigma_n}(z, \nabla\phi_n(z)) dz. \quad (2.3.58)$$

It follows that ϕ_n is bounded in $C^{1,\alpha}(\mathbb{T}^n)$, for all $\alpha \in [0, 1]$:

$$\|\phi_n\|_{C^{1,\alpha}(\mathbb{T}^n)} \leq C. \quad (2.3.59)$$

But this implies, in turn, that the functions $F_n(x)$ in (2.3.58) are also uniformly bounded in $C^\alpha(\mathbb{T}^n)$, hence ϕ_n are uniformly bounded in $C^{2,\alpha}(\mathbb{T}^n)$:

$$\|\phi_n\|_{C^{2,\alpha}(\mathbb{T}^n)} \leq C. \quad (2.3.60)$$

Now, the Arzela-Ascoli theorem implies that a subsequence ϕ_{n_k} will converge in $C^2(\mathbb{T}^n)$ to a function $\bar{\phi}$, which is a solution to (2.3.19) with $H = H_{\bar{\sigma}}$. Thus, $\bar{\sigma} \in \Sigma$, and Σ is closed. This finishes the proof of the existence part of the theorem.

2.3.2 Long time convergence and uniqueness of the wave solutions

We will now prove simultaneously the claim of the uniqueness of the speed c and of the profile $\phi(x)$ in Theorem 2.3.1, and the long time convergence for the solutions to the Cauchy problem stated in Theorem 2.3.3.

Let $u(t, x)$ be the solution to (2.3.1)

$$u_t = \Delta u + H(x, \nabla u), \quad t > 0, \quad x \in \mathbb{T}^n, \quad (2.3.61)$$

with $u(0, x) = u_0(x) \in C(\mathbb{T}^n)$. We also take a speed $c \in \mathbb{R}$ and a solution $\phi(x)$ to

$$\Delta\phi + H(x, \nabla\phi) = c, \quad (2.3.62)$$

without assuming that either c or ϕ is unique.

We wish to prove that there exists $\bar{k} \in \mathbb{R}$ so that $u(t, x) - ct$ is attracted exponentially fast in time to $\phi(x) + \bar{k}$:

$$|u(t, x) - ct - \bar{k} - \phi(x)| \leq C_0 e^{-\omega t}, \quad (2.3.63)$$

with some $C_0 > 0$ and $\omega > 0$, such that C_0 depends on the initial condition u_0 but ω does not. The idea is the same as in the proof of Theorem 1.8.2 for the Allen-Cahn equation: squeeze the solution between two different wave solutions, and show that the difference between the squeezers tends to zero as $t \rightarrow +\infty$. However, the situation here is much simpler: we do not have any tail as $|x| \rightarrow +\infty$ to control, because we are now considering the problem for $x \in \mathbb{T}^n$. Actually, the present setting realizes what would be the dream scenario for the Allen-Cahn equation.

As a simple remark, we may assume that $c = 0$, just by setting

$$\tilde{H}(x, p) = H(x, p) - c,$$

and dropping the tilde, and this is what we will do. In other words, $\phi(x)$ is the solution to

$$\Delta\phi + H(x, \nabla\phi) = 0. \quad (2.3.64)$$

Let ϕ be any solution to (2.3.64), and set

$$k_0^- = \sup\{k : u(0, x) \geq \phi(x) + k \text{ for all } x \in \mathbb{T}^n\},$$

and

$$k_0^+ = \inf\{k : u(0, x) \leq \phi(x) + k \text{ for all } x \in \mathbb{T}^n\}.$$

Because $\phi(x) + k_0^\pm$ solve (2.3.64) with $c = 0$, and $u(t, x)$ solves (2.3.61), we have, by the maximum principle:

$$\phi(x) + k_0^- \leq u(t, x) \leq \phi(x) + k_0^+, \text{ for all } t \geq 0 \text{ and } x \in \mathbb{T}^n. \quad (2.3.65)$$

Now, for all $q \in \mathbb{N}$, let us set

$$k_q^- = \sup\{k : u(t = q, x) \geq \phi(x) + k \text{ for all } x \in \mathbb{T}^n\} = \inf_{x \in \mathbb{T}^n} [u(t = q, x) - \phi(x)], \quad (2.3.66)$$

and

$$k_q^+ = \inf\{k : u(t = q, x) \leq \phi(x) + k \text{ for all } x \in \mathbb{T}^n\} = \sup_{x \in \mathbb{T}^n} [u(t = q, x) - \phi(x)]. \quad (2.3.67)$$

The strong maximum principle implies that the sequence k_q^- is increasing, whereas k_q^+ is decreasing, and that, as in (2.3.65), we have

$$\phi(x) + k_q^- \leq u(t, x) \leq \phi(x) + k_q^+, \text{ for all } t \geq q \text{ and } x \in \mathbb{T}^n. \quad (2.3.68)$$

Hence, the theorem will be proved if we manage to show that

$$0 \leq k_q^+ - k_q^- \leq Ca^q, \quad \text{for all } q \geq 0, \quad (2.3.69)$$

with some $C \in \mathbb{R}$ that may depend on the initial condition u_0 and $a \in (0, 1)$ that does not depend on u_0 . In order to prove (2.3.69), it suffices to show that

$$k_{q+1}^+ - k_{q+1}^- \leq (1 - r_0)(k_q^+ - k_q^-), \quad (2.3.70)$$

with some $r_0 \in (0, 1)$. This is a quantification of the strong maximum principle: by the time $t = q + 1$ $u(x)$ has to detach "by a fixed amount" from the respective lower and upper bounds $\phi(x) + k_q^\pm$ that hold at $t = q$. Such estimates typically rely on the Harnack inequality, and this is what we will use.

To bring the Harnack inequality in, note that the function

$$w(t, x) = u(t, x) - \phi(x) - k_q^-$$

is nonnegative for $t \geq q$, and solves an equation of the form

$$\partial_t w - \Delta w + b_j(t, x) \partial_{x_j} w = 0, \quad t > q, \quad x \in \mathbb{T}^n, \quad (2.3.71)$$

with a bounded drift $b(t, x)$ given by

$$b(t, x) = \int_0^1 \nabla_p H(x, (1-s)\nabla\phi(x) + s\nabla u(t, x)) ds, \quad (2.3.72)$$

so that

$$b(t, x) \cdot [\nabla u(t, x) - \nabla\phi(x)] = H(x, \nabla u(t, x)) - H(x, \nabla\phi(x)),$$

and

$$|b_j(t, x)| \leq C_L, \quad \text{for all } t \geq q \text{ and } x \in \mathbb{T}^n. \quad (2.3.73)$$

The Harnack inequality in Theorem ?? and (2.3.73) imply that there exists $r_0 > 0$ that depends only on C_L such that

$$\inf_{x \in \mathbb{T}^n} w(q+1, x) \geq r_0 \sup_{x \in \mathbb{T}^n} w(q, x). \quad (2.3.74)$$

Using (2.3.66) and (2.3.67), together with (2.3.74), we may write

$$\begin{aligned} r_0 \sup_{x \in \mathbb{T}^n} w(q, x) &= r_0 \sup_{x \in \mathbb{T}^n} [u(q, x) - \phi(x) - k_q^-] = r_0 [k_q^+ - k_q^-] \leq \inf_{x \in \mathbb{T}^n} w(q+1, x) \\ &= \inf_{x \in \mathbb{T}^n} [u(q+1, x) - \phi(x) - k_q^-] = k_{q+1}^- - k_q^-, \end{aligned} \quad (2.3.75)$$

so that

$$k_{q+1}^- \geq k_q^- + r_0 [k_q^+ - k_q^-]. \quad (2.3.76)$$

As $k_{q+1}^+ \leq k_q^+$, it follows that

$$k_{q+1}^+ - k_{q+1}^- \leq k_q^+ - k_q^- - r_0 (k_q^+ - k_q^-) \leq (1 - r_0)(k_q^+ - k_q^-), \quad (2.3.77)$$

which is (2.3.70). This implies the geometric decay as in (2.3.69), hence the theorem, because of (2.3.68) and (2.3.69). Note that the constant

$$a = 1 - r_0$$

comes from the Harnack inequality and does not depend on the initial condition u_0 but only on the Lipschitz constant C_L of $H(x, p)$. \square

- Exercise 2.3.11** (i) Why does the uniqueness of c and of the profile $\phi(x)$ follow?
(ii) How is the constant ω in Theorem 2.3.3 related to the constant a in the above proof?

Exercise 2.3.12 Consider a modified equation, not quite of the Hamilton-Jacobi form:

$$u_t - \Delta u = R(x, u)\sqrt{1 + |\nabla u|^2}, \quad (2.3.78)$$

where $R(x, u)$ is a smooth, positive function, that is 1-periodic in x and 1-periodic in u .

- (i) Let $u_0 \in C(\mathbb{T}^N)$, and show that the Cauchy problem for (2.3.78) with $u(0, x) = u_0(x)$ is well posed.
(ii) Prove the existence of a unique $T > 0$ such that equation (2.3.78) has solutions of the form

$$u(t, x) = \frac{t}{T} + \phi(t, x), \quad (2.3.79)$$

where ϕ is T -periodic in t and 1-periodic in x . We will call such a solution a wave solution. Why is it not reasonable to expect that under the above assumptions (2.3.78) has a wave solution of the form $u(t, x) = ct + \psi(x)$ with a 1-periodic function $\psi(x)$?

- (iii) Show that every solution of the Cauchy problem which is initially 1-periodic in x converges, exponentially fast in time, to a wave solution of the form (2.3.79).

If in doubt, you may consult [?]. Note that the topological degree argument used in that reference can be replaced by a more elementary implicit function theorem argument we have used in the existence proof here.

2.4 A glimpse of the classical solutions to the Hamilton-Jacobi equations

2.4.1 Smooth solutions and their limitations

We now turn our attention to first order inviscid Hamilton-Jacobi equations of the form

$$u_t + H(x, \nabla u) = 0. \quad (2.4.1)$$

The standard philosophy of the construction of a solution to a first order equation is to find its values on special curves, known as characteristics, that will, hopefully, fill the whole space. This is the strategy that is also classically used to solve (2.4.1). Consider a time $t > 0$ and a point $x \in \mathbb{R}^n$. In order to assign a value to $u(t, x)$ we consider a curve $\gamma(s)$, with $s \in [0, t]$, such that $\gamma(t) = x$, and set

$$p(s) = \nabla u(s, \gamma(s)).$$

Here, $u(t, x)$ is the sought for solution to (2.4.1). Assuming that everything is smooth we have, using the dot to denote the differentiation in s :

$$\begin{aligned} \dot{p}_k(s) &= \partial_{x_k} u_t(s, \gamma(s)) + \frac{\partial^2 u(s, \gamma(s))}{\partial x_k \partial x_m} \dot{\gamma}_m(s) \\ &= -\frac{\partial H(\gamma(s), p(s))}{\partial x_k} - \frac{\partial H(\gamma(s), p(s))}{\partial p_m} \frac{\partial^2 u(s, \gamma(s))}{\partial x_k \partial x_m} + \frac{\partial^2 u(s, \gamma(s))}{\partial x_k \partial x_m} \dot{\gamma}_m(s). \end{aligned} \quad (2.4.2)$$

We see that it is convenient to choose $\gamma(s)$ that satisfies the following system of ODEs:

$$\begin{aligned}\dot{\gamma}(s) &= \nabla_p H(\gamma(s), p(s)) \\ \dot{p}(s) &= -\nabla_x H(\gamma(s), p(s))\end{aligned}\tag{2.4.3}$$

for $0 \leq s \leq t$. This dynamical system is to be complemented by the boundary conditions at $s = 0$ and $s = t$:

$$p(0) = \nabla u_0(\gamma(0)), \quad \gamma(t) = x.\tag{2.4.4}$$

The system (2.4.3) has the form of a Hamiltonian system with the Hamiltonian $H(x, p)$, and the curves $(\gamma(s), p(s))$ are called the characteristic curves. In order to solve (2.1.1), we need to find a solution to (2.4.3)-(2.4.4), and it would be excellent to prove that such solution is unique. The trouble is that there is no good reason, in general, for existence and uniqueness of a solution to this boundary value problem.

Exercise 2.4.1 Consider $x_0 \in \mathbb{R}^n$ and $t > 0$ and assume that $u(t, x)$ is smooth in a ball around x_0 . Prove, for instance, with the help of the implicit function theorem, that the boundary value problem (2.4.3)-(2.4.4) has a unique solution $(\gamma(s), p(s))$ as soon as t is small enough and x is in the vicinity of x_0 , and that this solution is smooth in t and x .

Once $\gamma(s)$ and $p(s)$ are constructed, we may assign a value to $u(t, x)$ as follows. The function

$$\varphi(s) = u(s, \gamma(s))$$

satisfies

$$\dot{\varphi}(s) = u_t(s, \gamma(s)) + \dot{\gamma}(s) \cdot p(s) = -H(\gamma(s), p(s)) + \dot{\gamma}(s) \cdot p(s).\tag{2.4.5}$$

Integrating (2.4.5) from $s = 0$ to $s = t$ gives an expression for $u(t, x)$ in terms of the curves $\gamma(s)$ and $p(s)$, $0 \leq s \leq t$:

$$u(t, x) = u_0(\gamma(0)) + \int_0^t (-H(\gamma(s), p(s)) + \dot{\gamma}(s) \cdot p(s)) ds.\tag{2.4.6}$$

Exercise 2.4.2 Check that (2.4.6) indeed gives a solution to (2.4.1) such that $u(0, x) = u_0(x)$.

To see that this strategy can not always lead to smooth solutions for all times, just consider the simplest nonlinear equation in one space dimension

$$u_t + \frac{u_x^2}{2} = 0 \text{ for } t > 0 \text{ and } x \in \mathbb{R}, \quad u(0, x) = u_0(x).\tag{2.4.7}$$

The solution to the boundary value problem (2.4.3)-(2.4.4) amounts (this is very easily checked) to finding $\gamma(0)$ solving the equation

$$x = \gamma(0) + tu_0'(\gamma(0)),$$

for a given $t > 0$ and $x \in \mathbb{R}$. The issue is that this equation may, or may not have a unique solution $\gamma(0)$. If $u_0'' > 0$, solution is unique and we are on the safe side. But if $u_0''(x_0) < 0$ at some point x_0 , uniqueness fails as soon as

$$t \geq \frac{1}{\sup(-u_0'')}.$$

Thus, we need a more elaborate theory. Nevertheless, in the rest of this section, we wish to show the reader one interesting situation where smooth solutions can be constructed.

Before we end this short section, let us mention, in the form of an exercise (this will be revisited in the context of viscosity solutions), a very strong form of uniqueness.

Exercise 2.4.3 (*Finite speed of propagation*). Let H be uniformly Lipschitz with respect to its second variable, as well as $\nabla_p H$. Let u_0 and v_0 be two smooth, compactly supported initial conditions, and assume that each generates a smooth solution to the Cauchy problem for (2.4.1), on a common time interval $[0, T]$. Compute, in terms of H_p , a constant K such that, if

$$\text{dist}(x, \text{supp}(u_0 - v_0)) > Kt,$$

then $u(t, x) = v(t, x)$. Hint: it may be helpful to solve, first, the following question: let $b(t, x)$ be smooth and uniformly Lipschitz in its second variable. Let u_0 be a smooth compactly supported function, and $u(t, x)$ the solution to

$$\begin{aligned} u_t + b(t, x) \cdot \nabla u &= 0, & t > 0, x \in \mathbb{R}^n \\ u(0, x) &= u_0(x) \end{aligned} \tag{2.4.8}$$

If

$$\text{dist}(x, \text{supp}(u_0)) > t\|b\|_\infty,$$

then $u(t, x) = 0$.

2.4.2 An example of classical global solutions

We now discuss a situation when classical smooth solutions do exist. Consider solutions to the equation

$$u_t + \frac{1}{2}|\nabla u|^2 - R(x) = 0, \tag{2.4.9}$$

with an initial condition $u(0, x) = u_0(x)$. We assume that both u_0 and R are strictly convex smooth functions on \mathbb{R}^n , such that there is $\alpha \in (0, 1)$ so that, for all $x \in \mathbb{R}^n$ and all $\xi \in \mathbb{R}^n$ we have:

$$\alpha|\xi|^2 \leq (D^2u_0(x)\xi \cdot \xi) \leq \alpha^{-1}|\xi|^2, \quad \alpha|\xi|^2 \leq (D^2R(x)\xi \cdot \xi) \leq \alpha^{-1}|\xi|^2. \tag{2.4.10}$$

Exercise 2.4.4 First, consider the case $R = 0$. Argue informally, just by looking at the equation and using pictures that if $u_0(x)$ is strictly convex but its Hessian is uniformly bounded then the graph of $u(t, x)$ should not form a corner, and if $u_0(x)$ is strictly concave but its Hessian is uniformly bounded then it is plausible that the graph of $u(t, x)$ will form a corner. It may be helpful to start by looking at $u(0, x) = |x|^2$ and $u(0, x) = -|x|^2$. Explain, also informally, how the comparison principle should may come into play.

We now use the approach via the characteristic curves to show that a smooth solution exists under the above assumptions. Note that the Hamiltonian for (2.4.9) is

$$H(x, p) = \frac{1}{2}|p|^2 - R(x),$$

and the characteristic system (2.4.3)-(2.4.4) reduces to

$$\dot{\gamma}(s) = p(s), \quad \dot{p}(s) = \nabla R(\gamma(s)),$$

which can be written as

$$-\gamma'' + \nabla R(\gamma) = 0, \tag{2.4.11}$$

with the boundary conditions

$$\gamma'(0) - \nabla u_0(\gamma(0)) = 0, \quad \gamma(t) = x. \tag{2.4.12}$$

To establish uniqueness and smoothness of the solution $u(t, x)$ to (2.4.9) with the initial condition $u(0, x) = u_0(x)$, we need to prove that (2.4.11)-(2.4.12) has a unique solution $\gamma(s)$ that depends smoothly on t and x . Then, $u(t, x)$ will be given by (2.4.6):

$$u(t, x) = u_0(\gamma(0)) + \int_0^t (-H(\gamma(s), p(s)) + \dot{\gamma}(s) \cdot p(s)) ds. \tag{2.4.13}$$

Existence of the characteristic curves

To construct a solution to (2.4.11)-(2.4.12), we observe that (2.4.11) is the Euler-Lagrange equation for the energy functional

$$J_{t,x}(\gamma) = u_0(\gamma(0)) + \int_0^t \left(\frac{|\gamma'(s)|^2}{2} + R(\gamma(s)) \right) ds, \tag{2.4.14}$$

over $H^1([0, t])$, with the constraint $\gamma(t) = x$.

Exercise 2.4.5 Verify that claim: show that if the minimizer of $J_{t,x}(\gamma)$ over the set

$$S = \{\gamma \in H^1[0, t] : \gamma(t) = x\}$$

exists and is smooth then it satisfies both (2.4.11) and the boundary condition at $s = 0$ in (2.4.12). Next, define what it means for $\gamma \in H^1[0, t]$ (without assuming γ is smooth) to be a weak solution to (2.4.11)-(2.4.12) and show that a minimizer of $J_{t,x}$ over S (if it exists) is a weak solution.

As both $u_0(x)$ and $R(x)$ are strictly convex, they are bounded from below, and it is easy to see that the functional $J_{t,x}$ is bounded from below over S . Let us set

$$\bar{J}_{t,x} = \inf_{\gamma \in S} J_{t,x}(\gamma),$$

and let $\gamma_n \in S$ be a minimizing sequence, so that $J_{t,x}(\gamma_n)$ decreases to $\bar{J}_{t,x}$. Once again, as u_0 and R are bounded from below, there exists $C > 0$ so that

$$\int_0^t |\gamma_n'(s)|^2 ds \leq C,$$

for all n . As, in addition, $\gamma_n(t) = x$ for all n , there is a subsequence, that we will still denote by γ_n that converges uniformly over $[0, t]$, and weakly in $H^1([0, t])$ to a limit $\bar{\gamma}_{t,x} \in S$.

To prove that $J_{t,x}(\bar{\gamma}_{t,x}) = \bar{J}_{t,x}$ we simply observe that by the weak convergence we have

$$\|\bar{\gamma}'_{t,x}\|_{L^2}^2 \leq \liminf_{n \rightarrow +\infty} \|\gamma'_n\|_{L^2}^2,$$

which, combined with the uniform convergence of γ_n to $\bar{\gamma}_{t,x}$ on $[0, t]$ implies that

$$J_{t,x}(\bar{\gamma}_{t,x}) \leq \lim_{n \rightarrow +\infty} J_{t,x}(\gamma_n) = \bar{J}_{t,x},$$

and thus

$$J_{t,x}(\bar{\gamma}_{t,x}) = \bar{J}_{t,x}.$$

Uniqueness of the characteristic curve

To prove the uniqueness of the minimizer, we will use the convexity of $u_0(x)$ and $R(x)$ and not just their boundedness from below. Let γ_1 and γ_2 be two solutions to (2.4.11)-(2.4.12). The difference

$$\tilde{\gamma} = \gamma_2 - \gamma_1.$$

satisfies

$$-\tilde{\gamma}_k'' + A_{kj}(s)\tilde{\gamma}_j = 0, \quad 1 \leq k \leq n, \quad (2.4.15)$$

with the boundary conditions

$$\tilde{\gamma}'_k(0) - B_{kj}\tilde{\gamma}_j(0) = 0, \quad \tilde{\gamma}_k(t) = 0, \quad 1 \leq k \leq n. \quad (2.4.16)$$

The matrices A and B are given by

$$A_{kj}(s) = \int_0^1 \frac{\partial^2 R(\gamma_1(s) + \sigma(\gamma_2(s) - \gamma_1(s)))}{\partial x_k \partial x_j} d\sigma,$$

and

$$B_{kj} = \int_0^1 \frac{\partial^2 u_0(\gamma_1(0) + \sigma(\gamma_2(0) - \gamma_1(0)))}{\partial x_k \partial x_j} d\sigma.$$

Let us take the inner product of (2.4.15) with $\tilde{\gamma}$, and integrate. This gives

$$\int_0^t (|\tilde{\gamma}'(s)|^2 + (A\tilde{\gamma}(s) \cdot \tilde{\gamma}(s))) ds + (B\tilde{\gamma}(0) \cdot \tilde{\gamma}(0)) = 0. \quad (2.4.17)$$

Using (2.4.10), we deduce that the matrices A and B are strictly positive definite. Thus, (2.4.17) implies that $\tilde{\gamma}(s) \equiv 0$, so that the minimizer is unique. Hence, $u(t, x)$ is well-defined by (2.4.6):

$$\begin{aligned} u(t, x) &= u_0(\gamma(0)) + \int_0^t \left(-\frac{|p(s)|^2}{2} + R(\gamma(s)) + \dot{\gamma}(s) \cdot p(s) \right) ds \\ &= u_0(\gamma(0)) + \int_0^t \left(\frac{|\dot{\gamma}(s)|^2}{2} + R(\gamma(s)) \right) ds. \end{aligned} \quad (2.4.18)$$

This may be rephrased as

$$u(t, x) = \inf_{\gamma(\ell)=x} \left(u_0(\gamma(0)) + \int_0^t \left(\frac{|\dot{\gamma}(s)|^2}{2} + R(\gamma(s)) \right) ds \right). \quad (2.4.19)$$

This formula, known as the Lax-Oleinik formula, is the starting point of the Lagrangian theory of Hamilton-Jacobi equations, and has immense implications. We will spend some time with this aspect of Hamilton-Jacobi equations later in this chapter. We will see that we can take it as a good definition of a solution to the Cauchy problem, at least when the Hamiltonian is strictly convex in p .

Smoothness of the solution

Let us quickly examine the smoothness of $u(t, x)$ in x in the set-up of the present section. We see from (2.4.13) that it is equivalent to the smoothness of the minimizer γ in x . If $h \in \mathbb{R}$ and $i \in \{1, \dots, n\}$, consider the partial difference

$$\gamma_h^i(s) = \frac{\gamma_{t,x+he_i}(s) - \gamma_{t,x}(s)}{h}.$$

It solves a system similar to (2.4.15), except for the boundary condition at $s = t$ that is now $\gamma_h^i(t) = e_i$. The exact same integration by parts argument yields the uniform boundedness of $\|\gamma_h^i\|_{H^1}$, hence the uniform boundedness of γ_h^i . Sending h to 0 and repeating the analysis shows that γ_h^i converges to the unique solution of an equation of the type (2.4.15), with

$$A(s) = D^2R(\gamma_{t,x}(s)), \quad B = D^2u_0(\gamma_{t,x}(0)).$$

This argument may be repeated over and over again, to yield the C^∞ smoothness of $\gamma_{t,x}$ in t and x , as long as u_0 and $R(x)$ are infinitely differentiable. Finally, using (2.4.6) we can conclude that

$$u(t, x) = \bar{J}_{t,x},$$

is infinitely differentiable as well.

Exercise 2.4.6 Show that u is convex in x , for all $t > 0$, in two ways. First, fix $\xi \in \mathbb{R}^n$ and get a differential equation for $Q(t, x) = (D^2u(t, x)\xi \cdot \xi)$. Use a maximum principle type argument to conclude that $Q(t, x) > 0$ for all $t > 0$ and $x \in \mathbb{R}^n$. An alternative and more elegant way is to proceed as follows.

- (i) Assume the existence of $\kappa > 0$ such that, for all $(x, y) \in \mathbb{R}^n \times \mathbb{R}^n$ and $\sigma \in [0, 1]$, we have:

$$u(t, \sigma x + (1 - \sigma)y) \leq \sigma u(t, x) + (1 - \sigma)u(t, y) - \kappa\sigma(1 - \sigma)|x - y|^2. \quad (2.4.20)$$

Show that then the function $u(t, x)$ is strictly convex.

- (ii) Show that there exists $\lambda > 0$ such that if $\gamma_{t,x}$ and $\gamma_{t,y}$ are, respectively, the minimizing curves for $u(t, x)$ and $u(t, y)$, then

$$u(t, \sigma x + (1 - \sigma)y) \leq \sigma u(t, x) + (1 - \sigma)u(t, y) - \lambda\sigma(1 - \sigma) \left(|\gamma_{t,x}(0) - \gamma_{t,y}(0)|^2 + \|\gamma_{t,x} - \gamma_{t,y}\|_{H^1([0,t])}^2 \right). \quad (2.4.21)$$

Hint: use the test curve $\gamma_\sigma = \sigma\gamma_{t,x} + (1 - \sigma)\gamma_{t,y}$ in the Lax-Oleinik formula (2.4.19) for $u(t, \sigma x + (1 - \sigma)y)$, together with the convexity of the functions $u_0(x)$ and $R(x)$.

(iii) Finish the proof of (2.4.20), by noticing that

$$|\gamma_{t,x}(0) - \gamma_{t,y}(0)|^2 = |x - y|^2 - \int_0^t \frac{d}{ds} |\gamma_{t,x}(s) - \gamma_{t,y}(s)|^2 ds.$$

The qualitative behavior of $u(t, x)$ can be investigated further, implying the large time stabilization of the whole solution. We will come back to this class of questions later, when we study the large time behavior of viscosity solutions on the torus. For the time being, we leave the classical theory.

2.5 Viscosity solutions

We have just seen that, in order to find reasonable solutions to an inviscid Hamilton-Jacobi equation

$$u_t + H(x, \nabla u) = 0, \tag{2.5.1}$$

we should relax the constraint that " u is continuously differentiable". The first idea would be to replace it by " u is Lipschitz", and require (2.5.1) to hold almost everywhere. Alas, there are, in general, several such solutions to the Cauchy problem for (2.5.1) with a Lipschitz (or even smooth) initial condition. This parallels the fact that the weak solutions to the conservation laws are not unique – for uniqueness, one must require that the weak solution satisfies the entropy condition. See, for instance, [?] for a discussion of these issues. A simple illustration of this phenomenon is to consider the Hamilton-Jacobi equation

$$u_t + \frac{1}{2}u_x^2 = 0, \tag{2.5.2}$$

in one dimension, with the Lipschitz continuous initial condition

$$u_0(x) = 0 \text{ for } x \leq 0 \text{ and } u_0(x) = x \text{ for } x > 0. \tag{2.5.3}$$

It is easy to check that one Lipschitz solution to (2.5.2) that satisfies this equation almost everywhere and obeys the initial condition (2.5.3) is

$$u^{(1)}(t, x) = 0 \text{ for } x < t/2 \text{ and } u^{(1)}(t, x) = x - t/2 \text{ for } x > t/2.$$

However, another solution to (2.5.4)-(2.5.3) is given by

$$u^{(2)}(t, x) = 0 \text{ for } x < 0, u^{(2)}(t, x) = \frac{x^2}{2t} \text{ for } 0 < x < t \text{ and } u^{(2)}(t, x) = x - \frac{t}{2} \text{ for } x > t.$$

Exercise 2.5.1 Consider the solution $u^\varepsilon(t, x)$ to a viscous version of (2.5.4):

$$u_t^\varepsilon + \frac{1}{2}(u_x^\varepsilon)^2 = \varepsilon u_{xx}^\varepsilon, \tag{2.5.4}$$

also with the initial condition $u^\varepsilon(0, x) = u_0(x)$, as in (2.5.3). Use the Hopf-Cole transform

$$v^\varepsilon(t, x) = \exp\left(-\frac{u^\varepsilon(t/\varepsilon, x)}{2\varepsilon}\right),$$

to show that v^ε satisfies the standard heat equation

$$v_t^\varepsilon = v_{xx}^\varepsilon.$$

Find $v^\varepsilon(t, x)$ explicitly and use this to show that

$$u^\varepsilon(t, x) \rightarrow u^{(2)}(t, x) \text{ as } \varepsilon \rightarrow 0.$$

A natural question is, therefore, to know if an additional condition, less stringent than the C^1 -regularity, but stronger than the mere Lipschitz regularity, enables us to select a unique solution to the Cauchy problem – as the notion of the entropy solutions does for the conservation laws. Exercise 2.5.1 suggests that regularizing the inviscid Hamilton-Jacobi equation with a small diffusion can provide one such approach, but for more general Hamilton-Jacobi equations than (2.5.4), for which the Hopf-Cole transform is not available, this procedure would be much less explicit.

The above considerations have motivated the introduction, by Crandall and Lions [?], at the beginning of the 1980's, of the notion of a *viscosity solution* to (2.1.1). The idea is to select, among all the solutions of (2.1.1), “the one that has a physical meaning”, intrinsically, without directly appealing to the small diffusion regularization, – though understanding the connection to physics may require some additional thought. Being weaker than the notion of a classical solution, it introduces new difficulties to the existence, regularity and uniqueness issues, as well as into getting insight into the qualitative properties of solutions.

Finally, looking ahead, we mention that even if there is a unique viscosity solution to the Cauchy problem associated to (2.1.1), there will be no clear reason for the stationary equation (2.1.6) to have a unique steady viscosity solution – unlike what we have seen in the diffusive situation.

As a concluding remark to this introduction, we must mention that we will by no means do justice to a very rich subject in this short section and provide just a brief glance of a still developing subject. The reader interested to learn more may enjoy reading Barles [?], or Lions [?] as a starting point.

2.5.1 The definition and the basic properties of the viscosity solutions

The definition of a viscosity solution

Let us begin with more general equations than (2.1.1) – we will restrict the assumptions as the theory develops. Consider the Cauchy problem

$$u_t + F(x, u, \nabla u) = 0, \quad t > 0, \quad x \in \mathbb{T}^n, \tag{2.5.5}$$

with a continuous initial condition $u(0, x) = u_0(x)$, and $F \in C(\mathbb{T}^n \times \mathbb{R} \times \mathbb{R}^n; \mathbb{R})$.

In order to motivate the notion of a viscosity solution, one takes the point of view that the smooth solutions to the regularized problem

$$u_t^\varepsilon + F(x, u^\varepsilon, \nabla u^\varepsilon) = \varepsilon \Delta u^\varepsilon \tag{2.5.6}$$

are a good approximation to $u(t, x)$. Existence of the solution to the Cauchy problem for (2.5.6) for $\varepsilon > 0$ is not really an issue – we have already seen how it can be proved. Hence, a natural attempt would be to pass to the limit $\varepsilon \downarrow 0$ in (2.5.6). It is possible to prove that there is a unique limiting solution and that one actually ends up with a nonlinear semigroup. In particular, one may show that, if we take this notion of solution as a definition, there are uniqueness and contraction properties analogous to what we will see below – see [?] for further details. Taking this limit as a definition, however, raises an important issue: there is always the danger that the solution depends on the underlying regularization – why regularize with the Laplacian? What if we were to regularize differently? For instance, what if we would consider a dispersive regularization in one dimension

$$u_t^\varepsilon + F(x, u^\varepsilon, u_x^\varepsilon) = \varepsilon u_{xxx}^\varepsilon, \quad x \in \mathbb{R}, \quad (2.5.7)$$

which is a generalized Korteweg-de Vries equation, and let $\varepsilon \rightarrow 0$ in (2.5.7) instead?

We now describe an alternative and more intrinsic approach, instead of using (2.5.6) in this very direct fashion of passing to the limit $\varepsilon \downarrow 0$. The idea is that the key property that should be inherited from the diffusive regularization is the maximum principle, as it is usually inherent in the origins of such models in the corresponding applications, be it physics, such as motion of interfaces, or optimal control problems. There is an interesting separate question of what happens as $\varepsilon \rightarrow 0$ to the solutions coming from regularizations that do not admit the maximum principle, such as (2.5.7). The situation is not quite trivial, especially for non-convex fluxes F – we refer an interested reader to [?].

Our approach will be to use the comparison principle idea to extend the notions of a sub-solution and a super-solution to (2.5.5) and then simply say that a function $u(t, x)$ is a solution to (2.5.5) if it is both a sub-solution and a super-solution. To understand the upcoming definition of a viscosity sub-solution to (2.5.5), consider first a smooth sub-solution $u(t, x)$ to the regularized problem (2.5.6):

$$u_t + F(x, u, \nabla u) \leq \varepsilon \Delta u. \quad (2.5.8)$$

Let us take a smooth function $\phi(t, x)$ such that the difference $\phi - u$ attains its minimum at a point (t_0, x_0) . One may simply think of the case when $\phi(t_0, x_0) = u(t_0, x_0)$ and $\phi(t, x) \geq u(t, x)$ elsewhere. Then, at this point we have

$$u_t(t_0, x_0) = \phi_t(t_0, x_0), \quad \nabla \phi(t_0, x_0) = \nabla u(t_0, x_0),$$

and

$$D^2 \phi(t_0, x_0) \geq D^2 u(t_0, x_0),$$

in the sense of the quadratic forms. It follows that

$$\begin{aligned} & \phi_t(t_0, x_0) + F(x_0, u(t_0, x_0), \nabla \phi(t_0, x_0)) - \varepsilon \Delta \phi(t_0, x_0) \\ & \leq u_t(t_0, x_0) + F(x_0, u(t_0, x_0), \nabla u(t_0, x_0)) - \varepsilon \Delta u(t_0, x_0) \leq 0. \end{aligned} \quad (2.5.9)$$

In other words, if u is a smooth sub-solution to (2.5.6), and ϕ is a smooth function that touches u at the point (t_0, x_0) from above, then ϕ is also a sub-solution to (2.5.6) at this point.

In a similar vein, if $u(t, x)$ is a smooth super-solution to the regularized problem:

$$u_t + F(x, u, \nabla u) \geq \varepsilon \Delta u, \quad (2.5.10)$$

we consider a smooth function $\phi(t, x)$ such that the difference $\phi - u$ attains its maximum at a point (t_0, x_0) . Again, we may assume without loss of generality that $\phi(t_0, x_0) = u(t_0, x_0)$ and $\phi(t, x) \leq u(t, x)$ elsewhere. Then, at this point we have

$$\phi_t(t_0, x_0) + F(x_0, u(t_0, x_0), \nabla \phi(t_0, x_0)) - \varepsilon \Delta \phi(t_0, x_0) \geq 0, \quad (2.5.11)$$

by a computation similar to (2.5.9). That is, if u is a smooth super-solution to (2.5.6), and ϕ is a smooth function that touches u at (t_0, x_0) from below, then ϕ is also a super-solution to (2.5.6) at this point.

These two observations lead to the following definition, where we simply drop the requirement that u is smooth, only use the regularity of the test function that touches it from above or below, and send $\varepsilon \rightarrow 0$ in (2.5.9) and (2.5.11).

Definition 2.5.2 *A continuous function $u(t, x)$ is a viscosity sub-solution to*

$$u_t + F(x, u, \nabla u) = 0, \quad (2.5.12)$$

if, for all test functions $\phi \in C^1([0, +\infty) \times \mathbb{T}^n)$ and all $(t_0, x_0) \in (0, +\infty) \times \mathbb{T}^n$ such that (t_0, x_0) is a local minimum for $\phi - u$, we have:

$$\phi_t(t_0, x_0) + F(x_0, u(t_0, x_0), \nabla \phi(t_0, x_0)) \leq 0. \quad (2.5.13)$$

Furthermore, a continuous function $u(t, x)$ is a viscosity super-solution to (2.5.12) if, for all test functions $\phi \in C^1((0, +\infty) \times \mathbb{T}^n)$ and all $(t_0, x_0) \in (0, +\infty) \times \mathbb{T}^n$ such that the point (t_0, x_0) is a local maximum for the difference $\phi - u$, we have:

$$\phi_t(t_0, x_0) + F(x_0, u(t_0, x_0), \nabla \phi(t_0, x_0)) \geq 0. \quad (2.5.14)$$

Finally, $u(t, x)$ is a viscosity solution to (2.5.12) if it is both a viscosity sub-solution and a viscosity super-solution to (2.5.12).

Definition 2.5.2 extends to steady equations of the type

$$F(x, u, \nabla u) = 0 \text{ on } \mathbb{T}^n,$$

by requiring that $u(x)$ is a viscosity sub-solution (respectively, super-solution) to

$$u_t + F(x, u, \nabla u) = 0,$$

that happens to be time-independent.

This definition was introduced by Crandall and Lions in their seminal paper [?]. The name “viscosity solution” comes out of the diffusive regularization we have discussed above. Definition 2.5.2 is intrinsic and bypasses the philosophical question we have mentioned above: “Why regularize with the Laplacian?” much like the notion of an entropy solution does this for the conservation laws. We stress, however, that it does make the assumption that the

underlying model must respect the comparison principle. Let us also note that the notion of a viscosity solution has turned out to be also very much relevant to the second order elliptic and parabolic equations – especially those fully nonlinear with respect to the Hessian of the solution. There have been spectacular developments, which are out of the scope of this chapter.

The main issue we will need to face soon is whether such a seemingly weak definition has any selective power – can it possibly ensure uniqueness of the solution? The expectation is that it should, due to the general principle that "the comparison principle implies uniqueness".

First, the following exercises may help the reader gain some intuition.

Exercise 2.5.3 Show that a C^1 solution to

$$u_t + F(x, u, \nabla u) = 0, \quad t > 0, \quad x \in \mathbb{T}^n, \quad (2.5.15)$$

is a viscosity solution.

Exercise 2.5.4 Consider the Hamilton-Jacobi equation

$$u_t + u_x^2 = 0, \quad x \in \mathbb{R}. \quad (2.5.16)$$

(i) Which of the following two functions is a viscosity solution to (2.5.16): $v(t, x) = |x| - t$ or $w(t, x) = -t - |x|$? Hint: pay attention to the fact that at the point $x = 0$ a smooth function $\phi(t, x)$ can only touch $v(t, x)$ from the bottom, and $w(t, x)$ from the top. This will tell you something about $|\phi_x(t, 0)|$ and determine the answer to this question.

(ii) Consider (2.5.16) with a zigzag initial condition $u_0(x) = u(0, x)$:

$$u_0(x) = \begin{cases} x, & 0 \leq x \leq 1/2, \\ 1 - x, & 1/2 \leq x \leq 1, \end{cases} \quad (2.5.17)$$

extended periodically to \mathbb{R} . How will the viscosity solution $u(t, x)$ to the Cauchy problem look like? Where will it be smooth, and where will it be just Lipschitz? Hint: it may help to do this in at least two ways: (1) use the definition of the viscosity solution, (2) use the notion of the entropy solution for the Burgers' equation for $v(t, x) = u_x(t, x)$ if you are familiar with the basic theory of one-dimensional conservation laws.

Exercise 2.5.5 (*Intermezzo: a Laplace asymptotics of integrals*). Let $\varphi : \mathbb{R}^n \rightarrow \mathbb{R}$ be a real-valued smooth function such that there are two positive constants α and β such that

$$\varphi(x) \geq \alpha|x|^2 - \beta.$$

For $\varepsilon > 0$, consider the integral

$$I_\varepsilon = \int_{\mathbb{R}^n} e^{-\varphi(x)/\varepsilon} dx.$$

The goal of this exercise is to show that

$$\lim_{\varepsilon \rightarrow 0} (-\varepsilon \log I_\varepsilon) = \min_{x \in \mathbb{R}^n} \varphi(x). \quad (2.5.18)$$

Note that it suffices to assume that

$$\min_{x \in \mathbb{R}^n} \varphi(x) = 0, \quad (2.5.19)$$

and show that

$$\lim_{\varepsilon \rightarrow 0} (-\varepsilon \log I_\varepsilon) = 0. \quad (2.5.20)$$

Exercise 2.5.6 Let us add the term εu_{xx} to the right side of (2.5.16), which produces a solution $u_\varepsilon(t, x)$. Use the Hopf-Cole transformation $z_\varepsilon(t, x) = \exp(u_\varepsilon(t, x)/\varepsilon)$, solve the linear problem for $z(t, x)$ and then pass to the limit $\varepsilon \rightarrow 0$ using Exercise 2.5.5. Study what happens when $u'_0(x)$ has limits at $\pm\infty$.

Basic properties of the viscosity solutions

We now describe some basic corollaries of the definition of a viscosity solution.

Exercise 2.5.7 Show that the maximum of two viscosity subsolutions to (2.5.15) is a viscosity subsolution, and the minimum of two viscosity supersolutions is a viscosity supersolution.

Exercise 2.5.8 (Stability) Let $F_j(x, u, p)$ be a sequence of functions in $C(\mathbb{T}^n \times \mathbb{R} \times \mathbb{R}^n)$, which converges locally uniformly to $F \in C(\mathbb{T}^n \times \mathbb{R} \times \mathbb{R}^n)$. Let $u_j(t, x)$ be a sequence of viscosity solutions to (2.5.5) with $F = F_j$:

$$\partial_t u_j + F_j(x, u_j, \nabla u_j) = 0, \quad t > 0, \quad x \in \mathbb{T}^n, \quad (2.5.21)$$

and assume that u_j converges locally uniformly to $u \in C([0, +\infty), \mathbb{T}^n)$. Show that then u is a viscosity solution to the limiting problem

$$u_t + F(x, u, \nabla u) = 0, \quad t > 0, \quad x \in \mathbb{T}^n, \quad (2.5.22)$$

Hint: if (t_0, x_0) is, for instance, a local minimum of the difference $\phi - u$, one can turn it into a strict minimum by changing $\phi(t, x)$ into $\phi(x) + M((t - t_0)^2 + |x - x_0|^2)$, without changing $\phi_t(t_0, x_0)$ and $\nabla\phi(t_0, x_0)$. In this situation, show that there is a sequence (t_j, x_j) of minima of $\phi - j$ converging to (t_0, x_0) , and use the fact that each u_j is a viscosity solution to (2.5.21) to conclude.

The above exercise is extremely important: it shows that, somewhat similar to the weak solutions, we do not need to check the convergence of the derivatives of u_j to the derivatives of u to know that u is a viscosity solution – this is an extremely useful property to have. Exercise 2.5.8 asserts that one may safely “pass to the limit” in equation (2.5.5), as soon as estimates on the moduli of continuity of the solutions are available rather than on the derivatives.

The next proposition says that viscosity solutions that are Lipschitz continuous do satisfy the equation in the classical sense almost everywhere.

Proposition 2.5.9 *Let u be a locally Lipschitz viscosity solution to*

$$u_t + F(x, u, \nabla u) = 0, \quad t > 0, \quad x \in \mathbb{T}^n. \quad (2.5.23)$$

Then it satisfies (2.5.23) almost everywhere.

This property relies on the following lemma [?].

Lemma 2.5.10 *At a point of differentiability (t_0, x_0) of u , one may construct a C^1 test function $\phi(t, x)$ such that (t_0, x_0) is a local maximum (respectively, a local minimum) of $\phi - u$.*

Proof. For simplicity, we do not take the t -dependence into account, leaving this to the reader, so we assume that $u(x)$ is a function of x that is differentiable at x_0 . Without loss of generality, we assume that $x_0 = 0$, $u(0) = 0$, and that $\nabla u(0) = 0$, so that $u(x)$ satisfies, in the vicinity of $x = 0$:

$$u(x) = |x|\varepsilon(x), \quad \lim_{|x| \rightarrow 0} \varepsilon(x) = 0. \quad (2.5.24)$$

Our goal is to construct a C^1 function $\phi(x)$ such that $\phi(x) \leq u(x)$ and $\phi(0) = 0$. Note that this could be impossible for $u(x)$ that is merely Lipschitz and not differentiable – the simple counterexample is $u(x) = -|x|$. We look for a radially symmetric function $\phi(x)$ in the form $\phi(x) = |x|\zeta(|x|)$ with a C^1 -function $\zeta(r)$ such that

$$\zeta(|x|) \leq \varepsilon(x), \quad \lim_{r \rightarrow 0} \zeta(r) = 0. \quad (2.5.25)$$

To this end, consider the decreasing sequence

$$\varepsilon_n = \inf_{2^{-n-1} \leq |r| < 0} \varepsilon(r),$$

and produce the function $\zeta(r) \leq \varepsilon(r)$ by smoothing the piecewise constant function

$$\sum_{n=0}^{+\infty} \varepsilon_n \mathbf{1}_{2^{-n-1} \leq r < 2^{-n}}.$$

As the sequence $\varepsilon_n \rightarrow 0$ as $n \rightarrow +\infty$, and we have chosen the dyadic intervals in the above sum, we may ensure that

$$|\zeta'(r)| \leq \frac{\alpha(r)}{r},$$

with $\alpha(r) \rightarrow 0$ as $r \rightarrow 0$. It follows that $\phi(x) = |x|\zeta(|x|)$ is the sought C^1 -function. \square

Proof of Proposition 2.5.9. The conclusion of this proposition follows essentially immediately from Lemma 2.5.10 and the Rademacher theorem. The latter says that a Lipschitz function is differentiable a.e., see for instance [?]. At any differentiability point we can construct a C^1 -function $\phi(t, x)$ such that the difference $\phi - u$ attains its minimum at (t_0, x_0) , so that

$$\phi_t(t_0, x_0) = u_t(t_0, x_0) \text{ and } \nabla \phi(t_0, x_0) = \nabla u(t_0, x_0). \quad (2.5.26)$$

The definition of a viscosity sub-solution together with (2.5.26) implies that

$$u_t(t_0, x_0) + H(x, u(t_0, x_0), \nabla u(t_0, x_0)) = \phi_t(t_0, x_0) + H(x, u(t_0, x_0), \nabla \phi(t_0, x_0)) \leq 0.$$

Similarly, we can show that

$$u_t(t_0, x_0) + H(x, u(t_0, x_0), \nabla u(t_0, x_0)) \geq 0,$$

using the fact that $u(t, x)$ is a viscosity super-solution. This finishes the proof. \square

Warning. For the rest of this section, a solution of (2.1.1) or (2.1.6) will always be meant in the viscosity sense.

2.5.2 Uniqueness of the viscosity solutions

Let us first give the simplest uniqueness result, that we will prove by the method of doubling of variables. This argument appears in almost all uniqueness proofs, in more or less elaborate forms.

Proposition 2.5.11 *Assume that the Hamiltonian $H(x, p)$ is continuous in all its variables, and satisfies the coercivity assumption*

$$\lim_{|p| \rightarrow +\infty} H(x, p) = +\infty, \quad \text{uniformly in } x \in \mathbb{T}^n. \quad (2.5.27)$$

Consider the equation

$$H(x, \nabla u) + u = 0, \quad x \in \mathbb{T}^n. \quad (2.5.28)$$

Let \underline{u} and \bar{u} be, respectively, a viscosity sub- and a super-solution to (2.5.28), then $\underline{u} \leq \bar{u}$.

Proof. Assume for a moment that both \underline{u} and \bar{u} are C^1 -functions, so that we can use each of them as a test function in the definition of the viscosity sub- and super-solutions. First, we use the fact that \bar{u} is a super-solution to (2.5.28) and take \underline{u} as a test function. Let x_0 be a maximum of $\underline{u} - \bar{u}$, then we deduce from the definition of a viscosity super-solution to (2.5.28) that

$$H(x_0, \nabla \underline{u}(x_0)) + \bar{u}(x_0) \geq 0. \quad (2.5.29)$$

On the other hand, $\bar{u} - \underline{u}$ attains its minimum at the same point x_0 , and, as \underline{u} is a viscosity sub-solution to (2.5.28), and \bar{u} can serve as a test function, we have

$$H(x_0, \nabla \bar{u}(x_0)) + \underline{u}(x_0) \leq 0. \quad (2.5.30)$$

As x_0 is a minimum of $\bar{u} - \underline{u}$, and \underline{u} and \bar{u} are differentiable, we have $\nabla \bar{u}(x_0) = \nabla \underline{u}(x_0)$, whence (2.5.29) and (2.5.30) imply

$$\underline{u}(x_0) \leq \bar{u}(x_0).$$

Once again, as $\bar{u} - \underline{u}$ attains its minimum at x_0 , we conclude that $\bar{u}(x) \geq \underline{u}(x)$ for all $x \in \mathbb{T}^n$ if both of these functions are in $C^1(\mathbb{T}^n)$. Unfortunately, we only know that \underline{u} and \bar{u} are continuous, so we can not use the elegant argument above unless we know, in addition, that they are both C^1 -functions.

In the general case, the method of doubling the variables gives a way around the technical difficulty that $\underline{u}(x)$ and $\bar{u}(x)$ are merely continuous and not necessarily differentiable. Let us define, for $\varepsilon > 0$, the penalization

$$u_\varepsilon(x, y) = \bar{u}(x) - \underline{u}(y) + \frac{|x - y|^2}{2\varepsilon^2}$$

and let $(x_\varepsilon, y_\varepsilon) \in \mathbb{T}^{2n}$ be a minimum for $u_\varepsilon(x, y)$.

Exercise 2.5.12 Show that

$$\lim_{\varepsilon \rightarrow 0} |x_\varepsilon - y_\varepsilon| = 0. \quad (2.5.31)$$

and that the family $(x_\varepsilon, y_\varepsilon)$ converges, as $\varepsilon \rightarrow 0$, up to a subsequence, to a point (x_0, x_0) , where x_0 is a minimum for $\bar{u}(x) - \underline{u}(x)$.

Consider the function

$$\phi(x) = \underline{u}(y_\varepsilon) - \frac{|x - y_\varepsilon|^2}{2\varepsilon^2},$$

as a (smooth) quadratic function of the variable x . The difference

$$\phi(x) - \bar{u}(x) = -u_\varepsilon(x, y_\varepsilon)$$

attains its maximum, as a function of x , at the point $x = x_\varepsilon$. As $\bar{u}(x)$ is a viscosity super-solution to (2.5.28), we have

$$H(x_\varepsilon, \frac{y_\varepsilon - x_\varepsilon}{\varepsilon^2}) + \bar{u}(x_\varepsilon) \geq 0. \quad (2.5.32)$$

Next, we apply the viscosity sub-solution part of Definition 2.5.13 to the quadratic test function

$$\psi(y) = \bar{u}(x_\varepsilon) + \frac{|x_\varepsilon - y|^2}{2\varepsilon^2}.$$

The difference

$$\psi(y) - \underline{u}(y) = \bar{u}(x_\varepsilon) + \frac{|x_\varepsilon - y|^2}{2\varepsilon^2} - \underline{u}(y) = u_\varepsilon(x_\varepsilon, y)$$

attains its minimum at $y = y_\varepsilon$. As $\underline{u}(y)$ is a viscosity sub-solution to (2.5.28), we obtain

$$H(y_\varepsilon, \frac{y_\varepsilon - x_\varepsilon}{\varepsilon^2}) + \underline{u}(y_\varepsilon) \leq 0. \quad (2.5.33)$$

The coercivity of the Hamiltonian and (2.5.33), together with the boundedness of $\underline{u}_\varepsilon$, imply that $|x_\varepsilon - y_\varepsilon|/\varepsilon^2$ is bounded, uniformly in ε : there exists R so that

$$\frac{|x_\varepsilon - y_\varepsilon|}{\varepsilon^2} \leq R.$$

The uniform continuity of $H(x, p)$ on the set $\{(x, p) : x \in \mathbb{T}^n, p \in B(0, R)\}$ implies that, as consequence, we have

$$H(y_\varepsilon, \frac{y_\varepsilon - x_\varepsilon}{\varepsilon^2}) = H(x_\varepsilon, \frac{y_\varepsilon - x_\varepsilon}{\varepsilon^2}) + o(1), \text{ as } \varepsilon \rightarrow 0.$$

Subtracting (2.5.33) from (2.5.32), we obtain

$$\bar{u}(x_\varepsilon) - \underline{u}(y_\varepsilon) \geq o(1), \text{ as } \varepsilon \rightarrow 0.$$

Sending $\varepsilon \rightarrow 0$ with the help of the result of Exercise 2.5.12 implies

$$\bar{u}(x_0) - \underline{u}(x_0) \geq 0,$$

and, as x_0 is the minimum of $\bar{u} - \underline{u}$, the proof is complete. \square

An immediate consequence of Proposition 2.5.11 is that (2.5.28) has at most one solution.

The comparison principle and weak contraction

The proof of Proposition 2.5.11 can be adapted to establish two fundamental properties for the viscosity solutions to the Cauchy problem: the comparison principle and the weak contraction property.

Exercise 2.5.13 (The comparison principle) Assume that $H(x, p)$, is a continuous function that satisfies the coercivity property (2.5.27). Let $u_1(t, x)$ and $u_2(t, x)$ be, respectively, a viscosity sub-solution, and a viscosity super-solution to

$$u_t + H(x, \nabla u) = 0, \quad t > 0, \quad x \in \mathbb{T}^n, \quad (2.5.34)$$

with the initial conditions u_{10} and u_{20} such that $u_{10}(x) \leq u_{20}(x)$ for all $x \in \mathbb{T}^n$. Modify the proof of Proposition 2.5.11 to show that then $u_1(t, x) \leq u_2(t, x)$ for all $t \geq 0$ and $x \in \mathbb{T}^n$. This proves the uniqueness of the viscosity solutions.

Exercise 2.5.14 (Weak contraction) Let $H(x, p)$ be a continuous function that satisfies the coercivity property (2.5.27), and u_1 and u_2 be two solutions to (2.5.34) with continuous initial conditions u_{10} and u_{20} , respectively. Show that then we have

$$\|u_1(t, \cdot) - u_2(t, \cdot)\|_{L^\infty} \leq \|u_{10} - u_{20}\|_{L^\infty}.$$

Hint: notice that if $u(t, x)$ solves (2.5.34) then so does $u(t, x) + k$ for any $k \in \mathbb{R}$, and use Exercise 2.5.13.

2.5.3 Finite speed of propagation

We are now going to prove a finite speed of propagation property, partly to acquire some further familiarity with the notion of a viscosity solution, and partly to emphasize the sharp contrast with viscous models: if the equation carried a Laplacian, an initially nonnegative solution would instantly become positive everywhere. As this property makes better sense in \mathbb{R}^n and not on the torus, this is the case we will consider.

Proposition 2.5.15 *Let H be uniformly Lipschitz with respect to its second variable:*

$$|H(x, p_1) - H(x, p_2)| \leq C_L |p_1 - p_2| \quad \text{for all } x \in \mathbb{R}^n \text{ and } p_1, p_2 \in \mathbb{R}^n. \quad (2.5.35)$$

Let u_0 and v_0 be two continuous, compactly supported initial conditions, and assume that each generates a globally Lipschitz solution, respectively denoted by $u(t, x)$ and $v(t, x)$ to the Cauchy problem

$$u_t + H(x, \nabla u) = 0, \quad v_t + H(x, \nabla v) = 0, \quad 0 < t \leq T, \quad x \in \mathbb{R}^n, \quad (2.5.36)$$

with $u(0, x) = u_0(x)$ and $v(0, x) = v_0(x)$ for all $x \in \mathbb{R}^n$. Then, if $x_0 \in \mathbb{R}^n$ and $t_0 \in [0, T]$ satisfy

$$\text{dist}(x_0, \text{supp}(u_0 - v_0)) > t_0 C_L,$$

then $u(t_0, x_0) = v(t_0, x_0)$.

Proof. The idea is simple: assuming that everything is smooth, the function $w = u - v$ satisfies the inequalities

$$w_t \leq C_L |\nabla w|, \quad (2.5.37)$$

and

$$w_t \geq -C_L |\nabla w|. \quad (2.5.38)$$

Exercise 2.5.16 Use the method of characteristics to show that if w is a smooth function that satisfies (2.5.37) and

$$\text{dist}(x_0, \text{supp}(w(0, \cdot))) > C_L t_0, \quad (2.5.39)$$

then $w(t_0, x_0) \leq 0$, and if a smooth function w satisfies (2.5.38)-(2.5.39), then $w(t_0, x_0) \geq 0$.

Thus, the conclusion of this proposition follows from Exercise 2.5.16 if u and v are smooth. Unfortunately, if u and v are not smooth, then we can not use the characteristics but only the definition of a viscosity solution. Let us fix a point $x_0 \in \mathbb{R}^n$ and $t_0 > 0$ so that

$$\text{dist}(x_0, \text{supp}(u_0 - v_0)) > C_L t_0, \quad (2.5.40)$$

take $\varepsilon > 0$ sufficiently small, so that

$$\varepsilon < \frac{1}{2} (\text{dist}(x_0, \text{supp}(u_0 - v_0)) - C_L t_0), \quad (2.5.41)$$

and let $\phi_0(r)$ be a C^1 -function equal to 1 outside of the ball $B_{C_L t_0 + \varepsilon}(0)$, and to 0 in the ball $B_{C_L t_0}(0)$. The function

$$\bar{w}(t, x) = \|u_0 - v_0\|_{L^\infty} \phi_0(|x - x_0| + C_L t) \quad (2.5.42)$$

is a smooth solution to

$$\partial_t \bar{w} - C_L |\nabla \bar{w}| = 0, \quad t > 0, \quad x \in \mathbb{R}^n, \quad (2.5.43)$$

such that $\bar{w}(t, x_0) = 0$ for $t \leq t_0$. Moreover, because of (2.5.41), at $t = 0$ we have

$$\bar{w}(0, x) = \|u_0 - v_0\|_{L^\infty} \phi_0(|x - x_0|) \geq |u_0(x) - v_0(x)| \text{ for all } x \in \mathbb{R}^n. \quad (2.5.44)$$

Our goal is to show this inequality persists until the time t_0 :

$$|u(t, x) - v(t, x)| \leq \bar{w}(t, x) \text{ for all } 0 \leq t \leq t_0 \text{ and } x \in \mathbb{R}^n. \quad (2.5.45)$$

Indeed, using (2.5.45) at $x = x_0$ and $t = t_0$ would give

$$|u(t_0, x_0) - v(t_0, x_0)| \leq \|u_0 - v_0\|_{L^\infty} \phi_0(C_L t_0) = 0, \quad (2.5.46)$$

which is what we need.

The comparison principle in Exercise 2.5.13 together with (2.5.44) implies that (2.5.45) would follow if we show that $\bar{v}(t, x) = u(t, x) + \bar{w}(t, x)$ is a viscosity super-solution to (2.5.36):

$$\partial_t \bar{v} + H(x, \nabla \bar{v}) \geq 0. \quad (2.5.47)$$

Observe that (2.5.47) and (2.5.44) together would imply

$$v(t, x) \leq \bar{v}(t, x) = u(t, x) + \bar{w}(t, x) \text{ for all } 0 \leq t \leq t_0 \text{ and } x \in \mathbb{R}^n. \quad (2.5.48)$$

As the roles of u and v can be reversed, we would deduce (2.5.45).

Thus, we need to prove the viscosity super-solution property for $\bar{v}(t, x)$. Let $\varphi(t, x)$ be a smooth test function, and (t_1, x_1) be a minimum point for

$$\bar{v}(t, x) - \varphi(t, x) = u(t, x) + \bar{w}(t, x) - \varphi(t, x) = u(t, x) - \psi(t, x), \quad (2.5.49)$$

with a C^1 -function

$$\psi(t, x) = \varphi(t, x) - \bar{w}(t, x).$$

In other words, (t_1, x_1) is a minimum point for $u(t, x) - \psi(t, x)$. As u is a viscosity solution to (2.5.36), it follows that

$$\partial_t \psi(t_1, x_1) + H(x_1, \nabla \psi(t_1, x_1)) \geq 0, \quad (2.5.50)$$

which is nothing but

$$\partial_t \varphi(t_1, x_1) - \partial_t \bar{w}(t_1, x_1) + H(x_1, \nabla \varphi(t_1, x_1) - \nabla \bar{w}(t_1, x_1)) \geq 0, \quad (2.5.51)$$

Using the inequality

$$H(\bar{x}, \nabla \varphi - \nabla \bar{w}) \leq H(\bar{x}, \nabla \varphi) + C_L |\nabla \bar{w}|.$$

in (2.5.51) gives

$$\partial_t \varphi(t_1, x_1) - \partial_t \bar{w}(t_1, x_1) + H(x_1, \nabla \varphi(t_1, x_1)) + C_L |\nabla \bar{w}(t_1, x_1)| \geq 0. \quad (2.5.52)$$

Recalling (2.5.43), we obtain

$$\partial_t \varphi(t_1, x_1) + H(x_1, \nabla \varphi(t_1, x_1)) \geq 0. \quad (2.5.53)$$

We conclude that $\bar{v}(t, x)$ is a viscosity super-solution to (2.5.36), finishing the proof. \square

Exercise 2.5.17 (*Hole filling property*). Let $u(t, x)$ be a viscosity solution to

$$u_t = R(t, x) |\nabla u|, \quad t > 0, \quad x \in \mathbb{R}^n,$$

with $R(t, x) \geq R_0 > 0$. Assume that (i) $u(0, x) = u_0(x) \geq \delta_0 > 0$ outside a ball $B(0, R)$, and (ii) the set $\mathbb{R}^n \setminus (\text{supp}(u_0))$ is compact. Prove that there exists $T_0 > 0$ such that $u(t, x) > 0$ for all $t \geq T_0$, and all $x \in \mathbb{R}^n$.

2.6 Construction of solutions

So far, we have set up a beautiful notion of viscosity solutions, and we have also proved that this works well in settling our worries about uniqueness, distinguishing them from the Lipschitz solutions. Now, we have to prove that, as far as existence is concerned, this new notion does better than the classical solutions, in the sense that solutions to the Cauchy

problem exist for all $t > 0$ under reasonable assumptions. In this section, we will show how these solutions can be constructed. First, we will produce wave solutions to the time-dependent problem

$$\partial_t u + H(x, \nabla u) = 0, \quad x \in \mathbb{T}^n. \quad (2.6.1)$$

Next, we are going to prove that the Cauchy problem for (2.6.1) is well-posed as soon as a continuous initial condition is specified. Eventually, we will show that the wave solutions describe the long time behavior of the solutions to the Cauchy problem.

2.6.1 Existence of waves, and the Lions-Papanicolaou-Varadhan theorem

Wave solutions for (2.6.1) will be sought in the same form as viscous waves, that is

$$v(t, x) = -ct + u(x), \quad (2.6.2)$$

with a constant $c \in \mathbb{R}$. A function $v(t, x)$ of this form is a solution to (2.6.1) if $u(x)$ solves a time-independent problem

$$H(x, \nabla u) = c, \quad x \in \mathbb{T}^n. \quad (2.6.3)$$

Exercise 2.6.1 Show that a function $v(t, x)$ of the form (2.6.2) is a viscosity solution to (2.6.1) if and only if $u(x)$ is a viscosity solution to (2.6.3).

As before, we will think of $v(t, x)$ as the height of an interface, and refer to the constant c as the speed of the wave, and to $u(x)$ as its shape. Let us point out that the speed is unique: (2.6.3) may have viscosity solutions for at most one c . Indeed, assume there exist $c_1 \neq c_2$, such that (2.6.3) has a viscosity solution u_1 for $c = c_1$ and another viscosity solution u_2 for $c = c_2$. Let $K > 0$ be such that

$$u_1(x) - K \leq u_2(x) \leq u_1(x) + K, \quad \text{for all } x \in \mathbb{T}^n.$$

By Exercise 2.6.1 the functions

$$v_{1,\pm}(t, x) = -c_1 t + u_1(x) \pm K$$

and

$$v_2(t, x) = -c_2 t + u_2(x)$$

are the viscosity solutions to the Cauchy problem (2.1.1) with the respective initial conditions

$$v_{1,\pm}(x) = u_1(x) \pm K, \quad v_2(0, x) = u_2(x).$$

By the comparison principle (Exercise 2.5.13), we have

$$-c_1 t + u_1(x) - K \leq -c_2 t + u_2(x) \leq -c_1 t + u_1(x) + K, \quad \text{for all } t \geq 0 \text{ and } x \in \mathbb{T}^n.$$

This is a contradiction since $c_1 \neq c_2$, and the functions u_1 and u_2 are bounded. Therefore, the wave speed c is unique. Note that this does not address the question of uniqueness of the shape $u(x)$ – we leave this question for later.

The main result of this section is the following theorem, due to Lions, Papanicolaou, Varadhan [?], that asserts the existence of a constant c for which (2.6.3) has a solution.

Theorem 2.6.2 *Assume that $H(x, p)$ is continuous, uniformly Lipschitz:*

$$|H(x, p_1) - H(x, p_2)| \leq C_L |p_1 - p_2|, \text{ for all } x \in \mathbb{T}^n, \text{ and } p_1, p_2 \in \mathbb{R}^n, \quad (2.6.4)$$

the coercivity condition

$$\lim_{|p| \rightarrow +\infty} H(x, p) = +\infty, \quad \text{uniformly in } x \in \mathbb{T}^n. \quad (2.6.5)$$

holds, and

$$|\nabla_x H(x, p)| \leq K_0(1 + |p|), \text{ for all } x \in \mathbb{T}^n, \text{ and } p \in \mathbb{R}^n. \quad (2.6.6)$$

Then there is a unique $c \in \mathbb{R}$ for which

$$H(x, \nabla u) = c, \quad x \in \mathbb{T}^n. \quad (2.6.7)$$

has a solution.

It is important to point out that the periodicity assumption in x on the Hamiltonian is indispensable – for instance, when $H(x, p)$ is a random function (in x) on $\mathbb{R}^n \times \mathbb{R}^n$, the situation is much more complicated – an interested reader should consult the literature on stochastic homogenization of the Hamilton-Jacobi equations, a research area that is active and evolving at the moment of this writing. We also mention that the only assumption made in [?] is that $H(x, p)$ is continuous and coercive. The Lipschitz condition (2.6.4) in p and (2.6.6) in x have been added here for convenience.

The homogenization connection

Before proceeding with the proof of the Lions-Papanicolaou-Varadhan theorem, let us explain how the steady equation (2.6.7) appears in the context of periodic homogenization, which was probably the main motivation behind this theorem. We can not possibly do justice to the area of homogenization here – an interested reader should explore the huge literature on the subject, with the book [?] by G. Pavliotis and A. Stuart providing a good starting point. Let us just briefly illustrate the general setting on the example of the periodic Hamilton-Jacobi equations. Consider the Cauchy problem

$$u_t^\varepsilon + H(x, \nabla u^\varepsilon) = 0, \quad t > 0, \quad x \in \mathbb{R}^n, \quad (2.6.8)$$

in the whole space \mathbb{R}^n (and not on the torus), with the Hamiltonian $H(x, p)$ that is 1-periodic in all coordinates x_j , $j = 1, \dots, n$. We are interested in the regime where the initial condition is slowly varying and large:

$$u^\varepsilon(0, x) = \varepsilon^{-1} u_0(\varepsilon x). \quad (2.6.9)$$

Let us note that if one thinks of the solution to (2.6.8) as the height of some interface at the position $x \in \mathbb{R}^n$ at a time $t > 0$, then it is very natural that if $u^\varepsilon(0, x)$ varies on a scale ε^{-1} in the x -variable, then its height should also be of the order ε^{-1} , which is exactly what we see in (2.6.9).

The general question of homogenization is how the "microscopic" variations in the Hamiltonian that varies on the scale $O(1)$ affect the evolution of the initial condition that varies on

the "macroscopic" scale $O(\varepsilon^{-1})$. The goal is to describe the evolution in purely "macroscopic" terms, and extract an effective macroscopic problem that approximates the full microscopic problem well. This allows to avoid, say, in numerical simulations, modeling the microscopic variations of the Hamiltonian, and do the simulations on the macroscopic scale – a huge advantage in engineering problems. It also happens that from the purely mathematical view point, homogenization is also an extremely rich subject.

This general philosophy translates into the following strategy. As the initial condition in (2.6.9) is slowly varying, one should observe the solution on a macroscopic spatial scale, in the "slow" variable $y = \varepsilon x$. Since $u^\varepsilon(0, x)$ is also very large itself, of the size $O(\varepsilon^{-1})$, it is appropriate to rescale it down. In other words, instead of looking at $u^\varepsilon(t, x)$ directly, we would represent it as

$$u^\varepsilon(t, x) = \varepsilon^{-1} w^\varepsilon(t, \varepsilon x),$$

and consider the evolution of $w^\varepsilon(t, y)$, which satisfies

$$w_t^\varepsilon + \varepsilon H\left(\frac{y}{\varepsilon}, \nabla w^\varepsilon\right) = 0, \quad (2.6.10)$$

with the initial condition $w^\varepsilon(0, y) = u_0(y)$ that is now independent of ε . However, we see that w^ε evolves very slowly in t – its time derivative is of the size $O(\varepsilon)$. Hence, we need to wait a long time until it changes. To remedy this, we introduce a long time scale of the size $t = O(\varepsilon^{-1})$. In other words, we write

$$w^\varepsilon(t, y) = v^\varepsilon(\varepsilon t, y).$$

In the new variables the problem takes the form

$$v_s^\varepsilon + H\left(\frac{y}{\varepsilon}, \nabla v^\varepsilon\right) = 0, \quad y \in \mathbb{R}^n, \quad s > 0, \quad (2.6.11)$$

with the initial condition $v^\varepsilon(0, y) = u_0(y)$.

It seems that we have merely shifted the difficulty – we used to have ε in the initial condition in (2.6.9) while now we have it appear in the equation itself – the Hamiltonian depends on y/ε . However, it turns out that we may now find an ε -independent problem that has a spatially uniform Hamiltonian that provides a good approximation to (2.6.11). The reason this is possible is that we have chosen the correct temporal and spatial scales to track the evolution of the solution.

Here is an informal way to find the approximating problem. Let us seek the solution to (2.6.11) in the form of an asymptotic expansion

$$v^\varepsilon(s, y) = \bar{v}(s, y) + \varepsilon v_1(s, y, \frac{y}{\varepsilon}) + \varepsilon^2 v_2(s, y, \frac{y}{\varepsilon}) + \dots \quad (2.6.12)$$

The functions $v_j(s, y, z)$ are assumed to be periodic in the "fast" variable z but not in the "slow" variables s and y . Inserting this expansion into (2.6.11), and collecting the terms with various powers of ε , we obtain in the leading order

$$\bar{v}_s(s, y) + H\left(\frac{y}{\varepsilon}, \nabla_y \bar{v}(s, y) + \nabla_z v_1(s, y, \frac{y}{\varepsilon})\right) = 0. \quad (2.6.13)$$

As is standard in such multiple scale expansions, we consider (2.6.13) as

$$\bar{v}_s(s, y) + H(z, \nabla_y \bar{v}(s, y) + \nabla_z v_1(s, y, z)) = 0, \quad z \in \mathbb{T}^n, \quad (2.6.14)$$

an equation for v_1 as a function of the fast variable $z \in \mathbb{T}^n$, for each $s > 0$ and $y \in \mathbb{R}^n$ fixed. In other words, for each pair of the "macroscopic" variables s and y we consider a microscopic problem in the z -variable. In the area of numerical analysis, one would call this "sub-grid modeling": the variables s and y live on the macroscopic grid, and the z -variable lives on the microscopic sub-grid.

The function $\bar{v}(s, y)$ will then be found from the solvability condition for (2.6.13). Indeed, the terms $\bar{v}_s(s, y)$ and $\nabla_y \bar{v}(s, y)$ in (2.6.14) do not depend on the fast variable z and should be treated as constants – we solve (2.6.14) independently for each s and y . Let us then, for each fixed $p \in \mathbb{R}^n$, consider the problem

$$H(z, p + \nabla_z w) = c, \quad z \in \mathbb{T}^n. \quad (2.6.15)$$

The case of interest is $p = \nabla_y \bar{v}(s, y)$ and $c = -\bar{v}_s(s, y)$ but one needs to momentarily look at (2.6.15) for an arbitrary choice of $p \in \mathbb{R}^n$ and $c \in \mathbb{R}$. The Lions-Papanicolaou-Varadhan theorem says that for each $p \in \mathbb{R}^n$ there is a unique c that we will denote by $\bar{H}(p)$ such that (2.6.15) has a solution. We then write (2.6.15) as

$$H(z, p + \nabla_z w) = \bar{H}(p), \quad z \in \mathbb{T}^n. \quad (2.6.16)$$

Hence, the solvability condition for (2.6.14) is that the function $\bar{v}(s, y)$ satisfies the homogenized (also known as "effective") equation

$$\bar{v}_s + \bar{H}(\nabla_y \bar{v}) = 0, \quad \bar{v}(0, y) = u_0(y), \quad s > 0, \quad y \in \mathbb{R}^n, \quad (2.6.17)$$

and the function $\bar{H}(p)$ is called the effective, or homogenized Hamiltonian. Note that the effective Hamiltonian does not depend on the spatial variable – the "small scale" variations are averaged out via the above homogenization procedure. The point is that the solution $v^\varepsilon(s, y)$ of (2.6.11), an equation with highly oscillatory coefficients is well approximated by $\bar{v}(s, y)$, the solution of (2.6.17), an equation with spatially uniform coefficients, that is much simpler to study analytically or solve numerically.

Thus, the existence and uniqueness of the constant c for which solution of the steady equation (2.6.15) exists, is directly related to the homogenization (long time behavior) of the solutions to the Cauchy problem (2.6.8) with slowly varying initial conditions, as it provides the corresponding effective Hamiltonian. Unfortunately, there is a catch: not so much is known in general on how the effective Hamiltonian $\bar{H}(p)$ depends on the original Hamiltonian $H(x, p)$, except for some very generic properties. Estimating and computing numerically the effective Hamiltonian $\bar{H}(p)$ is a separate interesting line of research.

Exercise 2.6.3 (*The one-dimensional case*) Compute the effective Hamiltonian $\bar{H}(p)$ for

$$H(x, p) = R(x)\sqrt{1 + p^2}, \quad x \in \mathbb{T}^1, p \in \mathbb{R},$$

where $R(x)$ is a smooth 1-periodic function.

Exercise 2.6.4 Show that for every $p \in \mathbb{R}^n$ one can find a periodic in x function $u(x; p)$, $x \in \mathbb{T}^n$, $p \in \mathbb{R}^n$ such that the function

$$v(t, x; p) = p \cdot x + u(x; p) - t\bar{H}(p)$$

is a solution to

$$v_t + H(x; \nabla v) = 0.$$

What is the function $u(x; p)$ in terms of the approximate expansion (2.6.12)? Explain why it is natural that the function $u(x; p)$ appears when we try to approximate the solution to

$$u_t^\varepsilon + H(x, \nabla u^\varepsilon) = 0,$$

with an initial condition of the form $u^\varepsilon(0, x) = \varepsilon^{-1}u_0(\varepsilon x)$.

The proof of the Lions-Papanicolaou-Varadhan theorem

Recall that our goal is to construct a solution to (2.6.7):

$$H(x, \nabla u) = c, \quad x \in \mathbb{T}^n. \tag{2.6.18}$$

As we have already proved uniqueness of the constant c , we only need to prove its existence, and, of course, construct the solution $u(x)$. We will make use of the viscosity solution to the auxiliary problem

$$H(x, \nabla u^\varepsilon) + \varepsilon u^\varepsilon = 0, \quad x \in \mathbb{T}^n, \tag{2.6.19}$$

with $\varepsilon > 0$. Note that the regularization parameter $\varepsilon > 0$ in (2.6.19) has nothing to do with the small parameter $\varepsilon > 0$ that we have used in the discussion of the periodic homogenization theory, where it referred to the separation of scales between the scale of variation of the initial condition and that of the periodic Hamiltonian. Unfortunately, it is common to use the notation ε in both of these settings. We hope that the reader will find it not too confusing.

We have already shown that (2.6.19) has at most one solution. Let us for the moment accept that the solution to the regularized problem (2.6.19) exists and show how one can finish the proof of Theorem 2.6.2 from here. Then, we will come back to the construction of a solution to (2.6.19). Our task is to pass to the limit $\varepsilon \downarrow 0$ in (2.6.19).

Exercise 2.6.5 Show that for all $\varepsilon > 0$, the solution $u^\varepsilon(x)$ of (2.6.19) satisfies

$$-\frac{\|H(\cdot, 0)\|_{L^\infty}}{\varepsilon} \leq u^\varepsilon(x) \leq \frac{\|H(\cdot, 0)\|_{L^\infty}}{\varepsilon}, \tag{2.6.20}$$

for all $x \in \mathbb{T}^n$. Hint: use the comparison principle.

Note that the fact that $u^\varepsilon(x)$ is of the size ε^{-1} is not a fluke of the estimate. For instance, if the function $H(x, p)$ is bounded from below by a positive constant c_0 , then the solution to (2.6.19) will clearly satisfy $|u^\varepsilon(x)| \geq c_0/\varepsilon$ for all $x \in \mathbb{T}^n$. Therefore, one can not expect that the solution to (2.6.19) converges as $\varepsilon \rightarrow 0$ to a solution to (2.6.18). One can, however, hope that the solution becomes large but its gradient stays bounded, so if we subtract the

large mean the difference will be bounded. Accordingly, we will decompose u^ε into its mean and oscillation:

$$u^\varepsilon(x) = \langle u^\varepsilon \rangle + v^\varepsilon(x), \quad (2.6.21)$$

where

$$\langle u^\varepsilon \rangle = \int_{\mathbb{T}^n} u^\varepsilon(y) dy. \quad (2.6.22)$$

Recall that the torus \mathbb{T}^n is normalized so that $\text{Vol}(\mathbb{T}^n) = 1$. We will then show that there is a sequence $\varepsilon_k \rightarrow 0$ so that the limit

$$c = - \lim_{\varepsilon_k \rightarrow 0} \varepsilon_k \langle u^{\varepsilon_k} \rangle \quad (2.6.23)$$

exists, and $v^{\varepsilon_k}(x)$ also converge uniformly on \mathbb{T}^n to a limit u that satisfies (2.6.18) with c given by (2.6.23).

In order to pass to the limit $\varepsilon \downarrow 0$ in (2.6.19), we need a modulus of continuity estimate on u^ε (and hence v^ε) that does not depend on $\varepsilon \in (0, 1)$.

Lemma 2.6.6 *There is $C > 0$ independent of ε such that $|\text{Lip } u^\varepsilon| \leq C$.*

Proof. Again, we use the doubling of the independent variables. Fix $x \in \mathbb{T}^n$ and, for $K > 0$, consider the function

$$\zeta(y) = u^\varepsilon(y) - u^\varepsilon(x) - K|y - x|. \quad (2.6.24)$$

Let \hat{x} be a maximum of $\zeta(y)$ (the point \hat{x} depends on x). If $\hat{x} = x$ for all $x \in \mathbb{T}^n$, then, as $\zeta(x) = 0$, we obtain

$$u^\varepsilon(y) - u^\varepsilon(x) \leq K|x - y|, \quad (2.6.25)$$

for all $x, y \in \mathbb{T}^n$, which implies that u^ε is Lipschitz with the constant K . If there exists some x such that $\hat{x} \neq x$, then the function

$$\psi(y) = u^\varepsilon(x) + K|y - x|$$

is, in a vicinity of the point $y = \hat{x}$, an admissible test function, as a function of y . Moreover, the difference

$$\psi(y) - u^\varepsilon(y) = -\zeta(y)$$

attains its minimum at $y = \hat{x}$. As $u^\varepsilon(y)$ is a viscosity solution to (2.6.19), and

$$\nabla\psi(\hat{x}) = K \frac{\hat{x} - x}{|\hat{x} - x|},$$

it follows that

$$H\left(\hat{x}, K \frac{\hat{x} - x}{|\hat{x} - x|}\right) + \varepsilon u^\varepsilon(\hat{x}) \leq 0. \quad (2.6.26)$$

Since $\varepsilon u^\varepsilon(x)$ is bounded by $\|H(\cdot, 0)\|_{L^\infty}$, as in (2.6.20), we deduce that

$$H\left(\hat{x}, K \frac{\hat{x} - x}{|\hat{x} - x|}\right) \leq \|H(\cdot, 0)\|_{L^\infty}. \quad (2.6.27)$$

On the other hand, the coercivity condition (2.6.5) implies that we can take K sufficiently large, so that

$$\|H(\cdot, 0)\|_{L^\infty} < \inf_{x \in \mathbb{T}^n, |p|=K} H(x, p). \quad (2.6.28)$$

Hence, if we take K as in (2.6.28), then (2.6.27) can not hold. As a consequence, for such K we must have $\hat{x} = x$ for all $x \in \mathbb{T}^n$. It follows that for such K the inequality (2.6.25) holds for all $x, y \in \mathbb{T}^n$. This finishes the proof. \square

To finish the proof of Theorem 2.6.2, we go back to the decomposition (2.6.21)-(2.6.22). The function

$$v^\varepsilon = u^\varepsilon - \langle u^\varepsilon \rangle$$

satisfies

$$H(x, \nabla v^\varepsilon) + \varepsilon \langle u^\varepsilon \rangle + \varepsilon v^\varepsilon = 0. \quad (2.6.29)$$

As

$$\int_{\mathbb{T}^n} v^\varepsilon(x) dx = 0,$$

and because of Lemma 2.6.6, the family v^ε is both uniformly bounded in L^∞ and is uniformly Lipschitz. As a consequence, it converges uniformly, up to extraction of a subsequence, to a function $v \in C(\mathbb{T}^n)$, and $\varepsilon v^\varepsilon \rightarrow 0$. The bound (2.6.20) implies that the family $\varepsilon \langle u^\varepsilon \rangle$ is bounded. We may, therefore, assume its convergence (along a subsequence) to a constant denoted by $-c$, as in (2.6.23). By the stability result in Exercise 2.5.8, we deduce that v is a viscosity solution of

$$H(x, \nabla v) = c. \quad (2.6.30)$$

This finishes the proof of Theorem 2.6.2 except for the construction of a solution to (2.6.19).

Existence of the solution to the auxiliary problem

Let us now construct a solution to (2.6.19).

Proposition 2.6.7 *If $H(x, p)$ satisfies the assumptions of Theorem 2.6.2, then for all $\varepsilon > 0$ the problem*

$$H(x, \nabla u) + \varepsilon u = 0, \quad x \in \mathbb{T}^n, \quad (2.6.31)$$

has a viscosity solution.

We will treat a solution to (2.6.31) as a fixed point of the map $\mathcal{S}[v] = u$ defined via

$$H(x, \nabla u) + Mu = (M - \varepsilon)v, \quad x \in \mathbb{T}^n, \quad (2.6.32)$$

with $M > 0$ to be chosen appropriately. The point is that if M is sufficiently large, we will be able to prove that this map is a contraction on $C(\mathbb{T}^n)$, hence has a fixed point. Any such fixed point is a solution to (2.6.31). Our first task is to prove the following lemma.

Lemma 2.6.8 *There exists $M_0 > 0$ so that for all $M > M_0$ and all $f \in C(\mathbb{T}^n)$ there exists a solution to*

$$H(x, \nabla u) + Mu = f, \quad x \in \mathbb{T}^n. \quad (2.6.33)$$

This lemma shows that the map \mathcal{S} is well-defined for $M > M_0$. Its proof will use an explicit construction of the solutions via a limiting procedure that will give us sufficiently strong a priori bounds that will allow us to deduce that \mathcal{S} is a contraction.

The proof of Lemma 2.6.8

We take a function $f \in C(\mathbb{T}^n)$, and consider a regularized problem

$$-\delta \Delta u^{\gamma, \delta} + H(x, \nabla u^{\gamma, \delta}) + M u^{\gamma, \delta} = f_\gamma(x), \quad x \in \mathbb{T}^n, \quad (2.6.34)$$

with $\delta > 0$ and $\gamma > 0$, and

$$f_\gamma = G_\gamma \star f. \quad (2.6.35)$$

Here, G_γ is a compactly supported smooth approximation of identity:

$$G_\gamma(x) = \gamma^{-n} G\left(\frac{x}{\gamma}\right), \quad G(x) \geq 0, \quad \int_{\mathbb{R}^n} G(x) dx = 1,$$

so that $f_\gamma(x)$ is smooth, and $f_\gamma \rightarrow f$ in $C(\mathbb{T}^n)$. In particular, there exists K_γ that depends on $\gamma \in (0, 1)$ so that

$$\|f_\gamma\|_{L^\infty} \leq \|f\|_{L^\infty}, \quad \|f_\gamma\|_{C^1} \leq K_\gamma \|f\|_{L^\infty}. \quad (2.6.36)$$

It is straightforward to adapt what we have done in Section 1.5.2 for the time-dependent problems with a positive diffusion coefficient to show that (2.6.34) admits a smooth solution $u^{\gamma, \delta}$ for each $\gamma > 0$ and $\delta > 0$. The difficulty is to pass to the limit $\delta \downarrow 0$, followed by $\gamma \downarrow 0$ to construct in the limit a viscosity solution to (2.6.33). This will require a priori bounds on $u^{\gamma, \delta}$ summarized in the following lemma.

Lemma 2.6.9 *There exists $M_0 > 0$ so that if $M > M_0$ then the solution $u^{\gamma, \delta}$ to (2.6.34) obeys the following gradient bound, for all $\delta \in (0, 1)$:*

$$|\nabla u^{\gamma, \delta}(x)| \leq C_\gamma (1 + \|f\|_{L^\infty}) \text{ for all } x \in \mathbb{T}^n. \quad (2.6.37)$$

Here, the constant C_γ may depend on $\gamma \in (0, 1)$ but not on $\delta \in (0, 1)$. There also exists a constant $C > 0$ that does not depend on $\gamma \in (0, 1)$ or $\delta \in (0, 1)$ so that

$$|u^{\gamma, \delta}(x)| \leq \frac{C}{M} (1 + \|f\|_{L^\infty}) \text{ for all } x \in \mathbb{T}^n. \quad (2.6.38)$$

Proof. Let us look at the point x_0 where $|\nabla u^{\gamma, \delta}(x)|^2$ attains its maximum. Note that (we drop the super-scripts γ and δ for the moment)

$$\frac{\partial}{\partial x_i} (|\nabla u|^2) = 2 \frac{\partial u}{\partial x_j} \frac{\partial^2 u}{\partial x_i \partial x_j},$$

so that, using (2.6.34), we compute

$$\begin{aligned} \Delta(|\nabla u|^2) &= 2 \sum_{i,j=1}^n \left(\frac{\partial^2 u}{\partial x_i \partial x_j} \right)^2 + 2 \sum_{j=1}^n \frac{\partial u}{\partial x_j} \frac{\partial \Delta u}{\partial x_j} = 2 \sum_{i,j=1}^n \left(\frac{\partial^2 u}{\partial x_i \partial x_j} \right)^2 + \frac{2M}{\delta} |\nabla u|^2 \\ &+ \frac{2}{\delta} \sum_{j=1}^n \frac{\partial u}{\partial x_j} \frac{\partial H(x, \nabla u)}{\partial x_j} + \frac{2}{\delta} \sum_{k,j=1}^n \frac{\partial u}{\partial x_j} \frac{\partial H(x, \nabla u)}{\partial p_k} \frac{\partial^2 u}{\partial x_j \partial x_k} - \frac{2}{\delta} \sum_{j=1}^n \frac{\partial u}{\partial x_j} \frac{\partial f_\gamma}{\partial x_j} \\ &= 2 \sum_{i,j=1}^n \left(\frac{\partial^2 u}{\partial x_i \partial x_j} \right)^2 + \frac{2M}{\delta} |\nabla u|^2 + \frac{2}{\delta} \sum_{j=1}^n \frac{\partial u}{\partial x_j} \frac{\partial H(x, \nabla u)}{\partial x_j} + \frac{1}{\delta} \sum_{k=1}^n \frac{\partial H(x, \nabla u)}{\partial p_k} \frac{\partial |\nabla u|^2}{\partial x_k} \\ &- \frac{2}{\delta} \sum_{j=1}^n \frac{\partial u}{\partial x_j} \frac{\partial f_\gamma}{\partial x_j}. \end{aligned}$$

Thus, at the maximum x_0 of $|\nabla u|^2$ we have

$$0 \geq \Delta(|\nabla u|^2)(x_0) = 2 \sum_{i,j=1}^n \left(\frac{\partial^2 u}{\partial x_i \partial x_j} \right)^2 + \frac{2M}{\delta} |\nabla u|^2 + \frac{2}{\delta} \sum_{j=1}^n \frac{\partial u}{\partial x_j} \frac{\partial H(x, \nabla u)}{\partial x_j} - \frac{2}{\delta} \sum_{j=1}^n \frac{\partial u}{\partial x_j} \frac{\partial f_\gamma}{\partial x_j}. \quad (2.6.39)$$

Let us recall the gradient bound (2.6.6) on $H(x, p)$:

$$|\nabla_x H(x, p)| \leq K_0(1 + |p|). \quad (2.6.40)$$

We see from (2.6.39) and (2.6.40) that

$$Q = |\nabla u(x_0)| = \sup_{x \in \mathbb{T}^n} |\nabla u(x)|$$

satisfies

$$MQ^2 \leq K_0 Q(1 + Q) + Q \|f_\gamma\|_{C^1} \leq 5K_0(1 + Q^2) + K_\gamma Q \|f\|_{L^\infty}. \quad (2.6.41)$$

We used (2.6.36) above. It follows from (2.6.41) that there exist $M_0 > 0$ and C_1 that depend on K_0 but not on $\gamma \in (0, 1)$ and C_γ that depends on $\gamma \in (0, 1)$ so that for all $M > M_0$ we have

$$Q \leq C_1 + C_\gamma \|f\|_{L^\infty}. \quad (2.6.42)$$

This proves (2.6.37).

To prove (2.6.38) we look at the point x_M where u attains its maximum over \mathbb{T}^n . At this point we have

$$Mu(x_M) = f_\gamma(x_M) + \delta \Delta u(x_M) - H(x_M, 0) \leq \|f_\gamma\|_{L^\infty} + \|H(\cdot, 0)\|_{L^\infty}, \quad (2.6.43)$$

hence

$$u(x_M) \leq \frac{C}{M}(1 + \|f\|_{L^\infty}).$$

A similar estimate holds at the minimum of u , proving (2.6.38). \square

The Lipschitz bound (2.6.37) and (2.6.38) show that if $M > M_0$, after passing to a subsequence $\delta_k \downarrow 0$, the family $u^{\gamma, \delta_k}(x)$ converges uniformly in $x \in \mathbb{T}^n$, to a function $u^\gamma(x)$.

Exercise 2.6.10 Show that $u^\gamma(x)$ is the viscosity solution to

$$H(x, \nabla u^\gamma) + Mu^\gamma = f_\gamma(x), \quad x \in \mathbb{T}^n. \quad (2.6.44)$$

Hint: Exercise 2.5.8 and its solution should be helpful here.

The next step is to send $\gamma \rightarrow 0$.

Exercise 2.6.11 Mimic the proof of Lemma 2.6.6 to show that $u^\gamma(x)$ are uniformly Lipschitz: there exists a constant $C_f > 0$ that may depend on $\|f\|_{L^\infty}$ but is independent of $\gamma \in (0, 1)$ and of $M > M_0$ such that

$$|\text{Lip } u^\gamma| \leq C_f. \quad (2.6.45)$$

Also show that

$$\|u^\gamma\|_{L^\infty} \leq \frac{1}{M} (\|H(\cdot, 0)\|_{L^\infty} + \|f\|_{L^\infty}). \quad (2.6.46)$$

This exercise shows that as long as $M \geq M_0$, the family u^{γ_k} converges, along as subsequence $\gamma_k \downarrow 0$, uniformly in $x \in \mathbb{T}^n$, to a limit $u(x) \in C(\mathbb{T}^n)$ that obeys the same uniform Lipschitz and L^∞ -bounds in Exercise 2.6.11. Invoking again the stability result of Exercise 2.5.8 shows that $u(x)$ is the unique viscosity solution to

$$H(x, \nabla u) + Mu = f(x), \quad x \in \mathbb{T}^n. \quad (2.6.47)$$

This finishes the proof of Lemma 2.6.8. \square

The end of the proof of Proposition 2.6.7

We now explain how this construction implies the conclusion of Proposition 2.6.7. Let us take $\varepsilon < M$, and re-write equation (2.6.31)

$$H(x, \nabla u) + \varepsilon u = 0, \quad x \in \mathbb{T}^n. \quad (2.6.48)$$

for which we need to find a solution, as

$$H(x, \nabla u) + Mu = (M - \varepsilon)u, \quad x \in \mathbb{T}^n. \quad (2.6.49)$$

As we have mentioned, we define the map $\mathcal{S} : C(\mathbb{T}^n) \rightarrow C(\mathbb{T}^n)$ as follows: given $v \in C(\mathbb{T}^n)$, let $u = \mathcal{S}[v]$ be the unique viscosity solution to

$$H(x, \nabla u) + Mu = (M - \varepsilon)v, \quad x \in \mathbb{T}^n. \quad (2.6.50)$$

We claim that \mathcal{S} is a contraction in $C(\mathbb{T}^n)$. We have shown that $u = \mathcal{S}[v]$ can be constructed via the above procedure of passing to the limit $\delta \rightarrow 0$, followed by $\gamma \rightarrow 0$ in the regularized problem

$$-\delta \Delta u^{\gamma, \delta} + H(x, \nabla u^{\gamma, \delta}) + Mu^{\gamma, \delta} = (M - \varepsilon)v_\gamma, \quad x \in \mathbb{T}^n. \quad (2.6.51)$$

Given $v_1, v_2 \in C(\mathbb{T}^n)$, consider the corresponding solutions to the regularized problems (2.6.51):

$$-\delta \Delta u_1^{\gamma, \delta} + H(x, \nabla u_1^{\gamma, \delta}) + Mu_1^{\gamma, \delta} = (M - \varepsilon)v_{1, \gamma}, \quad x \in \mathbb{T}^n, \quad (2.6.52)$$

and

$$-\delta \Delta u_2^{\gamma, \delta} + H(x, \nabla u_2^{\gamma, \delta}) + Mu_2^{\gamma, \delta} = (M - \varepsilon)v_{2, \gamma}, \quad x \in \mathbb{T}^n. \quad (2.6.53)$$

Assume that the difference

$$w = u_1^{\gamma, \delta} - u_2^{\gamma, \delta}$$

attains its maximum at a point x_0 . The function w satisfies

$$-\delta \Delta w + H(x, \nabla u_1^{\gamma, \delta}) - H(x, \nabla u_2^{\gamma, \delta}) + Mw = (M - \varepsilon)(v_{1, \gamma} - v_{2, \gamma}), \quad x \in \mathbb{T}^n. \quad (2.6.54)$$

Evaluating this at $x = x_0$, as $\nabla u_1^{\gamma, \delta}(x_0) = \nabla u_2^{\gamma, \delta}(x_0)$, we see that

$$-\delta \Delta w(x_0) + Mw(x_0) = (M - \varepsilon)(v_{1, \gamma}(x_0) - v_{2, \gamma}(x_0)), \quad x \in \mathbb{T}^n. \quad (2.6.55)$$

As x_0 is the maximum of w , we deduce that

$$w(x_0) \leq \frac{M - \varepsilon}{M} \|v_{1, \gamma} - v_{2, \gamma}\|_{C(\mathbb{T}^n)}.$$

Using a nearly identical computation for the minimum, we conclude that

$$\|u_1^{\gamma,\delta} - u_2^{\gamma,\delta}\|_{C(\mathbb{T}^n)} \leq \frac{M - \varepsilon}{M} \|v_{1,\gamma} - v_{2,\gamma}\|_{C(\mathbb{T}^n)}. \quad (2.6.56)$$

Passing to the limit $\delta \downarrow 0$ and $\gamma \downarrow 0$, we obtain

$$\|u_1 - u_2\|_{C(\mathbb{T}^n)} \leq \frac{M - \varepsilon}{M} \|v_1 - v_2\|_{C(\mathbb{T}^n)}, \quad (2.6.57)$$

hence \mathcal{S} is a contraction on $C(\mathbb{T}^n)$, as claimed. Thus, this map has a fixed point, which is the viscosity solution to

$$H(x, \nabla u) + \varepsilon u = 0, \quad x \in \mathbb{T}^n. \quad (2.6.58)$$

This completes the proof of Proposition 2.6.7. \square

2.6.2 Existence of the solution to the Cauchy problem

We will now construct the viscosity solution to the Cauchy problem

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \quad t > 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x), \quad x \in \mathbb{T}^n, \end{aligned} \quad (2.6.59)$$

with a continuous initial condition $u_0(x)$. Recall that Exercise 2.5.13 implies the uniqueness of the solution with a given initial condition, so we do not need to address that issue. We make the same assumptions as in Theorem 2.6.2: there exists $C_L > 0$ so that

$$|H(x, p_1) - H(x, p_2)| \leq C_L |p_1 - p_2|, \quad \text{for all } x, p_1, p_2 \in \mathbb{R}^n, \quad (2.6.60)$$

and

$$\lim_{|p| \rightarrow +\infty} H(x, p) = +\infty, \quad \text{uniformly in } x \in \mathbb{T}^n. \quad (2.6.61)$$

We will again assume the gradient bound (2.6.6):

$$|\nabla_x H(x, p)| \leq K_0(1 + |p|), \quad \text{for all } x \in \mathbb{T}^n, \text{ and } p \in \mathbb{R}^n. \quad (2.6.62)$$

Theorem 2.6.12 *The Cauchy problem (2.6.59) has a unique viscosity solution $u(t, x)$. Moreover, the weak contraction property holds: if $u(t, x)$ and $v(t, x)$ are two solutions to (2.6.59) with the corresponding initial conditions $u_0 \in C(\mathbb{T}^n)$ and $v_0 \in C(\mathbb{T}^n)$, then*

$$\|u(t, \cdot) - v(t, \cdot)\|_{L^\infty} \leq \|u_0 - v_0\|_{L^\infty}. \quad (2.6.63)$$

The weak contraction property is recorded here simply for the sake of completeness: we have seen in Exercise 2.5.14 that it follows immediately from the comparison principle. Therefore, we will focus on the existence of the solutions.

An important consequence of the weak contraction principle is that we may restrict ourselves to initial conditions that are smooth. Indeed, suppose that we managed to prove the theorem for smooth initial conditions, and consider $u_0 \in C(\mathbb{T}^n)$. Let $u_0^{(k)}$ be a sequence of

smooth functions converging to u_0 in $C(\mathbb{T}^n)$ as $k \rightarrow +\infty$, and $u^{(k)}(t, x)$ be the corresponding sequence of solutions to (2.6.59), with the initial conditions $u_0^{(k)}$. It follows from the weak contraction principle that

$$\|u^{(k)} - u^{(m)}\|_{L^\infty(\mathbb{R}_+ \times \mathbb{T}^n)} \leq \|u_0^{(k)} - u_0^{(m)}\|_{L^\infty},$$

ensuring that $u^{(k)}$ is a uniformly Cauchy sequence on $C([0, +\infty) \times \mathbb{T}^n)$. Hence, it converges uniformly to a continuous function $u \in C(\mathbb{R}_+ \times \mathbb{T}^n)$. The stability result in Exercise 2.5.8 implies that u is a solution to the Cauchy problem (2.6.59) with the initial condition $u_0(x)$.

We are now left with the actual construction of a solution to (2.6.59), with the assumption that u_0 is smooth. We are going to use the most pedestrian way to do it: a time discretization. Take a family of time steps $\Delta t \rightarrow 0$. For a fixed $\Delta t > 0$, consider the sequence $u_{\Delta t}^n(x)$ defined by setting $u^0(x) := u_0(x)$ and the recursion relation:

$$\frac{u_{\Delta t}^{n+1} - u_{\Delta t}^n}{\Delta t} + H(x, \nabla u_{\Delta t}^{n+1}) = 0, \quad x \in \mathbb{T}^n, \quad (2.6.64)$$

that is an implicit time discretization of (2.6.59). Given $u_{\Delta t}^n(x)$, we look at (2.6.64) as a time-independent Hamilton-Jacobi equation

$$H(x, \nabla u_{\Delta t}^{n+1}) + \frac{1}{\Delta t} u_{\Delta t}^{n+1} = \frac{1}{\Delta t} u_{\Delta t}^n, \quad x \in \mathbb{T}^n. \quad (2.6.65)$$

It is of the type, for which Proposition 2.6.7 guarantees existence of a unique continuous solution $u_{\Delta t}^{n+1}$, as long as $u_{\Delta t}^n$ is continuous. This produces the sequence $u_{\Delta t}^n(x)$, for $n \geq 0$. An approximate solution $u_{\Delta t}$ to the Cauchy problem (2.6.59) is then constructed by interpolating linearly between the times $n\Delta t$ and $(n+1)\Delta t$:

$$u_{\Delta t}(t, x) = u_{\Delta t}^n(x) + \frac{t - n\Delta t}{\Delta t} (u_{\Delta t}^{n+1}(x) - u_{\Delta t}^n(x)), \quad t \in [n\Delta t, (n+1)\Delta t). \quad (2.6.66)$$

The help provided by the smoothness assumption on u_0 manifests itself in the next proposition.

Proposition 2.6.13 *There is $C > 0$, depending on $\|u_0\|_\infty$ and $Lip(u_0)$ but not on $\Delta t \in (0, 1)$, such that the function $u_{\Delta t}(t, x)$ is uniformly Lipschitz continuous in t and x on $[0, +\infty) \times \mathbb{T}^n$, and the Lipschitz constant $Lip(u_{\Delta t})$ of $u_{\Delta t}$ both in t and x , over the set $[0, +\infty) \times \mathbb{T}^n$, satisfies*

$$Lip(u_{\Delta t}) \leq C. \quad (2.6.67)$$

This ensures that there exists a sequence $\Delta t_n \rightarrow 0$, such that the corresponding sequence $u_{\Delta t_n}$ converges as $n \rightarrow \infty$ to a Lipschitz function $u(t, x)$ with the Lipschitz constant $Lip(u) \leq C$. The next step will be to prove

Proposition 2.6.14 *The function $u(t, x)$ is a viscosity solution to the Cauchy problem (2.6.59):*

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \quad t > 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x), \quad x \in \mathbb{T}^n. \end{aligned} \quad (2.6.68)$$

Proof. Let us prove this claim first, assuming the conclusion of Proposition 2.6.13. Note that the initial condition $u(0, x) = u_0(x)$ is satisfied by construction, so we only need to check that u is a viscosity solution to (2.6.68). We will only prove that u is a super-solution, the sub-solution property of u can be proved identically. Let $\varphi(t, x)$ be a C^1 -test function and (t_0, x_0) be a minimum point for the difference $u - \varphi$. As we have seen in the hint to Exercise 2.5.8, we may assume, possibly after subtracting a quadratic polynomial in t and x from the function φ , that the minimum is strict. Consider the linearly interpolated time discretization $\varphi_{\Delta t}$ of φ : set $\varphi^n(x) = \varphi(n\Delta t, x)$, for $n \geq 0$, and

$$\varphi_{\Delta t}(t, x) = \varphi^n(x) + \frac{t - n\Delta t}{\Delta t}(\varphi^{n+1}(x) - \varphi^n(x)), \quad \text{for } t \in [n\Delta t, (n+1)\Delta t).$$

Note a slight abuse of notation: the function $\varphi_{\Delta t}$ is a linear interpolation of the function φ , while $u_{\Delta t}$ is not the linear interpolation of the function u but rather the linear interpolation of the solution to the time-discretized problem (2.6.64), with the time step Δt . Nevertheless, as the minimum (t_0, x_0) of $u - \varphi$ is strict, and $u_{\Delta t}$ converges to u uniformly, for Δt sufficiently small, there exists a minimum point $(t_{\Delta t}, x_{\Delta t})$ for $u_{\Delta t} - \varphi_{\Delta t}$, such that

$$\lim_{\Delta t \rightarrow 0} (t_{\Delta t}, x_{\Delta t}) = (t_0, x_0).$$

In addition, because both $u_{\Delta t}$ and $\varphi_{\Delta t}$ are piecewise linear in t , we have $t_{\Delta t} = (n+1)\Delta t$ for some $n \geq 0$. Then we have, again, because $(t_{\Delta t}, x_{\Delta t})$ is a minimum for $u_{\Delta t} - \varphi_{\Delta t}$:

$$\frac{u_{\Delta t}^{n+1}(x_{\Delta t}) - u_{\Delta t}^n(x_{\Delta t})}{\Delta t} = \partial_t^- u_{\Delta t}((n+1)\Delta t, x_{\Delta t}) \leq \partial_t^- \varphi_{\Delta t}((n+1)\Delta t, x_{\Delta t}) = \partial_t \varphi(t_0, x_0) + o(1), \quad (2.6.69)$$

as $\Delta t \rightarrow 0$. We also have, in the vicinity of (t_0, x_0) :

$$\varphi(t, x) - \varphi_{\Delta t}(t, x) = O(\Delta t^2), \quad \partial_t \varphi(t, x) - \partial_t \varphi_{\Delta t}(t, x) = O(\Delta t), \quad \text{as } \Delta t \rightarrow 0, \quad (2.6.70)$$

with the slight catch here that we have to speak of the left and right derivatives of $\varphi_{\Delta t}$ at the discrete times $n\Delta t$. On the other hand, the point $x_{\Delta t}$ is a minimum of

$$u_{\Delta t}^{n+1}(x) - \varphi_{\Delta t}((n+1)\Delta t, x)$$

in the x -variable. Since $u_{\Delta t}^{n+1}$ is a viscosity solution to (2.6.64), we have

$$\frac{u_{\Delta t}^{n+1}(x_{\Delta t}) - u_{\Delta t}^n(x_{\Delta t})}{\Delta t} \geq -H(x_{\Delta t}, \nabla \varphi_{\Delta t}((n+1)\Delta t, x_{\Delta t})) = -H(x_0, \nabla \varphi(t_0, x_0)) + o(1), \quad (2.6.71)$$

as $\Delta t \rightarrow 0$. Putting together (2.6.69)-(2.6.71) and sending Δt to 0, we obtain

$$\partial_t \varphi(t_0, x_0) + H(x_0, \nabla \varphi(t_0, x_0)) \geq 0,$$

hence u is a super-solution to (2.6.68). This proves Proposition 2.6.14. \square

Proof of Proposition 2.6.13

The reason behind this proposition is quite simple: if u is a smooth solution to

$$u_t + H(x, \nabla u) = 0, \quad (2.6.72)$$

then the function $v(t, x) = u_t(t, x)$ solves

$$v_t + \nabla_p H(x, \nabla u) \cdot \nabla v = 0, \quad (2.6.73)$$

with the initial condition $v(0, x) = -H(x, \nabla u_0(x))$. It follows from the maximum principle, or the method of characteristics for smooth solutions, that

$$\|v(t, \cdot)\|_{L^\infty} \leq \|H(\cdot, \nabla u_0(\cdot))\|_{L^\infty}. \quad (2.6.74)$$

Moreover, (2.6.72) and (2.6.74) together with the coercivity of $H(x, p)$ yield the uniform boundedness of ∇u . The proof of the proposition consists in making this idea rigorous.

Let us recall that $u_{\Delta t}^n$ is the solution to the recursive equation (2.6.64)

$$\frac{u_{\Delta t}^{n+1} - u_{\Delta t}^n}{\Delta t} + H(x, \nabla u_{\Delta t}^{n+1}) = 0, \quad x \in \mathbb{T}^n, \quad (2.6.75)$$

interpolated between the times of the form $n\Delta t$ as in (2.6.66):

$$u_{\Delta t}(t, x) = u_{\Delta t}^n(x) + \frac{t - n\Delta t}{\Delta t} (u_{\Delta t}^{n+1}(x) - u_{\Delta t}^n(x)), \quad t \in [n\Delta t, (n+1)\Delta t). \quad (2.6.76)$$

The viscosity solution $u_{\Delta t}^{n+1}$ to (2.6.75) can be constructed using the by now familiar idea of a diffusive regularization:

$$-\delta \Delta u_{\Delta t}^{n+1, \delta} + H(x, \nabla u_{\Delta t}^{n+1, \delta}) + \frac{u_{\Delta t}^{n+1, \delta} - u_{\Delta t}^{n, \delta}}{\Delta t} = 0, \quad x \in \mathbb{T}^n, \quad (2.6.77)$$

with $\delta > 0$, and then sending $\delta \downarrow 0$. As we have assumed that $u_0(x)$ is smooth, all $u_{\Delta t}^{n, \delta}(x)$ are also smooth, for all $\delta > 0$.

Exercise 2.6.15 Show that

$$\|u_{\Delta t}^{n+1, \delta}\|_{L^\infty} \leq \|u_{\Delta t}^{n, \delta}\|_{L^\infty} + (\Delta t) \|H(\cdot, 0)\|_{L^\infty}. \quad (2.6.78)$$

Hint: look at the maximum x_0 of the smooth function $u_{\Delta t}^{n+1, \delta}$ over \mathbb{T}^n .

Exercise 2.6.16 Use the argument in the proof of Lemma 2.6.9 and Exercise 2.6.15 to show that there exists a constant $C_{n, \Delta t}$ that may depend on n and Δt but not on $\delta > 0$, so that

$$\|\nabla u_{\Delta t}^{n, \delta}\|_{L^\infty} \leq C_{n, \Delta t}. \quad (2.6.79)$$

The bound (2.6.79) is quite poor as we did not track the dependence of $C_{n,\Delta t}$ on n or Δt , but we have extra help. The differential quotient

$$v_{\Delta t}^{n,\delta} = \frac{u_{\Delta t}^{n+1,\delta} - u_{\Delta t}^{n,\delta}}{\Delta t}$$

satisfies

$$-\delta \Delta v_{\Delta t}^{n+1,\delta} + \frac{v_{\Delta t}^{n+1,\delta}}{\Delta t} + \frac{1}{\Delta t} \left(H(x, \nabla u_{\Delta t}^{n+1,\delta}) - H(x, \nabla u_{\Delta t}^{n,\delta}) \right) = \frac{v_{\Delta t}^{n,\delta}}{\Delta t}, \quad (2.6.80)$$

for all $n \geq 0$. At the maximum x_M and minimum x_m of the smooth function $v_{\Delta t}^{n,\delta}$ we have

$$\nabla u_{\Delta t}^{n+1,\delta}(x_M) = \nabla u_{\Delta t}^{n,\delta}(x_M), \quad \nabla u_{\Delta t}^{n+1,\delta}(x_m) = \nabla u_{\Delta t}^{n,\delta}(x_m).$$

Using this in (2.6.80) we obtain

$$\|v_{\Delta t}^{n+1,\delta}\|_{L^\infty} \leq \|v_{\Delta t}^{n,\delta}\|_{L^\infty} \leq \cdots \leq \|v_{\Delta t}^{0,\delta}\|_{L^\infty}. \quad (2.6.81)$$

For the last term in the right side we observe that

$$v_{\Delta t}^{0,\delta} = \frac{u_{\Delta t}^{1,\delta} - u_0}{\Delta t}$$

satisfies, instead of (2.6.80), the equation

$$-\delta \Delta v_{\Delta t}^{0,\delta} + \frac{v_{\Delta t}^{0,\delta}}{\Delta t} + \frac{1}{\Delta t} H(x, \nabla u_{\Delta t}^{1,\delta}) = \frac{\delta}{\Delta t} \Delta u_0. \quad (2.6.82)$$

Again, the maximum principle implies

$$\|v_{\Delta t}^{0,\delta}\|_{L^\infty} \leq \|H(\cdot, \nabla u_0)\|_{L^\infty} + \delta \|\Delta u_0\|_{L^\infty}. \quad (2.6.83)$$

Using this in (2.6.81), we conclude that

$$\|v_{\Delta t}^{n,\delta}\|_{L^\infty} \leq \|H(\cdot, \nabla u_0)\|_{L^\infty} + \delta \|\Delta u_0\|_{L^\infty}, \quad (2.6.84)$$

for all $n \geq 0$. This bound is the reason why we have assumed that u_0 is smooth.

We may now pass to the limit $\delta \rightarrow 0$ in (2.6.84) and recall the convergence of $u_{\Delta t}^{n,\delta}$ to $u_{\Delta t}^n$, to conclude that

$$v_{\Delta t}^{n,\delta} = \frac{u_{\Delta t}^{n+1,\delta} - u_{\Delta t}^{n,\delta}}{\Delta t} \rightarrow v_{\Delta t}^n := \frac{u_{\Delta t}^{n+1} - u_{\Delta t}^n}{\Delta t} \quad \text{as } \delta \downarrow 0. \quad (2.6.85)$$

Combining this with the uniform bound (2.6.84), we conclude that

$$\left\| \frac{u_{\Delta t}^{n+1} - u_{\Delta t}^n}{\Delta t} \right\|_{L^\infty} \leq \|H(\cdot, \nabla u_0)\|_{L^\infty}, \quad (2.6.86)$$

which is a uniform Lipschitz bound on $u_{\Delta t}$ in the t -variable that we need. The reader should compare it to the bound (2.6.74) that we have obtained easily for smooth solutions.

The Lipschitz bound for $u_{\Delta t}^n$ in the x -variable follows easily. Recall that the functions $u_{\Delta t}^n$ satisfy (2.6.64):

$$H(x, \nabla u_{\Delta t}^{n+1}) + \frac{1}{\Delta t} u_{\Delta t}^{n+1} = \frac{1}{\Delta t} u_{\Delta t}^n, \quad x \in \mathbb{T}^n. \quad (2.6.87)$$

We know from Exercise 2.6.16 that $u_{\Delta t}^n$ are Lipschitz – even though we do not know if they have a Lipschitz constant that does not depend on n or Δt . However, this already tells us that $u_{\Delta t}^n$ satisfy (2.6.87) almost everywhere. We write this equation in the form

$$H(x, \nabla u_{\Delta t}^{n+1}) = -v_{\Delta t}^n(x), \quad x \in \mathbb{T}^n. \quad (2.6.88)$$

The uniform bound on $v_{\Delta t}^n$ in (2.6.86) together with the coercivity of $H(x, p)$ imply that there exists a constant $K > 0$ that does not depend on n or Δt so that

$$\|\nabla u_{\Delta t}^{n+1}\|_{L^\infty} \leq K. \quad (2.6.89)$$

This finishes the proof of Proposition 2.6.13. \square

Exercise 2.6.17 Prove the following elementary fact that we used in the very last step in the above proof: if $u(x)$ is a Lipschitz function then $\text{Lip}(u) = \|\nabla u\|_{L^\infty}$.

Exercise 2.6.18 (*Hamiltonians that are coercive in u*). So far, we have been remarkably silent about Hamilton-Jacobi equations of the form

$$u_t + H(x, u, \nabla u) = 0, \quad t > 0, x \in \mathbb{T}^n, \quad (2.6.90)$$

with the Hamiltonian that depends also on the function u itself. There is one case when the above theory can be developed without any real input of new ideas: assume that $H(x, u, p)$ is non-decreasing in u , and that there exists $C_0 > 0$ so that for all $R > 0$, there exists $\delta_{1,2}(R)$ such that

$$0 < \delta_1(R) \leq \delta_2(R) < C_0,$$

and, for all $u \in [-R, R]$, we have

$$\delta_1(R)(|p| - 1) \leq H(x, u, p) \leq \delta_2(R)(|p| + 1) \text{ for all } |u| \leq R, x \in \mathbb{T}^n \text{ and } p \in \mathbb{R}^n.$$

Prove a well-posedness theorem analogous to Theorem 2.6.12. How far can one stretch the assumptions on $H(x, u, p)$? Hint: coercivity is really something one has to assume, one way or another.

2.7 When the Hamiltonian is strictly convex: the Lagrangian theory

Let us recall that in Section 2.4 we considered the Cauchy problem

$$u_t + \frac{1}{2} |\nabla u|^2 - R(x) = 0, \quad (2.7.1)$$

with an initial condition $u(0, x) = u_0(x)$. We have shown that when both $R(x)$ and $u_0(x)$ are convex, this problem has a smooth solution given by the (at first sight) strange looking expression (2.4.19)

$$u(t, x) = \inf_{\gamma(t)=x} \left(u_0(\gamma(0)) + \int_0^t \left(\frac{|\gamma'(s)|^2}{2} + R(\gamma(s)) \right) ds \right). \quad (2.7.2)$$

Moreover, this expression is well-defined even if the boundary value problem for the characteristic curves may be not well-posed. Hence, a natural idea is to generalize this formula to other Hamiltonians and take this generalization as the definition of a solution. On the other hand, we already have the notion of a viscosity solution, so an issue is if these objects agree. In this section, we investigate when the variational approach is possible and whether the solution you construct in this way is, indeed, a viscosity solution. We also discuss how the strict convexity of the Hamiltonian gives an improved regularity of the solution.

2.7.1 The Lax-Oleinik formula and viscosity solutions

In the construction of the viscosity solutions, we assumed very little about the Hamiltonian H : all we really needed was coercivity and continuity. The other regularity assumptions we have made are mostly of the technical nature and can be avoided. From now on, we will adopt an even stronger technical assumption that $H(x, p)$ is $C^\infty(\mathbb{T}^n \times \mathbb{R}^n)$ smooth but more crucially we will assume that $H(x, p)$ is uniformly strictly convex in its second variable: there exists $\alpha > 0$ so that

$$D_p^2 H(x, p) \geq \alpha I, \quad [D_p^2 H(x, p)]_{ij} = \frac{\partial^2 H(x, p)}{\partial p_i \partial p_j}, \quad (2.7.3)$$

in the sense of quadratic forms, for all $x \in \mathbb{T}^n$ and $p \in \mathbb{R}^n$. Unlike the regularity assumptions, the convexity of $H(x, p)$ in p is essential not only for this section, but also for many results on the Hamilton-Jacobi equations.

Exercise 2.7.1 The reader may be naturally concerned that in the construction of the viscosity solutions we have assumed that $H(x, p)$ is uniformly Lipschitz:

$$|H(x, p_1) - H(x, p_2)| \leq C_L |p_1 - p_2| \quad \text{for all } x \in \mathbb{T}^n \text{ and } p_1, p_2 \in \mathbb{R}^n, \quad (2.7.4)$$

and differentiable in x :

$$|\nabla_x H(x, p)| \leq C_0(1 + |p|) \quad \text{for all } x \in \mathbb{T}^n \text{ and } p \in \mathbb{R}^n, \quad (2.7.5)$$

These assumptions are, of course, incompatible with the strict convexity assumption on $H(x, p)$ in (2.7.3). Go through the proofs of existence and uniqueness of the viscosity solutions and show that the coercivity assumption

$$\lim_{|p| \rightarrow +\infty} H(x, p) = +\infty \quad (2.7.6)$$

together with the assumption that (2.7.4) and (2.7.5) hold locally in p , in the sense that for ever compact set $\mathcal{K} \subset \mathbb{R}^n$ there exist two constants $C_L(\mathcal{K})$ and $C_0(\mathcal{K})$ such that

$$\begin{aligned} |H(x, p_1) - H(x, p_2)| &\leq C_L |p_1 - p_2| \quad \text{for all } x \in \mathbb{T}^n \text{ and } p_1, p_2 \in \mathcal{K}, \\ |\nabla_x H(x, p)| &\leq C_0(1 + |p|) \quad \text{for all } x \in \mathbb{T}^n \text{ and } p \in \mathcal{K}, \end{aligned} \quad (2.7.7)$$

are sufficient to prove existence and uniqueness of the viscosity solutions both in the Lions-Papanicolaou-Varadhan Theorem 2.6.2 and in Theorem 2.6.12 for the solutions to the Cauchy problem.

The Legendre transform and extremal paths

Recall that in Section 2.2 we have informally argued as follows: given a path $\gamma(s)$, $t \leq s \leq T$, with the starting point $\gamma(t) = x$, we can define its cost as

$$\mathcal{C}(\gamma)(t) = \int_t^T \tilde{L}(\dot{\gamma}(s)) ds + f(x(T)). \quad (2.7.8)$$

Here, the function $\tilde{L}(v)$ represents the running cost, and the function $f(x)$ is the terminal cost. The corresponding value function is

$$\tilde{u}(t, x) = \inf_{\gamma: \gamma(t)=x} \mathcal{C}(\gamma)(t), \quad (2.7.9)$$

with the infimum taken over all curves $\gamma \in C^1$ such that $\gamma(t) = x$. We have shown, albeit very informally, that $\tilde{u}(t, x)$ satisfies the Hamilton-Jacobi equation

$$\tilde{u}_t + \tilde{H}(\nabla \tilde{u}) = 0, \quad (2.7.10)$$

with the terminal condition $u(T, x) = f(x)$. The Hamiltonian $\tilde{H}(p)$ is given in terms of the running cost $\tilde{L}(v)$ by (2.2.9):

$$\tilde{H}(p) = \inf_{v \in \mathcal{A}} [\tilde{L}(v) + v \cdot p]. \quad (2.7.11)$$

It is convenient to reverse the direction of time and set

$$u(t, x) = \tilde{u}(T - t, x). \quad (2.7.12)$$

This function satisfies the forward Cauchy problem

$$u_t + H(\nabla u) = 0, \quad (2.7.13)$$

with the initial condition $u(0, x) = f(x)$ and the Hamiltonian given by

$$H(p) = -\tilde{H}(p) = - \inf_{v \in \mathbb{R}^n} [\tilde{L}(v) + v \cdot p] = \sup_{v \in \mathbb{R}^n} [-p \cdot v - \tilde{L}(v)] = \sup_{v \in \mathbb{R}^n} [p \cdot v - L(v)], \quad (2.7.14)$$

with the time-reversed cost function

$$L(v) = \tilde{L}(-v). \quad (2.7.15)$$

The natural questions are, first, if the above construction, using the minimizer in (2.7.9), indeed, produces a solution to the initial value problem for (2.7.13) – so far, our arguments were rather informal, and, second, how it is related to the notion of the viscosity solution.

This brings us to the terminology of the Legendre transforms. One of the standard references for the basic properties of the Legendre transform is [?], where an interested reader may find much more information on this beautiful subject. Given a function $L(v)$, known as the Lagrangian, we define its Legendre transform as in (2.7.14)

$$H(p) = \sup_{v \in \mathbb{R}^n} (p \cdot v - L(v)). \quad (2.7.16)$$

Exercise 2.7.2 Show the function $H(p)$ defined by (2.7.16) is convex. Hint: use the fact that $H(p)$ is the supremum of a family of linear functions in p .

This shows that if we hope to connect the Hamilton-Jacobi equations to the above optimal control problem, this can only be done for convex Hamiltonians. Hence, our assumption (2.7.3) that the Hamiltonian $H(x, p)$ is convex in p .

If the function $L(v)$ is smooth and strictly convex, then, for a given $p \in \mathbb{R}^n$, the maximizer $\bar{v}(p)$ in (2.7.16) is explicit: it is the unique solution to

$$p = \nabla L(\bar{v}). \quad (2.7.17)$$

Exercise 2.7.3 Show that if $L(v)$ is strictly convex, and $H(p)$ is its Legendre transform given by (2.7.16), then we have the duality

$$L(v) = \sup_{p \in \mathbb{R}^n} (p \cdot v - H(p)),$$

so that the Lagrangian L is the Legendre transform of the Hamiltonian H . Hint: this is easier to verify if $L(v)$ is smooth, in addition to being convex.

As a consequence, if a function $H(p)$ is strictly convex, then we can define the Lagrangian L as the Legendre transform of H . If the Hamiltonian $H(x, p)$ depends, in addition, on a variable $x \in \mathbb{T}^n$ as a parameter, then the Lagrangian $L(x, v)$ is defined as the Legendre transform of $H(x, p)$ in the variable p :

$$L(x, v) = \sup_{p \in \mathbb{R}^n} (p \cdot v - H(x, p)), \quad (2.7.18)$$

with the dual relation

$$H(x, p) = \sup_{v \in \mathbb{R}^n} (p \cdot v - L(x, v)). \quad (2.7.19)$$

We usually refer to x as the spatial variable, and to p as the momentum variable.

Exercise 2.7.4 Compute the Lagrangian $L(x, v)$ for the classical mechanics Hamiltonian

$$H(x, p) = \frac{|p|^2}{2m} + U(x),$$

with a given $m > 0$. Why is it called the classical mechanics Hamiltonian? What is the meaning of the two terms in its definition? Hint: consider the characteristic curves for this Hamiltonian.

Exercise 2.7.5 Consider a sequence of smooth strictly convex Hamiltonians $H_\varepsilon(p)$ that converges locally uniformly, as $\varepsilon \rightarrow 0$, to $H(p) = |p|$. What happens to the corresponding Lagrangians $L_\varepsilon(v)$ as $\varepsilon \rightarrow 0$?

In the context of the forward in time Hamilton-Jacobi equations, with the Hamiltonian that depends on the spatial variable as well, the variational problem (2.7.8)-(2.7.9) is defined as follows. For $t > 0$, and two points $x \in \mathbb{T}^n$ and $y \in \mathbb{T}^n$, we define the function

$$h_t(y, x) = \inf_{\gamma(0)=y, \gamma(t)=x} \int_0^t L(\gamma(s), \dot{\gamma}(s)) ds. \quad (2.7.20)$$

Here, the infimum is taken over all paths γ on \mathbb{T}^n , that are piecewise $C^1[0, t]$, and $L(x, v)$ is the Lagrangian given by (2.7.18). The quantity

$$A(\gamma) = \int_0^t L(\gamma(s), \dot{\gamma}(s)) ds$$

is usually referred to as the Lagrangian action, or simply the action. This is a classical minimization problem, which admits the following result (Tonelli's theorem).

Proposition 2.7.6 *Given any $(t, x, y) \in \mathbb{R}_+^* \times \mathbb{T}^n \times \mathbb{T}^n$, there exists at least one minimizing path $\gamma(s) \in C^2([0, t]; \mathbb{T}^n)$, such that*

$$h_t(y, x) = \int_0^t L(\gamma(s), \dot{\gamma}(s)) ds.$$

Moreover there is $C(t, |x - y|) > 0$ such that

$$\|\dot{\gamma}\|_{L^\infty([0, t])} + \|\ddot{\gamma}\|_{L^\infty([0, t])} \leq C(t, |x - y|). \quad (2.7.21)$$

The function C tends to $+\infty$ as $t \rightarrow 0$ - keeping $|x - y|$ fixed. The function $\gamma(s)$ solves the Euler-Lagrange equation

$$\frac{d}{ds} \nabla_v L(\gamma(s), \dot{\gamma}(s)) - \nabla_x L(\gamma(s), \dot{\gamma}(s)) = 0. \quad (2.7.22)$$

We leave the proof as an exercise but give a hint for the proof. Think of how we proceeded in Section 2.4.2 as blueprint. Consider a minimizing sequence γ_n . First, use the strict convexity of L to obtain the H^1 -estimates for γ_n , thus ensuring compactness in the space of continuous paths and weak convergence to $\gamma \in H^1([0, t])$ with fixed ends. Next, show that the convexity of L implies that γ is, indeed, a minimizer. Finally, derive the Euler-Lagrange equation and show that γ is actually C^∞ . Such a curve γ is called an *extremal*.

The Lax-Oleinik semigroup and viscosity solutions

We now relate the solutions to the Cauchy problem for the Hamilton-Jacobi equations

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \quad t > 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x), \end{aligned} \quad (2.7.23)$$

with a strictly convex Hamiltonian $H(x, p)$, to the minimization problem. We let $L(x, v)$ be the Legendre transform of $H(x, p)$, and define the corresponding function $h_t(y, x)$. Given the initial condition $u_0 \in C(\mathbb{T}^n)$, we define the function

$$u(t, x) = \mathcal{T}(t)u_0(x) = \inf_{y \in \mathbb{T}^n} (u_0(y) + h_t(y, x)). \quad (2.7.24)$$

The following exercise gives the dynamic programming principle, the continuous in time analog of relation (2.2.5) in the time-discrete case we have considered in Section 2.2 .

Exercise 2.7.7 Show that the infimum in (2.7.24) is attained. Also show that $(\mathcal{T}(t))_{t>0}$ is a semi-group: for all $u_0 \in C(\mathbb{T}^n)$ one has

$$\mathcal{T}(t+s)u_0 = \mathcal{T}(t)\mathcal{T}(s)u_0, \quad \text{for all } t \geq 0 \text{ and } s \geq 0,$$

that is,

$$u(t, x) = \inf_{y \in \mathbb{T}^n} (u(s, y) + h_{t-s}(y, x)), \quad (2.7.25)$$

for all $0 \leq s \leq t$, and $\mathcal{T}(0) = I$.

This semigroup is sometimes referred to as the *Lax-Oleinik semigroup*. Here is its link to the Hamilton-Jacobi equations and the viscosity solutions.

Theorem 2.7.8 Given $u_0 \in C(\mathbb{T}^n)$, the function $u(t, x) := \mathcal{T}(t)u_0(x)$ is the unique viscosity solution to the Cauchy problem

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \\ u(0, x) &= u_0(x). \end{aligned} \quad (2.7.26)$$

Proof. The initial condition for $u(t, x)$ holds essentially automatically so we only need to check that u is the viscosity solution. We first show the super-solution property: take $t_0 > 0$ and $x_0 \in \mathbb{T}^n$ and let ϕ be a test function such that (t_0, x_0) is a minimum for $u - \phi$. As usual, without loss of generality, we may assume that $u(t_0, x_0) = \phi(t_0, x_0)$. Consider the minimizing point y_0 such that

$$u(t_0, x_0) = u_0(y_0) + h_{t_0}(y_0, x_0).$$

Let also γ be an extremal of the action between the times $t = 0$ and $t = t_0$, going from y_0 to x_0 : $\gamma(0) = y_0$, $\gamma(t_0) = x_0$. We have, for all $0 \leq t \leq t_0$:

$$\phi(t, \gamma(t)) \leq u(t, \gamma(t)) \leq u_0(y_0) + \int_0^t L(\gamma(s), \dot{\gamma}(s)) ds. \quad (2.7.27)$$

The first inequality above holds because (t_0, x_0) is a minimum of $u - \phi$ and $u(t_0, x_0) = \phi(t_0, x_0)$, and the second follows from the definition of $u(t, \gamma(t))$ in terms of the Lax-Oleinik semigroup. Note that at $t = t_0$ both inequalities in (2.7.27) become equalities: the first one because $u(t_0, x_0) = \phi(t_0, x_0)$, and the second because the curve γ is a minimizer for $u(t_0, x_0)$. This implies

$$\left. \frac{d}{dt} \left(u_0(y_0) + \int_0^t L(\gamma(s), \dot{\gamma}(s)) ds - \phi(t, \gamma(t)) \right) \right|_{t=t_0} \leq 0, \quad (2.7.28)$$

or, in other words

$$\phi_t(t_0, x_0) + \dot{\gamma}(t_0) \cdot \nabla \phi(t_0, x_0) - L(\gamma(t_0), \dot{\gamma}(t_0)) \geq 0. \quad (2.7.29)$$

Using the test point $v = \dot{\gamma}(t_0)$ in the definition (2.7.19) of $H(x, p)$, we then obtain

$$\phi_t(t_0, x_0) + H(x_0, \nabla \phi(t_0, x_0)) \geq 0. \quad (2.7.30)$$

Hence, $u(t, x)$ is a viscosity super-solution to (2.7.26).

To show the sub-solution property, consider a test function $\phi(t, x)$, as well as $t_0 > 0$ and $x_0 \in \mathbb{T}^n$, such that the difference $u - \phi$ attains its maximum at (t_0, x_0) , and assume,

once again, that $u(t_0, x_0) = \phi(t_0, x_0)$. Using the semigroup property (2.7.25), we obtain, for all $t \leq t_0$ and any curve $\gamma(t)$ such that $\gamma(t_0) = x_0$:

$$u(t_0, x_0) \leq u(t, \gamma(t)) + h_{t_0-t}(\gamma(t), x_0) \leq \phi(t, \gamma(t)) + h_{t_0-t}(\gamma(t), x_0). \quad (2.7.31)$$

Given $v \in \mathbb{R}^n$, we take the test curve

$$\gamma(s) = x_0 - (t_0 - s)v$$

in (2.7.31), so that

$$\gamma(t) = x_0 - (t_0 - t)v.$$

Note that the curve

$$\gamma_1(s) = x_0 - (t_0 - t)v + sv,$$

can be used as a test curve in the definition of $h_{t_0-t}(\gamma(t), x_0)$ because we have $\gamma_1(0) = \gamma(t)$, and $\gamma_1(t_0 - t) = x_0$. Using this in (2.7.31) gives

$$\begin{aligned} u(t_0, x_0) &\leq \phi(t, x_0 - (t_0 - t)v) + \int_0^{t_0-t} L(x_0 - (t_0 - t)v + sv, v) ds \\ &= \phi(t, x_0 - (t_0 - t)v) + \int_0^{t_0-t} L(x_0 - sv, v) ds, \end{aligned} \quad (2.7.32)$$

and, once again, this inequality becomes an equality at $t = t_0$, since $u(t_0, x_0) = \phi(t_0, x_0)$. Just as before, differentiating in t at $t = t_0$ gives

$$\phi_t(t_0, x_0) + v \cdot \nabla \phi(t_0, x_0) - L(x_0, v) \leq 0. \quad (2.7.33)$$

As (2.7.33) holds for all $v \in \mathbb{R}^n$, it follows that

$$\phi_t(t_0, x_0) + H(x_0, \nabla \phi(t_0, x_0)) \leq 0. \quad (2.7.34)$$

Therefore, u is also a viscosity sub-solution to (2.7.26), and the proof is complete. \square

Exercise 2.7.9 Show the weak contraction and the finite speed of propagation properties, directly from the Lax-Oleinik formula.

Instant regularization to Lipschitz

We conclude this section with a remarkable result on instant smoothing. We will show that if the initial condition u_0 is continuous on \mathbb{T}^n , then the solution to the Cauchy problem

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \quad t > 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x) \end{aligned} \quad (2.7.35)$$

becomes instantaneously Lipschitz. The improved regularity comes from the strict convexity of the Hamiltonian: indeed, nothing of that sort is true without this assumption, as can be seen from the following exercise.

Exercise 2.7.10 Consider the initial value problem

$$\begin{aligned} u_t + |u_x| &= 0, \quad t > 0, \quad x \in \mathbb{T}^1, \\ u(0, x) &= u_0(x). \end{aligned} \tag{2.7.36}$$

(i) Show that the solution to (2.7.36) is given by

$$u(t, x) = \inf_{|x-y| \leq t} u_0(y). \tag{2.7.37}$$

Hint: one may do this directly but also by considering a family of strictly convex Hamiltonians $H_\varepsilon(p)$ that converges to $H(p) = |p|$ as $\varepsilon \rightarrow 0$, and using the Lax-Oleinik semi-group for

$$\begin{aligned} u_t^\varepsilon + H_\varepsilon(u_x^\varepsilon) &= 0, \quad t > 0, \quad x \in \mathbb{T}^1, \\ u^\varepsilon(0, x) &= u_0(x). \end{aligned} \tag{2.7.38}$$

Exercise 2.7.5 may be useful here.

(ii) Given an example of a continuous initial condition $u_0(x)$ such that the viscosity solution to (2.7.36) is not Lipschitz.

On the other hand, if the Hamiltonian is strictly convex we have the following result.

Theorem 2.7.11 *Let $H(x, p)$ be strictly convex, and $u(t, x)$ be the unique solution to the Cauchy problem*

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \\ u(0, x) &= u_0(x), \end{aligned} \tag{2.7.39}$$

with $u_0 \in C(\mathbb{T}^n)$. Then, the function $u(t, x)$ is Lipschitz in t and x for all $t > 0$.

Let us point the key difference with Proposition 2.6.13: as can be seen from the proof of that proposition, we used the Lipschitz property of the initial condition u_0 , and showed that the solution remains Lipschitz at $t > 0$. Here, the initial condition is not assumed to be Lipschitz but only continuous, and the improved regularity comes from the convexity of the Hamiltonian.

Proof. It is sufficient to consider time intervals of length one, and repeat the argument on the subsequent intervals. Given $0 < t \leq 1$, and $x \in \mathbb{T}^n$, consider the extremal curve $\gamma(s)$ such that $\gamma(t) = x$, and

$$u(t, x) = u_0(\gamma(0)) + \int_0^t L(\gamma(s), \dot{\gamma}(s)) ds. \tag{2.7.40}$$

As $0 \leq s \leq 1$, both $\gamma(s)$ and $\dot{\gamma}(s)$ are uniformly bounded. Of course, on the torus $\gamma(s)$ is always bounded but it would also be bounded for $0 \leq s \leq 1$ if we were considering the problem on \mathbb{R}^n . Take $h \in \mathbb{R}^n$, and define the curve

$$\gamma_1(s) = \gamma(s) + \frac{s}{t}h, \quad 0 \leq s \leq t,$$

so that

$$\gamma_1(0) = \gamma(0), \quad \gamma_1(t) = x + h. \tag{2.7.41}$$

We may use the Lax-Oleinik formula for $u(t, x + h)$ and (2.7.40) for $u(t, x)$, as well as (2.7.41), to write

$$\begin{aligned} u(t, x + h) &= u(t, \gamma_1(t)) \leq u(\gamma_1(0)) + \int_0^t L(\gamma_1(s), \dot{\gamma}_1(s)) ds \\ &= u(t, x) + \int_0^t (L(\gamma_1(s), \dot{\gamma}_1(s)) - L(\gamma(s), \dot{\gamma}(s))) ds. \end{aligned} \quad (2.7.42)$$

The integral in the right side can be estimated as

$$\begin{aligned} \int_0^t (L(\gamma_1(s), \dot{\gamma}_1(s)) - L(\gamma(s), \dot{\gamma}(s))) ds &= \int_0^t [L(\gamma(s) + \frac{s}{t}h, \dot{\gamma}(s) + \frac{1}{t}h) - L(\gamma(s), \dot{\gamma}(s))] ds \\ &\leq \int_0^t \frac{1}{t} \left(sh \cdot \nabla_x L(\gamma(s), \dot{\gamma}(s)) + h \cdot \nabla_v L(\gamma(s), \dot{\gamma}(s)) \right) ds + C_t |h|^2, \end{aligned} \quad (2.7.43)$$

with a constant $C_t > 0$ that may blow up as $t \downarrow 0$. We may now use the Euler-Lagrange equation

$$\frac{d}{ds} \nabla_v L(\gamma(s), \dot{\gamma}(s)) - \nabla_x L(\gamma(s), \dot{\gamma}(s)) = 0$$

to rewrite (2.7.43) as

$$\begin{aligned} &\int_0^t (L(\gamma_1(s), \dot{\gamma}_1(s)) - L(\gamma(s), \dot{\gamma}(s))) ds \\ &\leq \frac{1}{t} \int_0^t h \cdot \left(s \frac{d}{ds} \nabla_v L(\gamma(s), \dot{\gamma}(s)) + \nabla_v L(\gamma(s), \dot{\gamma}(s)) \right) ds + C_t |h|^2 \\ &= h \cdot \nabla_v L(\gamma(t), \dot{\gamma}(t)) + C_t |h|^2. \end{aligned} \quad (2.7.44)$$

Using (2.7.44) in (2.7.42), we obtain

$$u(t, x + h) - u(t, x) \leq h \cdot \nabla_v L(\gamma(t), \dot{\gamma}(t)) + C_t |h|^2, \quad (2.7.45)$$

which proves the Lipschitz regularity in the spatial variable for all $0 < t \leq 1$, because both $\gamma(t)$ and $\dot{\gamma}(t)$ are bounded. Again, the boundedness of $\gamma(t)$ would only play a role if we considered the problem on \mathbb{R}^n , of course. Here, we use the fact that (2.7.45) holds for arbitrary x and $y = x + h$ so that the role of x and y can be switched.

In order to prove the Lipschitz regularity in time, let us examine a small variation of t , denoted by $t + \tau$ with $t + \tau > 0$. Perturbing the extremal curve γ into

$$\gamma_2(s) = \gamma\left(\frac{t}{t + \tau}s\right),$$

we still have

$$\gamma_2(0) = \gamma(0), \quad \gamma_2(t + \tau) = \gamma(t) = x.$$

The same computation as above gives

$$\begin{aligned} u(t + \tau, x) &= u(t + \tau, \gamma_2(t + \tau)) \leq u(\gamma_2(0)) + \int_0^{t + \tau} L(\gamma_2(s), \dot{\gamma}_2(s)) ds \\ &= u(t, x) + \int_0^t (L(\gamma_2(s), \dot{\gamma}_2(s)) - L(\gamma(s), \dot{\gamma}(s))) ds + \int_t^{t + \tau} L(\gamma_2(s), \dot{\gamma}_2(s)) ds. \end{aligned} \quad (2.7.46)$$

It is now straightforward to see that there exists $C'_t > 0$ that depends on t so that

$$u(t + \tau, x) - u(t, x) \leq C'_t |\tau|.$$

Once again, the role of t and $t' = t + \tau$ can be switched, hence $u(t, x)$ is Lipschitz in t as well, for any $t > 0$, finishing the proof. \square

Exercise 2.7.12 (i) Where did we use the strict convexity of the Hamiltonian in the above proof?

(ii) Consider again the initial value problem (2.7.36) with the convex but non strictly convex Hamiltonian $H(p) = |p|$ and a continuous initial condition $u_0(x)$ that is not Lipschitz continuous. Consider a sequence of smooth strictly convex Hamiltonians $H_\varepsilon(p)$ such that $H_\varepsilon(p) \rightarrow H(p)$ as $\varepsilon \rightarrow 0$, locally uniformly on \mathbb{R} . Review the above proof and see what will happen to the Lipschitz constant of the corresponding solution $u^\varepsilon(t, x)$ to the Cauchy problem

$$\begin{aligned} u_t^\varepsilon + H(u_x^\varepsilon) &= 0, \quad t > 0, \quad x \in \mathbb{T}^1, \\ u^\varepsilon(0, x) &= u_0(x), \end{aligned} \tag{2.7.47}$$

constructed by the Lax-Oleinik formula. Hint: again, Exercise 2.7.5 may be useful here.

Exercise 2.7.13 Take $t > 0$ and $\gamma(s)$ an extremal such that u is differentiable at $x = \gamma(t)$. Show that

$$\nabla u(t, x) = \nabla_v L(x, \dot{\gamma}(t)). \tag{2.7.48}$$

and

$$u_t(t, x) = -H(x, \nabla u(t, x)). \tag{2.7.49}$$

2.7.2 Semi-concavity and $C^{1,1}$ regularity

As we have mentioned, the Cauchy problem for a Hamilton-Jacobi equation

$$u_t + H(x, \nabla u) = 0, \tag{2.7.50}$$

with a prescribed initial condition $u(0, x) = u_0(x)$, may have more than one Lipschitz solution, so it is worth asking whether the unique viscosity solution has some additional regularity when the Hamiltonian is strictly convex, so that the solution can be constructed by the Lax-Oleinik semigroup. A relevant notion is that of semi-concavity. Most of the material of this section comes from [?].

Semi-concavity

We begin with the following definition.

Definition 2.7.14 *If B is an open ball in \mathbb{R}^n , F a closed subset of B and K a positive constant, we say that $u \in C(B)$ is K -semi-concave on F if for all $x \in F$, there is $l_x \in \mathbb{R}^n$ such that for all $h \in \mathbb{R}^n$ satisfying $x + h \in B$, we have:*

$$u(x + h) \leq u(x) + l_x \cdot h + K|h|^2. \tag{2.7.51}$$

The function u is said to be K -semi convex on F if $-u$ is K -semi-concave on F .

Exercise 2.7.15 Examine the proof of Theorem 2.7.11 and check that it actually proves that for any $t > 0$ there exists $C_t > 0$ so that $u(t, x)$ is C_t -semi-concave in x .

The next theorem is crucial for the sequel. If u is continuous in an open ball B in \mathbb{R}^n , and F is a closed subset of B , we say that $u \in C^{1,1}(F)$ if u is differentiable in F and ∇u is Lipschitz over F .

Theorem 2.7.16 *Let B be an open ball of \mathbb{R}^n and F closed in B . If $u \in C(B)$ is K -semi-concave and K -semi-convex in F , then $u \in C^{1,1}(F)$.*

Proof. As u is both K semi-concave and K -semi-convex, for all $x \in F$, there are two vectors l_x and m_x such that for all h such that $x + h \in B$ we have

$$\begin{aligned} u(x+h) &\leq u(x) + l_x \cdot h + K|h|^2, \\ u(x+h) &\geq u(x) + m_x \cdot h - K|h|^2 \end{aligned} \tag{2.7.52}$$

which yields

$$(m_x - l_x) \cdot h \leq 2K|h|^2.$$

As this is true for all h sufficiently small, we conclude that $l_x = m_x$ and, therefore, u is differentiable at x , and

$$l_x = m_x = \nabla u(x).$$

Next, we show that ∇u is Lipschitz over F . Given $(x, y, h) \in F \times F \times \mathbb{R}^n$, such that both $x + h \in B$ and $y + h \in B$, the semi-convexity and semi-concavity inequalities, written, respectively, between $x + h$ and x , x and y , and $x + h$ and y , give:

$$\begin{aligned} |u(x+h) - u(x) - \nabla u(x) \cdot h| &\leq K|h|^2 \\ |u(x) - u(y) - \nabla u(y) \cdot (x-y)| &\leq K|x-y|^2 \\ |u(y) - u(x+h) + \nabla u(y) \cdot (x+h-y)| &\leq K|x+h-y|^2. \end{aligned}$$

Adding the three inequalities above, we obtain:

$$|(\nabla u(x) - \nabla u(y)) \cdot h| \leq 3K(|h|^2 + |x-y|^2). \tag{2.7.53}$$

Taking

$$h = |x-y| \frac{\nabla u(x) - \nabla u(y)}{|\nabla u(x) - \nabla u(y)|},$$

in the inequality (2.7.53) gives

$$|\nabla u(x) - \nabla u(y)| \leq 6K|x-y|,$$

which is the Lipschitz property of ∇u that we sought. \square

Improved regularity of the viscosity solutions

Let us come back to the solution $u(t, x)$ to the Cauchy problem

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \\ u(0, x) &= u_0(x). \end{aligned} \tag{2.7.54}$$

We first prove that if γ is a minimizing curve for $u(t, x)$, with $\gamma(t) = x$, then it is also a minimizer for $u(s, \gamma(s))$ for all $0 \leq s \leq t$.

Proposition 2.7.17 *Fix $t > 0$ and $x \in \mathbb{T}^n$, and a minimizing curve γ such that $\gamma(t) = x$, and*

$$u(t, x) = u_0(\gamma(0)) + \int_0^t L(\gamma(s), \dot{\gamma}(s)) ds. \tag{2.7.55}$$

Then for all $0 \leq s \leq s' \leq t$ we have

$$u(s', \gamma(s')) = u_0(\gamma(0)) + \int_0^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma = u(s, \gamma(s)) + \int_s^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma. \tag{2.7.56}$$

Exercise 2.7.18 Relate the result of this proposition to the dynamic programming principle.

Proof. The Lax-Oleinik formula implies that for all $0 < s < t$ we have

$$u(s, \gamma(s)) \leq u_0(\gamma(0)) + \int_0^s L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma.$$

Assume that for some $0 < s < t$, we have a strict inequality

$$u(s, \gamma(s)) < u_0(\gamma(0)) + \int_0^s L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma. \tag{2.7.57}$$

Then, there exists a curve $\gamma_1(s')$, $0 \leq s' \leq s$, such that $\gamma_1(s) = \gamma(s)$, and

$$u_0(\gamma_1(0)) + \int_0^s L(\gamma_1(\sigma), \dot{\gamma}_1(\sigma)) d\sigma < u_0(\gamma(0)) + \int_0^s L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma.$$

Then, we can consider the concatenated curve $\gamma_2(s)$ so that $\gamma_2(s') = \gamma_1(s')$ for $0 \leq s' \leq s$, and $\gamma_2(s') = \gamma(s')$ for $s \leq s' \leq t$. The resulting curve is piece-wise $C^1[0, t]$, hence is an allowed trajectory. This would give

$$\begin{aligned} u(t, \gamma(t)) &= u_0(\gamma(0)) + \int_0^s L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma + \int_s^t L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma \\ &> u_0(\gamma_2(0)) + \int_0^t L(\gamma_2(s), \dot{\gamma}_2(s)) ds, \end{aligned} \tag{2.7.58}$$

which would contradict the extremal property of the curve γ between the times 0 and t . Therefore, (2.7.57) can not hold, and for all $0 \leq s \leq s' \leq t$ we have:

$$u(s', \gamma(s')) = u_0(\gamma(0)) + \int_0^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma = u(s, \gamma(s)) + \int_s^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma. \tag{2.7.59}$$

This finishes the proof of Proposition 2.7.17. \square

Definition 2.7.19 We say that $\gamma : [0, t] \rightarrow \mathbb{T}^n$ is calibrated by u if (2.7.55) holds.

Let us define the conjugate semigroup of the Lax-Oleinik semigroup by:

$$\tilde{\mathcal{T}}(t)u_0(x) = \sup_{y \in \mathbb{T}^n} (u_0(y) - h_t(x, y)), \quad \forall u_0 \in C(\mathbb{T}^n), \quad t > 0. \quad (2.7.60)$$

We will denote $\tilde{u}(t, x) = \tilde{\mathcal{T}}(t)u_0(x)$. The following lemma is proved exactly as Theorem 2.7.11.

Lemma 2.7.20 Let $u_0 \in C(\mathbb{T}^n)$ and $\sigma > 0$. There is $K(\sigma) > 0$ such that $\tilde{\mathcal{T}}(\sigma)u_0$ is $K(\sigma)$ -semi-convex. The constant $K(\sigma)$ blows up as $\sigma \rightarrow 0$.

Given $0 < s < s'$, we define the set $\Gamma_{s, s'}[u_0]$ as the union of all points $(s_1, x) \in [s, s'] \times \mathbb{T}^n$, so that the extremal calibrated by u , which passes through the point x at the time s_1 can be continued forward in time until the time s' , and backward in time until the time s .

Corollary 2.7.21 Let $u_0 \in C(\mathbb{T}^n)$ and $u(t, x) = \mathcal{T}(t)u_0(x)$, and $0 < s_1 < s_2$, then for any $\varepsilon > 0$, the function $u \in C^{1,1}(\Gamma_{s_1, s_2+\varepsilon} \cap ([s_1, s_2] \times \mathbb{T}^n))$.

Proof. Let us take $(s, x_0) \in \Gamma_{s_1, s_2+\varepsilon}$, with $s_1 \leq s \leq s_2$, so that that the extremal γ such that $x_0 = \gamma(s)$ can be continued past the time s , until the time $s_2 + \varepsilon$.

Let us first deal with the spatial regularity. As we have mentioned in Exercise 2.7.15, there is $K > 0$ depending on s_1 such that the function $u(s, x)$ is K -semi-concave at all $x \in \mathbb{T}^n$ for all $s \geq s_1$, in particular, at x_0 . Hence, we only need to argue that u is semi-convex at x_0 , and here we are going to use the fact that $(s, x_0) \in \Gamma_{s_1, s_2+\varepsilon}$. Note that for all $y \in \mathbb{R}^n$ we have, by the Lax-Oleinik formula,

$$u(s_2 + \varepsilon, y) \leq u(s, x_0) + h_{s_2+\varepsilon-s}(x_0, y). \quad (2.7.61)$$

In addition, the calibration relation (2.7.59) implies that if $x_0 = \gamma(s)$ and $(s, x_0) \in \Gamma_{s_1, s_2+\varepsilon}$, then equality is attained when $y = \gamma(s_2 + \varepsilon)$. We conclude that in this case we have

$$u(s, x_0) = \sup_{y \in \mathbb{T}^n} (u(s_2 + \varepsilon, y) - h_{s_2+\varepsilon-s}(x_0, y)) = \tilde{\mathcal{T}}(s_2 + \varepsilon - s)[u(s_2 + \varepsilon, \cdot)](x_0).$$

It follows from Lemma 2.7.20 that there is a constant \tilde{K} depending on ε , such that $u(s, \cdot)$ is \tilde{K} -semi-convex in x on $\Gamma_{s_1, s_2+\varepsilon} \cap ([s_1, s_2] \times \mathbb{T}^n)$.

Theorem 2.7.16 now implies that the function $u(s, \cdot)$ is $C^{1,1}$ in x on the set $\Gamma_{s_1, s_2+\varepsilon}$ for all $s_1 \leq s \leq s_2$. To end the proof, one just has to invoke relation (2.7.49) in Exercise 2.7.13 to obtain the corresponding regularity in the time variable. \square

This corollary may not, at first sight, look so striking. To enjoy its scope, let us specialize it to the solutions to the stationary equation

$$H(x, \nabla u) = 0, \quad (2.7.62)$$

assuming that they exist. Corollary 2.7.21 allows us to discover the following

Corollary 2.7.22 Consider a solution u of (2.7.62), and let F be the set of all points $x \in \mathbb{T}^n$ such that there exists $\varepsilon_x > 0$ and a C^1 curve $\gamma : (-\varepsilon_x, \varepsilon_x) \rightarrow \mathbb{T}^n$ such that $\gamma(0) = x$ and

$$u(\gamma(\varepsilon_x)) - u(\gamma(-\varepsilon_x)) = \int_{-\varepsilon_x}^{\varepsilon_x} L(\gamma(s), \dot{\gamma}(s)) ds. \quad (2.7.63)$$

Then $u \in C^{1,1}(F)$.

In other words, u is $C^{1,1}$ at every point through which an extremal of the Lagrangian passes, as opposed to ending at this point.

Let us examine some further consequences of this fact, in the form of a few exercises, just to give a glimpse of how far reaching these considerations can be. Their solution does not need more tools or ideas than the ones already presented, but they are fairly elaborate. We begin with an application of the finite speed of propagation property.

Exercise 2.7.23 Let $u(x)$ be a Lipschitz viscosity solution of

$$H(x, \nabla u) = 0$$

in a bounded open subset Ω of \mathbb{R}^n . Show that, for every open subset Ω_1 of Ω such that $\bar{\Omega}_1$ is compactly embedded in Ω , there is $\varepsilon > 0$ such that, for all $t \in [0, \varepsilon]$ and $x \in \Omega_1$ we have

$$u(x) = \mathcal{T}(t)u(x).$$

We continue with a statement that looks surprisingly elementary. However its solution is not.

Exercise 2.7.24 Let Ω be an open subset of \mathbb{R}^n and u_p a sequence in $C^1(\Omega)$, such that

$$|\nabla u_p| = 1 \text{ for all } p.$$

Show that all uniform limits of u_p are C^1 functions. Hint: if $x_0 \in \Omega$, then, for small $\varepsilon > 0$, the function $u_p(x)$ coincides, in a small neighborhood of x , with both $\mathcal{T}(\varepsilon)u_p$ and $\tilde{\mathcal{T}}(\varepsilon)u_p$. Note that the Hamiltonian is not strictly convex, so some care needs to be given to the definition of the Lax-Oleinik semigroup and its adjoint. If in doubt, look at (2.7.64) below.

We end the section with two regularity properties of the distance function. Recall that, if S is a subset of \mathbb{R}^n , the distance function to S is given by

$$d_S(x) = \inf_{v \in S} |x - v|.$$

It is, obviously, a Lipschitz function with Lipschitz constant 1. We can say much more, just recalling the age-old fact that the shortest path between two points is the line joining these two points: this makes d_S a viscosity solution of $|\nabla d| = 1$, or, even better:

$$|\nabla d|^2 = 1. \quad (2.7.64)$$

We may use the previous theory for the following results.

Exercise 2.7.25 If S is a compact set, $x_0 \notin S$ and v is such that

$$|x - v| = d_S(x),$$

then d_S is $C^{1,1}$ on the line segment $[v, x]$.

Exercise 2.7.26 If S is a convex set, then d_S is $C^{1,1}$ outside S .

If you are stuck with any of the above three exercises, see [?].

2.8 Large time behavior in particular cases

For the rest of this chapter, we go back to the long term behavior of the solutions to the Hamilton-Jacobi equations but unlike in Section 2.3, we now consider the inviscid case. In this initial section, we will focus on two examples. First, we will consider equations of the form

$$u_t + \frac{1}{2}|\nabla u|^2 = f(x), \quad (2.8.1)$$

with the classical Hamiltonian

$$H(x, p) = \frac{|p|^2}{2} - f(x). \quad (2.8.2)$$

This equation arises naturally in the context of classical mechanics. The strict convexity of the classical Hamiltonian (2.8.2) will allow us to use the Lax-Oleinik formula to understand the long time behavior for the solutions to (2.8.1), in a straightforward and elegant way.

Then, we will consider the Hamilton-Jacobi equation

$$u_t + R(x)\sqrt{1 + |\nabla u|^2} = 0, \quad (2.8.3)$$

with the Hamiltonian

$$H(x, p) = R(x)\sqrt{1 + |p|^2}. \quad (2.8.4)$$

The Hamiltonian in (2.8.4) is locally strictly convex in its second variable but not uniformly strictly convex. We could also attack the problem via the Lax-Oleinik formula, with a little extra technical argument due to the lack of the global strict convexity. We will not, however, rely on the strict convexity in any form in the analysis of the long time behavior for the solutions to (2.8.3). The separate arguments that we are going to display for this problem will work, at almost no additional cost, for the important class of Hamiltonians of the form

$$H(x, p) = |\nabla u| - f(x), \quad (2.8.5)$$

which are not strictly convex even locally. The proof is inspired by the arguments in [?].

Let us mention, looking ahead, that despite the difference in the approaches to the two cases, we will see some strong similarities in the underlying dynamics that will allow us to address the general case in the next section. We chose to start with these examples as the proofs here are much more concrete.

On the technical side, we will assume for (2.8.1) that the function $f(x)$ is smooth, and that the function $R(x)$ in (2.8.3) is smooth and positive: there exists $R_0 > 0$ so that

$$R(x) \geq R_0 > 0 \text{ for all } x \in \mathbb{T}^n, \quad (2.8.6)$$

and will use the notation

$$\bar{R} = \|R\|_{L^\infty}. \quad (2.8.7)$$

Note that the assumptions for the Hamiltonian $H(x, p) = R(x)\sqrt{1 + |p|^2}$ fall in line with those made in Section 2.3 on the convergence to the viscous waves, and in Section 2.6 on the existence of the inviscid waves and of the solutions to the inviscid Cauchy problem. As usual, the smoothness assumptions on the function $f(x)$ and $R(x)$ can be greatly relaxed.

Let us first explain how equation (2.8.3) comes up from simple geometric considerations. Consider a family of hypersurfaces $\Sigma(t)$ of \mathbb{R}^{n+1} , moving according to an imposed normal velocity $R(y)$:

$$V_n = R(y), \quad y \in \mathbb{R}^{n+1}, \quad (2.8.8)$$

the function $R(y)$ being given and positive. Assume that, at each time $t \geq 0$, the surface $\Sigma(t)$ is the level set of a function $v(t, y)$:

$$\Sigma(t) = \{y \in \mathbb{R}^{n+1} : v(t, y) = 0\}.$$

It is easy to see that the normal velocity V_n at the point y , at time t , is given by

$$V_n(t, y) = \frac{v_t(t, y)}{|\nabla v(t, y)|},$$

so that the evolution equation for the function $v(t, y)$ is

$$v_t = R(y)|\nabla v| \quad \text{on } \Sigma(t). \quad (2.8.9)$$

This evolution equation is interesting in itself, and is known in the literature on the mathematical theory of combustion as the G-equation. It also appears in many computational methods where it is often called the level sets equation. In particular, it allows to model coalescence of objects in digital animation.

We are going to consider a special situation when $\Sigma(t)$ is given in the form of a graph of a periodic function $u(t, x)$, $x \in \mathbb{T}^n$, that is, writing $y = (x, y_{n+1})$, with $x \in \mathbb{T}^n$ and $y_{n+1} \in \mathbb{R}$, we have

$$v(t, y) = y_{n+1} - u(t, x), \quad x \in \mathbb{T}^n,$$

and also that $R(y)$ is actually a function of the form $R(x)$ – it depends only on the first n coordinates of y . Then we obtain from (2.8.9)

$$u_t + R(x)\sqrt{1 + |\nabla_x u|^2} = 0, \quad x \in \mathbb{T}^n, \quad (2.8.10)$$

which is (2.8.3).

We will begin with the analysis of the wave solutions to (2.8.1) and (2.8.3) – as we will soon see, this study is essentially identical for both problems. Then, we will consider the long time convergence to the wave solutions, and there the two analyses will diverge.

2.8.1 Counting the waves

The first step is to understand the wave solutions to (2.8.1) and (2.8.3). Note that a wave solution to (2.8.3) satisfies

$$R(x)\sqrt{1 + |\nabla u|^2} = c, \quad x \in \mathbb{T}^n, \quad (2.8.11)$$

an equation that can be alternatively stated as

$$|\nabla u(x)|^2 = g(x), \quad x \in \mathbb{T}^n, \quad (2.8.12)$$

with

$$g(x) = \frac{c^2}{R^2(x)} - 1. \quad (2.8.13)$$

On the other hand, a wave solution to (2.8.1) solves

$$\frac{1}{2}|\nabla u(x)|^2 = f(x) + c, \quad x \in \mathbb{T}^n, \quad (2.8.14)$$

that can also be re-stated as (2.8.13), but now with

$$g(x) = 2(f(x) + c). \quad (2.8.15)$$

Thus, in both cases, existence of the wave solutions is equivalent to the question of existence of steady solutions to (2.8.12).

Identification of the speed

We begin with the following.

Proposition 2.8.1 *A solution to an equation of the form*

$$|\nabla u(x)|^2 = f(x) + \gamma, \quad x \in \mathbb{T}^n, \quad (2.8.16)$$

with a smooth function f exists if and only if

$$\gamma = -\min_{x \in \mathbb{T}^n} f(x). \quad (2.8.17)$$

In other words, a solution to

$$|\nabla u(x)|^2 = f(x), \quad x \in \mathbb{T}^n, \quad (2.8.18)$$

exists if and only if

$$\min_{x \in \mathbb{T}^n} f(x) = 0. \quad (2.8.19)$$

A consequence of this proposition is that the only c such that equation

$$R(x)\sqrt{1 + |\nabla u_\infty|^2} = c, \quad x \in \mathbb{T}^n, \quad (2.8.20)$$

has a solution $u_\infty(x)$ is $c = \bar{R}$, as seen from (2.8.12)-(2.8.13).

To understand the main idea of the proof, note that the unique γ for which (2.8.16) has a solution, can be alternatively defined as the only value of γ such that each solution to the Cauchy problem

$$\begin{aligned} u_t + |\nabla u|^2 &= f(x) + \gamma, \quad t > 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x), \end{aligned} \quad (2.8.21)$$

is uniformly bounded in time. This is an immediate consequence of the comparison principle for the viscosity solutions.

Exercise 2.8.2 (i) Explain this point: show that if $\gamma \neq c$ then the solution to the Cauchy problem (2.8.21) can not remain bounded as $t \rightarrow +\infty$, and, conversely, if $\gamma = c$ then it remains bounded as $t \rightarrow +\infty$.

(ii) Show also that c is the unique value γ such that there exists both a sub-solution \underline{u} and a super-solution \bar{u} to

$$|\nabla u|^2 = f(x) + \gamma, \quad x \in \mathbb{T}^n. \quad (2.8.22)$$

Hint: solutions to (2.8.21) may be helpful here.

Proof of Proposition 2.8.1. We know from the Lions-Papanicolaou-Varadhan theorem that for each $f \in C(\mathbb{T}^n)$ there exists some $\gamma \in \mathbb{R}$ such that a solution to

$$|\nabla u(x)|^2 = f(x) + \gamma, \quad x \in \mathbb{T}^n, \quad (2.8.23)$$

exists. We need to show that

$$\gamma = -\min_{x \in \mathbb{T}^n} f(x). \quad (2.8.24)$$

As in Exercise 2.8.2(ii), we only need to construct a sub-solution and a super-solution to (2.8.23) for γ as in (2.8.24). First, observe that if

$$\gamma + \min_{x \in \mathbb{T}^n} f(x) \geq 0, \quad (2.8.25)$$

then all constants are sub-solutions to (2.8.23).

On the other hand, a quadratic function of the form

$$\bar{u}(x) = \frac{\alpha}{2}|x - x_0|^2, \quad (2.8.26)$$

with some $x_0 \in \mathbb{T}^n$, is a super-solution to (2.8.23) if

$$\alpha^2|x - x_0|^2 \geq f(x) + \gamma, \quad \text{for all } x \in \mathbb{T}^n. \quad (2.8.27)$$

It follows that, in particular,

$$f(x_0) + \gamma \leq 0,$$

hence such super-solution can exist only if

$$\gamma + \min_{x \in \mathbb{T}^n} f(x) \leq 0. \quad (2.8.28)$$

On the other hand, if (2.8.28) does hold, x_0 is a minimum of $f(x)$, and f is smooth, as we assume here, then (2.8.27) does hold if we choose $\alpha > 0$ to be sufficiently large.

Thus, if $\gamma = -\min_{x \in \mathbb{T}^n} f(x)$ then we can find both a sub-solution and a super-solution to (2.8.23), finishing the proof. \square

Exercise 2.8.3 Note that the super-solution we have constructed in (2.8.26) is not periodic. Explain why this is not an issue.

Exercise 2.8.4 We did use the assumption that $f(x)$ is smooth in the construction of the super-solution in the above proof. Show that nevertheless the conclusion of Proposition 2.8.1 holds for $f \in C(\mathbb{T}^n)$. Hint: approximate $f \in C(\mathbb{T}^n)$ by a sequence of smooth functions f_k that converges uniformly to f and obtain a uniform Lipschitz bound for the solutions to

$$|\nabla u_k|^2 = f_k(x) + \gamma_k, \quad \gamma_k := -\min_{x \in \mathbb{T}^n} f_k(x),$$

such that $u_k(0) = 0$. Finally, use the stability property of the viscosity solutions to show that u_k converges, along a subsequence, to a viscosity solution to

$$|\nabla u|^2 = f(x) + \gamma, \quad \gamma := -\min_{x \in \mathbb{T}^n} f(x). \quad (2.8.29)$$

A simple example of the non-uniqueness of the waves

Before proceeding with the description of the set of the solutions to

$$|\nabla u(x)|^2 = f(x), \quad x \in \mathbb{T}^n, \quad (2.8.30)$$

under the assumption that

$$\min_{x \in \mathbb{T}^n} f(x) = 0, \quad (2.8.31)$$

let us explain why the solutions may be not unique. This is a big difference with the viscous case

$$-\Delta u + H(x, \nabla u) = c, \quad (2.8.32)$$

described in Theorem 2.3.1, where both the speed c and the solution u are unique.

We consider a very simple example in one dimension:

$$|u'| = f(x), \quad x \in \mathbb{T}^1. \quad (2.8.33)$$

Assume that $f \in C^1(\mathbb{T}^1)$ is 1/2-periodic, satisfies

$$f(x) > 0 \text{ on } (0, 1/2) \cup (1/2, 1), \text{ and } f(0) = f(1/2) = f(1) = 0.$$

and is symmetric with respect to $x = 1/4$ (and thus $x = 3/4$). Let u_1 and u_2 be 1-periodic and be defined, over a period, as follows:

$$u_1(x) = \begin{cases} \int_0^x f(y) dy, & 0 \leq x \leq \frac{1}{2}, \\ \int_x^1 f(y) dy, & \frac{1}{2} \leq x \leq 1, \end{cases} \quad u_2(x) = \begin{cases} \int_0^x f(y) dy, & 0 \leq x \leq \frac{1}{4}, \\ \int_x^{1/2} f(y) dy, & \frac{1}{4} \leq x \leq \frac{1}{2}, \\ u_2 \text{ is } \frac{1}{2}\text{-periodic.} \end{cases}$$

Note that $u_1(x)$ is continuously differentiable but $u_2(x)$ is only Lipschitz: its graph has corners at $x = 1/4$ and $x = 3/4$.

Exercise 2.8.5 Verify that both u_1 and u_2 are viscosity solutions of (2.8.33), and u_2 cannot be obtained from u_1 by the addition a constant. Pay attention to what happens at $x = 1/4$ and $x = 3/4$ with $u_2(x)$. Why can't you construct a solution that would have a corner at a minimum rather than the maximum?

Trajectories at very negative times

The above example of non-uniqueness inspires a more systematic study of the steady solutions to

$$|\nabla u|^2 = f(x), \quad x \in \mathbb{T}^n, \quad (2.8.34)$$

in order to understand how many steady solutions this problem may have. We assume that the function f is smooth and non-negative:

$$f(x) \geq 0 \quad \text{for all } x \in \mathbb{T}^1, \quad (2.8.35)$$

This ensures existence of a solution to (2.8.34), via Proposition 2.8.1. The smoothness assumption on the function f is adopted merely for convenience, continuity of f would certainly suffice.

A non-technical assumption is that the function $f(x)$ has finitely many zeroes x_1, \dots, x_N . We will see that an absolutely crucial role in the analysis will be played by the set

$$\mathcal{Z} = \{x : f(x) = 0\} = \{x_1, \dots, x_N\}. \quad (2.8.36)$$

What follows is a (much simplified) adaptation of the last chapter of the book of Fathi [?].

As we have mentioned, the viscosity solutions to (2.8.34) exist by Proposition 2.8.1 and our assumptions on f . They satisfy the Lax-Oleinik formula: for any $t < 0$ we have

$$u(x) = \inf_{\gamma(0)=x} \left(u(\gamma(t)) + \int_t^0 \left(\frac{|\dot{\gamma}(s)|^2}{4} + f(\gamma(s)) \right) ds \right). \quad (2.8.37)$$

We know from the preceding section that the infimum is, in fact, a minimum, attained at an extremal of the Lagrangian, that we denote $\gamma_t(s)$, $t \leq s \leq 0$. The Lagrangian associated to the Hamiltonian $H(x, p) = |p|^2 - f(x)$ is

$$L(x, v) = \frac{|v|^2}{4} + f(x). \quad (2.8.38)$$

This means, in particular, that $L(x, v)$ is nonnegative and vanishes only at the points of the form $(x, v) = (x_i, 0)$, $i \in \{1, \dots, N\}$. Hence, we expect that the minimizers in (2.8.37) should prefer to stay near the points where f vanishes, and move very slowly around those points. To formalize this idea, we would like to send the starting time $t \rightarrow -\infty$ and say that each minimizing curve $\gamma_t(s)$ is near one of $x_i \in \mathcal{Z}$, for s sufficiently large and negative.

Proposition 2.8.6 *The function $u(x)$ can be written as*

$$u(x) = \inf_{x_i \in \mathcal{Z}} \inf_{\gamma(-\infty)=x_i, \gamma(0)=x} \left(u(x_i) + \int_{-\infty}^0 \left(\frac{|\dot{\gamma}(s)|^2}{4} + f(\gamma(s)) \right) ds \right), \quad (2.8.39)$$

with the infimum taken over all curves $\gamma(s)$ such that $\gamma(0) = x$ and $\gamma(s) \rightarrow x_i$ as $s \rightarrow -\infty$.

Proof. First, note that $u(x)$ is bounded by the right side of (2.8.39), as follows immediately from the Lax-Oleinik formula (2.8.37). We need to show that equality is actually attained. Let

us fix $x \in \mathbb{T}^n$, take $t < 0$ large and negative, and consider the corresponding minimizer $\gamma_t(s)$, calibrated by u , so that

$$u(x) = u(\gamma_t(s)) + \int_s^0 \left(\frac{|\dot{\gamma}_t(\sigma)|^2}{4} + f(\gamma_t(\sigma)) \right) d\sigma, \quad \text{for all } t \leq s \leq 0. \quad (2.8.40)$$

The uniform bounds on $\gamma_t(s)$ and $\dot{\gamma}_t(s)$ imply that there is a sequence $t_n \rightarrow -\infty$ such that $\gamma_{t_n}(s)$ converges, locally uniformly, to a limit $\gamma(s)$ that is defined for all $s < 0$. Passing to the limit $t_n \rightarrow -\infty$ in (2.8.40) we see that $\gamma(s)$ is also calibrated by u : for all $s < 0$ we have

$$u(x) = u(\gamma(s)) + \int_s^0 \left(\frac{|\dot{\gamma}(\sigma)|^2}{4} + f(\gamma(\sigma)) \right) d\sigma. \quad (2.8.41)$$

We claim that there exists $x_k \in \mathcal{Z}$ so that

$$\lim_{s \rightarrow -\infty} \gamma(s) = x_k. \quad (2.8.42)$$

To see that (2.8.42) holds, take $\varepsilon > 0$ and consider the set

$$D_\varepsilon = \{y \in \mathbb{T}^n : |y - x_i| \leq \varepsilon \text{ for some } x_i \in \mathcal{Z}\}.$$

If $\varepsilon > 0$ is sufficiently small, then D_ε is a union of N pairwise disjoint balls

$$B_\varepsilon^{(k)} = \{y \in \mathbb{T}^n : |y - x_k| \leq \varepsilon\}.$$

The function $f(y)$ is strictly positive outside of D_ε : there exists $\lambda_\varepsilon > 0$ so that $f(y) > \lambda_\varepsilon$ for all $y \notin D_\varepsilon$. It follows from (2.8.41) that the total time that $\gamma(s)$ spends outside of D_ε is also bounded:

$$|\{s < 0 : \gamma(s) \notin D_\varepsilon\}| \leq \frac{2\|u\|_{L^\infty}}{\lambda_\varepsilon}. \quad (2.8.43)$$

Exercise 2.8.7 Show that there exists $\mu_\varepsilon > 0$ such that if $s_1 < s_2 < 0$, and $\gamma(s_1) \in B_\varepsilon^{(k)}$ while $\gamma(s_2) \in B_\varepsilon^{(k')}$ with $k \neq k'$, then

$$\int_{s_1}^{s_2} \left(\frac{|\dot{\gamma}(s)|^2}{4} + f(\gamma(s)) \right) ds \geq \mu_\varepsilon. \quad (2.8.44)$$

Hint: show that if the switch from $B_\varepsilon^{(k)}$ to $B_\varepsilon^{(k')}$ happens "quickly" then the contribution of the first term inside the integral is bounded from below, and if this switch happens "slowly", then the contribution of the second term inside the integral is bounded from below.

A consequence of (2.8.43) and Exercise 2.8.7 is that there exists T_ε and $1 \leq k \leq N$ such that $\gamma(s) \in B_\varepsilon(x_k)$ for all $t < T_\varepsilon$. This implies (2.8.42).

Now, we may let $s \rightarrow -\infty$ in (2.8.41) to obtain

$$u(x) = u(x_k) + \int_{-\infty}^0 \left(\frac{|\dot{\gamma}(\sigma)|^2}{4} + f(\gamma(\sigma)) \right) d\sigma. \quad (2.8.45)$$

It follows that $u(x)$ is bounded from below by the right side of (2.8.39), and the proof of Proposition 2.8.6 is complete. \square

Classification of steady solutions

We can now classify all solutions to

$$|\nabla u|^2 = f(x), \quad x \in \mathbb{T}^n. \quad (2.8.46)$$

The reader may remember that the proof of uniqueness of the waves in Theorems 2.3.1 and the long time behavior in Theorem 2.3.3 in the viscous case relied crucially on the strong maximum principle and the Harnack inequality for parabolic equations. It is exactly the lack of these properties for the inviscid Hamilton-Jacobi equations that leads to the non-uniqueness of the solutions to (2.8.87), and to different possible long time behaviors of the solutions to the corresponding Cauchy problem.

Let us set

$$S(x_i, x) = \inf_{\gamma(-\infty)=x_i} \int_{-\infty}^0 \left(\frac{|\dot{\gamma}(s)|^2}{4} + f(\gamma(s)) \right) ds. \quad (2.8.47)$$

It may be seen as the energy of a connection between x_i and x , or, in a more mathematically precise way, as a sort of distance between x_i and x . This fruitful point of view, developed in [?], will not be pushed further here. The next theorem classifies all solutions to (2.8.46).

Theorem 2.8.8 *Let $\{x_1, \dots, x_N\}$ be the set of zeros of a smooth non-negative function $f(x)$. Given a collection of numbers $\{a_1, \dots, a_N\}$ there is a unique solution $u(x)$ to (2.8.46), such that*

$$u(x_i) = a_i \text{ for all } 1 \leq i \leq N, \quad (2.8.48)$$

if and only if

$$a_j \leq a_i + S(x_i, x_j), \quad \text{for all } 1 \leq i, j \leq N. \quad (2.8.49)$$

Condition (2.8.49) has a simple interpretation: in order to be able to assign a value a_j at the zero x_j , the trajectory $\gamma(t) \equiv x_j$ for all $t < 0$, should be a minimizer.

Proof. Proposition 2.8.6 already shows that the values of $u(x_i)$ determine the value of $u(x)$ for all $x \in \mathbb{T}^n$, and that if a solution exists and (2.8.48) holds, then

$$a_j = \inf_{i \in \{1, \dots, N\}} (a_i + S(x_i, x_j)). \quad (2.8.50)$$

This implies (2.8.49).

To prove existence of a solution to (2.8.46) such that $u(x_i) = a_i$ for all $1 \leq i \leq N$, for given a_i , $i = 1, \dots, N$, that satisfy (2.8.49), set

$$u(x) = \inf_{i \in \{1, \dots, N\}} (a_i + S(x_i, x)). \quad (2.8.51)$$

Using the by now familiar arguments, it is easy to see that u is a solution to (2.8.46). Moreover, we have

$$u(x_j) = \inf_{i \in \{1, \dots, N\}} (a_i + S(x_i, x_j)).$$

This, together with (2.8.49) implies $u(x_j) = a_j$. \square

Exercise 2.8.9 Apply the above theorem to the equation $|u'| = f(x)$ on \mathbb{T}^1 , with a non-negative function $f(x)$ vanishing at 2 or 3 distinct points. Find out how many different solutions one may have.

Exercise 2.8.10 Let Ω be a smooth bounded subset of \mathbb{R}^n . Assume that f is nonnegative and vanishes only at a finite number of points and $u_0 \in C(\partial\Omega)$. Find a necessary and sufficient condition on the values of u_0 so that the boundary value problem

$$\begin{aligned} |\nabla u|^2 &= f(x), & x \in \Omega, \\ u(x) &= u_0(x), & x \in \partial\Omega, \end{aligned} \tag{2.8.52}$$

is well-posed. Count its solutions. If you have difficulty, we recommend that you read the very remarkable study of the non-uniqueness in Lions [?].

2.8.2 The large time behavior: a strictly convex example

The above analysis for the classification of the wave solutions can be adapted to understand the long time behavior of the solutions to the Cauchy problem

$$\begin{aligned} u_t + |\nabla u|^2 &= f(x), & t > 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x). \end{aligned} \tag{2.8.53}$$

The next theorem gives an (almost) explicit form of the asymptotic limit of the solution to (2.8.53), and exhibits again the role of the set \mathcal{Z} in the dynamics.

Theorem 2.8.11 *Let $u(t, x)$ be the solution to (2.8.53) with a smooth non-negative function $f(x)$ that vanishes on a finite set $\mathcal{Z} = \{x_1, \dots, x_N\}$, and $u_0 \in C(\mathbb{T}^n)$. Then, the function $u(t, x)$ is non-increasing in t on the set \mathcal{Z} , so that for each $x_k \in \mathcal{Z}$ the limit*

$$a_k := \lim_{t \rightarrow +\infty} u(t, x_k) \tag{2.8.54}$$

exists. Moreover, for all $x \in \mathbb{T}^n$ we have

$$\lim_{t \rightarrow +\infty} u(t, x) = \inf_{x_k \in \mathcal{Z}} (a_k + S(x_k, x)), \tag{2.8.55}$$

with $S(x_i, x)$ as in (2.8.47):

$$S(x_i, x) = \inf_{\gamma(-\infty)=x_i} \int_{-\infty}^0 \left(\frac{|\dot{\gamma}(s)|^2}{4} + f(\gamma(s)) \right) ds. \tag{2.8.56}$$

We will use throughout the proof the fact that the unique viscosity solution to (2.8.53) is uniformly bounded and is uniformly Lipschitz: there exists $C > 0$ so that for all $t \geq 1$ we have

$$\|u(t, \cdot)\|_{L^\infty} \leq C, \quad \|u_t(t, \cdot)\|_{L^\infty} + \|\nabla u(t, \cdot)\|_{L^\infty} \leq C. \tag{2.8.57}$$

The Lipschitz bound in (2.8.57) follows from Theorem 2.7.11, while the uniform bound on $u(t, x)$ is a simple consequence of the fact that steady solutions to (2.8.53) exist under our assumptions on $f(x)$. These estimates already tell us that there exists a sequence $t_n \rightarrow +\infty$ such that the sequence of functions $v_n(t, x) = u(t + t_n, x)$ converges in $L^\infty(\mathbb{T}^n)$ and locally uniformly in t , to a limit $\tilde{u}(t, x)$. However, we do not know that the limit is unique, nor that it is time-independent, nor that it is a solution to (2.8.53).

Monotonicity on \mathcal{Z}

We first prove that $u(t, x)$ is non-increasing in t for $x \in \mathcal{Z}$. If $u(t, x)$ were actually smooth at $x_k \in \mathcal{Z}$, then, as $f(x_k) = 0$ for $x_k \in \mathcal{Z}$, we would have

$$u_t(t, x_k) = -|\nabla u(t, x_k)|^2 \leq 0, \quad (2.8.58)$$

as desired. However, we only know that $u(t, x)$ is a viscosity solution, hence we can not use (2.8.58) directly. Instead, we fix $t_0 \geq 0$ and consider the function

$$\bar{u}(t, x) = u(t_0, x) + (t - t_0)f(x). \quad (2.8.59)$$

We claim that $\bar{u}(t, x)$ is a viscosity super-solution to (2.8.53). Consider a test function φ such that the difference $\bar{u} - \varphi$ attains its minimum at (t_1, x_1) . As $\bar{u}(t, x)$ is smooth in t , we have, in particular, that

$$0 \leq \varphi_t(t_1, x_1) - \bar{u}_t(t_1, x_1),$$

which implies

$$0 \leq \varphi_t(t_1, x_1) - \bar{u}_t(t_1, x_1) = \varphi_t(t_1, x_1) - f(x_1) \leq \varphi_t(t_1, x_1) + |\nabla \varphi(t_1, x_1)|^2 - f(x_1). \quad (2.8.60)$$

We deduce that $\bar{u}(t, x)$ is a super-solution to (2.8.53). Moreover, at $t = t_0$ we have

$$\bar{u}(t_0, x) = u(t_0, x) \quad \text{for all } x \in \mathbb{T}^n.$$

As a consequence, it follows that $u(t, x) \leq \bar{u}(t, x)$ for all $t \geq t_0$ and $x \in \mathbb{T}^n$. Specifying this at $x_k \in \mathcal{Z}$ gives

$$u(t, x_k) \leq u(t_0, x_k) \quad \text{for all } t \geq t_0,$$

thus $u(t, x)$ is non-increasing in t on \mathcal{Z} , proving the first claim of Theorem 2.8.11: the limit

$$a_k := \lim_{t \rightarrow +\infty} u(t, x_k) \quad (2.8.61)$$

exists for all $x_k \in \mathcal{Z}$, $1 \leq k \leq N$.

Convergence on the whole torus

The proof of the second part of Theorem 2.8.11 is similar to that of Proposition 2.8.6 but some technical points are different. For a fixed $t > 0$ and $x \in \mathbb{T}^n$, consider the Lax-Oleinik formula written as

$$u(t, x) = \inf_{\gamma(t)=x} \left(u(s, \gamma(s)) + \int_s^t \left(\frac{|\dot{\gamma}(\sigma)|^2}{4} + f(\gamma(\sigma)) \right) d\sigma \right), \quad (2.8.62)$$

with any $0 \leq s \leq t$. Taking a test curve $\gamma_{s,t}(\sigma)$, $s \leq \sigma \leq t$ such that $\gamma_{s,t}(s) = x_k \in \mathcal{Z}$, with both s and t large, and passing to the limit $t, s \rightarrow +\infty$ with $t - s \rightarrow +\infty$, we deduce that for all $x \in \mathbb{T}^n$ we have

$$\limsup_{t \rightarrow +\infty} u(t, x) \leq \inf_{x_i \in \mathcal{Z}} (a_i + S(x_i, x)), \quad (2.8.63)$$

with $S(x_i, x)$ defined in (2.8.56). We used here the existence of the limit in (2.8.61).

The longer step is to show the reverse inequality to (2.8.63). Let $\gamma_t(s)$, $0 \leq s \leq t$, be a minimizer in the Lax-Oleinik formula

$$u(t, x) = \inf_{\gamma(t)=x} \left(u_0(\gamma(0)) + \int_0^t \left(\frac{|\dot{\gamma}(\sigma)|^2}{4} + f(\gamma(\sigma)) \right) d\sigma \right). \quad (2.8.64)$$

As $\gamma_t(s)$ is calibrated by u , we have

$$\begin{aligned} u(t, x) &= u(t + s, \gamma_t(t + s)) + \int_{t+s}^t \left(\frac{|\dot{\gamma}_t(\sigma)|^2}{4} + f(\gamma_t(\sigma)) \right) d\sigma \\ &= u(t + s, \gamma_t(t + s)) + \int_s^0 \left(\frac{|\dot{\gamma}_t(t + \sigma)|^2}{4} + f(\gamma_t(t + \sigma)) \right) d\sigma, \quad \text{for all } -t \leq s \leq 0. \end{aligned} \quad (2.8.65)$$

Let us introduce the path $\eta_t(\sigma) = \gamma_t(t + \sigma)$, $-t \leq \sigma \leq 0$, and write (2.8.65) as

$$u(t, x) = u(t + s, \eta_t(s)) + \int_s^0 \left(\frac{|\dot{\eta}_t(\sigma)|^2}{4} + f(\eta_t(\sigma)) \right) d\sigma, \quad \text{for all } -t \leq s \leq 0. \quad (2.8.66)$$

We now pass to the limit $t \rightarrow +\infty$. The uniform a priori bounds on $\gamma_t(\sigma)$ and $\dot{\gamma}_t(\sigma)$ imply the corresponding bounds on $\eta_t(\sigma)$ and $\dot{\eta}_t(\sigma)$. Hence, there exists a sequence $t_n \rightarrow +\infty$ such that $\eta_{t_n}(\sigma)$ converges as $n \rightarrow +\infty$, locally uniformly in σ , to a limit $\eta(\sigma)$, $-\infty < \sigma \leq 0$. In addition, $\eta(s)$ inherits the minimizing property of η_t : for any $s \leq 0$, the curve $\eta(\sigma)$ is a minimizer of

$$\int_s^0 \left(\frac{|\dot{\gamma}(\sigma)|^2}{4} + f(\gamma(\sigma)) \right) d\sigma,$$

over all curves $\gamma(\sigma)$, $s \leq \sigma \leq 0$, that connect the point $\gamma(s) = \eta(s)$ to $x = \gamma(0) = \eta(0)$.

By the same token, the bounds (2.8.57)

$$\|u(t, \cdot)\| \leq C, \quad \|u_t(t, \cdot)\|_{L^\infty} + \|\nabla u(t, \cdot)\|_{L^\infty} \leq C \quad (2.8.67)$$

on the function $u(t, x)$ imply that the sequence

$$v_n(s, x) = u(t_n + s, x),$$

possibly after extracting a subsequence, converges in $L^\infty(\mathbb{T}^n)$ and locally uniformly in s , to a limit $v(s, x)$ such that

$$v(s, x_k) = a_k, \quad \text{for all } x_k \in \mathcal{Z} \text{ and } s \in \mathbb{R}. \quad (2.8.68)$$

The uniformity of the limits of $\eta_t(\sigma)$ and $v_n(s, x)$ and the uniform in t Lipschitz bounds on $u(t, x)$ allow us to pass to the limit $t_n \rightarrow +\infty$ in (2.8.66), giving

$$\lim_{t_n \rightarrow +\infty} u(t_n, x) = v(0, x) = v(s, \eta(s)) + \int_s^0 \left(\frac{|\dot{\eta}(\sigma)|^2}{4} + f(\eta(\sigma)) \right) d\sigma, \quad \text{for all } -\infty \leq s \leq 0. \quad (2.8.69)$$

As in the proof of Proposition 2.8.6, we deduce from (2.8.69) the boundedness of the integral

$$\int_{-\infty}^0 \left(\frac{|\dot{\eta}(\sigma)|^2}{4} + f(\eta(\sigma)) \right) d\sigma < +\infty.$$

This, in turn, as in that proof, implies that there exists $x_j \in \mathcal{Z}$ such that

$$\lim_{s \rightarrow -\infty} \eta(s) = x_j.$$

Using (2.8.68) together with the uniform in s Lipschitz bound on $v(s, x)$ we may now pass to the limit $s \rightarrow -\infty$ in the right side of (2.8.69) to conclude that

$$\lim_{t_n \rightarrow +\infty} u(t_n, x) = a_j + \int_{-\infty}^0 \left(\frac{|\dot{\eta}(\sigma)|^2}{4} + f(\eta(\sigma)) \right) d\sigma. \quad (2.8.70)$$

The minimizing property of $\eta(\sigma)$ implies that

$$\lim_{t_n \rightarrow +\infty} u(t_n, x) = a_j + S(x_j, x) \geq \inf_{x_k \in \mathcal{Z}} (a_k + S(x_k, x)). \quad (2.8.71)$$

Comparing to (2.8.63), we see that

$$\lim_{t_n \rightarrow +\infty} u(t_n, x) = \inf_{x_k \in \mathcal{Z}} (a_k + S(x_k, x)) := u_\infty(x). \quad (2.8.72)$$

On the other hand, as we have seen before, $u_\infty(x)$ is a solution to

$$|\nabla u_\infty|^2 = f(x).$$

The weak contraction property for the viscosity solutions implies that not only we have the limit along a sequence $t_n \rightarrow +\infty$ but actually

$$\lim_{t \rightarrow +\infty} u(t, x) = u_\infty(x). \quad (2.8.73)$$

This finishes the proof. \square

Exercise 2.8.12 Explain how the weak contraction property is used in the very last step of the proof.

An equation with a drift

The minimizers for the problem

$$u_t + |\nabla u|^2 = f(x),$$

that we have just considered, spend most of their time near one of the finitely many points in the zero set \mathcal{Z} of f . To illustrate a different possible behavior of the minimizers, consider the Cauchy problem

$$\begin{aligned} u_t + cu_x + u_x^2 &= 0, \quad t > 0, \quad x \in \mathbb{T}^1, \\ u(0, x) &= u_0(x). \end{aligned} \quad (2.8.74)$$

The Lagrangian corresponding to the Hamiltonian

$$H(p) = |p|^2 + cp \quad (2.8.75)$$

is

$$L(v) = \sup_{p \in \mathbb{R}} [pv - cp - p^2] = \frac{(v - c)^2}{4}, \quad (2.8.76)$$

and the solution to (2.8.74) is given by the Lax-Oleinik formula:

$$u(t, x) = \inf_{\gamma(0)=y, \gamma(t)=x} \left[u_0(y) + \frac{1}{4} \int_0^t (\dot{\gamma}(s) - c)^2 ds \right]. \quad (2.8.77)$$

It is easy to see that the minimizer $\gamma_t(s; x)$ for (2.8.77) is a straight line $\gamma(s) = x + c_t(s - t)$. The optimal speed c_t is given by

$$c_t = \operatorname{argmin}_{v \in \mathbb{R}} \left[u_0(x - vt) + \frac{t}{4} (v - c)^2 \right]. \quad (2.8.78)$$

This is a very different behavior from that in Theorem 2.8.11: the minimizers visit every point on the torus infinitely many times. An immediate consequence of (2.8.77) is that

$$u(t, x) \geq \min_{y \in \mathbb{T}^n} u_0(y). \quad (2.8.79)$$

On the other hand, if x_0 is a minimum of $u_0(y)$, we can take

$$v = \frac{x - x_0 - [x - x_0 - ct]}{t} \quad (2.8.80)$$

in (2.8.77). Here, $[\xi]$ is the integer part of $\xi \in \mathbb{R}$. This gives

$$x - vt = x_0 + [x - x_0 - ct], \quad u_0(x - vt) = u_0(x_0), \quad (2.8.81)$$

leading to an upper bound

$$u(t, x) \leq u_0(x - vt) + \frac{t(v - c)^2}{4} \leq u_0(x_0) + \frac{1}{4t} = \min_{y \in \mathbb{T}^n} u_0(y) + \frac{1}{4t}. \quad (2.8.82)$$

We deduce from (2.8.79) and (2.8.82) that

$$\lim_{t \rightarrow +\infty} u(t, x) = \min_{y \in \mathbb{T}^n} u_0(y), \quad (2.8.83)$$

uniformly in $x \in \mathbb{T}^n$. Note that (2.8.83) holds even though the minimizers do not spend any more time near the minima of $u_0(y)$ than at any other points. Thus, the specific behavior of the minimizers we have seen in Theorem 2.8.11 is helpful but is not needed for the long time limit of the solution to exist. We will revisit this issue in a more general setting in Section 2.9.

2.8.3 The large time behavior: without the Lax-Oleinik formula

We now turn to the long time behavior of the solutions to the Cauchy problem (2.8.3):

$$\begin{aligned} u_t + R(x) \sqrt{1 + |\nabla u|^2} &= 0, \quad t > 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x). \end{aligned} \quad (2.8.84)$$

Let us recall that we assume that the function $R(x)$ is smooth and non-negative:

$$R(x) \geq R_0 > 0 \text{ for all } x \in \mathbb{T}^n, \quad (2.8.85)$$

and we use the notation

$$\bar{R} = \|R\|_{L^\infty}. \quad (2.8.86)$$

We will assume for simplicity that the set \mathcal{Z} where $R(x)$ attains its maximum is finite, though this assumption may be very much relaxed. As we have seen in the discussion following Proposition 2.8.1, this problem admits wave solutions of the form $ct + u_\infty(x)$, moving with the speed $c = \bar{R}$. Our goal will be to prove the following long time behavior result.

Theorem 2.8.13 *Let $u(t, x)$ be the solution to (2.8.84) with $u_0 \in C(\mathbb{T}^n)$ and assume that $R(x)$ is smooth, satisfies (2.8.85), and attains its maximum on a finite set. There is a solution $u_\infty(x)$ to*

$$R(x)\sqrt{1 + |\nabla u_\infty|^2} = \bar{R}, \quad x \in \mathbb{T}^n, \quad (2.8.87)$$

such that we have, uniformly with respect to $x \in \mathbb{T}^n$:

$$\lim_{t \rightarrow +\infty} (u(t, x) + t\bar{R} - u_\infty(x)) = 0, \quad (2.8.88)$$

with \bar{R} defined in (2.8.86).

Note that there is no claim of uniqueness of the solutions to (2.8.87) in Theorem 2.8.13, even up to addition of a constant. Indeed, as we have seen, uniqueness need not hold, as soon as the function $R(x)$ attains its maximum at more than one point. Unlike in the strictly convex case considered in the previous section, we will not use the Lax-Oleinik formula to understand the long time behavior, to illustrate the fact that the strict convexity of the Hamiltonian is also not needed for the solutions to have a long time limit. Nevertheless, the set

$$\mathcal{Z} = \{x \in \mathbb{T}^n : R(x) = \bar{R}\} \quad (2.8.89)$$

will play an important role in the proof, and in the dynamics, very similar to that of the minima of the function $f(x)$ in the proofs of Theorems 2.8.8 and 2.8.11.

We start the proof of Theorem 2.8.13 by writing

$$u(t, x) = v(t, x) + t\bar{R},$$

which transforms (2.8.84) into

$$\begin{aligned} v_t + R(x)\sqrt{1 + |\nabla v|^2} - \bar{R} &= 0, \quad x \in \mathbb{T}^n \\ v(0, x) &= u_0(x). \end{aligned} \quad (2.8.90)$$

Our goal is to show that there is a solution $u_\infty(x)$ to

$$R(x)\sqrt{1 + |\nabla u_\infty|^2} = \bar{R}, \quad x \in \mathbb{T}^n, \quad (2.8.91)$$

such that

$$\lim_{t \rightarrow +\infty} v(t, x) = u_\infty(x), \quad \text{uniformly in } x \in \mathbb{T}^n. \quad (2.8.92)$$

It is easy to see from the weak contraction principle that we may assume without loss of generality that the initial condition $u_0 \in C^1(\mathbb{T}^n)$. As a technical remark, we have seen that

the unique viscosity solution to (2.8.90) is uniformly bounded and uniformly Lipschitz: there exists $C > 0$ so that for all $t \geq 1$ we have

$$\|v(t, \cdot)\| \leq C, \quad \|v_t(t, \cdot)\|_{L^\infty} + \|\nabla v(t, \cdot)\|_{L^\infty} \leq C. \quad (2.8.93)$$

The uniform bound on v in (2.8.93) follows from the existence of a steady solution to (2.8.91) and the comparison principle, and the Lipschitz bound is an implication of Theorem 2.7.11. These bounds will be useful again when we pass to the limit $t \rightarrow +\infty$.

Note that if we can show that $v(t, x)$ converges uniformly, as $t \rightarrow +\infty$, to a limit $u_\infty(x)$, as in (2.8.92), then the limit is a viscosity solution to (2.8.91). Indeed, in that case the functions $v_n(t, x) = v(t+n, x)$ are solutions to (2.8.90), and converge, as $n \rightarrow +\infty$, to $u_\infty(x)$, in $L^\infty(\mathbb{T}^n)$, and locally uniformly in t . The stability property of the viscosity solutions implies that $u_\infty(x)$ is a steady solution to (2.8.90), and thus solves (2.8.91). Thus, it suffices to prove that the limit in (2.8.92) exists. We will do this in two steps: first we will prove existence of the limit for $x \in \mathcal{Z}$, and then show that convergence on \mathcal{Z} implies convergence on $\mathbb{T}^n \setminus \mathcal{Z}$ as well. In other words, what happens on \mathcal{Z} controls the behavior off \mathcal{Z} . This is very similar to the dynamics in Theorem 2.8.11 even though unlike in that case we will not use the Lax-Oleinik minimizers.

Convergence on \mathcal{Z}

To show convergence on \mathcal{Z} , we are going to prove that $v(t, x)$ is non-increasing in t on \mathcal{Z} . This is intuitively obvious: if $v(t, x)$ is continuously differentiable at $x \in \mathcal{Z}$ at some time $t > 0$, so that (2.8.90) holds in the classical sense, then, as $R(x) = \bar{R}$ for $x \in \mathcal{Z}$, we have

$$v_t(t, x) = \bar{R}(1 - \sqrt{1 + |\nabla v(t, x)|^2}) \leq 0,$$

so that $v(t, x)$ is non-increasing in t . The familiar problem is that $v(t, x)$ is merely Lipschitz, and not necessarily differentiable, hence (2.8.90) holds only almost everywhere, and we have no guarantee that it holds at any given (t, x) .

To make the above simple reasoning rigorous, the argument is close to the corresponding step in the proof of Theorem 2.8.11: consider $t_0 > 0$ and $x_0 \in \mathcal{Z}$ and set

$$\bar{v}(t, x) = v(t_0, x) + (t - t_0)(\bar{R} - R(x)).$$

We claim that \bar{v} is a super-solution to (2.8.90) on $[t_0, +\infty) \times \mathbb{T}^n$, such that

$$\bar{v}(t_0, x) = v(t_0, x) \text{ for all } x \in \mathbb{T}^n. \quad (2.8.94)$$

The latter follows immediately from the definition of $\bar{v}(t, x)$. To see the super-solution property, consider a test function $\varphi(t, x)$, and let $(t_1, x_1) \in [t_0, +\infty)$ be a minimum point for $\bar{v} - \varphi$. Since $\bar{v}(t, x)$ is smooth in t , we have

$$0 \leq \varphi_t(t_1, x_1) - \bar{v}_t(t_1, x_1) = \varphi_t(t_1, x_1) + R(x_1) - \bar{R}. \quad (2.8.95)$$

Hence, we have

$$\varphi_t(t_1, x_1) + R(x_1)\sqrt{1 + |\nabla \varphi(t_1, x_1)|^2} - \bar{R} \geq \varphi_t(t_1, x_1) + R(x_1) - \bar{R} \geq 0. \quad (2.8.96)$$

This proves the super-solution property of $\bar{v}(t, x)$. Together with (2.8.94), this implies

$$v(t, x) \leq \bar{v}(t, x) \text{ for } (t, x) \in [t_0, +\infty) \times \mathbb{T}^n. \quad (2.8.97)$$

As $R(x_0) = \bar{R}$ for $x_0 \in \mathcal{Z}$, we obtain

$$v(t, x_0) \leq v(t_0, x_0), \quad \text{for all } t \geq t_0 \text{ and } x_0 \in \mathcal{Z}. \quad (2.8.98)$$

Since t_0 is arbitrary, it follows that $v(t, x_0)$ is non-increasing in t . As a consequence, for each $x \in \mathcal{Z}$ the limit

$$u_\infty(x) = \lim_{t \rightarrow +\infty} v(t, x)$$

exists.

Exercise 2.8.14 Show that for any $\delta > 0$ we can find t_δ such that, for all $x \in \mathcal{Z}$, $h > 0$ and $t \geq t_\delta$ we have

$$0 \leq v(t, x) - v(t + h, x) \leq \delta. \quad (2.8.99)$$

Convergence outside of \mathcal{Z}

The heart of the proof is to show that convergence of $v(t, x)$ as $t \rightarrow +\infty$ on the set \mathcal{Z} forces the convergence off \mathcal{Z} as well, without the use of the Lax-Oleinik minimizers. Instead, the large time convergence of $v(t, x)$ outside of \mathcal{Z} will follow from the fact that a transform of v solves a Hamilton-Jacobi equation that is more complex than (2.8.90), but that has the merit of carrying an absorption term. We will use the Krushkov transform:

$$w(t, x) = -e^{-v(t, x)}. \quad (2.8.100)$$

Because of the L^∞ and gradient bounds for the Lipschitz function v , the function w is also Lipschitz and satisfies L^∞ and gradient bounds of the same type, and, in particular, we have

$$w_t = |w|v_t = -wv_t, \quad \nabla w = |w|\nabla v = -w\nabla v.$$

Moreover, because the function $v \mapsto -e^{-v}$ is increasing in v , the function w is a viscosity solution to

$$w_t + R(x)\sqrt{w^2 + |\nabla w|^2} = -\bar{R}w, \quad (2.8.101)$$

which can be written as

$$w_t + R(x)\frac{|\nabla w|^2}{|w| + \sqrt{w^2 + |\nabla w|^2}} + (\bar{R} - R(x))w = 0, \quad t > 0, \quad x \in \mathbb{T}^n. \quad (2.8.102)$$

The last term in the left side of (2.8.102) is the aforementioned absorption that will eventually save the day.

Exercise 2.8.15 Show that if $z(t, x)$ is a viscosity solution to

$$z_t + H(x, \nabla z) = 0,$$

and the function $G(z)$ is increasing, then $\zeta = G(z)$ is a viscosity solution to

$$\zeta_t + \frac{1}{Q'(\zeta)}H(x, Q'(\zeta)\nabla\zeta) = 0.$$

Here, $Q(\zeta)$ is the inverse function of $G(z)$. Is this necessarily true if the function G is not monotonic?

Let \mathcal{Z}_δ be the closed set of all points that are at distance at most $\delta > 0$ from \mathcal{Z} . Under our simplifying assumption that the set \mathcal{Z} is finite, the set \mathcal{Z}_δ is a finite union of closed balls. The uniform bounds on ∇v , together with the result of Exercise 2.8.14 imply that there is $C > 0$ so that

$$|w(t, x) - w(t + h, x)| \leq C\delta \text{ for } t \geq t_\delta \text{ and } x \in \mathcal{Z}_\delta. \quad (2.8.103)$$

Our task is now to extend this inequality outside of \mathcal{Z}_δ . Note that there is $\rho_\delta > 0$ such that

$$\bar{R} - R(x) \geq \rho_\delta \text{ for } x \text{ outside } \mathcal{Z}_\delta,$$

meaning that the pre-factor in the last term in the left side of (2.8.102) is uniformly positive off \mathcal{Z}_δ . Intuitively this means that the dynamics for w outside of \mathcal{Z}_δ is "uniformly absorbing". Let us set

$$\underline{w}_\delta(t, x) = w(t + h, x) - C\delta - \|w(t_\delta, \cdot)\|_{L^\infty} e^{-\rho_\delta(t-t_\delta)}, \quad t \geq t_\delta, \quad x \notin \mathcal{Z}_\delta. \quad (2.8.104)$$

To show that (2.8.103) holds outside of \mathcal{Z}_δ , we are going to prove that $\underline{w}(t, x)$ is a sub-solution to (2.8.102) for $t \geq t_\delta$, and $x \notin \mathcal{Z}_\delta$, and, in addition,

$$\underline{w}_\delta(t, x) \leq w(t, x) \text{ for } (t, x) \in [t_\delta, +\infty) \times \mathcal{Z}_\delta, \text{ and } t = t_\delta, \quad x \in \mathbb{T}^n. \quad (2.8.105)$$

This will imply $w(t, x) \geq \underline{w}_\delta(t, x)$ for $t \geq t_\delta$ and $x \notin \mathcal{Z}_\delta$, which, in turn, entails

$$w(t, x) \geq w(t + h, x) - C(\delta + e^{-\rho_\delta(t-t_\delta)}), \quad \text{for } t \geq t_\delta, \quad x \in \mathbb{T}^n \text{ and } h > 0. \quad (2.8.106)$$

Since $\rho_\delta > 0$ is positive, and $\delta > 0$ is arbitrary, this implies the pointwise convergence of $w(t, x)$ to a limit $w_\infty(x)$ as $t \rightarrow +\infty$, and, consequently, its uniform convergence that follows from the Lipschitz bound on $w(t, x)$. Therefore, the function $v(t, x)$ also converges to a limit

$$v_\infty(x) = -\log(-w_\infty(x)),$$

as $t \rightarrow +\infty$. Note that the absorbing nature of the dynamics for w exhibits itself in the fact that $\rho_\delta > 0$ outside of \mathcal{Z}_δ – this is why the Kruzhkov transform is helpful here.

Thus, to finish the proof of Theorem 2.8.13, we only need to show that \underline{w}_δ is a sub-solution to (2.8.102) for $t \geq t_\delta$, and $x \notin \mathcal{Z}_\delta$, and check that (2.8.105) holds. We see from (2.8.103) that

$$w(t, x) \geq \underline{w}_\delta(t, x) \text{ for } t \geq t_\delta \text{ and } x \in \mathcal{Z}_\delta. \quad (2.8.107)$$

At the time $t = t_\delta$ we have

$$w(t_\delta, x) - \underline{w}_\delta(t_\delta, x) = w(t_\delta, x) + \|w(t_\delta, \cdot)\|_{L^\infty} + C\delta - w(t_\delta + h, x) \geq C\delta > 0, \quad \text{for all } x \in \mathbb{T}^n. \quad (2.8.108)$$

We used here the fact that $w(t, x) \leq 0$ for all $t > 0$ and $x \in \mathbb{T}^n$. Putting (2.8.107) and (2.8.108) together, we conclude that (2.8.105) does hold.

It remains to check the sub-solution property for \underline{w}_δ , outside of \mathcal{Z}_δ . Let φ be a test function and (t_1, x_1) a minimum point of $\varphi - \underline{w}_\delta$, with $x_1 \notin \mathcal{Z}_\delta$. In other words, (t_1, x_1) is a minimum point of the function

$$(\varphi(t, x) + C\delta + \|w(t_\delta, \cdot)\|_{L^\infty} e^{-\rho_\delta(t-t_\delta)}) - w(t, x + h).$$

As w is a viscosity solution to (2.8.102):

$$w_t + R(x) \frac{|\nabla w|^2}{|w| + \sqrt{w^2 + |\nabla w|^2}} + (\bar{R} - R(x))w = 0, \quad t > 0, \quad x \in \mathbb{T}^n, \quad (2.8.109)$$

we deduce that the following inequality holds at (t_1, x_1) :

$$\varphi_t - \|w(t_\delta, \cdot)\|_{L^\infty} \rho_\delta e^{-\rho_\delta(t-t_\delta)} + R(x_1) \frac{|\nabla \varphi|^2}{|\tau_h w| + \sqrt{(\tau_h w)^2 + |\nabla \varphi|^2}} + (\bar{R} - R(x_1))\tau_h w \leq 0. \quad (2.8.110)$$

Here, we have set $\tau_h w(t, x) = w(t + h, x)$. The definition of \underline{w}_δ implies that

$$\tau_h w(t, x) \geq \underline{w}_\delta(t, x) \quad \text{for } t \geq t_\delta \text{ and } x \notin \mathcal{Z}_\delta,$$

so that

$$|\tau_h w(t, x)| \leq |\underline{w}_\delta(t, x)|. \quad (2.8.111)$$

Also, as

$$\bar{R} - R(x) \geq \rho_\delta \text{ for } x \notin \mathcal{Z}_\delta,$$

we have

$$\begin{aligned} & (\bar{R} - R(x_1))\tau_h w(t, x_1) - \|w(t_\delta, \cdot)\|_{L^\infty} \rho_\delta e^{-\rho_\delta(t-t_\delta)} \\ & \geq (\bar{R} - R(x_1))w(\tau + h, x_1) - (\bar{R} - R(x_1))\|w(t_\delta, \cdot)\|_{L^\infty} e^{-\rho_\delta(t-t_\delta)} \\ & = (\bar{R} - R(x_1))[\underline{w}_\delta(t, x_1) + C\delta] \geq (\bar{R} - R(x_1))\underline{w}_\delta(t, x_1). \end{aligned} \quad (2.8.112)$$

Using the inequalities (2.8.111) and (2.8.112) in (2.8.110) leads to

$$\varphi_t + R(x_1) \frac{|\nabla \varphi|^2}{|\underline{w}_\delta| + \sqrt{(\underline{w}_\delta)^2 + |\nabla \varphi|^2}} + (\bar{R} - R(x_1))\underline{w}_\delta \leq 0, \quad (2.8.113)$$

at (t_1, x_1) . This is the desired viscosity sub-solution inequality for \underline{w}_δ . Thus, $\underline{w}_\delta(t, x)$ is, indeed, a sub-solution to (2.8.102) for $t \geq t_\delta$, and $x \notin \mathcal{Z}_\delta$. This finishes the proof of Theorem 2.8.13. \square

Exercise 2.8.16 Carry out the same analysis for the equation

$$u_t + |\nabla u| = f(x), \quad t > 0, \quad x \in \mathbb{T}^n,$$

where $f \in C(\mathbb{T}^n)$ satisfies the usual assumptions of this section: continuous, nonnegative, with a nontrivial zero set.

2.9 Convergence of the Lax-Oleinik semigroup

In this section, we prove that the solutions to the Cauchy problem

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x) \end{aligned} \quad (2.9.1)$$

converge to a wave solution as $t \rightarrow +\infty$, under the assumption of uniform strict convexity of the Hamiltonian $H(x, p)$. So far, we have seen a very particular mechanism for convergence: the dynamics on a special set dictates in turn the convergence in the area where the equation is coercive. This was the zero set of the function $f(x)$ in Theorem 2.8.11, and the set \mathcal{Z} where the function $R(x)$ attains its maximum in the proof of Theorem 2.8.13. In both cases, the key is a monotonicity property on this "controlling" set. This is also the set where, in the strictly convex situation of Theorem 2.8.11 characteristics spend most of the time. On the other hand, in the example following Theorem 2.8.11, we have seen a situation where the minima of the initial condition $u_0(y)$ dictate the long time behavior, even though the minimizing curves do not spend any extra time near these points.

It turns out that the existence of such "controlling" set is a general fact. For a general Hamilton-Jacobi equation of the type (2.9.1), there is a set where the extremals associated to the wave solutions accumulate, and which orchestrates the convergence to a steady solution. The reader interested to learn more may consult [?] or [?], where their general theory by Fathi is exposed. Our goal here is much more modest: we want to identify a set where, following the ideas of the preceding section, the dynamics of u will dictate the behavior on the whole torus. The following theorem is due to Fathi [?] but we present an alternative proof inspired by [?].

Theorem 2.9.1 *Let $H(x, p)$ be smooth and uniformly strictly convex in p :*

$$\alpha I \leq D_p^2 H(x, p), \text{ in the sense of quadratic forms,}$$

and c be the corresponding wave speed: there exists a solution $u_\infty(x)$ to

$$H(x, \nabla u) = c, \quad x \in \mathbb{T}^n. \tag{2.9.2}$$

Then, for any given $u_0 \in C(\mathbb{T}^n)$, there exists a solution $u_\infty(x)$ to (2.9.2) such that the solution to the Cauchy problem

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x) \end{aligned} \tag{2.9.3}$$

converges to $u_\infty(x)$ as $t \rightarrow +\infty$:

$$\lim_{t \rightarrow +\infty} \|u(t, \cdot) + ct - u_\infty(\cdot)\|_{L^\infty} = 0. \tag{2.9.4}$$

We will assume throughout the proof, without loss of generality, that $c = 0$. Otherwise, we would simply replace the Hamiltonian $H(x, p)$ by $H(x, p) - c$.

As usual, existence of the steady solutions implies, via the comparison principle, that there exists $C_0 > 0$ such that

$$|u(t, x)| \leq C_0, \quad \text{for all } t > 0 \text{ and } x \in \mathbb{T}^n. \tag{2.9.5}$$

In addition, we have the uniform Lipschitz bound:

$$\text{Lip}_{t,x}[u] \leq C, \quad \text{for all } t \geq 1 \text{ and } x \in \mathbb{T}^n. \tag{2.9.6}$$

Here, $\text{Lip}_{t,x}[u]$ is the Lipschitz constant of u both in the t and x variables. These uniform bounds show that, at least along a sequence $t_n \rightarrow +\infty$, the uniform limits

$$u_\infty(t, x) = \lim_{t_n \rightarrow +\infty} u(t + t_n, x) \quad (2.9.7)$$

exist but may exist on the sequence t_n . Our goal is to show that there is actually a limit in (2.9.7) that does not depend on the sub-sequence, this limit is time-independent, and is a solution to the steady equation

$$H(x, \nabla u_\infty) = 0.$$

Before going directly into the proof of Theorem 2.9.1, we would like to explain the construction of the set \mathcal{Z} , and what sort of monotonicity we can use for the proof, as we did in the proofs of Theorem 2.8.11 and Theorem 2.8.13. This will require the notion of the ω -limit set of a solution, and that is where we will start the discussion. After introducing these objects and discussing their basic properties we will turn to the bona fide proof of Theorem 2.9.1 that will be quite short once we have obtained the desired properties of the set \mathcal{Z} .

The ω -limit set

The ω -limit set for a given initial condition $u_0 \in C(\mathbb{T}^n)$ with respect to the Lax-Oleinik semi-group is denoted by $\omega(u_0) \subset C(\mathbb{T}^n)$, and is constructed as follows. The uniform bounds (2.9.5) and (2.9.6) imply that there is a sequence $t_n \rightarrow +\infty$ such that the family $v_n(t, x) = u(t + t_n, x)$ converges :

$$v_n(t, x) = u(t + t_n, x) \rightarrow v(t, x), \quad (2.9.8)$$

in $L^\infty(\mathbb{T}^n)$ and uniformly on compact intervals of $t \in \mathbb{R}$. Here and below, $u(t, x)$ is the solution to (2.9.3). The function $v(t, x)$ is a solution to

$$v_t + H(x, \nabla v) = 0, \quad (2.9.9)$$

defined for all $t \in \mathbb{R}$, and not just for $t > 0$. Sometimes such solutions are called "entire in time", to indicate that they are also defined for negative times. The set $\omega(u_0)$ consists of all "one-time" snapshots of the limits: the functions $\psi \in C(\mathbb{T}^n)$ such that there is a sequence $t_n \rightarrow +\infty$ and the corresponding limit $v(t, x)$ so that (2.9.8) holds, and

$$\psi(x) = v(0, x) = \lim_{n \rightarrow +\infty} u(t_n, x). \quad (2.9.10)$$

It will be convenient for us to consider, in addition, the set $\tilde{\omega}(u_0) \subset C(\mathbb{R} \times \mathbb{T}^n)$ of all functions $v \in C(\mathbb{R} \times \mathbb{T}^n)$ that can be obtained by the limiting procedure in (2.9.8). A simple observation is that if $v \in \tilde{\omega}(u_0)$ then for any $s \in \mathbb{R}$ fixed, the function $v(s, \cdot)$ is in $\omega(u_0)$. This can be seen simply by taking the sequence $t'_n = t_n + s$.

The claim of Theorem 2.9.1 can be now reformulated as saying that each $v \in \tilde{\omega}(u_0)$ does not depend on t , and that $\omega(u_0)$ contains exactly one function ψ . The following exercise gives a sufficient condition for this to be true.

Exercise 2.9.2 (i) Assume that $v \in \tilde{\omega}(u_0)$ does not depend on t . Show that then $v(x)$ is a viscosity solution to

$$H(x, \nabla v) = 0. \quad (2.9.11)$$

(ii) Show that if there exists $v \in \tilde{\omega}(u_0)$ that does not depend on t , then the limit

$$\lim_{t \rightarrow +\infty} u(t, x),$$

exists, is unique, and is a viscosity solution to (2.9.11). Hint: use the contraction property for the solutions to (2.9.3).

Exercise 2.9.2 gives us a blueprint for the proof of Theorem 2.9.1: it suffices to show that any $v \in \tilde{\omega}(u_0)$ does not depend on t . As in the proof of Theorem 2.8.13, we will first identify a set \mathcal{Z} where it is easier to show that $v(t, x)$ is time-independent, and then show this outside of \mathcal{Z} .

Monotonicity along the minimizers

Our first goal is to recycle the main idea of the proofs of Theorems 2.8.11 and 2.8.13, namely, to find a set where convergence will hold because of some monotonicity property. The following easy and very general remark can be made: let $\phi(x)$ be a steady solution to

$$H(x, \nabla \phi) = 0, \tag{2.9.12}$$

and $\gamma : [0, t] \rightarrow \mathbb{T}^n$ be an extremal path calibrated by ϕ . For all $0 \leq s \leq s' \leq t$ we have

$$\phi(\gamma(s')) = \phi(\gamma(s)) + \int_s^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma, \tag{2.9.13}$$

whereas, by the definition of the Lax-Oleinik semigroup we have

$$u(s', \gamma(s')) \leq u(s, \gamma(s)) + \int_s^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma. \tag{2.9.14}$$

Subtracting, we obtain

$$u(s', \gamma(s')) - \phi(\gamma(s')) \leq u(s, \gamma(s)) - \phi(\gamma(s)) \quad \text{if } s \leq s'. \tag{2.9.15}$$

Thus, the difference $u(s, \gamma(s)) - \phi(\gamma(s))$ is non-increasing in s along any extremal path calibrated by ϕ . This simple observation will bear a lot of fruit.

The ω -limits of paths and the set \mathcal{Z}

We now use the monotonicity property (2.9.15) to construct a candidate for the set \mathcal{Z} . It will contain paths that calibrate all steady solutions but it will also do more. Let us fix a steady solution ϕ . We define \mathcal{Z}_ϕ as the collection of all "eternal" extremal paths calibrated by ϕ , which is the set of all trajectories $\gamma : \mathbb{R} \rightarrow \mathbb{T}^n$ such that

$$\phi(\gamma(s')) = \phi(\gamma(s)) + \int_s^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma, \quad \text{for all } -\infty < s \leq s' < +\infty. \tag{2.9.16}$$

The next exercise shows that eternal extremal paths calibrated by ϕ exist, so the set \mathcal{Z}_ϕ is not empty.

Exercise 2.9.3 (i) Fix $x \in \mathbb{T}^n$ and consider a family of extremal paths $\gamma_t(s)$, $t \leq s \leq 0$, calibrated by ϕ . Show that there is a sequence $t_n \rightarrow -\infty$ such that $\gamma_{t_n}(s) \rightarrow \gamma(s)$, locally uniformly in $s \leq 0$, and $\gamma(0) = x$.

(ii) Let $\gamma(s)$ be constructed as in part (i). Show that there exists a sequence $s_n \rightarrow -\infty$ such that the paths $\gamma_n(s) = \gamma(s_n + s)$, $-\infty \leq s \leq -s_n$, converge, locally uniformly on \mathbb{R} , to a path $\bar{\gamma}(s)$, $s \in \mathbb{R}$.

(iii) Show that the path $\bar{\gamma}(s)$, $s \in \mathbb{R}$, is calibrated by ϕ .

The set \mathcal{Z} is then defined as follows: a point $x \in \mathbb{T}^n$ is in \mathcal{Z} if there is a path $\gamma_\infty : \mathbb{R} \rightarrow \mathbb{T}^n$ that passes through x , a path $\gamma \in \mathcal{Z}_\phi$, and a sequence $s_n \rightarrow +\infty$ such that

$$\gamma_\infty(s) = \lim_{n \rightarrow +\infty} \gamma(s + s_n), \quad (2.9.17)$$

with the limit uniform on every bounded interval of \mathbb{R} . In other words, \mathcal{Z} is the union of ω -limits of the paths in \mathcal{Z}_ϕ .

Exercise 2.9.4 (i) Find the set \mathcal{Z} for the Hamiltonian $H(x, p) = |p|^2 - f(x)$ with a smooth non-negative function $f(x)$, $x \in \mathbb{T}^n$. Does it depend on the steady solution $\phi(x)$ with which you start?

(ii) Find the set \mathcal{Z} for the Hamiltonian $H(p) = |p|^2 + cp$, with $c > 0$.

The following exercise will be important when we discuss the time monotonicity of the eternal solutions on \mathcal{Z} in the proof of Theorem 2.9.1.

Exercise 2.9.5 Show that if $\gamma_\infty(\sigma)$, $\sigma \in \mathbb{R}$ is in \mathcal{Z} , then so is the time-shifted path

$$\gamma_\infty^{(s)}(\sigma) = \gamma_\infty(\sigma + s), \quad \sigma \in \mathbb{R},$$

for any $s \in \mathbb{R}$ fixed. Hint: this is because the original solution ϕ , that we used to construct \mathcal{Z}_ϕ and then \mathcal{Z} , is time-independent, so that if a path $\gamma \in \mathcal{Z}_\phi$ calibrates ϕ , then any time-shifted path $\gamma_s(s') = \gamma(s + s')$ also calibrates ϕ , and is therefore in \mathcal{Z}_ϕ .

Calibration by paths in \mathcal{Z}

Let us now take a path $\gamma_\infty(s)$, $s \in \mathbb{R}$ that lies in \mathcal{Z} , obtained as the limit in (2.9.17), with a given $\gamma \in \mathcal{Z}_\phi$. Writing

$$\begin{aligned} \phi(\gamma(s' + s_n)) &= \phi(\gamma(s - s_n)) + \int_{s+s_n}^{s'+s_n} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma \\ &= \phi(\gamma(s + s_n)) + \int_s^{s'} L(\gamma(\sigma + s_n), \dot{\gamma}(\sigma + s_n)) d\sigma, \end{aligned} \quad (2.9.18)$$

and passing to the limit $s_n \rightarrow +\infty$, we see immediately that $\gamma_\infty(s)$ is also calibrated by ϕ . The miracle is that $\gamma_\infty(s)$ is also calibrated by every other steady solution $\psi(x)$ to (2.9.12). Indeed, it follows from the monotonicity property (2.9.15) used with $u(t, x) = \psi(x)$ that for any path $\gamma \in \mathcal{Z}_\phi$ the limit

$$\lim_{s \rightarrow -\infty} [\psi(\gamma(s)) - \phi(\gamma(s))] = K(\gamma),$$

exists. It follows that if $\gamma(s + s_n) \rightarrow \gamma_\infty(s)$ as $s_n \rightarrow +\infty$, then the two solutions differ by a constant on γ_∞ :

$$\psi(\gamma_\infty(s)) = \phi(\gamma_\infty(s)) + K(\gamma), \quad \text{for all } -\infty < s < +\infty. \quad (2.9.19)$$

As $\gamma_\infty(s)$ is calibrated by ϕ , we conclude from (2.9.19) that it is calibrated by ψ as well. We stress that it is not true that every path in \mathcal{Z}_ϕ is calibrated by any other steady solution – this is only true for their ω -limits that form the set \mathcal{Z} . The following proposition shows that we can say even more.

Proposition 2.9.6 *Let $\gamma_\infty(s)$, $s \in \mathbb{R}$ be a path in \mathcal{Z} and $v \in \tilde{\omega}(u_0)$. There exists $K(\gamma) \in \mathbb{R}$ such that*

$$v(t, \gamma_\infty(t)) - \phi(\gamma_\infty(t)) = K(\gamma), \quad \text{for all } t \in \mathbb{R}. \quad (2.9.20)$$

In particular, the path $\gamma_\infty(s)$, $s \in \mathbb{R}$, is calibrated by $v(t, x)$.

Proof. By definition of γ_∞ , there is a global extremal path γ calibrated by ϕ , and a sequence $s_n \rightarrow +\infty$ such that

$$\gamma_\infty(\sigma) = \lim_{n \rightarrow +\infty} \gamma(\sigma + s_n),$$

uniformly in every compact in $\sigma \in \mathbb{R}$. Observe that to prove (2.9.20) it suffices to find a subsequence s_{n_k} such that

$$\lim_{k \rightarrow +\infty} [v(s, \gamma(s + s_{n_k})) - \phi(\gamma(s + s_{n_k}))] = \text{const}, \quad \text{independent of } s \in \mathbb{R}. \quad (2.9.21)$$

Let $u(t, x)$ be the solution to (2.9.3).

$$\begin{aligned} u_t + H(x, \nabla u) &= 0, \quad t > 0, \quad x \in \mathbb{T}^n, \\ u(0, x) &= u_0(x). \end{aligned} \quad (2.9.22)$$

Since $v \in \tilde{\omega}(u_0)$, we may also find a sequence $t_n \rightarrow +\infty$ such that for all $s \in \mathbb{R}$ we have

$$v(s, x) = \lim_{n \rightarrow +\infty} u(t_n + s_n + s, x), \quad \text{in } L^\infty(\mathbb{T}^n), \quad (2.9.23)$$

locally uniformly in $s \in \mathbb{R}$. Thus, for every compact set $K \subset \mathbb{R}$ and $\varepsilon > 0$ there exists $N_{\varepsilon, K}$ such that for all $n > N_{\varepsilon, K}$ we have

$$|v(s, \gamma(s + s_n)) - u(t_n + s_n + s, \gamma(s + s_n))| < \varepsilon \quad \text{for all } s \in K. \quad (2.9.24)$$

Hence, (2.9.21) would follow if we can show that, along a sub-sequence $n_k \rightarrow +\infty$, we have

$$\lim_{k \rightarrow +\infty} [u(t_{n_k} + s_{n_k} + s, \gamma(s + s_{n_k})) - \phi(\gamma(s + s_{n_k}))] = \text{const}, \quad \text{independent of } s \in \mathbb{R}. \quad (2.9.25)$$

Note that there is a sub-sequence $n_k \rightarrow +\infty$ such that the limit

$$\tilde{v}(t, \cdot) = \lim_{k \rightarrow +\infty} u(t + t_{n_k}, \cdot), \quad \text{in } L^\infty(\mathbb{T}^n) \quad (2.9.26)$$

exists, locally uniformly in t , and set

$$\psi_1(\cdot) = \tilde{v}(0, \cdot) = \lim_{k \rightarrow +\infty} u(t_{n_k}, \cdot), \quad \text{in } L^\infty(\mathbb{T}^n). \quad (2.9.27)$$

By the weak contraction property, we have

$$\|u(t_{n_k} + t, \cdot) - \tilde{v}(t, \cdot)\|_{L^\infty} \leq \|u(t_{n_k}, \cdot) - \psi_1\|_{L^\infty} \rightarrow 0, \quad \text{as } k \rightarrow +\infty, \quad (2.9.28)$$

and, very importantly, this limit is uniform in $t \geq 0$. Thus, for any $\varepsilon > 0$ we can find N_ε so that for all $k > N_\varepsilon$ we have

$$|u(t_{n_k} + s_{n_k} + s, \gamma(s_{n_k} + s)) - \tilde{v}(s_{n_k} + s, \gamma(s_{n_k} + s))| \leq \varepsilon, \quad (2.9.29)$$

locally uniformly in $s \in \mathbb{R}$. Hence, (2.9.25) would follow if we show that

$$\lim_{k \rightarrow +\infty} [\tilde{v}(s_{n_k} + s, \gamma(s_{n_k} + s)) - \phi(\gamma(s_{n_k} + s))] = \text{const}, \quad \text{independent of } s \in \mathbb{R}. \quad (2.9.30)$$

However, the monotonicity property (2.9.15) along the extremals implies that the limit

$$\ell = \lim_{\xi \rightarrow +\infty} (\tilde{v}(\xi + s, \gamma(\xi + s)) - \phi(\gamma(\xi + s))).$$

exists and is independent of s . It follows that

$$\lim_{k \rightarrow +\infty} [\tilde{v}(s_{n_k} + s, \gamma(s_{n_k} + s)) - \phi(\gamma(s_{n_k} + s))] = \ell \quad \text{for all } s \in \mathbb{R}, \quad (2.9.31)$$

finishing the proof. \square

We will also need the following proposition which says that paths calibrated by solutions in $\tilde{\omega}(u_0)$ come arbitrarily close to the set \mathcal{Z} – this is what eventually leads to the fact that the behavior of the solutions on \mathcal{Z} controls the behavior outside of \mathcal{Z} as well.

Proposition 2.9.7 *Let $v \in \tilde{\omega}(u_0)$. Given any path $\gamma(s)$, $s \in \mathbb{R}$, that is calibrated by v , there exists a sequence $t_n \rightarrow +\infty$ such that $\text{dist}(\gamma(t_n), \mathcal{Z}) \rightarrow 0$ as $n \rightarrow +\infty$.*

Proof. Let $\gamma(s)$, $s \in \mathbb{R}$, be a path calibrated by $v(t, x)$, and $\phi(x)$ be the steady solution used to generate \mathcal{Z} . The Lax-Oleinik formula tells us that for any $s < s' \leq t$ we have

$$v(s', \gamma(s')) = v(s, \gamma(s)) + \int_s^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma,$$

and

$$\phi(\gamma(s')) \leq \phi(\gamma(s)) + \int_s^{s'} L(\gamma(\sigma), \dot{\gamma}(\sigma)) d\sigma.$$

Subtracting we get the monotonicity relation

$$v(s', \gamma(s')) - \phi(\gamma(s')) \geq v(s, \gamma(s)) - \phi(\gamma(s)), \quad \text{for all } s < s'. \quad (2.9.32)$$

As both $v(s, x)$ and $\phi(x)$ are uniformly bounded, it follows that the limit

$$\lim_{s \rightarrow +\infty} [v(s, \gamma(s)) - \phi(\gamma(s))] \quad (2.9.33)$$

exists. In addition, the uniform bounds on $\gamma(s)$ and $\dot{\gamma}(s)$ imply that there exists a sequence $s_n \rightarrow +\infty$ so that the paths $\gamma_n(s) = \gamma(s + s_n)$ converge, as $n \rightarrow +\infty$, to a limiting path $\gamma_\infty(s)$, $s \in \mathbb{R}$, locally uniformly in s .

Exercise 2.9.8 Show that the path $\gamma_\infty(s)$ is calibrated by any steady solution, in particular, by ϕ .

Exercise 2.9.8 shows that $\gamma_\infty(s)$, $s \in \mathbb{R}$, lies in \mathcal{Z}_ϕ but we do not yet know that the path $\gamma_\infty(s)$, $s \in \mathbb{R}$, is in \mathcal{Z} . Since $\gamma_\infty(s)$ inherits the uniform bounds obeyed by $\gamma(s)$ and $\dot{\gamma}(s)$, we can find a sequence $s'_n \rightarrow +\infty$ such that

$$\gamma_\infty^{(n)}(s) := \gamma_\infty(s + s'_n) \rightarrow \bar{\gamma}_\infty(s),$$

locally uniformly in $s \in \mathbb{R}$. As all $\gamma_\infty^{(n_k)}(s)$ are calibrated by ϕ , we know that the limiting path $\bar{\gamma}_\infty(s)$, $s \in \mathbb{R}$, lies in \mathcal{Z} , by the definition of the set \mathcal{Z} .

To finish the proof of the proposition, consider the points $\gamma(s_n + s'_m)$. First, we fix m and choose $n = N_m$ sufficiently large, so that

$$|\gamma(s_{N_m} + s'_m) - \gamma_\infty(s'_m)| < \frac{\varepsilon}{2}.$$

Next, we choose m sufficiently large, so that $|\gamma_\infty(s'_m) - \bar{\gamma}_\infty(0)| < \varepsilon/2$. This gives

$$|\gamma(s_{N_m} + s'_m) - \bar{\gamma}_\infty(0)| < \varepsilon.$$

Since $\bar{\gamma}_\infty(0)$ is in \mathcal{Z} , the proof is complete. \square

Convergence on $\bar{\mathcal{Z}}$

After setting up the required objects, we turn to the proof of Theorem 2.9.1. The strategy comes from Exercise 2.9.2: we need to show that any solution to

$$v_t + H(x, \nabla v) = 0, \quad t \in \mathbb{R}, \quad x \in \mathbb{T}^n, \quad (2.9.34)$$

such that $v \in \tilde{\omega}(u_0)$, is time-independent. The reader has certainly guessed what will happen: the set \mathcal{Z} will play the same role as the zero set of the function $f(x)$ in Theorem 2.8.11, and the set where the function $R(x)$ attains its maximum in the proof of Theorem 2.8.13. This is confirmed by the following proposition, showing that such $v(t, x)$ is independent of $t \in \mathbb{R}$ on the closure $\bar{\mathcal{Z}}$ of the set \mathcal{Z} , though we do not know yet that this happens everywhere.

Proposition 2.9.9 *Any $v \in \tilde{\omega}_0(u_0)$ does not depend on $t \in \mathbb{R}$ for all $x \in \bar{\mathcal{Z}}$.*

Proof. Consider an eternal extremal path $\gamma_\infty(s)$, $s \in \mathbb{R}$, in \mathcal{Z} . We are going to show that

$$\partial_t v(t, \gamma_\infty(t)) = 0, \quad \text{for all } t \in \mathbb{R}. \quad (2.9.35)$$

We have shown in Proposition 2.9.6 that $\gamma_\infty(s)$, $s \in \mathbb{R}$, is calibrated both by ϕ and by v . It follows then from Corollary 2.7.21 that both ϕ and v are $C^{1,1}$ on γ_∞ , and we have

$$\nabla v(t, \gamma_\infty(t)) = \nabla_v L(\gamma_\infty(t), \dot{\gamma}_\infty(t)), \quad \nabla \phi(\gamma_\infty(t)) = \nabla_v L(\gamma_\infty(t), \dot{\gamma}_\infty(t)),$$

for all $t \in \mathbb{R}$, as in (2.7.48) in Exercise 2.7.13. This gives

$$\nabla v(t, \gamma_\infty(t)) = \nabla \phi(\gamma_\infty(t)), \quad \text{for all } t \in \mathbb{R}. \quad (2.9.36)$$

This relation holds in the classical sense, as both v and ϕ are $C^{1,1}$ on $\gamma_\infty(t)$. Since ϕ is a solution to the steady equation (2.9.12):

$$H(x, \nabla\phi) = 0, \tag{2.9.37}$$

we deduce that

$$H(\gamma_\infty(t), \nabla v(t, \gamma_\infty(t))) = 0.$$

As v is $C^{1,1}$ at $(t, \gamma_\infty(t))$, this entails (2.9.35):

$$\partial_t v(t, \gamma_\infty(t)) = 0, \quad \text{for all } t \in \mathbb{R}. \tag{2.9.38}$$

Consider $x \in \mathcal{Z}$ and an eternal extremal path $\gamma_\infty(\sigma)$, $\sigma \in \mathbb{R}$, in \mathcal{Z} that passes through x , so that $\gamma_\infty(t) = x$, with some $t \in \mathbb{R}$. Given any $s \in \mathbb{R}$, Exercise 2.9.5 allows us to use (2.9.38) with the shifted path

$$\gamma_\infty^{(s)}(\sigma) = \gamma_\infty(\sigma + t - s).$$

Note that

$$x = \gamma_\infty(t) = \gamma_\infty^{(s)}(s),$$

and (2.9.38) applied to $\gamma_\infty^{(s)}(\sigma)$ at $\sigma = s$ gives

$$0 = \partial_t v(s, \gamma_\infty^{(s)}(s)) = \partial_t v(s, \gamma_\infty(t)) = \partial_t v(s, x). \tag{2.9.39}$$

Since s is arbitrary, we conclude that $v(t, x)$ does not depend on t for all $x \in \mathcal{Z}$. The continuity of $v(t, x)$ implies that the same is true for $\overline{\mathcal{Z}}$ as well. \square

Convergence away from \mathcal{Z}

To finish the proof of Theorem 2.9.1 we now show that the claim of Proposition 2.9.9 holds also outside of $\overline{\mathcal{Z}}$. This step relies crucially on Proposition 2.9.7.

Proposition 2.9.10 *Any $v \in \tilde{\omega}_0(u_0)$ does not depend on $t \in \mathbb{R}$ for all $x \in \mathbb{T}^n$.*

Proof. Proposition 2.9.9 shows that we only need to consider $x \notin \overline{\mathcal{Z}}$. Using the by now familiar arguments based on the uniform bounds on minimizers, for any $x \in \mathbb{T}^n$ and $t > 0$ fixed, we can find a path $\gamma_t(\sigma)$, $\sigma \leq t$, calibrated by v , so that for any $s < t$ we have

$$v(t, x) = v(s, \gamma_t(s)) + \int_s^t L(\gamma_t(\sigma), \dot{\gamma}_t(\sigma)) d\sigma.$$

Our first goal is to show that $\gamma_t(s)$ has visited positions $x \in \mathbb{T}^n$ that are arbitrarily close to the set \mathcal{Z} at some times $s \leq t$. As usual, we know that there exists a sequence $s_n \rightarrow +\infty$ so that the time-shifted paths $\gamma_n(s) = \gamma_t(s + t - s_n)$, $s \leq s_n$, converge, as $n \rightarrow +\infty$, to a limiting path $\gamma_\infty(s)$, $s \in \mathbb{R}$. The path $\gamma_\infty(s)$ is also calibrated by $v(t, x)$, hence Proposition 2.9.7 shows that for any $\varepsilon > 0$ there exists $t_\varepsilon > 0$ so that

$$\text{dist}(\gamma_\infty(t_\varepsilon), \mathcal{Z}) < \frac{\varepsilon}{2}.$$

Let us fix $\varepsilon > 0$ and the corresponding t_ε , and choose $n > N_\varepsilon$ sufficiently large, so that

$$|\gamma_t(t_\varepsilon + t - s_n) - \gamma_\infty(t_\varepsilon)| = |\gamma_n(t_\varepsilon) - \gamma_\infty(t_\varepsilon)| < \frac{\varepsilon}{2}.$$

We deduce that

$$\text{dist}(\gamma_t(t_\varepsilon + t - s_n), \mathcal{Z}) < \varepsilon.$$

Therefore, we can find a sequence $\tau_n \rightarrow +\infty$ such that

$$\text{dist}(\gamma_t(-\tau_n), \mathcal{Z}) \rightarrow 0, \quad \text{as } n \rightarrow +\infty.$$

That is, as we have claimed, the path $\gamma_t(s)$ did get arbitrarily close to \mathcal{Z} in the past. Hence, there exists $z \in \overline{\mathcal{Z}}$ and a subsequence $\tau_{n_k} \rightarrow +\infty$ such that

$$\gamma_t(-\tau_{n_k}) \rightarrow z \quad \text{as } k \rightarrow +\infty.$$

Since the function $v(t, x)$ is Lipschitz in x , uniformly in $t \in \mathbb{R}$ and in $x \in \mathbb{T}^n$, we know that

$$\Delta_k(t, x) = v(-\tau_{n_k}, \gamma_t(-\tau_{n_k})) - v(-\tau_{n_k}, z) \rightarrow 0 \quad \text{as } k \rightarrow +\infty. \quad (2.9.40)$$

Let us then write

$$v(t, x) = v(-\tau_{n_k}, z) + \int_{-\tau_{n_k}}^t L(\gamma_t(\sigma), \dot{\gamma}_t(\sigma)) d\sigma + \Delta_k(t, x). \quad (2.9.41)$$

As $z \in \overline{\mathcal{Z}}$, by Proposition 2.9.9 we know that $v(s, z)$ does not depend on s , so that for any $s \in \mathbb{R}$ we have

$$v(-\tau_{n_k}, z) = v(-\tau_{n_k} - t + s, z),$$

and (2.9.41) can be written as

$$\begin{aligned} v(t, x) &= v(-\tau_{n_k} - t + s, z) + \int_{-\tau_{n_k}}^t L(\gamma_t(\sigma), \dot{\gamma}_t(\sigma)) d\sigma + \Delta_k(t, x) \\ &= v(-\tau_{n_k} - t + s, z) + \int_{-\tau_{n_k} - t + s}^s L(\gamma_t(\sigma + t - s), \dot{\gamma}_t(\sigma + t - s)) d\sigma + \Delta_k(t, x) \\ &= v(-\tau_{n_k} - t + s, z) + \int_{-\tau_{n_k} - t + s}^s L(\tilde{\gamma}_t(\sigma), \tilde{\dot{\gamma}}_t(\sigma)) d\sigma + \Delta_k(t, x), \end{aligned} \quad (2.9.42)$$

with the path

$$\tilde{\gamma}_t(\sigma) = \gamma_t(\sigma + t - s), \quad \sigma \leq s.$$

Note that $\tilde{\gamma}_t(s) = \gamma_t(t) = x$, and

$$\tilde{\gamma}_t(-\tau_{n_k} - t + s) = \gamma_t(-\tau_{n_k}) = z + \delta_k, \quad \delta_k \rightarrow 0 \quad \text{as } k \rightarrow +\infty. \quad (2.9.43)$$

Therefore, the Lax-Oleinik formula, together with the uniform Lipschitz property of the function $v(t, x)$, tells us that

$$\begin{aligned} v(s, x) &\leq v(-\tau_{n_k} - t + s, \tilde{\gamma}_t(-\tau_{n_k} - t + s)) + \int_{-\tau_{n_k} - t + s}^s L(\tilde{\gamma}_t(\sigma), \tilde{\dot{\gamma}}_t(\sigma)) d\sigma \\ &\leq v(-\tau_{n_k} - t + s, z) + \int_{-\tau_{n_k} - t + s}^s L(\tilde{\gamma}_t(\sigma), \tilde{\dot{\gamma}}_t(\sigma)) d\sigma + \Delta_k(t, x) + C\delta_k. \end{aligned} \quad (2.9.44)$$

Comparing to (2.9.42) and passing to the limit $k \rightarrow +\infty$ we conclude that

$$v(s, x) \leq v(t, x) \quad \text{for all } t, s \in \mathbb{R}. \quad (2.9.45)$$

As t and s are arbitrary, it follows that $v(t, x)$ does not depend on t , finishing the proof of Proposition 2.9.10, and thus that of Theorem 2.9.1 as well. \square

Our tour of the Hamilton-Jacobi equations ends here. One could say much more on the organization of the steady solutions, and the reader should consult [?]. They would be, however, outside the scope of this book. Let us just notice that the results of the present section provide a complete parallel with the large time behavior of the solutions to viscous Hamilton-Jacobi equations, which was the goal we wanted to achieve: the viscosity solutions of the inviscid problem still converge to waves, although their organization, that we have largely uncovered, is much more complicated.

Chapter 3

Adaptive dynamics models in biology

We mostly follow the material from the books by Benoit Perthame, with some research articles in the latter sections of this chapter. The plan for this chapter is as follows. We will start with a simple ODE system that models evolution of a population with various traits that correspond to different birth rates and compete for common resources. We will show, quite easily that the trait with the highest birth rate dominates in the long time limit. Next, we will allow both the birth and death rates to depend on the trait, and show the selection mechanism in that case. Mutations will be introduced next: the off-spring may have a different trait than the parent. This leads to a nonlinear integral equation for the population density. Existence of its solutions will be quite straightforward. In the limit of small mutations, we will obtain a Hamilton-Jacobi equation with a constraint. The zero set of the solution corresponds to the traits that survive in the long time limit. If such point is unique, the constrained Hamilton-Jacobi equation describes the evolution of a monomorphic population. The last step will be to investigate the Hamilton-Jacobi equations with such constraint, bringing us to the results obtained in the last 5-7 years.

3.1 Adaptive dynamics

Adaptive dynamics studies two biological effects: (i) the selection principle, which favors population with the best adapted trait, and (ii) mutations which allow the off-spring to have a slightly different trait from the parent. Here, we look at simple ODE models and study how the selection principle arises as the long time limit of small mutations.

A simple of example of a structured population and the selection principle

We begin with a very simple example of the logistic ordinary differential equation

$$\frac{dn}{dt} = bn - n^2,$$

modified to take into account the traits. We assumed that the population is structured by a trait $x \in \mathbb{R}$ and assume that (i) the reproduction rate b depends only on the trait, and (ii)

the death rate depends on the total population. This is reasonable because of the competition between various traits. The resulting system of ODEs is

$$\frac{\partial n(t, x)}{\partial t} = b(x)n(t, x) - \rho(t)n(t, x), \quad (3.1.1)$$

where $\rho(t)$ is the total population:

$$\rho(t) = \int_{\mathbb{R}} n(t, x) dx.$$

The ODEs in (3.1.1) for various $x \in \mathbb{R}$ are coupled via the term $\rho(t)$ that depends on $n(t, x)$ at all $x \in \mathbb{R}$.

We assume that the initial condition for (3.1.1) is $n(0, x) = n_0(x)$ such that $n_0(x) > 0$ for all $x \in (x_m, x_M)$, and $n_0(x) = 0$ otherwise. It follows that $n(t, x) = 0$ for all x outside of (x_m, x_M) for $t > 0$ as well. Note that there are no mutations in this model, and any state of the form

$$n(t, x) = b(y)\delta(x - y)$$

is a steady solution, for all $y \in \mathbb{R}$. The question is which of these states will be selected in the long time limit for the time-dependent problem (3.1.1). We have the selection principle – the best adapted population will be selected.

Theorem 3.1.1 *Assume that $b(x)$ is continuous, $b(x) \geq \underline{b} > 0$ for all $x \in [x_m, x_M]$, and that $b(x)$ attains its maximum over the interval $[x_m, x_M]$ at a single point $\bar{x} \in (x_m, x_M)$. Then the solution to (3.1.1) satisfies*

$$\lim_{t \rightarrow +\infty} \rho(t) = \bar{\rho} = b(\bar{x}), \quad n(t, x) \rightarrow b(\bar{x})\delta(x - \bar{x}), \quad \text{as } t \rightarrow +\infty, \quad (3.1.2)$$

the last convergence in the sense of distributions.

Proof. In this simple case, we give a computational proof. The function

$$N(t, x) = n(t, x) \exp \left\{ \int_0^t \rho(s) ds \right\}$$

satisfies

$$\frac{dN}{dt} = b(x)N(t, x),$$

hence

$$N(t, x) = n_0(x)e^{b(x)t}.$$

We also have

$$\frac{d}{dt} \left(\exp \left\{ \int_0^t \rho(s) ds \right\} \right) = \rho(t) \exp \left\{ \int_0^t \rho(s) ds \right\} = \int_{\mathbb{R}} N(t, x) dx = \int_{\mathbb{R}} n_0(x) e^{b(x)t} dx,$$

so that (recall that $b(x) > 0$ by assumption)

$$\exp \left\{ \int_0^t \rho(s) ds \right\} = \int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx + K, \quad K = 1 - \int_{\mathbb{R}} \frac{n_0(x)}{b(x)} dx,$$

whence

$$\int_0^t \rho(s) ds = \log \left(\int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx + K \right).$$

Differentiating, we get an explicit expression for the total population density:

$$\rho(t) = \left(\int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx + K \right)^{-1} \int_{\mathbb{R}} n_0(x) e^{b(x)t} dx.$$

To see what happens as $t \rightarrow +\infty$, we note that

$$\rho(t) \leq \left(\int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx + K \right)^{-1} b(\bar{x}) \int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx \rightarrow b(\bar{x}),$$

as $t \rightarrow +\infty$, because of the assumption that $b(x) \geq \underline{b} > 0$ for all $x \in [x_m, x_M]$. For the converse, we take $\varepsilon > 0$ and look at the set

$$I_\varepsilon = \{x : b(x) \geq b(\bar{x}) - \varepsilon\}.$$

Then, we have

$$\begin{aligned} \rho(t) &\geq \left(\int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx + K \right)^{-1} \int_{I_\varepsilon} n_0(x) e^{b(x)t} dx \\ &\geq \left(\int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx + K \right)^{-1} (b(\bar{x}) - \varepsilon) \int_{I_\varepsilon} \frac{n_0(x)}{b(x)} e^{b(x)t} dx = \frac{(b(\bar{x}) - \varepsilon)}{A_\varepsilon(t)}, \end{aligned}$$

where

$$\begin{aligned} A_\varepsilon(t) &= \left(\int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx + K \right) \left(\int_{I_\varepsilon} \frac{n_0(x)}{b(x)} e^{b(x)t} dx \right)^{-1} \\ &= \left(\int_{\mathbb{R}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx \right) \left(\int_{I_\varepsilon} \frac{n_0(x)}{b(x)} e^{b(x)t} dx \right)^{-1} + o(1). \end{aligned}$$

Note that

$$\int_{I_\varepsilon} \frac{n_0(x)}{b(x)} e^{b(x)t} dx \geq \int_{I_{\varepsilon/2}} \frac{n_0(x)}{b(x)} e^{b(x)t} dx \geq C e^{(b(\bar{x}) - \varepsilon/2)t},$$

while

$$\int_{\mathbb{R} \setminus I_\varepsilon} \frac{n_0(x)}{b(x)} e^{b(x)t} dx \leq C e^{(b(\bar{x}) - \varepsilon)t}.$$

It follows that $A_\varepsilon(t) \rightarrow 1$ as $t \rightarrow +\infty$ for each $\varepsilon > 0$ fixed, and therefore

$$\rho(t) \rightarrow b(\bar{x}) \text{ as } t \rightarrow +\infty. \quad (3.1.3)$$

Next, from the expression for $N(t, x)$ we know that

$$n(t, x) = n_0(x) e^{b(x)t} \exp \left\{ - \int_0^t \rho(s) ds \right\} \quad (3.1.4)$$

It is easy to see from (3.1.3) and (3.1.4) that for $x \neq \bar{x}$ we have $n(t, x) \rightarrow 0$, uniformly on each set of the form $|x - \bar{x}| > \varepsilon$. It then follows from (3.1.3) that

$$n(t, x) \rightarrow b(\bar{x}) \delta(x - \bar{x}) \text{ as } t \rightarrow +\infty,$$

in the sense of distributions. This finishes the proof. \square

A more general situation

A more general model than (3.1.1) may have the form

$$\frac{\partial n(t, x)}{\partial t} = b(x, \rho(t))n(t, x) - d(x, \rho(t))n(t, x), \quad (3.1.5)$$

so that the birth and death rates of the population with a trait $x \in \mathbb{R}$ depend both on x and the total population

$$\rho(t) = \int_{\mathbb{R}} n(t, x) dx.$$

We will assume for simplicity that the functions $b(x, \rho)$ and $d(x, \rho)$ factorize:

$$b(x, \rho) = b(x)Q_b(\rho), \quad d(x, \rho) = d(x)Q_d(\rho), \quad (3.1.6)$$

the functions b and d are smooth and positive, as are Q_b and Q_d , and the following uniform bounds hold:

$$0 < b_m \leq b(x) \leq b_M, \quad 0 < d_m \leq d(x) \leq d_M, \quad \text{for all } x \in \mathbb{R}.$$

Then (3.1.5) becomes

$$\frac{\partial n(t, x)}{\partial t} = [b(x)Q_b(\rho(t)) - d(x)Q_d(\rho(t))]n(t, x). \quad (3.1.7)$$

Hence, for a given population size $\rho(t)$, we may characterize x such that

$$b(x)Q_b(\rho(t)) - d(x)Q_d(\rho(t)) > 0 \quad (3.1.8)$$

as a favourable region, and x such that

$$b(x)Q_b(\rho(t)) - d(x)Q_d(\rho(t)) < 0 \quad (3.1.9)$$

as an unfavourable region.

The following bounds will ensure that the population does not explode or disappear completely: first, there exists $\rho_M > 0$ so that

$$\alpha_M = \max_{x \in \mathbb{R}} [b(x)Q_b(\rho_M) - d(x)Q_d(\rho_M)] < 0, \quad (3.1.10)$$

and, second, there exists $\rho_m \in (0, \rho_M)$ such that

$$\alpha_m = \min_{x \in \mathbb{R}} [b(x)Q_b(\rho_m) - d(x)Q_d(\rho_m)] > 0. \quad (3.1.11)$$

In other words, if the population is too large then all of \mathbb{R} becomes unfavourable, and if it is too small, then all of \mathbb{R} becomes favourable.

Proposition 3.1.2 *Assume that $n_0(x) \geq 0$ and $\rho_m \leq \rho(t=0) \leq \rho_M$, then $\rho_m \leq \rho(t) \leq \rho_M$ for all $t > 0$.*

Proof. We will just show that $\rho(t) \geq \rho_m$. This is a consequence of the maximum principle. Indeed, assume that τ_0 is the first time such that $\rho(\tau_0) = \rho_m$, then

$$\left. \frac{d\rho}{dt} \right|_{t=\tau_0} = \int_{\mathbb{R}} [b(x)Q_b(\rho_m) - d(x)Q_d(\rho_m)]n(\tau_0, x)dx \geq \alpha_m \rho_m > 0.$$

It follows that $\rho(t) \geq \rho_m$ for all $t > 0$. \square

Next, we will assume that

$$Q'_b(\rho) < 0, \quad Q'_d(\rho) > 0, \quad \text{for all } \rho \in \mathbb{R}, \quad (3.1.12)$$

that is, the growth rate decreases, and the death rate increases, as the population grows. We will make an additional assumption that there exists a unique $\rho = \bar{\rho}$ such that

$$\max_{\mathbb{R}} [b(x)Q_b(\rho) - d(x)Q_d(\rho)] = 0, \quad (3.1.13)$$

and that there exists a unique \bar{x} such that

$$b(\bar{x})Q_b(\bar{\rho}) - d(\bar{x})Q_d(\bar{\rho}) = 0. \quad (3.1.14)$$

Therefore, if $\rho = \bar{\rho}$, then

$$\frac{\partial n(t, x)}{\partial t} < 0, \text{ for all } x \neq \bar{x}. \quad (3.1.15)$$

The last assumption we need is that there exists $\delta_0 > 0$ and $R > 0$ so that for all $|\rho - \bar{\rho}| < \delta_0$ we have

$$\beta_R = \max_{|x| \geq R} [b(x)Q_b(\rho) - d(x)Q_d(x)] < 0. \quad (3.1.16)$$

Under these assumptions, we still have the selection principle.

Theorem 3.1.3 *With the above assumptions, if $n_0(x) > 0$ and $\rho_m \leq \rho(t=0) \leq \rho_M$, then*

$$\rho(t) \rightarrow \bar{\rho}, \quad n(t, x) \rightarrow \bar{\rho} \delta(x - \bar{x}), \quad \text{as } t \rightarrow +\infty. \quad (3.1.17)$$

Proof. Consider an auxiliary function $Z(r)$ that satisfies

$$Z'(r) = Q(r), \quad Q(r) = \frac{Q_d(r)}{Q_b(r)}, \quad \rho_m \leq r \leq \rho_M. \quad (3.1.18)$$

This function is defined only up to constant: $Z(r) = Z_0(r) + \lambda$, where $Z_0(r)$ is a particular solution and $\lambda \in \mathbb{R}$ is a constant we are free to choose but, regardless, $Z(r)$ is uniformly bounded on $[\rho_m, \rho_M]$. Let us also define the Lyapunov functional

$$L(t) = \int_{\mathbb{R}} \frac{b(x)}{d(x)} n(t, x) dx - Z(\rho(t)). \quad (3.1.19)$$

Our first goal will be to show that $L(t)$ has a limit as $t \rightarrow +\infty$. Note that $L(t)$ is uniformly bounded – this follows from our assumptions and Proposition 3.1.2. We compute:

$$\begin{aligned}
\frac{dL(t)}{dt} &= \int_{\mathbb{R}} \frac{b(x)}{d(x)} \left(b(x)Q_b(\rho(t)) - d(x)Q_d(\rho(t)) \right) n(t, x) dx - Z'(\rho(t)) \frac{d\rho(t)}{dt} \\
&= \int_{\mathbb{R}} \frac{b(x)}{d(x)} \left(b(x)Q_b(\rho(t)) - d(x)Q_d(\rho(t)) \right) n(t, x) dx \\
&= \int_{\mathbb{R}} \frac{Q_d(\rho(t))}{Q_b(\rho(t))} \left(b(x)Q_b(\rho(t)) - d(x)Q_d(\rho(t)) \right) n(t, x) dx \\
&= \int_{\mathbb{R}} \left(\frac{b(x)}{d(x)} - \frac{Q_d(\rho(t))}{Q_b(\rho(t))} \right) \left(b(x)Q_b(\rho(t)) - d(x)Q_d(\rho(t)) \right) n(t, x) dx \\
&= \int_{\mathbb{R}} d(x)Q_b(\rho(t)) \left(\frac{b(x)}{d(x)} - Q(\rho(t)) \right)^2 n(t, x) dx \geq d_m Q_b(\rho_M) D(t).
\end{aligned} \tag{3.1.20}$$

We used the uniform bounds on $b(x)$ and $Q_d(\rho(t))$ in the last inequality. Here, the dissipation rate is

$$D(t) = \int_{\mathbb{R}} \left(\frac{b(x)}{d(x)} - Q(\rho(t)) \right)^2 n(t, x) dx. \tag{3.1.21}$$

We used above the fact that $Q_b(\rho)$ is decreasing in ρ , and $\rho(t) \leq \rho_M$ for all $t > 0$. The above computation shows that $L(t)$ is increasing. As $Z(r)$ is bounded, this function is also uniformly bounded:

$$L(t) \leq \frac{b_M}{d_m} \rho_M + \max_{r \in [\rho_m, \rho_M]} |Z(r)|, \quad \text{for all } t \geq 0,$$

thus it approaches a limit as $t \rightarrow +\infty$:

$$L(t) \rightarrow \bar{L}, \quad \text{as } t \rightarrow +\infty. \tag{3.1.22}$$

Now, we need to boost (3.1.22) to show that $\rho(t)$ has a limit. If $L(t)$ were a function of $\rho(t)$, that would be simple but as it is not, we need to do some work.

Our next goal is to show that $D(t) \rightarrow 0$ as $t \rightarrow +\infty$. We deduce from (3.1.20) a bound

$$\int_0^\infty D(t) dt < +\infty. \tag{3.1.23}$$

Let us now compute

$$\begin{aligned}
\frac{dD(t)}{dt} &= \int_{\mathbb{R}} \left(\frac{b(x)}{d(x)} - Q(\rho(t)) \right)^2 \left(b(x)Q_b(\rho(t)) - d(x)Q_d(\rho(t)) \right) n(t, x) dx \\
&\quad - 2Q'(\rho(t)) \int_{\mathbb{R}} \left(\frac{b(x)}{d(x)} - Q(\rho(t)) \right) n(t, x) dx \int_{\mathbb{R}} \left(b(y)Q_b(\rho(t)) - d(y)Q_d(\rho(t)) \right) n(t, y) dy = I + II.
\end{aligned} \tag{3.1.24}$$

As $\rho(t)$ is a priori bounded, we have

$$|I| \leq C_1 D(t),$$

with a constant $C > 0$ that depends on $\rho_m, \rho_M, b_m, b_M, d_m$ and d_M . The second term can be bounded using the bound on ρ and the Cauchy-Schwartz inequality as

$$\begin{aligned} |II| &\leq C \left(\int_{\mathbb{R}} \left(\frac{b(x)}{d(x)} - Q(\rho(t)) \right)^2 n(t, x) dx \right)^{1/2} \rho(t)^{1/2} \\ &\times \left(\int_{\mathbb{R}} (b(y)Q_b(\rho(t)) - d(y)Q_d(\rho(t)))^2 n(t, y) dy \right)^{1/2} \rho(t)^{1/2} \leq C_2 D(t), \end{aligned}$$

once again, with the constant $C_2 > 0$ that depends on the same quantities as C_1 . We conclude that

$$\int_0^\infty \left| \frac{dD(t)}{dt} \right| dt \leq C \int_0^\infty D(t) dt < +\infty. \quad (3.1.25)$$

Therefore, $D(t)$ has a limit as $t \rightarrow +\infty$. In addition, as $D(t)$ is integrable, we conclude that

$$D(t) \rightarrow 0 \text{ as } t \rightarrow +\infty. \quad (3.1.26)$$

The Cauchy-Schwartz inequality implies that

$$\int_{\mathbb{R}} \left| \frac{b(x)}{d(x)} - Q(\rho(t)) \right| n(t, x) dx \leq (D(t))^{1/2} \rho(t)^{1/2} \rightarrow 0 \text{ as } t \rightarrow +\infty. \quad (3.1.27)$$

Note that, by the definition of $Z(t)$, we have

$$\rho(t)Q(\rho(t)) - Z(\rho(t)) = L(t) + \int_{\mathbb{R}} \left(Q(\rho(t)) - \frac{b(x)}{d(x)} \right) n(t, x) dx,$$

thus (3.1.22) and (3.1.27) together imply that

$$G(\rho(t)) := \rho(t)Q(\rho(t)) - Z(\rho(t)) \rightarrow \bar{L}, \text{ as } t \rightarrow +\infty.$$

The function $G(\rho(t))$ not only has a limit as $t \rightarrow +\infty$, but unlike $L(t)$ is a function of ρ . Note that $G(\rho)$ is increasing:

$$G'(r) = (rQ(r) - Z(r))' = rQ'(r) > 0,$$

because the function $Q(r)$ is increasing in r . It follows that $\rho(t)$ has a limit:

$$\rho(t) \rightarrow \rho^*, \quad (3.1.28)$$

where ρ^* is the unique $\rho \in [\rho_m, \rho_M]$ such that $G(\rho^*) = L$. Let us now show that $\rho^* = \bar{\rho}$. Indeed, if $\rho^* > \bar{\rho}$ then

$$\max_{x \in \mathbb{R}} [b(x)Q_b(\rho^*) - d(x)Q_d(\rho^*)] < 0,$$

which implies that $n(t, x) \rightarrow 0$ as $t \rightarrow +\infty$, uniformly in $x \in \mathbb{R}$., which is a contradiction since $\rho(t) \geq \rho_m$. On the other hand, if $\rho^* < \bar{\rho}$, then

$$\max_{x \in \mathbb{R}} [b(x)Q_b(\rho^*) - d(x)Q_d(\rho^*)] > 0,$$

which, in turn, implies that $\rho(t) \rightarrow +\infty$ as $t \rightarrow +\infty$ contradicting $\rho(t) \leq \rho_M$. Therefore, we have $\rho^* = \bar{\rho}$. It follows from assumption (3.1.16) that for t sufficiently large we have

$$\frac{d}{dt} \int_{|x| \geq R} n(t, x) dx \leq \beta_R \int_{|x| \geq R} n(t, x) dx, \quad (3.1.29)$$

thus

$$\int_{|x| \geq R} n(t, x) dx \rightarrow 0 \text{ as } t \rightarrow +\infty.$$

Hence, $n(t, x)$ has a weak limit $n^*(x)$ in the space of measures along a sequence $t_n \rightarrow +\infty$, and

$$\int_{\mathbb{R}} n^*(x) dx = \bar{\rho}.$$

Finally, we know from (3.1.27) that $n^*(x)$ has to be concentrated on the set where

$$\frac{b(x)}{d(x)} - \frac{Q_d(\bar{\rho})}{Q_b(\bar{\rho})} = 0,$$

which consists of one point \bar{x} . It follows that the limit is unique and

$$n^*(x) = \bar{\rho} \delta(x - \bar{x}).$$

The proof is complete.

Mutations and spatial heterogeneities

A slight variation of the previous system is dynamics of the form

$$\frac{\partial n(t, x)}{\partial t} = b(x) Q_b(\rho_b(t)) n(t, x) - d(x) Q_d(\rho_d(t)) n(t, x). \quad (3.1.30)$$

Here, we have set

$$\rho_b(t) = \int_{\mathbb{R}} \psi_b(x) n(t, x) dx, \quad \rho_d(t) = \int_{\mathbb{R}} \psi_d(x) n(t, x) dx.$$

The function $\psi_b(y)$ measures how much the presence of the species of the trait y hurts the ability of other species to reproduce, and the function $\psi_d(y)$ measures how much stronger the competition becomes if species with the trait y are present. It is natural to assume, as before, that the functions $b(x)$ and $d(x)$ are continuous, and

$$0 < b_m \leq b(x) \leq b_M, \quad 0 < d_m \leq d(x) \leq d_M, \text{ for all } x \in \mathbb{R}. \quad (3.1.31)$$

The functions Q_b and Q_d are $C^1(\mathbb{R}^+)$ and

$$Q'_b(\rho) \leq a_1 < 0, \quad Q'_d(\rho) \geq a_2 > 0 \text{ for all } \rho > 0. \quad (3.1.32)$$

For the birth and death competition rates we assume that

$$0 < \psi_m \leq \psi_d(x), \psi_b(x) \leq \psi_M \text{ for all } x \in \mathbb{R}. \quad (3.1.33)$$

These assumptions help prevent the blow-up of the total population in a finite time that potentially may be an issue because the equation is quadratic in $n(t, x)$. We also introduce a generalization of (3.1.10) and (3.1.11): first, there exists ρ_M such that:

$$\alpha_M = \max_{x \in \mathbb{R}} [b(x)Q_b(\psi_m \rho_M) - d(x)Q_d(\psi_m \rho_M)] < 0, \quad (3.1.34)$$

and, second, there exists $\rho_m \in (0, \rho_M)$ such that

$$\alpha_m = \min_{x \in \mathbb{R}} [b(x)Q_b(\psi_M \rho_m) - d(x)Q_d(\psi_M \rho_m)] > 0. \quad (3.1.35)$$

The second modification is to also add the possibility of mutations – an individual with a trait x may give birth to offspring with a trait y , and not just with trait x as we have so far assumed. This would lead to the following dynamics:

$$\frac{\partial n(t, x)}{\partial t} = Q_b(\rho_b(t)) \int_{\mathbb{R}} b(y)K(x - y)n(t, y)dy - d(x)Q_d(\rho_d(t))n(t, x). \quad (3.1.36)$$

Here, $K(x)$ is a non-negative probability density:

$$\int_{\mathbb{R}} K(x)dx = 1.$$

The case of no mutations corresponds to $K(x) = \delta(x)$. Note that only the birth rate is affected in (3.1.36).

Let us mention in passing that (3.1.36) is a very simple nonlinear kinetic model which is very natural as it comes from a branching process with jumps that results from the mutations.

Existence of the solutions

The first step is to prove existence of the solutions. This is reasonably standard but one needs to be careful that there is no possibility of a blow-up in the "quadratic-like" dynamics (3.1.36).

Theorem 3.1.4 *Assume that the non-negative initial condition $n_0(x) \in L^1(\mathbb{R})$, and*

$$\rho_m \leq \rho_0 \leq \rho_M.$$

Then (3.1.36) has a non-negative solution such that

$$n, \frac{\partial n}{\partial t} \in C(0, +\infty; L^1(\mathbb{R})),$$

and for all $t \geq 0$ we have

$$\rho_m \leq \rho(t) \leq \rho_M. \quad (3.1.37)$$

Proof. The proof is similar to that for the Cauchy-Kovalevskaya theorem.

An a priori bound. We first obtain the a priori bound (3.1.37) on the solution (assuming that it exists). Let us integrate (3.1.36). Note that

$$\int_{\mathbb{R} \times \mathbb{R}} b(y)K(x - y)n(t, y)dydx = \int_{\mathbb{R}} b(y)n(t, y)dy.$$

It follows that

$$\begin{aligned} \frac{d\rho(t)}{dt} &= Q_b(\rho_b(t)) \int_{\mathbb{R}} b(x)n(t,x)dx - Q_d(\rho_d(t)) \int_{\mathbb{R}} d(x)n(t,x)dx \\ &= \int [Q_b(\rho_b(t))b(x) - Q_d(\rho_d(t))d(x)]n(t,x)dx \leq \rho(t) \max_y [Q_b(\rho_b(t))b(y) - Q_d(\rho_d(t))d(y)]. \end{aligned} \quad (3.1.38)$$

Note that

$$\rho_b(t) \geq \psi_m \rho(t), \quad \rho_d(t) \geq \psi_m \rho(t),$$

thus

$$Q_b(\rho_b(t)) \leq Q_b(\psi_m \rho(t)),$$

and

$$Q_d(\rho_d(t)) \geq Q_d(\psi_m \rho(t)).$$

Using this in (3.1.38) gives

$$\frac{d\rho(t)}{dt} \leq \rho(t) \max_y [Q_b(\psi_m \rho(t))b(y) - Q_d(\psi_m \rho(t))d(y)].$$

Therefore, if $\rho(t) > \rho_M$ then

$$\frac{d\rho(t)}{dt} < 0,$$

and $\rho(t)$ decreases. Similarly, if $\rho(t) < \rho_m$, then

$$\frac{d\rho(t)}{dt} > 0,$$

and $\rho(t)$ increases. Hence, if initially we have $\rho_m \leq \rho_0 \leq \rho_M$, then for all $t > 0$ we still have $\rho_m \leq \rho(t) \leq \rho_M$.

Existence. We will use the fixed point theorem for the existence. The argument is applicable to other kinetic models as long as good a priori bounds are available, which is usually the case for linear models.

Consider the Banach space

$$X = C([0, T]; L^1(\mathbb{R})), \quad \|m\|_X = \sup_{0 \leq t \leq T} \|m(t)\|_{L^1(\mathbb{R})},$$

for some $T > 0$ to be chosen. The choice of this space is, again, very natural for branching kinetic models as it measures the total number of particles. Let us choose $C_0 = 2\rho_M$ and T sufficiently small so that

$$\rho_0 + Tb_M Q_b(0)C_0 \leq C_0,$$

and set

$$S = \{m \in X, m \geq 0, \|m\|_X \leq C_0\}.$$

Given a function $m \in S$, define

$$R_b(t) = \int_{\mathbb{R}} \psi_b(x)m(t,x)dx, \quad R_d(t) = \int_{\mathbb{R}} \psi_d(x)m(t,x)dx,$$

and let $n(t, x)$ be the solution of the ODE, that we solve x by x :

$$\frac{\partial n(t, x)}{\partial t} = Q_b(R_b(t)) \int_{\mathbb{R}} b(y)K(x-y)m(t, y)dy - d(x)Q_d(R_d(t))n(t, x),$$

with the initial condition $n(0, x) = n_0(x)$. We may then define the mapping $m \rightarrow \Phi(m) = n$, and the claim is that Φ has a unique fixed point in S if we choose a good C_0 and a sufficiently small T . We need to verify two conditions: (i) Φ maps S into S , and (ii) that Φ is a contraction for T sufficiently small. If we can verify these conditions then the Banach-Picard fixed point theorem implies that Φ has a fixed point in S , which is a solution we seek. We can then iterate this argument on the intervals $[T, 2T]$, $[2T, 3T]$, \dots . Note that on each time step the solution will satisfy $\rho_m \leq \rho(t) \leq \rho_M$, hence we can restart the argument each time.

To check (i) we simply write down the solution formula:

$$\begin{aligned} n(t, x) &= n_0(x) \exp\left(-d(x) \int_0^t Q_b(R_b(s))ds\right) \\ &+ \int_0^t Q_b(R_b(s)) \int_{\mathbb{R}} b(y)K(x-y)m(s, y)dy \exp\left\{-d(x) \int_s^t Q_d(R_d(s'))ds'\right\} ds. \end{aligned} \quad (3.1.39)$$

It follows that $n \geq 0$, and we also have

$$\frac{\partial n(t, x)}{\partial t} \leq Q_b(R_b(t)) \int_{\mathbb{R}} b(y)K(x-y)m(t, y)dy, \quad (3.1.40)$$

so that

$$\|n(t)\|_{L^1} \leq \rho_0 + TQ_b(0)b_M C_0 \leq C_0, \quad (3.1.41)$$

if T is sufficiently small. Thus, Φ maps S to S .

To check that Φ is a contraction, take $m_{1,2} \in S$, and write

$$\begin{aligned} \frac{\partial}{\partial t}(n_1 - n_2) &= Q_b(R_b^1(t)) \int_{\mathbb{R}} b(y)K(x-y)m_1(t, y)dy - d(x)Q_d(R_d^1(t))n_1(t, x) \\ &- Q_b(R_b^2(t)) \int_{\mathbb{R}} b(y)K(x-y)m_2(t, y)dy + d(x)Q_d(R_d^2(t))n_2(t, x) \\ &= [Q_b(R_b^1(t)) - Q_b(R_b^2(t))] \int_{\mathbb{R}} b(y)K(x-y)m_1(t, y)dy \\ &+ Q_b(R_b^2(t)) \int_{\mathbb{R}} b(y)K(x-y)[m_1^1(t, y) - m_2^2(t, y)]dy \\ &- d(x)Q_d(R_d^1(t))(n_1(t, x) - n_2(t, x)) + d(x)[Q_d(R_d^2(t)) - Q_d(R_d^1(t))]n_2(t, x). \end{aligned}$$

Integrating in x we obtain

$$\begin{aligned} \|n_1 - n_2\|_X &\leq \psi_M C_0 T b_M \|m_1 - m_2\|_X + Q_b(0)b_M T \|m_1 - m_2\|_X \\ &+ d_M Q_d(\rho_M) T \|n_1 - n_2\|_X + d_M \rho_M \psi_M T \|m_1 - m_2\|_X. \end{aligned}$$

Therefore, if T is sufficiently small, then

$$\|n_1 - n_2\|_X \leq c \|m_1 - m_2\|_X,$$

with $c < 1$. Thus, for such T the mapping $\Phi : S \rightarrow S$ is a contraction, and has a fixed point, which is the solution we seek. \square

Small mutations: the asymptotic limit

We now consider the situation when mutations are small: this is modeled by taking a smooth compactly supported kernel $K(x)$ of the form

$$K_\varepsilon(x) = \frac{1}{\varepsilon}K\left(\frac{x}{\varepsilon}\right), \quad k(x) \geq 0, \quad \int_{\mathbb{R}} K(z)dz = 1. \quad (3.1.42)$$

Of course, one would not expect small mutations to have a non-trivial effect on times of the order $t \sim O(1)$, because

$$\int_{\mathbb{R}} b(y)K_\varepsilon(x-y)n(t, y)dy = \int_{\mathbb{R}} b(x-\varepsilon z)K(z)n(t, x-\varepsilon z)dz \rightarrow \int_{\mathbb{R}} b(x)K(z)n(t, x)dz = b(x)n(t, x), \quad (3.1.43)$$

as $\varepsilon \rightarrow 0$. That is, the model with small mutations should be well-approximated by the model (3.1.30) with no mutations. In order for the small mutations to have a non-trivial effect, we need to wait for times of the order $t \sim O(\varepsilon^{-1})$. Accordingly, we consider the system in the rescaled time variable:

$$\varepsilon \frac{\partial n^\varepsilon(t, x)}{\partial t} = Q_b(\rho_b^\varepsilon(t)) \int_{\mathbb{R}} b(y)K_\varepsilon(x-y)n^\varepsilon(t, y)dy - d(x)Q_d(\rho_d^\varepsilon(t))n^\varepsilon(t, x), \quad (3.1.44)$$

with

$$\rho_b^\varepsilon(t) = \int \psi_b(x)n^\varepsilon(t, x)dx, \quad \rho_d^\varepsilon(t) = \int \psi_d(x)n^\varepsilon(t, x)dx.$$

We will show that in the limit $\varepsilon \rightarrow 0$ there is a selection principle, so that at every time t there is only one dominant trait $\bar{x}(t)$ but $\bar{x}(t)$ itself has a non-trivial dynamics, so that typically we will have

$$n_\varepsilon(t, x) \rightarrow \bar{n}(t, x) = \bar{\rho}(t)\delta(x - \bar{x}(t)). \quad (3.1.45)$$

Our goal will be to understand the dynamics of $\bar{x}(t)$ and $\bar{\rho}(t)$. Such limiting population is called monomorphic. It is also possible that the limit is a sum of several Dirac masses at $\bar{x}_1(t), \bar{x}_2(t), \dots, \bar{x}_N(t)$, and then the population is called polymorphic.

We will assume that the initial population is nearly monomorphic:

$$n_0^\varepsilon(x) = e^{\phi_0^\varepsilon(x)/\varepsilon}, \quad (3.1.46)$$

with a function $\phi_0^\varepsilon(x)$ such that

$$\phi_0^\varepsilon(x) \rightarrow \phi_0(x) \leq 0, \quad \text{uniformly in } \mathbb{R}, \quad (3.1.47)$$

and

$$\int_{\mathbb{R}} n_0^\varepsilon(x)dx \rightarrow M_0 > 0, \quad \varepsilon \rightarrow 0. \quad (3.1.48)$$

Note that $n_0^\varepsilon(x)$ is very small where $\phi_0^\varepsilon(x) \ll -\varepsilon$, which is, approximately, the region where $\phi_0(x) < 0$. Thus, in order to ensure we have initially a nearly monomorphic population, we will assume that

$$\max_{x \in \mathbb{R}} \phi_0(x) = 0 = \phi_0(\bar{x}_0) \text{ for a unique } \bar{x}_0 \in \mathbb{R}. \quad (3.1.49)$$

A typical example is the Gaussian family

$$n_0^\varepsilon(x) = \frac{1}{\sqrt{2\pi\varepsilon}} e^{-|x|^2/(2\varepsilon)}, \quad \phi_0^\varepsilon(x) = -\frac{|x|^2}{2} - \frac{\varepsilon}{2} \log(2\pi\varepsilon).$$

Let us write the equation for ϕ_ε :

$$\begin{aligned} \frac{\partial \phi^\varepsilon(t, x)}{\partial t} &= e^{-\phi^\varepsilon(t, x)/\varepsilon} Q_b(\rho_b^\varepsilon(t)) \int_{\mathbb{R}} b(y) K_\varepsilon(x - y) e^{\phi^\varepsilon(t, y)/\varepsilon} dy - d(x) Q_d(\rho_d^\varepsilon(t)) \quad (3.1.50) \\ &= Q_b(\rho_b^\varepsilon(t)) \int_{\mathbb{R}} b(x - \varepsilon y) K(y) e^{[\phi^\varepsilon(t, x - \varepsilon y) - \phi^\varepsilon(t, x)]/\varepsilon} dy - d(x) Q_d(\rho_d^\varepsilon(t)). \end{aligned}$$

It is convenient to assume that K is even: $K(y) = K(-y)$, then, expanding in ε we get the formal limit:

$$\frac{\partial \phi(t, x)}{\partial t} = Q_b(\rho_b(t)) b(x) \int_{\mathbb{R}} K(y) \exp \left\{ y \frac{\partial \phi(t, x)}{\partial x} \right\} dy - d(x) Q_d(\rho_d(t)). \quad (3.1.51)$$

Let us define

$$H(p) = \int_{\mathbb{R}} K(y) e^{py} dy.$$

The limiting constrained Hamilton-Jacobi problem should be understood as follows: the function $\phi(t, x)$ satisfies the Hamilton-Jacobi equation

$$\frac{\partial \phi(t, x)}{\partial t} = Q_b(\rho_b(t)) b(x) H \left(\frac{\partial \phi(t, x)}{\partial x} \right) - d(x) Q_d(\rho_d(t)). \quad (3.1.52)$$

This equation is not closed yet, since $\rho_b(t)$ and $\rho_d(t)$ are unknown. They are determined from an additional constraint:

$$\max_{x \in \mathbb{R}} \phi(t, x) = 0 \text{ for all } t \geq 0. \quad (3.1.53)$$

The total density $\bar{\rho}(t)$ is a Lagrange multiplier that ensures that the constraint (3.1.53) holds. If the maximum $\bar{x}(t)$ is unique then, thinking of

$$n(t, x) = \bar{\rho}(t) \delta(x - \bar{x}(t)),$$

we can postulate that

$$\bar{\rho}_b(t) = \psi_b(\bar{x}(t)) \bar{\rho}(t), \quad \bar{\rho}_d(t) = \psi_d(\bar{x}(t)) \bar{\rho}(t). \quad (3.1.54)$$

Therefore, the formal limit is as follows: find a function $\phi(t, x)$, and $\bar{\rho}(t)$ and $\bar{x}(t)$, so that $\phi(t, x)$ satisfies (3.1.52) with $\rho_b(t)$ and $\rho_d(t)$ given in terms of $\bar{\rho}(t)$ and $\bar{x}(t)$ by (3.1.54), the constraint (3.1.53) holds, and $\phi(t, x)$ attains its maximum at $\bar{x}(t)$, where

$$\phi(t, \bar{x}(t)) = 0. \quad (3.1.55)$$

An example of the constrained Hamilton-Jacobi problem

Let us explain the above scheme on a simple example. Let us assume that $Q_b \equiv 1$, $d \equiv 1$, $\psi_d \equiv 1$ and $Q_d(u) = u$, so that the starting problem is

$$\frac{\partial n^\varepsilon(t, x)}{\partial x} = \int_{\mathbb{R}} b(y)K_\varepsilon(x - y)n^\varepsilon(t, y)dy - \rho_\varepsilon(t)n^\varepsilon(t, x), \quad (3.1.56)$$

with

$$\rho^\varepsilon(t) = \int_{\mathbb{R}} n^\varepsilon(t, x)dx.$$

For short times this model reduces to the familiar simple problem

$$\frac{\partial n(t, x)}{\partial t} = b(x)n(t, x) - \rho(t)n(t, x),$$

with which we have started. The function $\phi^\varepsilon(t, x)$ satisfies

$$\frac{\partial \phi^\varepsilon(t, x)}{\partial t} = \int_{\mathbb{R}} b(x + \varepsilon y)K(y)e^{[\phi^\varepsilon(t, x + \varepsilon y) - \phi^\varepsilon(t, x)]/\varepsilon} dy - \rho^\varepsilon(t). \quad (3.1.57)$$

This gives the following constrained Hamilton-Jacobi problem (3.1.52):

$$\begin{aligned} \frac{\partial \phi(t, x)}{\partial t} &= b(x)H\left(\frac{\partial \phi(t, x)}{\partial x}\right) - \bar{\rho}(t), \\ \max_{x \in \mathbb{R}} \phi(t, x) &= 0 = \phi(t, \bar{x}(t)), \text{ for all } t \geq 0, \\ \phi(0, x) &= \phi_0(x). \end{aligned} \quad (3.1.58)$$

The Hamiltonian is, as before,

$$H(p) = \int_{\mathbb{R}} K(y)e^{py}dy.$$

In this simple example, we can use the following trick: set

$$R(t) = \int_0^t \bar{\rho}(s)ds, \quad \psi(t, x) = \phi(t, x) + R(t),$$

then we arrive at the unconstrained Hamilton-Jacobi equation for the function $\psi(t, x)$:

$$\begin{aligned} \frac{\partial \psi(t, x)}{\partial t} &= b(x)H\left(\frac{\partial \psi(t, x)}{\partial x}\right), \\ \psi(0, x) &= \phi_0(x). \end{aligned} \quad (3.1.59)$$

Then, after solving (3.1.59) we may simply set

$$R(t) = \max_{x \in \mathbb{R}} \psi(t, x),$$

enforcing the constraint on $\phi(t, x)$. Note that the point $\bar{x}(t)$ where the function $\phi(t, x)$ vanishes (and attains its maximum) is simply

$$\bar{x}(t) = \operatorname{argmax}_{\mathbb{R}} \psi(t, x).$$

Exercise 3.1.5 Is it true that $\bar{x}(t) \rightarrow \operatorname{argmax}_{\mathbb{R}} b(x)$ as $t \rightarrow +\infty$ under some reasonable assumptions on $b(x)$?

Theorem 3.1.6 *Under the above assumptions, assume, in addition, that*

$$\phi_0^\varepsilon(x) \leq C_0^\varepsilon - |x|,$$

then the function

$$\psi^\varepsilon(t, x) = \phi^\varepsilon(t, x) + R^\varepsilon(t), \quad R^\varepsilon(t) = \int_0^t \rho^\varepsilon(t) dt, \quad (3.1.60)$$

satisfies

$$\psi^\varepsilon(t, x) \rightarrow \psi(t, x), \text{ locally uniformly in } x.$$

Here, $\psi(t, x)$ is the viscosity solution of the Hamilton-Jacobi equation (3.1.59), and

$$\phi^\varepsilon(t, x) \rightarrow \phi(t, x) = \psi(t, x) - \max_{y \in \mathbb{R}} \psi(t, y).$$

The first step toward the proof are the following propositions.

Proposition 3.1.7 *We have, for all $t \geq 0$ the bound*

$$\min \left(\min_{y \in \mathbb{R}} b(y), \rho_0^\varepsilon \right) \leq \rho^\varepsilon(t) \leq \max \left(\max_{y \in \mathbb{R}} b(y), \rho_0^\varepsilon \right). \quad (3.1.61)$$

Proof. Indeed, integrating (3.1.56) in x gives

$$\frac{d\rho^\varepsilon(t)}{dt} = \int_{\mathbb{R}} b(y) n^\varepsilon(t, x) dx - (\rho^\varepsilon(t))^2. \quad (3.1.62)$$

It follows that

$$\frac{d\rho^\varepsilon(t)}{dt} \leq b_M \rho^\varepsilon(t) - (\rho^\varepsilon(t))^2, \quad \frac{d\rho^\varepsilon(t)}{dt} \geq b_m \rho^\varepsilon(t) - (\rho^\varepsilon(t))^2,$$

with

$$b_m = \min(b(y)), \quad b_M = \max b(y).$$

The maximum principle implies then (3.1.61). \square

Proposition 3.1.8 *If the initial condition satisfies*

$$\phi_0^\varepsilon(x) \leq C_0^\varepsilon - |x|,$$

then $\psi_\varepsilon(t, x)$ defined by (3.1.60) satisfies

$$\psi^\varepsilon(t, x) \leq C_0^\varepsilon - |x| + t \left(\max_{y \in \mathbb{R}} b(y) \right) \left(\max_{|p| \leq 1} H(p) \right), \quad (3.1.63)$$

and

$$\left| \frac{\partial \psi_\varepsilon(t, x)}{\partial t} \right| \leq 2 \left(\max_{y \in \mathbb{R}} b(y) \right) H(\|\nabla \phi_0^\varepsilon\|_{L^\infty}). \quad (3.1.64)$$

Proof. The function $\psi^\varepsilon(t, x)$ satisfies

$$\frac{\partial \psi^\varepsilon(t, x)}{\partial t} = \int_{\mathbb{R}} b(x + \varepsilon y) K(y) e^{[\psi^\varepsilon(t, x + \varepsilon y) - \psi^\varepsilon(t, x)]/\varepsilon} dy. \quad (3.1.65)$$

The function

$$\bar{\psi}(t, x) = C_0^\varepsilon - |x| + tB, \quad B = (\max_{x \in \mathbb{R}} b(x)) (\max_{|p| \leq 1} H(p))$$

is a super-solution to (3.1.65): indeed, we have

$$\frac{\partial \bar{\psi}_\varepsilon(t, x)}{\partial t} = B, \quad (3.1.66)$$

and

$$\begin{aligned} \int_{\mathbb{R}} b(x + \varepsilon y) K(y) e^{[\bar{\psi}(t, x + \varepsilon y) - \bar{\psi}(t, x)]/\varepsilon} dy &\leq (\max_{x \in \mathbb{R}} b(x)) \int_{\mathbb{R}} K(y) e^{(|x| - |x + \varepsilon y|)/\varepsilon} dy \\ &\leq (\max_{x \in \mathbb{R}} b(x)) \int_{\mathbb{R}} K(y) e^{|y|} dy \leq B. \end{aligned}$$

We used here the fact that $K(y)$ is non-negative and even in y . Now, (3.1.63) follows from the maximum principle in a slightly roundabout way: set

$$m_\varepsilon(t, x) = e^{\psi^\varepsilon(t, x)/\varepsilon}, \quad \bar{m}(t, x) = e^{\bar{\psi}(t, x)/\varepsilon},$$

then

$$\frac{\partial m_\varepsilon(t, x)}{\partial t} = \int b(y) K_\varepsilon(x - y) m_\varepsilon(t, y) dy, \quad (3.1.67)$$

and

$$\frac{\partial \bar{m}(t, x)}{\partial t} \geq \int b(y) K_\varepsilon(x - y) \bar{m}(t, y) dy. \quad (3.1.68)$$

It is easy to see that (3.1.67) and (3.1.68) together with the inequality $m_\varepsilon(0, x) \leq \bar{m}(0, x)$ imply that

$$m_\varepsilon(t, x) \leq \bar{m}(t, x), \quad (3.1.69)$$

and (3.1.63) follows.

Finally, to get (3.1.64) we define

$$\Psi^\varepsilon(t, x) = \frac{\partial \psi^\varepsilon(t, x)}{\partial t},$$

and differentiate (3.1.65) to get

$$\frac{\partial \Psi^\varepsilon(t, x)}{\partial t} = \frac{1}{\varepsilon} \int_{\mathbb{R}} b(x + \varepsilon y) K(y) e^{[\psi^\varepsilon(t, x + \varepsilon y) - \psi^\varepsilon(t, x)]/\varepsilon} [\Psi^\varepsilon(t, x + \varepsilon y) - \Psi^\varepsilon(t, x)] dy. \quad (3.1.70)$$

Therefore, at the point x_0 where $\Psi^\varepsilon(t, x)$ attains its maximum we have

$$\frac{\partial \Psi^\varepsilon(t, x_0)}{\partial t} = \frac{1}{\varepsilon} \int_{\mathbb{R}} b(x_0 + \varepsilon y) K(y) e^{[\psi^\varepsilon(t, x_0 + \varepsilon y) - \psi^\varepsilon(t, x_0)]/\varepsilon} [\Psi^\varepsilon(t, x_0 + \varepsilon y) - \Psi^\varepsilon(t, x_0)] dy \leq 0, \quad (3.1.71)$$

whence

$$\max_{x \in \mathbb{R}} \Psi^\varepsilon(t, x) \leq \max_{x \in \mathbb{R}} \Psi^\varepsilon(t = 0, x).$$

The same argument shows that

$$\min_{x \in \mathbb{R}} \Psi^\varepsilon(t, x) \geq \min_{x \in \mathbb{R}} \Psi^\varepsilon(t = 0, x).$$

Finally, we use (3.1.65) at $t = 0$ to observe that, with some intermediate point $\xi(y)$ we have

$$\begin{aligned} |\Psi_\varepsilon(t = 0, x)| &= \int_{\mathbb{R}} b(x + \varepsilon y) K(y) e^{[\phi_0^\varepsilon(x + \varepsilon y) - \phi_0^\varepsilon(x)]/\varepsilon} dy \\ &\leq (\max_{y \in \mathbb{R}} b(y)) \int_{\mathbb{R}} K(y) \exp\left\{y \frac{\partial \phi_0^\varepsilon(\xi(y))}{\partial x}\right\} dy \leq 2(\max_{y \in \mathbb{R}} b(y)) H\left(\|\nabla \phi_0^\varepsilon\|_{L^\infty}\right). \end{aligned} \quad (3.1.72)$$

In the last step we used the following inequality: if $|f(y)| \leq M$, then

$$\begin{aligned} \int K(y) e^{yf(y)} dy &\leq \int_{y < 0} K(y) e^{-My} dy + \int_{y > 0} K(y) e^{My} dy \\ &\leq \int_{y < 0} K(y) e^{-My} dy + \int_{y > 0} K(y) e^{-My} dy + \int_{y > 0} K(y) e^{My} dy + \int_{y < 0} K(y) e^{My} dy \\ &= 2 \int_{\mathbb{R}} K(y) e^{My} dy = 2H(M). \end{aligned}$$

Proof of Theorem 3.1.6

First, we would like to bound the spatial derivative of ψ^ε , in order to get some compactness. Fix a time $T > 0$ and let

$$\Phi^\varepsilon(t, x) = \frac{\partial \psi_\varepsilon(t, x)}{\partial x}.$$

In order not to get a large term from differentiating $b(x + \varepsilon y)$ in x , let us write

$$\frac{\partial}{\partial t} \left(\frac{\psi^\varepsilon(t, x)}{b(x)} \right) = \int_{\mathbb{R}} \frac{b(x + \varepsilon y)}{b(x)} K(y) e^{[\psi^\varepsilon(t, x + \varepsilon y) - \psi^\varepsilon(t, x)]/\varepsilon} dy, \quad (3.1.73)$$

and differentiate in x only now:

$$\begin{aligned} \frac{\partial}{\partial t} \left(\frac{\Phi^\varepsilon(t, x)}{b(x)} \right) &= \frac{1}{b(x)^2} \frac{\partial b(x)}{\partial x} \frac{\partial \psi_\varepsilon(t, x)}{\partial t} + \int_{\mathbb{R}} \frac{\partial}{\partial x} \left(\frac{b(x + \varepsilon y)}{b(x)} \right) K(y) e^{[\psi^\varepsilon(t, x + \varepsilon y) - \psi^\varepsilon(t, x)]/\varepsilon} dy \\ &+ \frac{1}{\varepsilon} \int_{\mathbb{R}} \frac{b(x + \varepsilon y)}{b(x)} K(y) e^{[\psi^\varepsilon(t, x + \varepsilon y) - \psi^\varepsilon(t, x)]/\varepsilon} (\Phi^\varepsilon(t, x + \varepsilon y) - \Phi^\varepsilon(t, x)) dy. \end{aligned} \quad (3.1.74)$$

Note that by dividing by $b(x)$ we achieved a better situation in the sense that

$$\left| \frac{\partial}{\partial x} \frac{b(x + \varepsilon y)}{b(x)} \right| \leq C\varepsilon.$$

We consider the maximal point of $|\Phi^\varepsilon|$:

$$Q_\varepsilon(t) = \max_{x \in \mathbb{R}} |\Phi^\varepsilon(t, x)|.$$

This maximum is attained either at a point where $\Phi^\varepsilon(t, x)$ has a positive maximum, or a negative minimum. The last term in the right side of (3.1.74) is non-positive where $\Phi^\varepsilon(t, x)$ attains its maximum in x and non-negative where Φ^ε attains its minimum, while the first term is bounded by Proposition 3.1.8. We obtain therefore

$$\frac{dQ_\varepsilon}{dt} \leq C + C\varepsilon \int_{\mathbb{R}} K(y) e^{[\psi^\varepsilon(t, x+\varepsilon y) - \psi^\varepsilon(t, x)]/\varepsilon} dy \leq C + C\varepsilon \int K(y) e^{|y|Q_\varepsilon(t)} dy.$$

As $K(y)$ is compactly supported, we deduce that $Q_\varepsilon(t)$ satisfies

$$\frac{dQ_\varepsilon}{dt} \leq C + C\varepsilon e^{CQ_\varepsilon(t)}.$$

It follows that for each $T > 0$ there exist $\varepsilon_0(T)$ and C_T so that we have $Q_\varepsilon(t) \leq C_T$ for all $\varepsilon < \varepsilon_0(T)$. Therefore, the family of functions $\psi_\varepsilon(t, x)$ is locally compact due to the Arzela-Ascoli theorem.

Thus, we may extract a subsequence $\varepsilon_k \rightarrow 0$, so that both $R_\varepsilon(t)$, which is Lipschitz continuous in time, and $\psi_\varepsilon(t, x)$ have local uniform limits. The limit $\psi(t, x)$ satisfies the Hamilton-Jacobi equation in the viscosity sense (this is a non-trivial step but part of the general theory of stability of the viscosity solutions that we have seen already in the lecture notes). The fact that the maximum of $\phi(t, x)$ has to be equal to zero follows from the upper and lower bounds on $\rho^\varepsilon(t)$ – if the maximum were different from zero, then $\rho(t)$ would either tend to zero or grow at a rate which is unbounded in ε .

Dynamics of the dominant trait: the monomorphic population

Let us now explain how the dominant trait $\bar{x}(t)$ can be recovered from the solution of the Hamilton-Jacobi equation in the general case, as long as the population is monomorphic, that is, the function $\phi(t, x)$ attains a single maximum $\bar{x}(t)$ where

$$\phi(t, \bar{x}(t)) = 0. \tag{3.1.75}$$

Let us recall that $\phi(t, x)$ satisfies

$$\frac{\partial \phi(t, x)}{\partial t} = Q_b(\rho_b(t))b(x)H\left(\frac{\partial \phi(t, x)}{\partial x}\right) - d(x)Q_d(\rho_d(t)). \tag{3.1.76}$$

We will assume that $\phi(t, x)$ is smooth even though at the moment we only know that $\phi(t, x)$ is merely a viscosity solution. Later, we will see a situation when it is actually smooth. Note that (3.1.75) implies that, in addition to

$$\frac{\partial \phi(t, \bar{x}(t))}{\partial x} = 0, \tag{3.1.77}$$

which holds simply because $\bar{x}(t)$ is the maximum of $\phi(t, x)$, we have

$$0 = \frac{d}{dt}\phi(t, \bar{x}(t)) = \frac{\partial \phi(t, \bar{x}(t))}{\partial t} + \frac{d\bar{x}(t)}{dt} \frac{\partial \phi(t, \bar{x}(t))}{\partial x} = \frac{\partial \phi(t, \bar{x}(t))}{\partial t}. \tag{3.1.78}$$

We deduce then from (3.1.76) that (again we assume that $\phi(t, x)$ is smooth at $x = \bar{x}(t)$)

$$Q_b(\rho_b(t))b(\bar{x})H(0) - d(\bar{x})Q_d(\rho_d(t)) = 0. \quad (3.1.79)$$

We know that

$$H(0) = \int K(y)dy = 1, \quad (3.1.80)$$

and we get an important relation

$$Q_b(\rho_b(t))b(\bar{x}) = d(\bar{x})Q_d(\rho_d(t)). \quad (3.1.81)$$

In order to get the evolution of $\bar{x}(t)$ let us differentiate (3.1.77) in t :

$$0 = \frac{d}{dt} \frac{\partial \phi(t, \bar{x}(t))}{\partial x} = \frac{\partial^2 \phi(t, \bar{x}(t))}{\partial t \partial x} + \frac{\partial^2 \phi(t, \bar{x}(t))}{\partial x^2} \frac{d\bar{x}(t)}{dt}. \quad (3.1.82)$$

On the other hand, differentiating (3.1.76) in x gives

$$\begin{aligned} \frac{\partial^2 \phi(t, \bar{x}(t))}{\partial t \partial x} &= Q_b(\rho_b(t)) \frac{\partial b(\bar{x})}{\partial x} H\left(\frac{\partial \phi(t, \bar{x})}{\partial x}\right) + b(\bar{x}) H_p\left(\frac{\partial \phi(t, \bar{x})}{\partial x}\right) \frac{\partial^2 \phi(t, \bar{x})}{\partial x^2} \\ &\quad - \frac{\partial d(\bar{x})}{\partial x} Q_d(\rho_d(t)). \end{aligned} \quad (3.1.83)$$

However, as $K(y)$ is even, we have $H(0) = 1$ and

$$H_p(0) = \int_{\mathbb{R}} yK(y)dy = 0,$$

thus we get

$$\frac{\partial^2 \phi(t, \bar{x}(t))}{\partial t \partial x} = Q_b(\rho_b(t)) \frac{\partial b(\bar{x})}{\partial x} - \frac{\partial d(\bar{x})}{\partial x} Q_d(\rho_d(t)).$$

Using this in (3.1.82) leads to an evolution equation for $\bar{x}(t)$:

$$\frac{d\bar{x}(t)}{dt} = - \left(\frac{\partial^2 \phi(t, \bar{x}(t))}{\partial x^2} \right)^{-1} \left[Q_b(\rho_b(t)) \frac{\partial b(\bar{x}(t))}{\partial x} - \frac{\partial d(\bar{x}(t))}{\partial x} Q_d(\rho_d(t)) \right]. \quad (3.1.84)$$

The pre-factor $(\phi_{xx}(t, \bar{x}(t)))^{-1}$ is very natural – if the second derivative is very small, the function $\phi(t, x)$ is very flat near $x = \bar{x}(t)$, so that it is easier for the maximum to move.

Equation (3.1.84) is still not closed – we need to find $\rho_b(t)$ and $\rho_d(t)$. To close it, observe that if the population is monomorphic, that is, $\phi(t, x)$ attains a unique maximum, then

$$n(t, x) = \bar{\rho}(t) \delta(x - \bar{x}(t)) \quad (3.1.85)$$

and

$$\rho_b(t) = \psi_b(\bar{x}(t)) \bar{\rho}(t), \quad \rho_d(t) = \psi_d(\bar{x}(t)) \bar{\rho}(t). \quad (3.1.86)$$

We may then re-write (3.1.81) as an equation for $\bar{\rho}(t)$ in terms of $\bar{x}(t)$:

$$Q_b(\psi_b(\bar{x}(t)) \bar{\rho}(t)) b(\bar{x}(t)) = d(\bar{x}(t)) Q_d(\psi_d(\bar{x}(t)) \bar{\rho}(t)). \quad (3.1.87)$$

Then we may use (3.1.86) and (3.1.87) in (3.1.84) to get a closed equation for $\bar{x}(t)$ as soon as the function $\phi(t, x)$ is known.

The above argument assumes that there is a unique maximum $\bar{x}(t)$, we will see below a case when this assumption may be rigorously justified.

The dimorphic case

Let us see what happens if the population is dimorphic: the density $n(t, x)$ has the form

$$n(t, x) = \bar{\rho}_1(t)\delta(x - \bar{x}_1(t)) + \bar{\rho}_2(t)\delta(x - \bar{x}_2(t)), \quad (3.1.88)$$

and

$$0 = \max_{x \in \mathbb{R}} \phi(t, x) = \phi(t, \bar{x}_1(t)) = \phi(t, \bar{x}_2(t)). \quad (3.1.89)$$

As before, we may derive (3.1.81) both at $\bar{x}_1(t)$ and $\bar{x}_2(t)$, so that

$$R(t) := \frac{Q_b(\rho_b(t))}{Q_d(\rho_d(t))} = \frac{d(\bar{x}_1(t))}{b(\bar{x}_1(t))} = \frac{d(\bar{x}_2(t))}{b(\bar{x}_2(t))}. \quad (3.1.90)$$

Thus, a necessary condition for dimorphism is that the function $s(x) = d(x)/b(x)$ is not one-to-one. If $s(x)$ has a "parabolic profile", so that for every y we can find two pre-images x_1 and x_2 so that

$$y = s(x_1) = s(x_2),$$

then $\bar{x}_1(t)$ and $\bar{x}_2(t)$ determine each other. The functions $\rho_b(t)$ and $\rho_d(t)$ are now given by

$$\begin{aligned} \rho_b(t) &= \psi_b(\bar{x}_1(t))\bar{\rho}_1(t) + \psi_b(\bar{x}_2(t))\bar{\rho}_2(t), \\ \rho_d(t) &= \psi_d(\bar{x}_1(t))\bar{\rho}_1(t) + \psi_d(\bar{x}_2(t))\bar{\rho}_2(t). \end{aligned} \quad (3.1.91)$$

Then, $\bar{\rho}_1(t)$ and $\bar{\rho}_2(t)$ are two Lagrange multipliers that are needed in the Hamilton-Jacobi equation to ensure that the solution $\phi(t, x)$ has exactly two maxima and it vanishes at both of them.

3.2 Hamilton-Jacobi equations with a constraint

This section is based on a paper by S. Mirrahimi and J.-M. Roquejoffre. Motivated by the models of adaptive dynamics we have discussed above, they consider equations of the form

$$\phi_t = |\nabla \phi|^2 + R(x, I), \quad (3.2.1)$$

for unknown functions $I(t)$ and $\phi(t, x)$, with the constraint

$$\max_{x \in \mathbb{R}^d} \phi(t, x) = 0, \quad (3.2.2)$$

and the initial conditions $I(0) = I_0 > 0$ and $\phi(0, x) = \phi_0(x)$. As in the mutation models, the function $R(x, I)$ is prescribed, and satisfies the familiar assumptions: first, that there exists $I_M > 0$ so that

$$\max_{x \in \mathbb{R}^d} R(x, I_M) = 0. \quad (3.2.3)$$

Second, we assume that $R(x, I)$ is strictly decreasing in I : there exist $K_1 > K_2 > 0$ so that

$$-K_1 \leq \frac{\partial R(x, I)}{\partial I} < -K_2. \quad (3.2.4)$$

One may think of $R(x, I)$ as a parabolic profile in x for each $I \in \mathbb{R}$ fixed, such that for each $I < I_M$ there is an interval $(\ell(I), r(I))$ so that $R(x, I) > 0$ for all $x \in (\ell(I), r(I))$, and $R(x, I) < 0$ for $x \notin [\ell(I), r(I)]$, with $\ell(I) - r(I) \rightarrow 0$ as $I \uparrow I_M$. In addition, we have $R(x, I) < 0$ for all $x \in \mathbb{R}$ for $I > I_M$.

As for the initial condition $\phi_0(x)$, we assume that it has a unique maximum \bar{x}_0 :

$$\max_{x \in \mathbb{R}^d} \phi_0(x) = \phi_0(\bar{x}_0) = 0, \quad (3.2.5)$$

and that the initial condition I_0 is consistent, in the sense that

$$R(\bar{x}_0, I_0) = 0. \quad (3.2.6)$$

The above assumptions are very natural and follow the reasoning of the mutation models. We now make some extra technical assumptions. We will strengthen the assumption that the maximum of $\phi_0(x)$ is unique by assuming that $\phi_0(x)$ is strictly concave and is bounded by two quadratics: there exist $K_1 > K_2 > 0$ and $K_0 > 0$ so that for all $x \in \mathbb{R}$ we have

$$-K_1 \leq D_x^2 \phi_0(x) \leq -K_2, \quad (3.2.7)$$

$$-K_0 - K_1|x|^2 \leq \phi_0(x) \leq K_0 - K_2|x|^2. \quad (3.2.8)$$

We will also make additional convexity assumptions on $R(x, I)$ in the x -variable, that hold uniformly for all $0 \leq I \leq I_M$: there exist $K_1 > K_2 > 0$ and $K_0 > 0$ so that

$$-K_1 \text{Id} \leq D_x^2 R(x, I) \leq -K_2 \text{Id}, \quad (3.2.9)$$

$$-K_1|x|^2 \leq R(x, I) \leq K_0 - K_2|x|^2, \quad (3.2.10)$$

so that, indeed, $R(x, I)$ looks like a parabolic profile in x for each I fixed (but is, of course, not necessarily exactly a parabola). The inequalities in (3.2.7) and (3.2.9) hold in the sense of positive-definite matrices. Our goal will be to prove existence and uniqueness of solutions to the constrained Hamilton-Jacobi equation (3.2.1)-(3.2.2) under these assumptions, together with the regularity assumptions

$$u_0 \in C_b^3(\mathbb{R}^d), \text{ and } R \in C_b^3(\mathbb{R}^d \times [0, I_M]). \quad (3.2.11)$$

Unlike the other assumptions made above, the regularity of $\phi_0(x)$ and $R(x, I)$ in (3.2.11) is needed in the proof but for purely technical reasons – the reader may think of both $\phi_0(x)$ and $R(x, I)$ as smooth functions.

Theorem 3.2.1 *Under the above assumptions, the Hamilton-Jacobi equation (3.2.1) with the constraint (3.2.2) and the initial conditions $\phi(0, x) = \phi_0(x)$ and $I(0) = I_0$, has a unique solution $\phi(t, x)$, $I(t)$ such that*

$$I \in W^{1, \infty}(\mathbb{R}), \text{ and } \phi \in L_{loc}^\infty((0, +\infty); W_{loc}^{3, \infty}(\mathbb{R}^d)) \cap W_{loc}^{1, \infty}((0, +\infty); L_{loc}^\infty(\mathbb{R}^d)). \quad (3.2.12)$$

3.2.1 The unconstrained problem

The first step is to look at the unconstrained problem.

$$\phi_t = |\nabla\phi|^2 + Q(t, x), \quad (3.2.13)$$

with a prescribed function $Q(t, x)$. We make the following assumptions on $Q(t, x)$ that mimic the assumptions on $R(x, I)$:

$$-K_1\text{Id} \leq D_x^2Q(t, x) \leq -K_2\text{Id}, \quad (3.2.14)$$

$$-K_1|x|^2 \leq Q(t, x) \leq K_0 - K_2|x|^2, \quad (3.2.15)$$

$$\|D_x^3Q(t, \cdot)\|_{L^\infty(\mathbb{R}^d)} \leq K_0, \quad \text{for all } t \geq 0. \quad (3.2.16)$$

Once again, assumption (3.2.16) is purely technical – one can think of $Q(t, x)$ as a smooth function. We also make the same assumptions on the initial condition $\phi_0(x)$ as in Theorem 3.2.1.

Theorem 3.2.2 *Equation (3.2.13) has a unique viscosity solution $\phi(t, x)$ that is bounded from above. Moreover, it is a classical solution, is strictly concave:*

$$-C_1\text{Id} \leq D_x^2\phi(t, x) \leq -C_2\text{Id}, \quad (3.2.17)$$

with the constants $C_1 > 0$ and $C_2 > 0$ that depend only on the constants K_0, K_1 and K_2 above, and, in addition, we have $\phi \in L_{loc}^\infty((0, +\infty); W_{loc}^{3,\infty}(\mathbb{R}^d)) \cap W_{loc}^{1,\infty}((0, +\infty); L_{loc}^\infty(\mathbb{R}^d))$, with $\|D^3\phi\|_{L^\infty} \leq C$, with C that depends only on the constants K_0, K_1 and K_2 above.

The assumption that $\phi(t, x)$ is bounded from above is natural in our context – after all, we are looking for solutions to the constrained problem that achieve their maximum at a point where they equal to zero.

The strict concavity of solutions will be essential for us later when we use this result for the constrained problem as it tells us that solutions may achieve their maximum only at a single point.

Uniqueness of the solution

Let us recall from the theory of viscosity solutions that comparison principle applies to viscosity solutions that satisfy

$$|\phi(t, x)| \leq C_T(1 + |x|^2), \quad \text{for all } x \in \mathbb{R}^d, 0 \leq t \leq T. \quad (3.2.18)$$

This is similar to the uniqueness of the solutions to the heat equation that grow at most as $\exp(c|x|^2)$ as $|x| \rightarrow +\infty$. As we have assumed that $u(t, x)$ is bounded from above, we only need to show that

$$\phi(t, x) \geq -C_T(1 + |x|^2), \quad \text{for all } x \in \mathbb{R}^d, 0 \leq t \leq T. \quad (3.2.19)$$

Taking

$$v(t, x) = \phi(t, x) + K_1t|x|^2,$$

we see that $v(t, x)$ satisfies

$$\frac{\partial v}{\partial t} = K_1|x|^2 + |\nabla\phi|^2 + Q(t, x) \geq K_1|x|^2 - K_1|x|^2 \geq 0. \quad (3.2.20)$$

This inequality only holds in the viscosity sense but it still implies that $v(t, x) \geq v(s, x)$ for all $0 \leq s \leq t$. Indeed, let us fix T and assume that there exists some $0 < s < s' \leq T$ and x_0 so that $v(s, x_0) > v(s', x_0)$. Consider the test function

$$w_\varepsilon(t, x) = v(s, x_0) - \frac{\varepsilon}{T-t} - \frac{|x-x_0|^2}{\varepsilon^2}.$$

If ε is sufficiently small, then the difference

$$z(t, x) = v(t, x) - w_\varepsilon(t, x) = v(t, x) - v(s, x_0) + \frac{\varepsilon}{T-t} + \frac{|x-x_0|^2}{\varepsilon^2}$$

satisfies $z(s', x_0) < 0$ while $z(t, x) > 0$ for $|x-x_0| = 1$ and $s \leq t \leq T$. It follows that the function $z(t, x)$ attains a local minimum over all $(s, x) \in [s, T] \times \{|x-x_0| \leq 1\}$ at a point $(t_\varepsilon, x_\varepsilon)$ with $s \leq t_\varepsilon < T$, and $x_\varepsilon \rightarrow x_0$ as $\varepsilon \rightarrow 0$. The viscosity inequality (3.2.20) then implies that at the minimum point we have

$$\frac{\partial w_\varepsilon(t_\varepsilon, x_\varepsilon)}{\partial t} \geq 0,$$

which translates simply into

$$0 \leq \frac{\partial}{\partial t} \left(-\frac{1}{T-t} \right) = -\frac{1}{(T-t)^2},$$

which is a contradiction. It follows that $v(t, x)$ is increasing in t . In addition, at $t = 0$ we have

$$v(0, x) = \phi_0(x) \geq -K_0 - K_1|x|^2,$$

thus

$$\phi(t, x) = v(t, x) - K_1t|x|^2 \geq v(0, x) - K_1t|x|^2 \geq -K_0 - (1+t)K_1|x|^2,$$

and (3.2.19) follows. Hence, comparison principle can be applied to any pair of solutions to the initial value problem that are uniformly bounded from above, which gives uniqueness.

Existence and regularity of the solution

A viscosity solution to (3.2.13) is given by the dynamic programming principle

$$\phi(t, x) = \sup \left\{ F(\gamma) : \gamma \in C^1([0, t]; \mathbb{R}^d), \gamma(t) = x \right\}, \quad (3.2.21)$$

where

$$F(\gamma) = \phi_0(\gamma(0)) + \int_0^t \left(-\frac{|\dot{\gamma}(s)|^2}{4} + Q(s, \gamma(s)) \right) ds. \quad (3.2.22)$$

It is bounded from above because $Q(t, x)$ and $u_0(x)$ are, thus the uniqueness argument above shows that this is the unique solution bounded from above. Let us show that the maximizing

trajectory exists and is unique for each (t, x) fixed. Let $\gamma_n \in C^1([0, t]; \mathbb{R}^d)$ be a sequence of trajectories such that $F(\gamma_n) \rightarrow u(t, x)$. As $F(\gamma_n) \geq F(\Gamma) - 1$, where $\Gamma(t) \equiv x$, we know that

$$\frac{1}{4} \int_0^t |\dot{\gamma}_n(s)|^2 ds \leq 1 + \phi_0(\gamma_n(0)) + \int_0^t Q(s, \gamma_n(s)) ds - \phi_0(x) - \int_0^t Q(s, x) ds \leq C(1 + |x|^2)(1 + t).$$

It follows that there exists $\bar{\gamma} \in W^{1,2}([0, t] \times \mathbb{R}^d)$ such that $\gamma_n \rightarrow \bar{\gamma}$ strongly in $C([0, t] \times \mathbb{R}^d)$ and weakly in $W^{1,2}([0, t] \times \mathbb{R}^d)$. Hence, we have

$$\phi_0(\gamma_n(0)) \rightarrow \phi_0(\bar{\gamma}(0)),$$

as well as

$$\int_0^t Q(s, \gamma_n(s)) ds \rightarrow \int_0^t Q(s, \bar{\gamma}(s)) ds,$$

and

$$\int_0^t |\dot{\bar{\gamma}}(s)|^2 ds \leq C_{t,x}.$$

Thus, we have

$$\phi(t, x) = \phi_0(\bar{\gamma}(0)) + \int_0^t \left(-\frac{|\dot{\bar{\gamma}}(s)|^2}{4} + Q(s, \bar{\gamma}(s)) \right) ds. \quad (3.2.23)$$

We claim that a trajectory $\bar{\gamma}(t)$ that realizes (3.2.23) is unique. Indeed, the concavity assumptions (3.2.7) on $\phi_0(x)$ and (3.2.14) on $Q(t, x)$ imply that $F(\gamma)$ is strictly concave in γ , hence the minimizer is unique. In addition, it has to satisfy the Euler-Lagrange equations

$$\frac{d^2 \bar{\gamma}(s)}{ds^2} = -2\nabla Q(s, \bar{\gamma}(s)), \quad (3.2.24)$$

with the boundary conditions

$$\frac{d\bar{\gamma}(0)}{ds} = -2\nabla \phi_0(\bar{\gamma}(0)), \quad \bar{\gamma}(t) = x. \quad (3.2.25)$$

Now, the regularity for $\phi(t, x)$ claimed in Theorem 3.2.2 follows simply from differentiating (3.2.23) and the regularity of the solution $\bar{\gamma}(t)$ of (3.2.24) in t and x that, in turn, follows from our regularity assumptions on $\phi_0(x)$ and $Q(t, x)$. We skip the details.

Concavity of the solution

We now prove that $\phi(t, x)$ is strictly concave by showing that there exists $\lambda > 0$ so that for all $\sigma \in [0, 1]$ and $x, y \in \mathbb{R}^d$ we have

$$\sigma \phi(t, x) + (1 - \sigma) \phi(t, y) + \lambda \sigma(1 - \sigma) |x - y|^2 \leq \phi(t, \sigma x + (1 - \sigma)y). \quad (3.2.26)$$

Let $\bar{\gamma}_x(t)$ and $\bar{\gamma}_y(t)$ be the optimal trajectories, solutions to (3.2.24) corresponding to x and y , respectively, so that

$$\phi(t, x) = \phi_0(\bar{\gamma}_x(0)) + \int_0^t \left(-\frac{|\dot{\bar{\gamma}}_x(s)|^2}{4} + Q(s, \bar{\gamma}_x(s)) \right) ds, \quad (3.2.27)$$

and

$$\phi(t, y) = \phi_0(\bar{\gamma}_y(0)) + \int_0^t \left(-\frac{|\dot{\bar{\gamma}}_y(s)|^2}{4} + Q(s, \bar{\gamma}_y(s)) \right) ds. \quad (3.2.28)$$

Taking $\gamma(s) = \sigma\bar{\gamma}_x(s) + (1 - \sigma)\bar{\gamma}_y(s)$ as a test trajectory for the point $\sigma x + (1 - \sigma)y$, which we can do because

$$\sigma\bar{\gamma}_x(t) + (1 - \sigma)\bar{\gamma}_y(t) = \sigma x + (1 - \sigma)y,$$

gives

$$\begin{aligned} \phi(t, \sigma x + (1 - \sigma)y) &\geq \phi_0(\sigma\bar{\gamma}_x(0) + (1 - \sigma)\bar{\gamma}_y(0)) \\ &+ \int_0^t \left(-\frac{|\sigma\dot{\bar{\gamma}}_x(s) + (1 - \sigma)\dot{\bar{\gamma}}_y(s)|^2}{4} + Q(s, \sigma\bar{\gamma}_x(s) + (1 - \sigma)\bar{\gamma}_y(s)) \right) ds. \end{aligned} \quad (3.2.29)$$

Strict concavity (3.2.7) of ϕ_0 implies that

$$\phi_0(\sigma\bar{\gamma}_x(0) + (1 - \sigma)\bar{\gamma}_y(0)) \geq \sigma\phi_0(\bar{\gamma}_x(0)) + (1 - \sigma)\phi_0(\bar{\gamma}_y(0)) + k\sigma(1 - \sigma)|\bar{\gamma}_x(0) - \bar{\gamma}_y(0)|^2, \quad (3.2.30)$$

and, similarly, the strict concavity (3.2.14) of $Q(t, x)$ implies that

$$\begin{aligned} \int_0^t Q(s, \sigma\bar{\gamma}_x(s) + (1 - \sigma)\bar{\gamma}_y(s)) ds &\geq \sigma \int_0^t Q(s, \bar{\gamma}_x(s)) ds + (1 - \sigma) \int_0^t Q(s, \bar{\gamma}_y(s)) ds \\ &+ k\sigma(1 - \sigma) \int_0^t |\bar{\gamma}_x(s) - \bar{\gamma}_y(s)|^2 ds. \end{aligned} \quad (3.2.31)$$

Finally, we also have

$$\begin{aligned} \int_0^t \left(-\frac{|\sigma\dot{\bar{\gamma}}_x(s) + (1 - \sigma)\dot{\bar{\gamma}}_y(s)|^2}{4} \right) ds &\geq \sigma \int_0^t \left(-\frac{|\dot{\bar{\gamma}}_x(s)|^2}{4} \right) ds + (1 - \sigma) \int_0^t \left(-\frac{|\dot{\bar{\gamma}}_y(s)|^2}{4} \right) ds \\ &+ k\sigma(1 - \sigma) \int_0^t |\dot{\bar{\gamma}}_x(s) - \dot{\bar{\gamma}}_y(s)|^2 ds, \end{aligned} \quad (3.2.32)$$

with k sufficiently small. Putting (3.2.29)-(3.2.32) together gives

$$\begin{aligned} \phi(t, \sigma x + (1 - \sigma)y) &\geq \sigma\phi(t, x) + (1 - \sigma)\phi(t, y) \\ &+ k\sigma(1 - \sigma) \left(|\bar{\gamma}_x(0) - \bar{\gamma}_y(0)|^2 + \int_0^t \left(|\bar{\gamma}_x(s) - \bar{\gamma}_y(s)|^2 + |\dot{\bar{\gamma}}_x(s) - \dot{\bar{\gamma}}_y(s)|^2 \right) ds \right). \end{aligned} \quad (3.2.33)$$

Now, to bound the integral in the right side, note that

$$\begin{aligned} |x - y|^2 &= |\bar{\gamma}_x(t) - \bar{\gamma}_y(t)|^2 = |\bar{\gamma}_x(0) - \bar{\gamma}_y(0)|^2 + \int_0^t \frac{d}{ds} \left(|\bar{\gamma}_x(s) - \bar{\gamma}_y(s)|^2 \right) ds \\ &\leq |\bar{\gamma}_x(0) - \bar{\gamma}_y(0)|^2 + \int_0^t \left(|\bar{\gamma}_x(s) - \bar{\gamma}_y(s)|^2 + |\dot{\bar{\gamma}}_x(s) - \dot{\bar{\gamma}}_y(s)|^2 \right) ds. \end{aligned} \quad (3.2.34)$$

Using this in (3.2.33) gives

$$\phi(t, \sigma x + (1 - \sigma)y) \geq \sigma\phi(t, x) + (1 - \sigma)\phi(t, y) + k\sigma(1 - \sigma)|x - y|^2, \quad (3.2.35)$$

which is the strict concavity in (3.2.26).

The constrained system

We now go back to the constrained system (3.2.1)-(3.2.2):

$$\phi_t = |\nabla\phi|^2 + R(x, I), \quad (3.2.36)$$

$$\max_{x \in \mathbb{R}^d} \phi(t, x) = 0, \quad (3.2.37)$$

with the initial conditions $I(0) = I_0 > 0$ and $\phi(0, x) = \phi_0(x)$. Recall that we assume that the initial condition $\phi_0(x)$ satisfies the concavity conditions (3.2.7)-(3.2.8), that $\phi_0(x)$ attains its unique maximum at a point \bar{x}_0 such that

$$\max_{x \in \mathbb{R}^d} \phi_0(x) = \phi_0(\bar{x}_0) = 0, \quad (3.2.38)$$

and that the initial condition I_0 is consistent, in the sense that

$$R(\bar{x}_0, I_0) = 0. \quad (3.2.39)$$

The reformulated constrained system

We now reformulate the constrained system. As the function $R(x, I)$ is concave in x by assumption (3.2.9), Theorem 3.2.2 implies that if the solution $\phi(t, x)$ to the constrained system exists, then it is strictly concave. In addition, assumptions (3.2.8) and (3.2.10) imply that $\phi(t, x)$ is sandwiched between two parabolas. Hence, in particular, $\phi(t, x)$ has a unique maximum $\bar{x}(t)$, and (3.2.37) implies that

$$\phi(t, \bar{x}(t)) = 0. \quad (3.2.40)$$

Note that

$$\nabla\phi(t, \bar{x}(t)) = 0, \quad (3.2.41)$$

because $\bar{x}(t)$ is the maximum of $\phi(t, x)$. Thus, differentiating (3.2.40) in t gives

$$\frac{\partial\phi(t, \bar{x}(t))}{\partial t} = 0. \quad (3.2.42)$$

Together, (3.2.41) and (3.2.42) imply that

$$R(t, \bar{x}(t)) = 0, \quad (3.2.43)$$

so that (3.2.39) propagates to later times.

Next, observe that differentiating (3.2.41) in t gives

$$(\partial_t \partial_k \phi)(t, \bar{x}(t)) + \partial_{k_j}^2 \phi(t, \bar{x}(t)) \frac{d\bar{x}_j}{dt} = 0, \quad (3.2.44)$$

for all fixed $k = 1, \dots, d$, with summation over the repeated index j . On the other hand, differentiating (3.2.36) in x_k gives

$$\partial_t \partial_k \phi(t, \bar{x}(t)) = 2\partial_j \phi(t, \bar{x}(t)) \partial_{k_j}^2 \phi(t, \bar{x}(t)) + \partial_k R(\bar{x}(t), I(t)) = \partial_k R(\bar{x}(t), I(t)), \quad (3.2.45)$$

because of (3.2.41). We conclude from (3.2.44) and (3.2.45) that $\bar{x}(t)$ satisfies

$$\frac{d\bar{x}(t)}{dt} = (-D^2\phi(t, \bar{x}(t)))^{-1} \nabla R(\bar{x}(t), I(t)). \quad (3.2.46)$$

Summarizing, we have shown that if $\phi(t, x)$ satisfies the constrained system, then $\phi(t, x)$, $\bar{x}(t)$ and $I(t)$ satisfy the following ODE-PDE system

$$\phi_t = |\nabla\phi|^2 + R(x, I(t)), \quad (3.2.47)$$

$$\frac{d\bar{x}(t)}{dt} = (-D^2\phi(t, \bar{x}(t)))^{-1} \nabla R(\bar{x}(t), I(t)), \quad (3.2.48)$$

$$R(\bar{x}(t), I(t)) = 0, \quad (3.2.49)$$

supplemented by the initial condition $\phi(0, x) = \phi_0(x)$, $I(0) = I_0$, and $\bar{x}(0) = x_0$, such that (3.2.38) holds, and I_0 satisfies (3.2.39). Note that $I(t)$ is determined in terms of $\bar{x}(t)$ by (3.2.49), hence we can think of (3.2.47)-(3.2.48) as a system for $\phi(t, x)$ and $\bar{x}(t)$, with $I(t)$ determined by (3.2.49).

In addition, if (3.2.47)-(3.2.49) holds, then Theorem 3.2.2 implies that $\phi(t, x)$ is concave and three times differentiable. It is straightforward to check that then if the initial conditions are consistent, so that $\bar{x}(0)$ is the unique maximum for $\phi_0(x)$, and $R(\bar{x}_0, I_0) = 0$, then $\bar{x}(t)$ is the unique maximum for $\phi(t, x)$ and, moreover, if $\phi_0(\bar{x}(0)) = 0$ then $\phi(t, \bar{x}(t)) = 0$ for all $t \geq 0$. Therefore, the system (3.2.47)-(3.2.49) is completely equivalent to the original constrained system.

Let us note that for all $t > 0$ the trajectory $\bar{x}(t)$ stays in the bounded region Ω_0 that consists of all $x \in \mathbb{R}^d$ for which $R(x, I_0) \geq 0$. The key observation to this end is that $I(t)$ is increasing. Indeed, differentiating (3.2.49) in t gives

$$\frac{d\bar{x}(t)}{dt} \cdot \nabla R(\bar{x}(t), I(t)) + \frac{\partial R(\bar{x}(t), I(t))}{\partial I} \frac{dI(t)}{dt} = 0. \quad (3.2.50)$$

In addition, multiplying (3.2.48) by $\nabla R(\bar{x}(t), I(t))$ gives

$$\frac{d\bar{x}(t)}{dt} \cdot \nabla R(\bar{x}(t), I(t)) = ((-D^2\phi(t, \bar{x}(t)))^{-1} \nabla R(\bar{x}(t), I(t)) \cdot \nabla R(\bar{x}(t), I(t))) \geq 0, \quad (3.2.51)$$

because the function $\phi(t, x)$ is strictly concave. As $R(x, I)$ is strictly decreasing in I , we deduce from (3.2.50)-(3.2.51) that

$$\frac{dI(t)}{dt} \geq 0. \quad (3.2.52)$$

Hence, $I(t)$ increasing in t , so that $I(t) \geq I_0$. It follows that

$$R(\bar{x}(t), I_0) \geq R(\bar{x}(t), I(t)) = 0,$$

thus $\bar{x}(t) \in \Omega_0$.

Existence and uniqueness for the reformulated system

The previous analysis shows that instead of proving existence and uniqueness of solutions for the constrained Hamilton-Jacobi equation, we may do the same for the reformulated system (3.2.47)-(3.2.49). This is what we will do here via a fixed point argument. The set-up is as follows. Let us fix a "nice" path $\bar{x}(t)$ with $\bar{x}(0) = \bar{x}_0$. Given $\bar{x}(t)$, we can find $I(t)$ as the solution to (3.2.49). With $I(t)$ in hand, we can solve the initial value problem

$$v_t = |\nabla v|^2 + R(x, I(t)), \quad (3.2.53)$$

with the initial condition $v(0, x) = \phi_0(x)$. Next, given $v(t, x)$, we may find the solution $y(t)$ to the ODE

$$\frac{dy}{dt} = (-D^2v(t, x(t)))^{-1} \nabla R(x(t), I(t)), \quad y(0) = \bar{x}_0, \quad (3.2.54)$$

simply by integrating the right side in time. This defines the map Φ that sends $x(t)$ to $y(t)$. A solution $\bar{x}(t)$ to (3.2.47)-(3.2.49) is a fixed point of this map, and vice versa. To prepare an application of the contraction mapping principle, take $\delta > 0$ and $r > 0$, to be determined later, and define

$$\mathcal{A} = \{x(\cdot) \in C([0, \delta]; B(x_0, r)) : x(0) = \bar{x}_0\}.$$

Our goal will be to show that if we choose δ sufficiently small, with an appropriate $r = r_\delta$, then Φ has a unique fixed point in \mathcal{A} . To this end, we need to show first that we can take δ and r so that Φ maps \mathcal{A} into \mathcal{A} . Note that Theorem 3.2.2 implies that for any $x(t)$ the solution $v(t, x)$ to (3.2.53) satisfies

$$-c_1 \text{Id} \leq D^2v(t, x) \leq -c_2 \text{Id}, \quad \|D^3v\|_{L^\infty} \leq c_3, \quad (3.2.55)$$

with the constants c_1, c_2 and c_3 that do not depend on $x(t)$. In addition, as $I(t)$ is increasing, and $R(\bar{x}_0, I(t)) \geq 0$ for all $t \geq 0$, it satisfies

$$I_0 \leq I(t) \leq I_M.$$

Let us take C so that $|\nabla R(x, I)| \leq C$ for all $|x - \bar{x}_0| \leq 1$ and $0 \leq I \leq I_M$. Then $y(t)$ satisfies

$$|y(t) - \bar{x}_0| \leq Cc_2^{-1}\delta. \quad (3.2.56)$$

Thus if we take δ sufficiently small so that $|Cc_2^{-1}\delta| < 1$ and then take $r = r_\delta := Cc_2^{-1}\delta$, then Φ maps \mathcal{A} to itself.

Our next goal is to show that Φ is a contraction on \mathcal{A} if δ is sufficiently small and $r = r_\delta$. First, given x_1 and x_2 , let I_1 and I_2 be the corresponding solutions to (3.2.49):

$$R(x_1, I_1) = R(x_2, I_2) = 0.$$

Then, we can write

$$R(x_1, I_1) - R(x_2, I_1) = R(x_2, I_2) - R(x_2, I_1),$$

so that there exist $c \in (0, 1)$ and $J \in (I_1, I_2)$, such that

$$\nabla R(cx_1 + (1-c)x_2, I_1) \cdot (x_2 - x_1) = \frac{\partial R(x_2, J)}{\partial I} (I_2 - I_1).$$

As R is strictly decreasing in I and $x(t)$ is in $B(\bar{x}_0, r_\delta)$, it follows that

$$|I(x_1(t)) - I(x_2(t))| \leq C\|x_2 - x_1\|_{C[0,\delta]}, \quad (3.2.57)$$

that is, the mapping $x(\cdot) \rightarrow I(\cdot)$ is Lipschitz.

Next, we look at how the solution $v(t, x)$ to (3.2.53) depends on $I(t)$.

Lemma 3.2.3 *Let $v^{(1)}(t, x)$ and $v^{(2)}(t, x)$ be the solutions to (3.2.53) with the initial conditions $v^{(k)}(0, x) = \phi_0(x)$ and $I(t) = I_k(t)$, $k = 1, 2$. Assume that $I_1, I_2 \in C([0, \delta]; [0, I_M])$, then*

$$\|v^{(1)} - v^{(2)}\|_{W^{2,\infty}([0,\delta] \times \mathbb{R}^d)} \leq C\delta\|I_1 - I_2\|_{L^\infty[0,\delta]}. \quad (3.2.58)$$

We do not prove this lemma, the proof is quite standard if a bit long due to the need to estimate the difference of the derivatives.

Let y_k , $k = 1, 2$ be the corresponding solutions to (3.2.54):

$$\frac{dy_k}{dt} = (-D^2v_k(t, x_k(t)))^{-1}\nabla R(x_k(t), I_k(t)), \quad y_k(0) = \bar{x}_0. \quad (3.2.59)$$

Then we have

$$\begin{aligned} & |y_1(t) - y_2(t)| \quad (3.2.60) \\ & \leq \delta \sup_{0 \leq s \leq \delta} |(-D^2v_1(t, x_1(s)))^{-1}\nabla R(x_1(s), I_1(s)) - (-D^2v_2(s, x_2(s)))^{-1}\nabla R(x_2(s), I_2(s))| \\ & \leq \delta \sup_{0 \leq s \leq \delta} |(D^2v_1(s, x_1(s)))^{-1}\nabla R(x_1(s), I_1(s)) - (D^2v_1(s, x_2(s)))^{-1}\nabla R(x_1(s), I_1(s))| \\ & \quad + \delta \sup_{0 \leq s \leq \delta} |(D^2v_1(s, x_2(s)))^{-1}\nabla R(x_1(s), I_1(s)) - (D^2v_2(s, x_2(s)))^{-1}\nabla R(x_2(s), I_2(s))| \\ & \leq \delta\|x_1 - x_2\|_{C[0,\delta]}\|D^3v_1\|_{L^\infty([0,\delta] \times B(\bar{x}_0, r_\delta))}\|\nabla R\|_{L^\infty(B(\bar{x}_0, r_\delta) \times [0, I_M])} \\ & \quad + \delta \sup_{0 \leq s \leq \delta} |(D^2v_1(s, x_2(s)))^{-1}\nabla R(x_1(s), I_1(s)) - (D^2v_2(s, x_2(s)))^{-1}\nabla R(x_1(s), I_1(s))| \\ & \quad + \delta \sup_{0 \leq s \leq \delta} |(D^2v_2(s, x_2(s)))^{-1}\nabla R(x_1(s), I_1(s)) - (D^2v_2(s, x_2(s)))^{-1}\nabla R(x_2(s), I_2(s))| \\ & \leq \delta\|x_1 - x_2\|_{C[0,\delta]}\|D^3v_1\|_{L^\infty([0,\delta] \times B(\bar{x}_0, r_\delta))}\|\nabla R\|_{L^\infty(B(\bar{x}_0, r_\delta) \times [0, I_M])} \\ & \quad + \delta\|v_1 - v_2\|_{W^{2,\infty}([0,\delta] \times \mathbb{R}^d)}\|\nabla R\|_{L^\infty(B(\bar{x}_0, r_\delta) \times [0, I_M])} \\ & \quad + \delta\|(D^2v_2)^{-1}\|_{L^\infty([0,\delta] \times \mathbb{R}^d)} \sup_{0 \leq s \leq \delta} |\nabla R(x_1(s), I_1(s)) - \nabla R(x_2(s), I_2(s))| \quad (3.2.61) \end{aligned}$$

$$\begin{aligned} & \leq \delta\|x_1 - x_2\|_{C[0,\delta]}\|D^3v_1\|_{L^\infty([0,\delta] \times B(\bar{x}_0, r_\delta))}\|\nabla R\|_{L^\infty(B(\bar{x}_0, r_\delta) \times [0, I_M])} \\ & \quad + \delta\|v_1 - v_2\|_{W^{2,\infty}([0,\delta] \times \mathbb{R}^d)}\|\nabla R\|_{L^\infty(B(\bar{x}_0, r_\delta) \times [0, I_M])} \\ & \quad + C\delta\|(D^2v_2)^{-1}\|_{L^\infty([0,\delta] \times \mathbb{R}^d)}\|D_{x,I}^2R\|_{\mathbb{R}^d \times [0, I_M]}(\|x_1 - x_2\|_{C[0,\delta]} + \|I_1 - I_2\|_{C[0,\delta]}) \quad (3.2.62) \\ & \leq C\delta\|x_1 - x_2\|_{C[0,\delta]}. \end{aligned}$$

We conclude that Φ is a contraction on $C([0, \delta]; B(\bar{x}_0, r_\delta))$ if δ is sufficiently small. Hence, the reformulated system has a unique solution for $0 \leq t \leq \delta$.

To extend this result to existence and uniqueness for all $t \in [0, T]$, one can use the standard argument, constructing the solution on $[0, \delta]$, then $[\delta, 2\delta]$ and so on. Note that δ does not depend on the initial condition, as long as \bar{x}_0 is such that there exists $\bar{I}_0 \in [0, I_M]$ such that $R(\bar{x}_0, \bar{I}_0) = 0$. This remains true as $\bar{x}(t)$ stays in Ω_0 because of the argument in (3.2.50)-(3.2.52). Therefore, we can extend the existence and uniqueness of the solution to all $t \in [0, T]$, for any $T > 0$ fixed, and the proof of Theorem 3.2.1 is complete.