



HAL
open science

On the Cauchy problem for microlocally symmetrizable hyperbolic systems with log-Lipschitz coefficients

Ferruccio Colombini, Daniele Del Santo, Francesco Fanelli, Guy Métivier

► To cite this version:

Ferruccio Colombini, Daniele Del Santo, Francesco Fanelli, Guy Métivier. On the Cauchy problem for microlocally symmetrizable hyperbolic systems with log-Lipschitz coefficients. *Indiana University Mathematics Journal*, 2020, 69 (3), 10.1512/iumj.2020.69.7963 . hal-01380371

HAL Id: hal-01380371

<https://hal.science/hal-01380371>

Submitted on 12 Oct 2016

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire **HAL**, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

ON THE CAUCHY PROBLEM FOR MICROLOCALLY SYMMETRIZABLE HYPERBOLIC SYSTEMS WITH LOG-LIPSCHITZ COEFFICIENTS

*Ferruccio Colombini*¹, *Daniele Del Santo*², *Francesco Fanelli*³, *Guy Métivier*⁴

¹ UNIVERSITÀ DI PISA
Dipartimento di Matematica
colombini@dm.unipi.it

² UNIVERSITÀ DI TRIESTE
Dipartimento di Matematica e Geoscienze
delsanto@units.it

³ UNIVERSITÉ DE LYON, UNIVERSITÉ CLAUDE BERNARD LYON 1
Institut Camille Jordan, UMR 5208
fanelli@math.univ-lyon1.fr

⁴ UNIVERSITÉ DE BORDEAUX
Institut de Mathématiques de Bordeaux, UMR 5251
guy.metivier@math.u-bordeaux1.fr

October 12, 2016

Abstract

The present paper concerns the well-posedness of the Cauchy problem for microlocally symmetrizable hyperbolic systems whose coefficients and symmetrizer are log-Lipschitz continuous, uniformly in time and space variables. For the global in space problem we establish energy estimates with finite loss of derivatives, which is linearly increasing in time. This implies well-posedness in H^∞ , if the coefficients enjoy enough smoothness in x . From this result, by standard arguments (i.e. extension and convexification) we deduce also local existence and uniqueness. A huge part of the analysis is devoted to give an appropriate sense to the Cauchy problem, which is not evident a priori in our setting, due to the very low regularity of coefficients and solutions.

2010 Mathematics Subject Classification: 35L45 (primary); 35B45, 35B65 (secondary).

Keywords: hyperbolic systems, microlocal symmetrizability, log-Lipschitz regularity, loss of derivatives, global and local Cauchy problem, well-posedness.

1 Introduction

In the present paper we study local and global questions related to the well-posedness of the Cauchy problem for $m \times m$ first order hyperbolic systems

intro_eq:L

 (1)
$$Lu(t, x) := A_0(t, x) \partial_t u(t, x) + \sum_{j=1}^n A_j(t, x) \partial_j u(t, x) + B(t, x) u(t, x)$$

under low regularity assumptions on its coefficients.

In [17], Ivrii and Petkov proved that a necessary condition for the well-posedness in the energy space $L^2(\mathbb{R}^n)$ is the existence of a bounded microlocal symmetrizer $S(t, x, \xi)$ for L (see Definition 2.1 below). We remark that this is equivalent to the hypothesis of strong hyperbolicity of the operator (see [20] and the references therein). Nonetheless, this condition is far from being sufficient, even for C^∞ well-posedness: see counterexamples in [24], [20] and [12].

On the other hand, in the simplest case when the first order coefficients A_j are symmetric matrices, i.e. the system is symmetric in the sense of Friedrichs, if they are also Lipschitz continuous over $[0, T] \times \mathbb{R}^n$, energy estimates (and then well-posedness) in L^2 are quite easy to obtain. Notice that, in this instance, the symmetrizer is simply the identity matrix, and it is constant (and then smooth) in (t, x, ξ) .

Keeping the $W^{1,\infty}([0, T] \times \mathbb{R}^n)$ regularity for the coefficients of the principal part, the previous result was extended by Métivier in [19] to the case of microlocally symmetrizable systems, under the assumption that the symmetrizer S is Lipschitz in (t, x) and smooth with respect to ξ . More recently, in [20] Métivier also proved that, under a stronger $W^{2,\infty}([0, T] \times \mathbb{R}^n)$ condition on the A_j 's, the existence of a Lipschitz continuous symmetrizer *in all the variables* (t, x, ξ) is sufficient for proving energy estimates in L^2 . Whether the additional regularity for the coefficients of L is merely a technical requirement or rather a necessary hypothesis is not clear at present.

Indeed, pathological phenomena may be produced by the lack of suitable regularity of the coefficients. To explain this assertion, let us make a brief parallel with scalar strictly hyperbolic operators of second order

intro_eq:wave

$$(2) \quad Wu(t, x) := \partial_t^2 u(t, x) - \sum_{j,k=1}^n \partial_j \left(a_{jk}(t, x) \partial_k u(t, x) \right).$$

It is well-known that, if the coefficients a_{jk} are Lipschitz continuous in t and just bounded in x , then the Cauchy problem for W is well-posed in the energy space $H^1 \times L^2$. On the contrary, whenever the Lipschitz regularity in t is not met, additional smoothness in x is required, and the Cauchy problem is no more well-posed in the usual sense. Namely, the solution loses regularity in the evolution, and energy estimates are recovered in weaker spaces $H^s \times H^{s-1}$, for $|s| < 1$, which deteriorates with the passing of time: actually, $s = s(t)$, with $s'(t) < 0$. Therefore, this justifies the additional regularity in x , which is not just a technical requirement to make energy estimates work, but which is really needed to give sense to the Cauchy problem.

Many are the relevant papers on this subject: see e.g. [4], [10], [11], [25], [5]. We refer to [6] and [7] for an overview and recent results. In passing, we mention that this loss of smoothness produces relevant effects also at the level of the control of waves. In addition, we have to point out that similar phenomena were proved to occur also for transport equations with non-regular coefficients (see e.g. [1], [13]): we refer e.g. to Chapter 3 of [2] for a review of previous results in this direction and for further references.

There is an important feature to point out about the wave operator (2). The work of Tarama [25] showed that Zygmund type conditions in time are well-adapted to this kind of analysis. These are smoothness assumptions which are made on the second variation of the function, i.e. on the symmetric difference $|f(t+\tau) + f(t-\tau) - 2f(t)|$, rather than on its modulus of continuity; besides, they can be related with special Besov type regularities. It turns out that these conditions are weaker than the corresponding ones made on the first variation, namely on $|f(t+\tau) - f(t)|$: in particular, one can recover well-posedness (with no loss) for the pure Zygmund condition

$$(Z) \quad |f(t+\tau) + f(t-\tau) - 2f(t)| \leq C|\tau|,$$

which is weaker than the Lipschitz one. In order to deal with this worse behaviour in hyperbolic Cauchy problems, one has to introduce a lower order corrector in the definition of the energy: this additional term is necessary to produce special algebraic cancellations in the estimates, erasing bad remainders arising in the time derivative of the energy. We refer to the above mentioned works [5], [6], [7] (where the coefficients depend also on x) for further progress in this direction.

Let us come back to the operator L , defined in (1). In the present paper we aim at investigating the problem of its well-posedness looking at minimal regularity conditions in time and space variables for coefficients and symmetrizers: namely, we will consider the case of non-Lipschitz dependence on (t, x) , keeping however the smoothness of $S(t, x, \xi)$ with respect to ξ (see the discussion above). Indeed, our analysis makes a broad use of paradifferential calculus, which requires smoothness of symbols in the dual variable.

In a preliminary study (see [9]), we focused on time-dependent coefficients: $A_j(t)$ and $B = 0$ for simplicity. Inspired by [25], it was natural to formulate Zygmund type hypotheses on them; on the other hand, in this context it is out of use to know a priori the existence of a microlocal symmetrizer $S(t, \xi)$ with the same regularity in time. As a matter of fact, analogously to the case of the wave equation, the time derivative of the energy $E(t) \sim (Su, u)_{L^2}$ produces bad remainders, and one needs to introduce correctors to cancel them out: the challenge there was to build up a *suitable* microlocal symmetrizer, for which energy estimates work well. Hence, we supposed the system to be hyperbolic with constant multiplicities, which means that all the eigenvalues of the principal symbol are real, semi-simple and with constant multiplicities in t and ξ ; this condition implies in particular strong hyperbolicity.

In the end, in [9] we proved well-posedness in any H^s for Zygmund type assumptions (even of integral type), while energy estimates with a finite time-increasing loss of derivatives for log-Zygmund type assumptions (in the right-hand side of (Z) an extra logarithmic factor $|\log \tau|$ is added), which entail well-posedness just in H^∞ .

We have to point out that the hypothesis of dependence of the A_j 's just on time was crucial in [9], in order to construct a good symmetrizer, and it is far to be clear at present how to deal with the more general case of dependence also on x .

As a first step in this direction, in the present paper we consider coefficients $A_j(t, x)$ with slightly better regularity conditions in t , and non-zero matrices $B(t, x)$: we will see that, now, the presence of 0-th order term makes some differences in the analysis. More precisely, we assume B to be $L^\infty([0, T]; \mathcal{C}^\gamma(\mathbb{R}^n))$, for some Hölder exponent $0 < \gamma < 1$, and the A_j 's to be bounded and uniformly log-Lipschitz continuous in their variables.

d:LL **Definition 1.1.** Let $\Omega \subset \mathbb{R}^N$ be an open domain. A function $f \in L^\infty(\Omega)$ is said to be *log-Lipschitz* in Ω , and we write $f \in LL(\Omega)$, if the quantity

$$|f|_{LL} := \sup_{y, z \in \Omega, |y-z| < 1} \left(\frac{|f(y) - f(z)|}{|y - z| \log(1 + 1/|y - z|)} \right) < +\infty.$$

We define $\|f\|_{LL} := \|f\|_{L^\infty} + |f|_{LL}$.

The same regularity hypothesis is assumed in (t, x) also for the microlocal symmetrizer (of course, it is taken smooth in ξ). Indeed, as we will see, corrector terms are no more needed in this case, and the energy can be defined in a classical way, in terms (roughly speaking) of the L^2 scalar product with respect to S . Then, the microlocal symmetrizability assumption will be still suitable for our purposes.

For operator L , supplemented with these additional hypotheses, in Theorem 2.4 we establish energy estimates on the whole $[0, T] \times \mathbb{R}^n$ with time-increasing loss of derivatives, which are the exact counterpart of similar inequalities for the wave operator (2) (see [10], [11]). These estimates, however, hold true in low regularity Sobolev spaces H^s , for $0 < s < \gamma$, due to the weak smoothness of the coefficients in x . For the same reason, the result is just local in time: if the loss of derivatives is too high, u ends up in very weak classes, for which the product with log-Lipschitz or Hölder functions is no more well-defined.

Let us explain better this point. One has to remark (see also Proposition 3.9 and Corollary 3.10 below) that multiplication by a log-Lipschitz function is a self-map of H^s if and only if $|s| < 1$, so that the first order part makes sense in H^{s-1} if and only if $-1 < s - 1 < 0$. On the other hand,

the 0-th order term will be treated as a remainder, and so we need it to be in H^s : this is true for $|s| < \gamma$. Notice here that the argument is symmetric for the conservative counterpart

$$L^* u = -\partial_t u - \sum_{j=1}^n \partial_j (A_j^* u) + B^* u,$$

and in fact, for L^* one gets a sort of “dual” version of the previous result, in H^s classes for $-\gamma < s < 0$ (but the loss is always increasing in time, of course).

Our result strongly relies on paradifferential calculus in logarithmic Sobolev classes, and especially on a fine analysis of paradifferential operators associated to symbols which are log-Lipschitz (or Hölder) continuous in x , for which we develop also a symbolic calculus. This analysis allows us to approximate the principal part \mathcal{A} of L by its paradifferential operator $T_{\mathcal{A}}$, and to define an approximate symmetrizer for $T_{\mathcal{A}}$ by taking, roughly speaking, the paradifferential operator associated to the original symmetrizer S (up to small modifications, required in order to deal with low frequencies).

On the other hand, the time dependence comes into play as well, and the weak smoothness in t causes some troubles in view of energy estimates. Then, following the approach initiated in [4] by Colombini, De Giorgi and Spagnolo for the scalar wave equation, we need to introduce a regularization in time, and to link the regularization parameter (say) ε with the dual variable ξ . More precisely, given a symbol $a = a(t, x, \xi)$, in a first moment we smooth it out by convolution with a smoothing kernel: we define the family $a_\varepsilon = \rho_\varepsilon *_t a$. Then, in view of closing energy estimates, we have to make the key choice $\varepsilon = 1/|\xi|$: this means that the approximation is different, depending on the size of the frequencies we are looking at. Therefore, the previous paradifferential calculus construction has to be completely revisited and adapted to treat new symbols $\tilde{a}(t, x, \xi)$ which are obtained by the family $(a_\varepsilon)_\varepsilon$ performing the choice $\varepsilon = |\xi|^{-1}$.

Let us point out that, in our analysis, the log-Lipschitz continuity in (t, x) of *both* the first order coefficients and the symmetrizer is exploited in a fundamental way. Also, the counterexamples established in [20] and [12] imply somehow the sharpness of our result. In particular, in [12] Colombini and Métivier were able to exhibit explicit examples of 2×2 microlocally symmetrizable systems with smooth time-dependent coefficients, for which the following phenomena occur: if the symmetrizer is ω -continuous for some modulus of continuity ω which is even slightly worse than Lipschitz, a loss of derivatives has to occur in the energy estimates (then ill-posedness in L^2); if ω is worse than log-Lipschitz, then the loss is in general infinite (which shows ill-posedness of the Cauchy problem in \mathcal{C}^∞).

This having been done, we turn our attention to local in space existence and uniqueness questions. Indeed, our regularity hypotheses are invariant under smooth change of variables, and thus they are suitable for local analysis. On the other hand, we will prove that also the microlocal symmetrizability assumption (reformulated in Definition 2.9 in a coordinate independent way) is invariant under change of variables. So, in Theorems 2.10 and 2.12 we show respectively local existence and uniqueness of solutions to the Cauchy problem for our operator L (written in a coordinate independent way, see formula (12) below). These results are the analogue of what established in [11] for wave operators W . In passing, we mention that local uniqueness will be derived from a result about propagation of zero across space-like manifolds, which also implies finite speed of propagation and a sharp description of the propagation of supports by standard arguments (see e.g. [15], [22]).

Of course, since we are in a low regularity framework, the sense of the local Cauchy problem is not clear a priori: therefore, the initial efforts (see Subsection 5.1) will be devoted to explain the setting, how one has to interpret the Cauchy problem under our assumptions and in which sense one can aim at solving it. Then, by a change of variables we will reconduct the analysis to operator L written in the form (1) above, and, by a classical convexification argument, we will be able to reduce the proof of the local statements to the global in space results, previously established.

Before going on, let us give a brief overview of the paper.

In the next section we collect our hypotheses, and we state our main results, first for the global Cauchy problem, and then for the local one. Section 3 is devoted to Littlewood-Paley analysis of log-Lipschitz and Sobolev classes, and to developing paradifferential calculus associated to low regularity symbols. This analysis will be the key to the proof of the global statement, which will be presented in Section 4. Finally, in Section 5 we discuss the local questions.

Notations

Before going on, let us introduce some notations.

First of all, given two vectors v and w in \mathbb{C}^m , we will denote by $v \cdot w$ the usual hermitian product in \mathbb{C}^m and by $|v|$ the usual norm of a vector in \mathbb{C}^m :

$$v \cdot w = \sum_{j=1}^m v_j \bar{w}_j \quad \text{and} \quad |v|^2 = v \cdot v.$$

On the contrary, given a infinite-dimensional Banach space X , we will denote by $\|\cdot\|_X$ its norm and, if it's Hilbert, by $(\cdot, \cdot)_X$ its scalar product. Typically, for us $X = L^2(\mathbb{R}^n; \mathbb{R}^m)$ or $H^s(\mathbb{R}^n; \mathbb{R}^m)$.

We will also set $\mathcal{M}_m(\mathbb{C})$ the set of all $m \times m$ matrices whose components are complex numbers, and we will denote by $|\cdot|_{\mathcal{M}}$ its norm:

$$|A|_{\mathcal{M}} := \sup_{|v|=1} |Av| \equiv \sup_{|v| \leq 1} |Av| \equiv \sup_{v \neq 0} \frac{|Av|}{|v|}.$$

Acknowledgements

The first two authors are members of the Gruppo Nazionale per l'Analisi Matematica, la Probabilità e le loro Applicazioni (GNAMPA) of the Istituto Nazionale di Alta Matematica (INdAM).

2 Basic definitions and main results

s:results

We state here our main results. Let us start by the global in space questions.

2.1 The global Cauchy problem

For $m \geq 1$, let us consider the $m \times m$ linear first order system

$$\text{def:L} \quad (3) \quad Lu(t, x) = \partial_t u(t, x) + \sum_{j=1}^n A_j(t, x) \partial_j u(t, x) + B(t, x) u(t, x)$$

defined on a strip $[0, T] \times \mathbb{R}^n$, for some time $T > 0$ and $n \geq 1$. We suppose $u(t, x) \in \mathbb{R}^m$ and, for all $1 \leq j \leq n$, the matrices $A_j(t, x) \in \mathcal{M}_m(\mathbb{C})$ as well as $B(t, x)$.

We define the principal symbol \mathcal{A} associated to the operator L : for all $(t, x, \xi) \in [0, T] \times \mathbb{R}^n \times \mathbb{R}^n$,

$$\text{def:symbol} \quad (4) \quad \mathcal{A}(t, x, \xi) := \sum_{j=1}^n \xi_j A_j(t, x).$$

Then, for all (t, x, ξ) fixed, $\mathcal{A}(t, x, \xi)$ is an $m \times m$ matrix which has complex-valued coefficients. We denote by $(\lambda_k(t, x, \xi))_{1 \leq k \leq m} \subset \mathbb{C}$ its eigenvalues at any point (t, x, ξ) .

Analogously, we consider also its conservative counterpart, i.e. the operator \tilde{L} defined by

$$\text{def:L*} \quad (5) \quad \tilde{L}u(t, x) = \partial_t u(t, x) + \sum_{j=1}^n \partial_j (A_j(t, x) u(t, x)) + B(t, x) u(t, x)$$

and whose principal symbol is still $\mathcal{A}(t, x, \xi)$, with a different quantization.

Let us state now our working hypothesis. First of all, in the sequel we will always assume that our operator L is *hyperbolic*, namely $(\lambda_k(t, x, \xi))_{1 \leq k \leq m} \subset \mathbb{R}$. As for the coefficients of L , we always suppose boundedness: we suppose that, for all $M \in \{A_1 \dots A_n, B\}$,

hyp:bound

$$(6) \quad \|M\|_{L^\infty([0, T] \times \mathbb{R}^n; \mathcal{M}_m(\mathbb{C}))} := \sup_{(t, x)} |M(t, x)|_{\mathcal{M}} \leq K_0.$$

Concerning their regularity, let us start by focusing on first order coefficients: we assume that, for all $1 \leq j \leq n$, the matrix-valued functions A_j are *log-Lipschitz* continuous (*LL* in brief) in their variables: there exists a constant $K_1 > 0$ such that, for all $\tau > 0$ and all $y \in \mathbb{R}^n \setminus \{0\}$, one has

hyp:LL

$$(7) \quad \sup_{(t, x)} |A_j(t + \tau, x + y) - A_j(t, x)|_{\mathcal{M}} \leq K_1 (|\tau| + |y|) \log \left(1 + \frac{1}{|\tau| + |y|} \right).$$

Notice that, since we are in finite dimension, this is equivalent to require the same condition on each component of A_j .

Concerning the coefficient of the lower order term, we will assume γ -Hölder continuity in space, uniformly in time. More precisely, we suppose that there exist a $\gamma \in]0, 1[$ and a constant $K_2 > 0$ such that, for all $y \in \mathbb{R}^n \setminus \{0\}$, one has

hyp:Holder

$$(8) \quad \sup_{(t, x)} |B(t, x + y) - B(t, x)|_{\mathcal{M}} \leq K_2 |y|^\gamma.$$

Finally, we will require that the system is *uniformly microlocally symmetrizable*, in the sense of Métivier (see [19], Chapter 7). The word *uniformly* here means with respect to $(t, x) \in [0, T] \times \mathbb{R}^n$ (see also Section 4 of [20]).

d:micro_symm

Definition 2.1. System (3) is *uniformly symmetrizable* if there exists a $m \times m$ matrix $S(t, x, \xi)$, homogeneous of degree 0 in ξ , such that:

- $\xi \mapsto S(t, x, \xi)$ is C^∞ for $\xi \neq 0$;
- for any point (t, x, ξ) , the matrix $S(t, x, \xi)$ is self-adjoint;
- there exist constants $0 < \lambda \leq \Lambda$ such that $\lambda \text{Id} \leq S(t, x, \xi) \leq \Lambda \text{Id}$ for any (t, x, ξ) ;
- for any point (t, x, ξ) , the matrix $S(t, x, \xi) \mathcal{A}(t, x, \xi)$ is self-adjoint.

The matrix valued function S is called a *microlocal symmetrizer* for system (1).

ex:symm

Example 2.2. Obviously, symmetric systems are microlocally symmetrizable hyperbolic systems, whose symmetrizer is simply the identity matrix.

ex:const

Example 2.3. Also hyperbolic systems with constant multiplicities (and in particular strictly hyperbolic systems, for which multiplicities are all equal to 1) are microlocally symmetrizable. Indeed, a symmetrizer can be easily constructed (see e.g. [19]) in terms of the eigenvalues and projection operators onto the eigenspaces related to $\mathcal{A}(t, x, \xi)$.

In addition, standard perturbation theory for linear operators (see e.g. [16], [20],[23]) entails that the eigenvalues and eigenprojectors inherit the same regularity in (t, x) as the coefficients of L . In particular, this implies that, in general, one cannot expect to find a symmetrizer having more smoothness than the one of the coefficients of L or \tilde{L} .

In what follows, we are going to consider the case when also the symmetrizer S is log-Lipschitz continuous in (t, x) , in the sense that it verifies an inequality of the same type as (7) at any $\xi \neq 0$ fixed. Such a regularity hypothesis for S will be exploited in a fundamental way in order to get

our result. Without loss of generality, we can assume that the constants K_0 in (6) and K_1 in (7) are large enough, to control also the corresponding quantities computed for the symmetrizer S .

Under the previous hypotheses, we can show an energy estimate with finite loss of derivatives for L and \tilde{L} . We point out that such a loss is linearly increasing in time, which implies that the solution becomes more and more irregular in the time evolution.

The estimates are stated for smooth enough u at this level. However, as a consequence of a “weak = strong” type result (see Theorem 4.10 below), they remain true for tempered distributions in a much broader class.

th:en_LL

Theorem 2.4. *Let us consider the first-order system (3), and assume it to be microlocally symmetrizable, in the sense of Definition 2.1. Suppose moreover that the coefficients $(A_j)_{1 \leq j \leq n}$ and the symmetrizer S satisfy the boundedness and log-Lipschitz conditions (6)-(7). Suppose also that the coefficient B verifies hypotheses (6)-(8), for some $\gamma \in]0, 1[$.*

Then, for all $s \in]0, \gamma[$, there exist positive constants C_1, C_2 (depending just on s, K_0 and K_1), a $\beta > 0$ (depending just on K_1) and a time $T_ \in]0, T[$, with $\beta T_* < s$, such that the estimate*

est:u_LL

$$(9) \quad \sup_{t \in [0, T_*]} \|u(t)\|_{H^{s-\beta t}} \leq C_1 e^{C_2 T} \left(\|u(0)\|_{H^s} + \int_0^{T_*} \|Lu(\tau)\|_{H^{s-\beta\tau}} d\tau \right)$$

holds true for any tempered distribution $u \in L^2([0, T_*]; H^1(\mathbb{R}^n; \mathbb{R}^m)) \cap H^1([0, T_*]; L^2(\mathbb{R}^n; \mathbb{R}^m))$.

An analogous estimate holds true also for operator \tilde{L} defined in (5), but for any $s \in]-\gamma, 0[$ and under the condition $\beta T_* < \gamma + s$.

Some remarks on the previous statement are in order.

r:no-lower

Remark 2.5. (i) The technical limitation $|s| < \gamma$ is dictated by product continuity properties (see Proposition 3.9 below). The same can be said about the conditions on the time T_* .

(ii) A careful but easy inspection of our proof reveals that, if $B \in L^\infty([0, T_*]; LL(\mathbb{R}^n))$, then Theorem 2.4 holds true replacing γ by 1 (see also Corollary 3.10). On the other hand, having additional regularity for B does not help to improve the result: in particular, this is the case when the operator is homogeneous of first order, i.e. if $B \equiv 0$.

(iii) In the case of operator \tilde{L} , the Hölder regularity of the 0-th order term imposes an additional limitation on the lifespan T_* . We notice that this is coherent with what is known for scalar wave equations (see e.g. [8]).

From the previous theorem, we can deduce the existence and uniqueness of a local in time solution to the Cauchy problem associated to L and \tilde{L} .

t:global_e

Theorem 2.6. *Let us consider the first-order system (3), and assume it to be microlocally symmetrizable, in the sense of Definition 2.1. Suppose moreover that the coefficients $(A_j)_{1 \leq j \leq n}$ and the symmetrizer S satisfy the boundedness and log-Lipschitz conditions (6)-(7). Suppose also that the coefficient B verifies hypotheses (6)-(8), for some $\gamma \in]0, 1[$.*

For $s \in]0, \gamma[$, let $\beta > 0$ and $T_ > 0$ respectively the loss parameter and the existence time given by Theorem 2.4. Set $s_0 := s - \beta T_*$.*

Then, for any fixed $u_0 \in H^{s_0}(\mathbb{R}^n; \mathbb{R}^m)$ and $f \in L^1([0, T_]; H^s(\mathbb{R}^n; \mathbb{R}^m))$, there exists a unique solution $u \in \mathcal{C}([0, T_*]; H^{s_0}(\mathbb{R}^n; \mathbb{R}^m))$ to the Cauchy problem*

eq:Cauchy

$$(10) \quad \begin{cases} Lu = f \\ u|_{t=0} = u_0, \end{cases}$$

which satisfies the energy inequality (9). In particular, for any $t \in [0, T_*]$, one has

$$u(t) \in H^{s-\beta t}(\mathbb{R}^n; \mathbb{R}^m),$$

and the map $t \mapsto u(t)$ is continuous between the respective functional spaces.

An analogous statement still holds true for the conservative operator \tilde{L} .

r:Cauchy

Remark 2.7. Let us denote by $C_b^\infty(\mathbb{R}^n)$ the space of $C^\infty(\mathbb{R}^n)$ functions which are uniformly bounded with all their derivatives.

Theorems 2.4 and 2.6 immediately imply the following statement: if the coefficients of L (respectively \tilde{L}) and the symmetrizer S are $L^\infty([0, T]; C_b^\infty(\mathbb{R}^n))$, with the first order coefficients and S log-Lipschitz continuous in time, then the Cauchy problem for L (respectively \tilde{L}) is well-posed in the space $H^\infty(\mathbb{R}^n; \mathbb{R}^m)$, with a finite loss of derivatives.

2.2 Local in space results

ss:local_th

We turn now our attention to the local in space problem. For simplicity of exposition, we will always consider *smooth* bounded domains and manifolds.

Hence, we fix a smooth open bounded domain $\Omega \subset \mathbb{R}^{1+n}$, and we suppose that the coefficients $A_j = A_j(z)$ are log-Lipschitz in Ω (keep in mind Definition 1.1).

We recall also that $H^s(\Omega)$ is defined (see Chapter 3 of [26], where the more general context of Besov spaces is treated) as

def:H~s

$$(11) \quad H^s(\Omega) := \left\{ u \in \mathcal{D}'(\Omega) \mid \text{there exists } \tilde{u} \in H^s(\mathbb{R}^{1+n}) \text{ such that } \tilde{u}|_\Omega \equiv u \right\},$$

endowed with the norm

$$\|u\|_{H^s(\Omega)} := \inf \left\{ \|\tilde{u}\|_{H^s(\mathbb{R}^{1+n})} \mid \tilde{u} \in H^s(\mathbb{R}^{1+n}) \text{ and } \tilde{u}|_\Omega \equiv u \right\}.$$

As for the boundary $\partial\Omega$, partition of unity leads to a similar definition for $H^s(\partial\Omega)$, by use of local charts and extension operator; the same can be said for smooth submanifolds Σ .

Let us consider the operator P , defined in Ω by the formula

eq:def_P

$$(12) \quad P(z, \partial_z)u := \sum_{j=0}^n A_j(z) \partial_{z_j} u + B(z)u,$$

where, for all $1 \leq j \leq n$, the A_j 's and B are real $m \times m$ matrices; we will specify later on their regularity. For the time being, let us introduce the principal symbol P_1 of P , identified by the formula

$$P_1(z, \zeta) = \sum_{j=0}^n i \zeta^j A_j(z),$$

and recall some basic definitions (see Section 4 of [20]). At this level, the dependence of the coefficients on the variable $z \in \Omega$ is not really important, so let us omit it from the notations here.

d:hyperbolic

Definition 2.8. The operator P is said to be *hyperbolic* in the direction $\nu \in \mathbb{R}^{1+d}$ if the following conditions are verified:

- (i) $\det(P_1(\nu)) \neq 0$;
- (ii) there exists a $\eta_0 > 0$ such that $\det(P_1(i\tau\nu + \zeta) + B) \neq 0$ for all $\zeta \in \mathbb{R}^{1+d}$ and all $\tau \in \mathbb{R}$ such that $|\tau| > \eta_0$.

The principal operator $P_1(\partial_z)$ is *strongly hyperbolic* in the direction ν if for all matrix $B \in \mathcal{M}_m(\mathbb{R})$, then $P_1 + B$ is hyperbolic in the direction ν .

Let us recall that Proposition 4.2 of [20] gives a characterization of the strong hyperbolicity. We do not enter into the details here; however, we will come back to this notion in a while.

Now, we are interested in considering the dependence on $z \in \Omega$: more precisely, for any $z \in \Omega$, we assign a direction $\nu(z) \in \mathbb{S}^d$ (where \mathbb{S}^d is the unitary sphere in \mathbb{R}^{1+d}) in a smooth way. In what follows, we are going to assume that

(H-1) $P_1(z, \cdot)$ is *uniformly strongly hyperbolic* in the direction $\nu(z)$, for all $z \in \Omega$.

By Theorem 4.10 of [20], hypothesis **(H-1)** is equivalent to the following conditions:

1. one has the property

$$\boxed{\text{eq:det_pos}} \quad (13) \quad C_\Omega := \inf_{z \in \Omega} \left| \det P_1(z, \nu(z)) \right| > 0;$$

2. P_1 admits a bounded family of *full symmetrizers* $\mathbf{S}(z, \cdot)$ which is *uniformly positive* in the direction $\nu(z)$.

Hence, let us recall Definition 4.7 of [20] here below.

Definition 2.9. A bounded family of *full symmetrizers* for $P_1(z, \zeta)$ is a family of $m \times m$ matrices $\mathbf{S}(z, \zeta)$, homogeneous of degree 0 in $\zeta \neq 0$, such that the following conditions are satisfied:

- uniform boundedness: there exists a constant $\Lambda > 0$ such that $\sup_{(z, \zeta)} |\mathbf{S}(z, \zeta)|_{\mathcal{M}} \leq \Lambda$;
- symmetrizability: for any (z, ζ) , the matrix $\mathbf{S}(z, \zeta) P_1(z, \zeta)$ is self-adjoint.

The symmetrizer $\mathbf{S}(z, \zeta)$ is *positive* in the direction $\nu \neq 0$ if there exists a constant $\lambda > 0$ such that, for all $\zeta \neq 0$, one has

$$v \in \text{Ker } P_1(z, \zeta) \quad \implies \quad \text{Re} \left(\mathbf{S}(z, \zeta) P_1(z, \nu) v \cdot v \right) \geq \lambda |v|^2.$$

Uniform positivity means positivity of $\mathbf{S}(z, \cdot)$ in the direction $\nu(z)$ for all $z \in \Omega$, for a constant λ independent of z .

We are going to need also the following assumptions concerning the family of symmetrizers:

(H-2) the map $\zeta \mapsto \mathbf{S}(z, \zeta)$ is \mathcal{C}^∞ for $\zeta \neq 0$ (smoothness in ζ);

(H-3) the map $z \mapsto \mathbf{S}(z, \zeta)$ is uniformly *LL* in Ω (log-Lipschitz regularity in z).

As for the coefficients of the operator P , defined in (12), we suppose instead that:

(H-4) the matrices A_j have coefficients in the *LL*(Ω) class;

(H-5) B has coefficients in the Hölder space $\mathcal{C}^\gamma(\Omega)$.

Let us now fix a smooth hypersurface $\Sigma \subset \Omega$. As in [11], up to shrink Ω , we can assume that Σ is defined by the equation $\varphi = 0$, for a smooth φ such that $d\varphi \neq 0$. Finally, we suppose that, for any $z \in \Sigma$, the vector $\nu(z)$ coincides with the normal to Σ in z , i.e. $d\varphi(z)$.

We introduce the notations $\Omega_{\geq} := \Omega \cap \{\varphi \geq 0\}$ and $\Omega_{>} := \Omega \cap \{\varphi > 0\}$. For $s \in \mathbb{R}$, we say that $u \in H_{loc}^s(\Omega_{\geq})$ if, for any open $\Omega' \subset \Omega$, relatively compact in Ω , the restriction of u to $\Omega' \cap \{\varphi > 0\}$ belongs to $H^s(\Omega_{>})$. In a similar way, we say that $v \in H_{comp}^s(\Omega_{\geq})$ if $v \in H^s(\Omega_{>})$ has compact support in Ω_{\geq} .

For a $z_0 \in \Sigma$, we are interested in solving the Cauchy problem for P in a neighborhood of z_0 . We have the following *local existence* result.

t:local_e

Theorem 2.10. Let $1/2 < \gamma < 1$ and $s \in]1 - \gamma, \gamma[$. Let P be the operator defined in (12), satisfying all hypotheses from (H-1) to (H-5).

Fix a neighborhood ω of z_0 in Σ . Then, there exist a $s_0 \in]1 - \gamma, s[$ and a neighborhood Ω_0 of z_0 in Ω such that, for any $u_0 \in H^{s_0}(\omega)$ and any $f \in H^s(\Omega_0 \cap \{\varphi > 0\})$, there exists a solution $u \in H^{s_0}(\Omega_0 \cap \{\varphi > 0\})$ to the Cauchy problem

$$(CP) \quad \begin{cases} Pu = f \\ u|_{\Sigma} = u_0. \end{cases}$$

r:sense

Remark 2.11. Notice that, due to the low regularity of the coefficients, the meaning of (CP) is not clear a priori: we will precise it in Subsection 5.1. In particular, one step of the proof is devoted exactly to giving sense to the trace operator in a weak smoothness framework.

For the problem (CP) we have also *local uniqueness* of a solution. From this statement, which establish propagation of zero across any space-like hypersurface, one can deduce also further local results, for instance about finite propagation speed and domain of dependence, by use of classical arguments (see e.g. [15] and [22]).

t:local_u

Theorem 2.12. Let $1/2 < \gamma < 1$ and $s \in]1 - \gamma, \gamma[$. Let P be the operator defined in (12), satisfying all hypotheses from (H-1) to (H-5).

If $u \in H^s(\Omega_{>})$ satisfies (CP) with $f = 0$ and Cauchy datum $u_0 = 0$, then $u \equiv 0$ on a neighborhood of z_0 in Ω_{\geq} .

In Section 5, after rigorously justify the good formulation of the Cauchy problem (CP), we will show the invariance of our hypotheses with respect to smooth changes of variables. This fact will enable us to pass in (t, x) coordinates: then, the proof of Theorems 2.10 and 2.12 will be deduced from the global in space results. Therefore, let us focus first on these latter properties.

3 Tools from Littlewood-Paley theory

s:tools

We collect here some notions and results which turn out to be useful in our proof. First, by use of the Littlewood-Paley decomposition, we describe some properties of Sobolev spaces and of log-Lipschitz functions. Then, we recall some notions of Paradifferential Calculus, focusing on operators whose symbol is in the log-Lipschitz class.

3.1 Dyadic analysis of Sobolev and log-Lipschitz classes

ss:L-P

Let us first define the so called *Littlewood-Paley decomposition* in \mathbb{R}^N (for any $N \geq 1$), based on a non-homogeneous dyadic partition of unity with respect to the Fourier variable. We refer to [2] (Chapter 2) and [19] (Chapters 4 and 5) for the details.

So, fix a smooth radial function χ supported in the ball $B(0, 2) \subset \mathbb{R}^N$, equal to 1 in a neighborhood of $B(0, 1)$ and such that $r \mapsto \chi(re)$ is nonincreasing over \mathbb{R}_+ for all unitary vectors $e \in \mathbb{R}^N$. Set $\varphi(\xi) = \chi(\xi) - \chi(2\xi)$ and $\varphi_j(\xi) := \varphi(2^{-j}\xi)$ for all $j \geq 0$.

The dyadic blocks $(\Delta_j)_{j \in \mathbb{Z}}$ are defined by¹

$$\Delta_j := 0 \text{ if } j \leq -1, \quad \Delta_0 := \chi(D) \quad \text{and} \quad \Delta_j := \varphi(2^{-j}D) \text{ if } j \geq 1.$$

We also introduce the following low frequency cut-off:

$$S_j u := \chi(2^{-j}D) u = \sum_{k \leq j} \Delta_k u \quad \text{for } j \geq 0.$$

By use of the previous operators, for any $u \in \mathcal{S}'$, we have the Littlewood-Paley decomposition of u : namely, the equality $u = \sum_j \Delta_j u$ holds true in \mathcal{S}' .

Let us recall the fundamental *Bernstein's inequalities*.

¹Throughout we agree that $f(D)$ stands for the pseudo-differential operator $u \mapsto \mathcal{F}^{-1}(f(\xi)\mathcal{F}u(\xi))$.

l:bern

Lemma 3.1. *Let $0 < r < R$. A constant C exists so that, for any nonnegative integer k , any couple (p, q) in $[1, +\infty]^2$ with $p \leq q$ and any function $u \in L^p$, we have, for all $\lambda > 0$,*

$$\begin{aligned} \text{supp } \widehat{u} \subset B(0, \lambda R) &\implies \|\nabla^k u\|_{L^q} \leq C^{k+1} \lambda^{k+N(\frac{1}{p}-\frac{1}{q})} \|u\|_{L^p}; \\ \text{supp } \widehat{u} \subset \{\xi \in \mathbb{R}^N \mid r\lambda \leq |\xi| \leq R\lambda\} &\implies C^{-k-1} \lambda^k \|u\|_{L^p} \leq \|\nabla^k u\|_{L^p} \leq C^{k+1} \lambda^k \|u\|_{L^p}. \end{aligned}$$

Let us now introduce the class of *logarithmic Sobolev spaces*, which naturally come into play in the study of hyperbolic operators with low regularity coefficients (see [11], [6] and [7]). Let us set $\Pi(D) := \log(2 + |D|)$, i.e. its symbol is $\pi(\xi) := \log(2 + |\xi|)$.

d:log-H~s

Definition 3.2. For all $(s, \alpha) \in \mathbb{R}^2$, we define the space $H^{s+\alpha \log}(\mathbb{R}^N)$ as $\Pi^{-\alpha} H^s(\mathbb{R}^N)$, i.e.

$$f \in H^{s+\alpha \log} \iff \Pi^\alpha f \in H^s \iff \pi^\alpha(\xi) (1 + |\xi|^2)^{s/2} \widehat{f}(\xi) \in L^2.$$

Obviously, for $\alpha = 0$ one recovers the classical Sobolev space H^s .

We have the following dyadic characterization of these classes (see [19], Proposition 4.1.11), which generalizes the classical property for the H^s scale.

p:log-H

Proposition 3.3. *Let $s, \alpha \in \mathbb{R}$. Then $u \in \mathcal{S}'$ belongs to the space $H^{s+\alpha \log}$ if and only if:*

(i) *for all $k \in \mathbb{N}$, $\Delta_k u \in L^2(\mathbb{R}^N)$;*

(ii) *set $\delta_k := 2^{ks} (1+k)^\alpha \|\Delta_k u\|_{L^2}$ for all $k \in \mathbb{N}$, the sequence $(\delta_k)_k$ belongs to $\ell^2(\mathbb{N})$.*

Moreover, $\|u\|_{H^{s+\alpha \log}} \sim \|(\delta_k)_k\|_{\ell^2}$.

The previous proposition can be summarized by the equivalence $H^{s+\alpha \log} \equiv B_{2,2}^{s+\alpha \log}$, where, for any $(s, \alpha) \in \mathbb{R}^2$ and $1 \leq p, r \leq +\infty$, the *non-homogeneous logarithmic Besov space* $B_{p,r}^{s+\alpha \log}$ is the subset of tempered distributions u for which

eq:log-Besov

$$(14) \quad \|u\|_{B_{p,r}^{s+\alpha \log}} := \left\| (2^{js} (1+j)^\alpha \|\Delta_j u\|_{L^p})_{j \in \mathbb{N}} \right\|_{\ell^r} < +\infty.$$

In addition, we point out that an analogous characterization holds true also for Hölder classes: namely, for any $\gamma \in]0, 1[$ one has $\mathcal{C}^\gamma \equiv B_{\infty,\infty}^\gamma = B_{\infty,\infty}^{\gamma+0 \log}$.

We recall that the previous definition and properties do not depend on the choice of the cut-off functions used in a Littlewood-Paley decomposition: in the logarithmic framework, this comes from Lemma 3.5 of [9], which we recall here.

l:log-B_ind

Lemma 3.4. *Let $\mathcal{C} \subset \mathbb{R}^d$ be a ring, $(s, \alpha) \in \mathbb{R}^2$ and $(p, r) \in [1, +\infty]^2$. Let $(u_j)_{j \in \mathbb{N}}$ be a sequence of smooth functions such that*

$$\text{supp } \widehat{u}_j \subset 2^j \mathcal{C} \quad \text{and} \quad \left\| (2^{js} (1+j)^\alpha \|u_j\|_{L^p})_{j \in \mathbb{N}} \right\|_{\ell^r} < +\infty.$$

Then $u := \sum_{j \in \mathbb{N}} u_j$ belongs to $B_{p,r}^{s+\alpha \log}$ and $\|u\|_{B_{p,r}^{s+\alpha \log}} \leq C_{s,\alpha} \left\| (2^{js} (1+j)^\alpha \|u_j\|_{L^p})_{j \in \mathbb{N}} \right\|_{\ell^r}$.

This fact will be used freely throughout the paper.

Now we mention a couple of results which will be useful in the sequel. They are classical (see [2], Chapter 2), and their extension to the logarithmic framework is proved in [9].

l:log-S_j

Lemma 3.5. *Fix $(s, \alpha) \in \mathbb{R}^2$ and let $u \in \mathcal{S}'$ given.*

(i) *If the sequence $(2^{js} (1+j)^\alpha \|S_j u\|_{L^2})_{j \in \mathbb{N}}$ belongs to ℓ^2 , then $u \in H^{s+\alpha \log}$ and*

$$\|u\|_{H^{s+\alpha \log}} \leq C \left\| (2^{js} (1+j)^\alpha \|S_j u\|_{L^2})_{j \in \mathbb{N}} \right\|_{\ell^2},$$

for some constant $C > 0$ depending only on s and α , but not on u .

(ii) Suppose $u \in H^{s+\alpha \log}$, with $s < 0$. Then the sequence $(2^{js} (1+j)^\alpha \|S_j u\|_{L^2})_{j \in \mathbb{N}} \in \ell^2$, and

$$\left\| (2^{js} (1+j)^\alpha \|S_j u\|_{L^2})_{j \in \mathbb{N}} \right\|_{\ell^2} \leq \tilde{C} \|u\|_{H^{s+\alpha \log}},$$

for some constant $\tilde{C} > 0$ depending only on s and α .

Observe that, in general, the second property fails in the endpoint case $s = 0$. Indeed, for $s = 0$ one can only infer, for any $\alpha \leq 0$,

$$\left\| \left((1+j)^\alpha \|S_j u\|_{L^2} \right)_{j \in \mathbb{N}} \right\|_{\ell^\infty} \leq \tilde{C} \|u\|_{B_{2,1}^{0+\alpha \log}}.$$

The second result we want to mention is a sort of “dual” version of the previous lemma.

Lemma 3.6. *Let \mathcal{B} be a ball of \mathbb{R}^N and take $s > 0$ and $\alpha \in \mathbb{R}$. Let $(u_j)_{j \in \mathbb{N}}$ be a sequence of smooth functions such that*

$$\text{supp } \hat{u}_j \subset 2^j \mathcal{B} \quad \text{and} \quad (2^{js} (1+j)^\alpha \|u_j\|_{L^2})_{j \in \mathbb{N}} \in \ell^2.$$

Then the function $u := \sum_{j \in \mathbb{N}} u_j$ belongs to the space $H^{s+\alpha \log}$, and there exists a constant C , depending only on s and α , such that

$$\|u\|_{H^{s+\alpha \log}} \leq C \left\| (2^{js} (1+j)^\alpha \|u_j\|_{L^2})_{j \in \mathbb{N}} \right\|_{\ell^2}.$$

Once again, the previous statement is not true in the endpoint case $s = 0$: then, one can just infer, for any $\alpha \geq 0$,

$$\|u\|_{B_{2,\infty}^{0+\alpha \log}} \leq C \left\| \left((1+j)^\alpha \|u_j\|_{L^p} \right)_{j \in \mathbb{N}} \right\|_{\ell^1}.$$

We now turn our attention to the study of the class of *log-Lipschitz functions*. We have given the general definition in Definition 1.1; now, we restrict to the case $\Omega = \mathbb{R}^N$, for some $N \geq 1$.

Let us recall some properties which can be deduced by use of dyadic decomposition (see [10] and [11] for the proof), and which are true in any dimension $N \geq 1$.

Proposition 3.7. *There exists a positive constant C such that, for all $a \in LL(\mathbb{R}^N)$ and all integers $k \geq 0$, we have*

$$\begin{aligned} \|\Delta_k a\|_{L^\infty} &\leq C (k+1) 2^{-k} \|a\|_{LL} \\ \|a - S_k a\|_{L^\infty} &\leq C (k+1) 2^{-k} \|a\|_{LL} \\ \|S_k a\|_{\text{Lip}} := \|S_k a\|_{L^\infty} + \|\nabla S_k a\|_{L^\infty} &\leq C (k+1) \|a\|_{LL}. \end{aligned}$$

Remark 3.8. By Proposition 3.3 of [10], the last property is a characterization of the space LL .

We conclude this part by showing continuity properties of multiplication of Sobolev distributions by Hölder-type functions.

Proposition 3.9. *Let $b \in B_{\infty,\infty}^{\gamma+\varrho \log}$, where $\gamma > 0$ and $\varrho \in \mathbb{R}$, or $\gamma = 0$ and $\varrho > 1$. Then the multiplication operator $u \mapsto bu$ is a continuous self-map of $H^{s+\alpha \log}(\mathbb{R}^N)$ if:*

- $|s| < \gamma$, no matter the value of $\alpha \in \mathbb{R}$;
- $s = \gamma$ and $\alpha < \varrho - 1/2$, or $s = -\gamma$ and $\alpha > 1/2 - \varrho$.

Proof. We use Bony's paraproduct decomposition (see [3], [2] and [19]) to write

$$bu = T_b u + T_u b + R(b, u) = T_b u + T'_u b,$$

where the previous operators are defined by the formulas

$$(15) \quad T_b u = \sum_j S_{j-3} b \Delta_j u, \quad R(b, u) = \sum_j \sum_{|j-k| \leq 3} \Delta_j b \Delta_k u \quad T'_u a = \sum_j S_{j+3} u \Delta_j b.$$

We remark that the conditions on γ and ϱ imply the chain of embeddings $B_{\infty, \infty}^{\gamma+\varrho \log} \hookrightarrow B_{\infty, 1}^0 \hookrightarrow L^\infty$. Then, by classical properties of paraproduct, we immediately have that $T_b u \in H^{s+\alpha \log}$, with the estimate $\|T_b u\|_{H^{s+\alpha \log}} \leq C \|b\|_{L^\infty} \|u\|_{H^{s+\alpha \log}}$.

Let us now focus on the case $s > 0$ (and then $\gamma > 0$), and let us consider the operator $T'_u b$. For any $j \in \mathbb{N}$, using Bernstein inequalities and definition (14), we deduce

$$\|S_{j+3} u \Delta_j b\|_{L^2} \leq \|S_{j+3} u\|_{L^2} \|\Delta_j b\|_{L^\infty} \leq C 2^{-j\gamma} (j+1)^{-\varrho} \|u\|_{H^{s+\alpha \log}} \|b\|_{B_{\infty, \infty}^{\gamma+\varrho \log}}.$$

Since the term $S_{j+3} u \Delta_j b$ is supported in dyadic balls $2^j \mathcal{B}$, we conclude by use of Lemma 3.6. We point out that, for $s = \gamma$, the condition $\alpha < \varrho - 1/2$ is needed to have the right-hand side of the previous inequality in ℓ^2 .

For $s \leq 0$ (and then γ can be taken even 0), instead, we employ the finer decomposition in $T_u b + R(b, u)$. First of all, by Lemma 3.5 and (14) again, we have the estimate

$$\|S_{j-3} u \Delta_j b\|_{L^2} \leq \|S_{j-3} u\|_{L^2} \|\Delta_j b\|_{L^\infty} \leq C \|u\|_{H^{s+\alpha \log}} \|b\|_{B_{\infty, \infty}^{\gamma+\delta \log}} 2^{-j(\gamma+s)} (1+j)^{-(\alpha+\varrho)} \zeta_j,$$

for some $(\zeta_j)_j \in \ell^2$ of unitary norm. Observe that, for $s = 0$, we have no more the presence of $(\zeta_j)_j$, but the right-hand side still belongs to ℓ^2 , thanks to our hypotheses on γ , α and δ . Since the generic term $S_{j-3} u \Delta_j b$ is supported in dyadic rings $2^j \mathcal{C}$, from Lemma 3.4 we infer $\|T_u b\|_{H^{s+\alpha \log}} \leq C \|u\|_{H^{s+\alpha \log}} \|b\|_{B_{\infty, \infty}^{\gamma+\varrho \log}}$.

For the remainder term, we use again Lemma 3.4: focusing just on the ‘‘diagonal’’ term $R_0(b, u) = \sum_j \Delta_j b \Delta_j u$ (the other ones being similar), we have to bound, for any $\nu \geq 0$, the quantity $2^{s\nu} (1+\nu)^\alpha \|\Delta_\nu R_0(b, u)\|_{L^2}$. By use of Proposition 3.3 we have

$$\|\Delta_\nu R_0(b, u)\|_{L^2} \leq \sum_{j \geq \nu-3} \|\Delta_j b \Delta_j u\|_{L^2} \leq C \|u\|_{H^{s+\alpha \log}} \|b\|_{B_{\infty, \infty}^{\gamma+\varrho \log}} \sum_{j \geq \nu-3} 2^{-(\gamma+s)j} (1+j)^{-(\alpha+\varrho)} \zeta_j,$$

where the sequence $(\zeta_j)_j$ is as above. Then, by Cauchy-Schwarz inequality we get

$$2^{s\nu} (1+\nu)^\alpha \|\Delta_\nu R_0(b, u)\|_{L^2} \leq C \|u\|_{H^{s+\alpha \log}} \|b\|_{B_{\infty, \infty}^{\gamma+\varrho \log}} \times \\ \times 2^{s\nu} (1+\nu)^\alpha \left(\sum_{j \geq \nu-3} 2^{-2(\gamma+s)j} (1+j)^{-2(\varrho+\alpha)} \right)^{1/2},$$

and this completes the proof of the case $s \leq 0$, and therefore of the proposition. \square

We point out that, by Proposition 3.7, one has $LL \hookrightarrow B_{\infty, \infty}^{1-\log}$. Hence, from the previous statement, we immediately infer continuity properties of multiplication by log-Lipschitz functions, which generalize Proposition 3.5 of [10] to the framework of logarithmic Sobolev spaces.

c:LL-H \tilde{s}

Corollary 3.10. *Let $a \in LL(\mathbb{R}^N)$. Then the multiplication operator $u \mapsto au$ is a continuous map of $H^{s+\alpha \log}(\mathbb{R}^N)$ into itself if:*

- $|s| < 1$, no matter the value of $\alpha \in \mathbb{R}$;
- $s = 1$ and $\alpha < -3/2$, or $s = -1$ and $\alpha > 3/2$.

3.2 Paradifferential Calculus in the log-Lipschitz class

In our study, we need to resort to tools from Paradifferential Calculus, as introduced by J.-M. Bony in the celebrated paper [3]. We refer to [2] (Chapter 2) for a complete treatment, and to papers [18]-[21] for a construction depending on parameters. Here, we follow the approach of [19] (see Chapters 4 and 5).

The first part of this section is devoted to recall basic properties: we adapt the classical construction to consider symbols having log-Lipschitz smoothness with respect to the space variable, and we define general paradifferential operators associated to them, for which we develop also a symbolic calculus. In the final part, we consider the case of time-dependent symbols, which are log-Lipschitz in t : at this point, time cannot be considered as a parameter anymore, and we need to establish some properties for paradifferential operators whose symbols belong to this class.

3.2.1 Symbols having log-Lipschitz regularity

Fix a cut-off function $\psi \in \mathcal{C}^\infty(\mathbb{R}^N \times \mathbb{R}^N)$ which verifies the following properties:

- there exist $0 < \varepsilon_1 < \varepsilon_2 < 1$ such that

$$\psi(\eta, \xi) = \begin{cases} 1 & \text{for } |\eta| \leq \varepsilon_1 (1 + |\xi|) \\ 0 & \text{for } |\eta| \geq \varepsilon_2 (1 + |\xi|) ; \end{cases}$$

- for all $(\beta, \alpha) \in \mathbb{N}^N \times \mathbb{N}^N$, there exists a constant $C_{\beta, \alpha} > 0$ such that

$$\left| \partial_\eta^\beta \partial_\xi^\alpha \psi(\eta, \xi) \right| \leq C_{\beta, \alpha} (1 + |\xi|)^{-|\alpha| - |\beta|} .$$

For instance, it is easy to verify (see Ex. 5.1.5 [19]) that

$$\psi(\eta, \xi) \equiv \psi_{-3}(\eta, \xi) := \sum_{k=0}^{+\infty} \chi_{k-3}(\eta) \varphi_k(\xi) ,$$

where χ and φ are the localization (in phase space) functions associated to a Littlewood-Paley decomposition, satisfies the previous requirements.

Define now G^ψ as the inverse Fourier transform of ψ with respect to the variable η :

$$G^\psi(x, \xi) := (\mathcal{F}_\eta^{-1} \psi)(x, \xi) .$$

We have the following result (see Lemma 5.1.7 of [19]).

1:G Lemma 3.11. *For all $(\beta, \alpha) \in \mathbb{N}^N \times \mathbb{N}^N$, there exist constants $C_{\beta, \alpha} > 0$ such that:*

$$\begin{aligned} \left\| \partial_x^\beta \partial_\xi^\alpha G^\psi(\cdot, \xi) \right\|_{L^1(\mathbb{R}_x^N)} &\leq C_{\beta, \alpha} (1 + |\xi|)^{-|\alpha| + |\beta|} \\ \left\| |\cdot| \log \left(1 + \frac{1}{|\cdot|} \right) \partial_x^\beta \partial_\xi^\alpha G^\psi(\cdot, \xi) \right\|_{L^1(\mathbb{R}_x^N)} &\leq C_{\beta, \alpha} (1 + |\xi|)^{-|\alpha| + |\beta| - 1} \log(2 + |\xi|) . \end{aligned}$$

Let us now take a symbol $a = a(x, \xi)$: thanks to G^ψ , we can smooth it out in the space variable, and then define the paradifferential operator associated to a as the pseudodifferential operator related to this smooth function.

First of all, let us specify the class of symbols we are interested in.

d:symbols Definition 3.12. Let $X \subset L^\infty(\mathbb{R}^N)$ a Banach space and fix $(m, \delta) \in \mathbb{R}^2$.

- (i) We denote by $\Gamma_X^{m+\delta \log}$ the space of functions $a(x, \xi)$ which are locally bounded over $\mathbb{R}^N \times \mathbb{R}^N$, of class \mathcal{C}^∞ with respect to ξ and which satisfy the following property: for all $\alpha \in \mathbb{N}^N$ and all $\xi \in \mathbb{R}^N$, the map $x \mapsto \partial_\xi^\alpha a(x, \xi)$ belongs to X , and, for some $C_\alpha > 0$,

$$\|\partial_\xi^\alpha a(\cdot, \xi)\|_X \leq C_\alpha (1 + |\xi|)^{m-|\alpha|} \log^\delta(2 + |\xi|).$$

- (ii) $\Sigma_X^{m+\delta \log}$ is the space of symbols $\sigma \in \Gamma_X^{m+\delta \log}$ which satisfy the following spectral condition: there exists a $0 < \epsilon < 1$ such that, for all $\xi \in \mathbb{R}^N$, the spectrum of the function $x \mapsto \sigma(x, \xi)$ is contained in the ball $\{|\eta| \leq \epsilon(1 + |\xi|)\}$.

In a quite natural way, we can equip $\Gamma_X^{m+\delta \log}$ with the family of seminorms

$$\boxed{\text{eq:L-inf-sem}} \quad (16) \quad \|a\|_{(X,k)}^{(m,\delta)} := \sup_{|\alpha| \leq k} \sup_{\mathbb{R}_\xi^N} \left((1 + |\xi|)^{-m+|\alpha|} \log^{-\delta}(2 + |\xi|) \|\partial_\xi^\alpha a(\cdot, \xi)\|_X \right).$$

Typically, $X = L^\infty(\mathbb{R}^N)$ or $X = LL(\mathbb{R}^N)$ for us. In the former case, for convenience we will use the notations $\Gamma_\infty^{m+\delta \log}$, $\Sigma_\infty^{m+\delta \log}$ and $\|\cdot\|_{(\infty,k)}^{(m,\delta)}$. In the final part of the present section (see the end of Paragraph 3.2.2), for the sake of generality we will consider also the case $X = B_{\infty,\infty}^{\gamma+\varrho \log}(\mathbb{R}^N)$ (recall (14) for its definition).

In the particular case $X = LL$, we explicitly notice the following fact: for $a \in \Gamma_{LL}^{m+\delta \log}$, there exists $K > 0$ such that, for all $\xi \in \mathbb{R}^N$ and all $y \in \mathbb{R}^N \setminus \{0\}$, one has

$$\boxed{\text{eq:LL-sem}} \quad (17) \quad \sup_{\mathbb{R}_x^N} |a(x+y, \xi) - a(x, \xi)| \leq K (1 + |\xi|)^m \log^\delta(2 + |\xi|) |y| \log\left(1 + \frac{1}{|y|}\right).$$

Hence we can set $|a|_{LL} = |a|_{(LL,0)}^{(m,\delta)}$ to be the smallest constant K such that the previous inequality holds true. In a quite natural way, we can also define the LL seminorms $|a|_{(LL,k)}^{(m,\delta)}$.

When $X = B_{\infty,\infty}^{\gamma+\varrho \log}$, instead, we set $\Gamma_{\gamma+\varrho \log}^{m+\delta \log} := \Gamma_X^{m+\delta \log}$. Moreover, introducing a Littlewood-Paley decomposition $(\Delta_\nu)_{\nu \geq 0}$ in the x -variable (not in ξ), we have

$$\boxed{\text{eq:B-sem}} \quad (18) \quad \sup_{\mathbb{R}_x^N} |\Delta_\nu a(\cdot, \xi)| \leq K (1 + |\xi|)^m \log^\delta(2 + |\xi|) 2^{-\gamma\nu} (1 + \nu)^{-\varrho},$$

for a constant $K > 0$, for all $\xi \in \mathbb{R}^N$ and all $\nu \geq 0$. Once again, in a natural way we can introduce the seminorms $\|a\|_{(\gamma+\varrho \log, k)}^{(m,\delta)}$.

Finally, we explicitly point out that, by spectral localization and Paley-Wiener Theorem, a symbol $\sigma \in \Sigma_X^{m+\delta \log}$ is smooth also in the x variable.

Now let us consider a symbol $a \in \Gamma_X^{m+\delta \log}$: we can associate to it a the classical symbol according to the formula

$$\boxed{\text{classical-symb}} \quad (19) \quad \sigma_a^\psi(x, \xi) := (\psi(D_x, \xi) a)(x, \xi) = \left(G^\psi(\cdot, \xi) *_x a(\cdot, \xi) \right)(x).$$

The following proposition holds true.

Proposition 3.13. *Let $(m, \delta) \in \mathbb{R}^2$.*

- (i) *For $X \subset L^\infty(\mathbb{R}^N)$ a Banach space, the smoothing operator $\mathbb{S} : a(x, \xi) \mapsto \sigma_a(x, \xi)$ maps continuously $\Gamma_X^{m+\delta \log}$ into $\Sigma_X^{m+\delta \log}$.*

- (ii) *For $a \in \Gamma_{LL}^{m+\delta \log}$, then the difference symbol $a - \sigma_a$ belongs to $\Gamma_\infty^{(m-1)+(\delta+1) \log}$.*

(iii) In particular, if ψ_1 and ψ_2 are two admissible cut-off functions, then the difference of the two smoothed symbols, $\sigma_a^{\psi_1} - \sigma_a^{\psi_2}$, belongs to $\Sigma_\infty^{(m-1)+(\delta+1)\log}$.

Proof. The first property is classical. The second one immediately follows from the log-Lipschitz continuity assumption and Lemma 3.11. The last statement is a straightforward consequence of the previous ones. \square

We conclude this part by noting that, at this level, the time variable can be treated as a parameter in the construction. In particular, the previous properties still hold true for symbols in $L^\infty([0, T]; \Gamma_X^{m+\delta\log})$.

3.2.2 Operators, symbolic calculus

As announced above, we can use the previous construction to associate to any symbol $a(x, \xi)$, which is non-regular in x , an operator: this will be the pseudodifferential operator associated to the smooth symbol σ_a^ψ .

Let us formalize the discussion: for $a \in \Gamma_\infty^{m+\delta\log}$, we define the paradifferential operator associated to a via the formula

$$T_a^\psi := \sigma_a^\psi(x, D_x) : u \mapsto T_a^\psi u(x) = \frac{1}{(2\pi)^N} \int_{\mathbb{R}_\xi^N} e^{ix \cdot \xi} \sigma_a^\psi(x, \xi) \widehat{u}(\xi) d\xi.$$

r:p-prod

Remark 3.14. Notice that, if $f = f(\xi)$ is a Fourier multiplier, then $T_f \equiv f(D_x)$ (see e.g. [19]).

Let us also point out that if $a = a(x) \in L^\infty$ and if we take the cut-off function ψ_{-3} , then T_a^ψ is actually the classical paraproduct operator, defined first in [3].

Let us recall some basic definitions and properties in Paradifferential Calculus. The corresponding proofs in the logarithmic setting are analogous to the classical case, and they are not detailed here.

d:op_order

Definition 3.15. We say that an operator P is of order $m + \delta \log$ if, for every $(s, \alpha) \in \mathbb{R}^2$, P maps $H^{s+\alpha\log}$ into $H^{(s-m)+(\alpha-\delta)\log}$ continuously.

With slight modifications to the proof of Proposition 2.9 of [21], stated for the classical Sobolev class, we get the next fundamental result.

l:action

Lemma 3.16. For all $\sigma \in \Sigma_\infty^{m+\delta\log}$, the corresponding operator $\sigma(\cdot, D_x)$ is of order $m + \delta \log$.

The following result is an immediate consequence of the previous lemma and Proposition 3.13.

t:action

Theorem 3.17. Given a symbol $a \in \Gamma_\infty^{m+\delta\log}$, for any admissible cut-off function ψ , the operator T_a^ψ is of order $m + \delta \log$.

Notice that, in Lemma 3.16 and Theorem 3.17, the LL hypothesis is not necessary: these results hold true if the symbol is even just L^∞ with respect to x . On the contrary, we are going to exploit the additional regularity in space in the next result. It states that the whole construction does not depend on the cut-off function ψ .

p:act-psi

Proposition 3.18. If ψ_1 and ψ_2 are two admissible cut-off functions and $a \in \Gamma_{LL}^{m+\delta\log}$, then the difference $T_a^{\psi_1} - T_a^{\psi_2}$ is of order $(m-1) + (\delta+1)\log$.

Therefore, changing the cut-off function ψ doesn't change the paradifferential operator associated to a , up to lower order terms. So, in what follows we fix the cut-off function $\psi = \psi_{-3}$ (defined in Paragraph 3.2.1) and we will miss out the dependence of σ_a and T_a on it.

We want now to develop symbolic calculus in the LL class. A preliminary result is in order: it can be viewed as a generalization of Proposition 3.7.

l:11-symb

Lemma 3.19. *Let $a \in \Gamma_{LL}^{m+\delta \log}$, and denote by σ_a the classical symbol associated to it via formula (19). Then the following estimates hold true:*

$$\begin{aligned} |\partial_\xi^\alpha \sigma_a| &\leq C_\alpha (1 + |\xi|)^{m-|\alpha|} \log^\delta (2 + |\xi|) \\ \left| \partial_x^\beta \partial_\xi^\alpha \sigma_a \right| &\leq C_{\beta,\alpha} (1 + |\xi|)^{m-|\alpha|+|\beta|-1} \log^{\delta+1} (2 + |\xi|). \end{aligned}$$

The constants C_α just depend on the quantities $\|a\|_{\infty,k}^{(m,\delta)}$ defined in (16), where $|\alpha| \leq k$.

The constants $C_{\beta,\alpha}$, instead, depend only on the quantities $\|a\|_{LL,k}^{(m,\delta)}$, where again $|\alpha| \leq k$.

The proof of the previous result is somehow classical, as it follows the same lines of Lemma 3.16 in [6] and Lemma 3.15 in [7]. Therefore we omit it.

From Lemma 3.19 we immediately deduce the following properties.

Theorem 3.20. (i) *Let us take two symbols $a \in \Gamma_{LL}^{m+\delta \log}$ and $b \in \Gamma_{LL}^{n+\varrho \log}$ and denote by T_a , T_b the respective associated paradifferential operators. Then one has*

$$T_a \circ T_b = T_{ab} + R_o.$$

The principal part T_{ab} is of order $(m+n) + (\delta + \varrho) \log$.

The remainder operator R_o has order $(m+n-1) + (\delta + \varrho + 1) \log$.

(ii) *Let $a \in \Gamma_{LL}^{m+\delta \log}$. The adjoint operator (over L^2) of T_a is given by the formula*

$$(T_a)^* = T_{\bar{a}} + R_*.$$

The order of $T_{\bar{a}}$ is still $m + \delta \log$.

The remainder operator R_* has order $(m-1) + (\delta + 1) \log$.

The last statement of this paragraph is a fundamental parilinearization result, in the general instance $X = B_{\infty,\infty}^{\gamma+\varrho \log}$, which will allow us to treat both the first and lower order terms of operator L in energy estimates. In order to give sense to all terms, we have to restrict to differential operators, and then operator T_a reduces to the classical paraproduct operator (keep in mind also Remark 3.14).

t:paralin

Theorem 3.21. *Let $m \in \mathbb{N}$, and $\eta \in \mathbb{N}^N$ of length $|\eta| = m$. Take a pair $(\gamma, \varrho) \in \mathbb{R}^2$ such that $\gamma \geq 0$, and consider a function $a \in B_{\infty,\infty}^{\gamma+\varrho \log}$, introduced in (14). Define the difference operator $\mathcal{D} : u \mapsto a \partial_x^\eta u - T_a \partial_x^\eta u = a \partial_x^\eta u - T_{a \xi^\eta} u$.*

(i) *If $s > m$, then \mathcal{D} maps continuously $H^{s+\alpha \log}$ into $H^{\gamma+(\varrho-h) \log}$, for any $h > 1/2$.*

(ii) *If $s = m$ and $\alpha \geq 0$, the previous statement remains true.*

(iii) *For any $s \in]m - \gamma, m[$ (and then $\gamma > 0$) and any $\alpha \in \mathbb{R}$, \mathcal{D} maps continuously $H^{s+\alpha \log}$ into $H^{\sigma+h \log}$, where $\sigma = s - m + \gamma$ and $h = \alpha - \delta + \varrho$.*

The norms of the operators just depend on the quantity $\|a\|_{B_{\infty,\infty}^{\gamma+\varrho \log}}$.

Proof. We start by noticing that, according to (15), $\mathcal{D}(u)$ can be rewritten as

$$\mathcal{D}(u) = \sum_{\nu \geq 0} \partial_x^\eta (S_{\nu+3} u) \Delta_\nu a = T'_{\partial_x^\eta u} a = T_{\partial_x^\eta u} a + R(\partial_x^\eta u, a).$$

First of all, we focus on the case of high regularity, i.e. $s > m$, or $s = m$ and $\alpha \geq 0$. This in particular implies that $\partial_x^\eta u \in L^2$, with $\|\partial_x^\eta u\|_{L^2} \leq C \|u\|_{H^{s+\alpha \log}}$. Therefore, we can estimate

$$\|S_{\nu-3} \partial_x^\eta u \Delta_\nu a\|_{L^2} \leq C \|S_{\nu-3} \partial_x^\eta u\|_{L^2} \|\Delta_\nu a\|_{L^\infty} \leq C \|u\|_{H^{s+\alpha \log}} \|a\|_{B_{\infty,\infty}^{\gamma+\varrho \log}} 2^{-\gamma \nu} (1 + \nu)^{-\varrho}.$$

Thanks to Lemma 3.4, we immediately deduce that $T_{\partial_x^\eta u} a \in H^{\gamma+(\varrho-h)\log}$ for all $h > 1/2$.

As for R , once again we can focus just on the diagonal terms $R_0(\partial_x^\eta u, a) = \sum_\nu \partial_x^\eta \Delta_\nu u \Delta_\nu a$. For any $k \in \mathbb{N}$, let us estimate

$$\begin{aligned} \|\Delta_k R_0(\partial_x^\eta u, a)\|_{L^2} &\leq C \sum_{\nu \geq k-3} \|\Delta_\nu \partial_x^\eta u\|_{L^2} \|\Delta_\nu a\|_{L^\infty} \\ &\leq C \|u\|_{H^{s+\alpha \log}} \|a\|_{B_{\infty, \infty}^{\gamma+\varrho \log}} \sum_{\nu \geq k-3} 2^{\nu(m-s-\gamma)} (1+\nu)^{-\alpha-\varrho} \zeta_\nu \\ &\leq C \|u\|_{H^{s+\alpha \log}} \|a\|_{B_{\infty, \infty}^{\gamma+\varrho \log}} \left(\sum_{\nu \geq k-3} 2^{2\nu(m-s-\gamma)} (1+\nu)^{-2(\alpha+\varrho)} \right)^{1/2}, \end{aligned}$$

where, as usual, $\|\zeta_\nu\|_{\ell^2} = 1$. Therefore, by Lemma 3.4 again, the remainder term $R(\partial_x^\eta u, a) \in H^{(s-m+\gamma)+(\alpha+\varrho-h)\log}$ for all $h > 1/2$, which is in particular included in $H^{\gamma+(\varrho-h)\log}$. When $s = m$, the same speech holds true for any $\alpha \geq 0$.

Now let us turn our attention to the case $m - \gamma < s < m$ (and then $\gamma > 0$): for this, we consider directly the operator $T'_{\partial_x^\eta u} a$. Since $s - m < 0$, we can apply Lemma 3.5 and estimate

$$\|S_{\nu+3} \partial_x^\eta u \Delta_\nu a\|_{L^2} \leq C \|u\|_{H^{s+\alpha \log}} 2^{\nu(m-s)} (1+\nu)^{-\alpha} \zeta_\nu \|a\|_{B_{\infty, \infty}^{\gamma+\varrho \log}} 2^{-\nu\gamma} (1+\nu)^{-\varrho},$$

where the sequence $(\zeta_\nu)_\nu \in \ell^2$ is as above. Thanks to this inequality and the fact that $s - m + \gamma > 0$ by hypothesis, Lemma 3.6 implies that $T'_{\partial_x^\eta u} a$ belongs to $H^{(s-m+\gamma)+(\alpha+\varrho)\log}$.

The proof of the theorem is now completed. \square

We notice the following fact: if $a \in W^{1, \infty}$ is Lipschitz, it is well-known that $au - T_a u \in H^1$ for any $u \in L^2$ (see e.g. Theorem 5.2.8 of [19]). On the contrary, if we applied our result with $m = s = \alpha = \varrho = 0$ and $\gamma = 1$, we would get just $au - T_a u \in H^{1-(1/2+\delta)\log}$ for any $\delta > 0$.

Motivated by this consideration, let us make a remark.

r:paralin

Remark 3.22. Attaining the limit case $h = 0$ in points (i), (ii), and $s = m - \gamma$ in point (iii), would require further technical extensions of the theory, in the same spirit of Paragraph 5.2.4 of [19], to functions a in logarithmic Hölder classes $B_{\infty, \infty}^{\gamma+\varrho \log}$, where $\varrho \neq 0$.

However, these adaptations go beyond the scopes of the present paper, and we decided not to address these issues here in order to keep the presentation as concise as possible. Indeed, due to product properties (see Proposition 3.9) and loss of derivatives in the energy estimates, we will always be far away from these endpoint cases, and the previous statement turns to be enough for our scopes.

Before going further, let us remark that, at this level, time can be treated once again as a parameter in the construction: for a symbol $a = a(t, x, \xi) \in L^\infty([0, T]; \Gamma_\infty^{m+\delta \log})$ and $u \in \mathcal{S}'([0, T] \times \mathbb{R}^d)$, we set

$$(T_a u)(t, \cdot) := T_{a(t, \cdot, \xi)} u(t, \cdot).$$

As a consequence, all the properties stated above still hold true for the time-dependent operator, at any time $t \in [0, T]$ fixed. Indeed, only regularity in space is used at this level.

3.2.3 The case of symbols which are log-Lipschitz in time

sss:LL-t

Now, we get closer to our hypotheses, and we introduce a new class of symbols, by imposing additional regularity in the time variable.

Definition 3.23. Let $Y \subset L^\infty([0, T])$ a Banach space. For $(m, \delta) \in \mathbb{R}^2$, we define $Y_T(\Gamma_X^{m+\delta \log})$ as the class of symbols $a(t, x, \xi) \in L^\infty([0, T]; \Gamma_X^{m+\delta \log})$ such that, for almost every $(x, \xi) \in \mathbb{R}^N \times \mathbb{R}^N$, the map $t \mapsto a(t, x, \xi)$ belongs to Y .

d:symbol_t

Analogously, we define the class $Y_T(\Sigma_X^{m+\delta \log})$ if, moreover, the spectral condition (in x) of Definition 3.12 is verified for almost every $(t, \xi) \in [0, T] \times \mathbb{R}^N$.

We omit the subscript T whenever $T = +\infty$.

For us Y_T will be always $L^\infty([0, T])$ or $LL([0, T])$.

In particular, when $Y = LL$, the previous definition implies that there exists a $C_0 > 0$ for which, for all $(x, \xi) \in \mathbb{R}^N \times \mathbb{R}^N$ and all $0 < \tau \leq T/2$, one has

$$|a(t + \tau, x, \xi) - a(t, x, \xi)| \leq C_0 (1 + |\xi|)^m \log^\delta(2 + |\xi|) \tau \log\left(1 + \frac{1}{\tau}\right).$$

As before, for such a symbol a we can then define the seminorm $|a|_{LL_t}$. Of course, higher order seminorms (related to derivatives in ξ) can be defined, but they will not be used in our study, so that we prefer to limit the presentation to this case.

r:LL **Remark 3.24.** We remark that the LL continuity of the symbol, separately with respect to time and space variables, is enough to our scopes. This will be evident from our computations.

In order to perform energy estimates, for a symbol a as above we need to introduce a regularization in time. So, take an even function $\rho \in C_0^\infty(\mathbb{R})$, $0 \leq \rho \leq 1$, whose support is contained in the interval $[-1, 1]$ and such that $\int \rho(t) dt = 1$, and define the mollifier kernel

$$\rho_\varepsilon(t) := \frac{1}{\varepsilon} \rho\left(\frac{t}{\varepsilon}\right) \quad \forall \varepsilon \in]0, 1].$$

Let us fix a symbol $a \in LL_T(\Gamma_X^{m+\delta \log})$; if $T < +\infty$, we extend this symbol out of $[0, T]$ (for instance, by taking the constant values at the extremities of the interval), in such a way to get a new symbol (which we will still denote by a) in the class $LL(\Gamma_X^{m+\delta \log})$. Now, we smooth a out setting, for all $\varepsilon \in]0, 1]$,

eq:a~e (20)
$$a_\varepsilon(t, x, \xi) := (\rho_\varepsilon *_t a(\cdot, x, \xi))(t) = \int_{\mathbb{R}_s} \rho_\varepsilon(t - s) a(s, x, \xi) ds.$$

Then, we have the following estimates (see e.g. [4], [10]).

l:LL-reg **Lemma 3.25.** *Let $a \in LL(\Gamma_X^{m+\delta \log})$. Then $(a_\varepsilon)_\varepsilon$ is a bounded family in the same space. In particular, Proposition 3.13 still holds true for a_ε , uniformly in $\varepsilon \in]0, 1]$.*

In addition, there exist constants $C > 0$ such that, for all $\varepsilon \in]0, 1]$ and for all $(t, x, \xi) \in \mathbb{R} \times \mathbb{R}^N \times \mathbb{R}^N$, one has

$$\begin{aligned} |a_\varepsilon(t, x, \xi) - a(t, x, \xi)| &\leq C |a|_{LL_t} \varepsilon \log\left(1 + \frac{1}{\varepsilon}\right) (1 + |\xi|)^m \log^\delta(2 + |\xi|) \\ |\partial_t a_\varepsilon(t, x, \xi)| &\leq C |a|_{LL_t} \log\left(1 + \frac{1}{\varepsilon}\right) (1 + |\xi|)^m \log^\delta(2 + |\xi|). \end{aligned}$$

In the course of the proof, it will be fundamental to link the approximation parameter ε with the dual variable, following the original idea of [4] (see also [6], [7], [8], [9], [10], [11], [25]). More precisely, in [4] the authors took $\varepsilon = 1/|\xi|$: here, we make an analogous choice, but replacing $|\xi|$ by $\langle \xi \rangle := (1 + |\xi|^2)^{1/2}$ (we need the new symbol to be defined for all $\xi \in \mathbb{R}^N$).

p:Y-tilde **Proposition 3.26.** *For $a \in Y_T(\Gamma_X^{m+\delta \log})$, where $Y \subset L^\infty([0, T])$ is a Banach space, we define the new function*

$$\tilde{a}(t, x, \xi) := a_{1/\langle \xi \rangle}(t, x, \xi).$$

Then \tilde{a} is still a symbol in the class $Y_T(\Gamma_X^{m+\delta \log})$, which is actually smooth in time.

Proof. Proposition 3.22 of [9] proves that the convolution acts as an operator of order $0 + 0 \log$ in the time variable. We are going to show that this is true also with respect to ξ .

First of all, we notice that it is enough to consider the case $a \in L_T^\infty(\Gamma_\infty^{0+0 \log})$: regularity in time comes from Lemma 3.25, while the x variable does not play any role at this level.

For $\alpha = 0$, the estimate of Definition 3.12 is immediate. Hence, let us consider the case $|\alpha| = 1$: by definitions we have

$$\begin{aligned} \partial_\xi^\alpha \tilde{a}(t, x, \xi) &= \partial_\xi^\alpha \langle \xi \rangle \int \rho((t-s)\langle \xi \rangle) a(s, x, \xi) ds + \\ &+ \langle \xi \rangle \partial_\xi^\alpha \langle \xi \rangle \int \rho'((t-s)\langle \xi \rangle) (t-s) a(s, x, \xi) ds + \langle \xi \rangle \int \rho((t-s)\langle \xi \rangle) \partial_\xi^\alpha a(s, x, \xi) ds. \end{aligned}$$

The estimates for the first and last term are straightforward: thanks to the uniform bounds

$$|\partial_\xi^\alpha a(t, x, \xi)| \leq C_\alpha (1 + |\xi|)^{-|\alpha|},$$

(recall that we are taking $a \in L_T^\infty(\Gamma_\infty^{0+0 \log})$ only), we have to control

$$\int \rho(t\langle \xi \rangle) dt \leq C / \langle \xi \rangle,$$

where we used the change of variable $\tau = \langle \xi \rangle t$ and fact that $\int \rho \equiv 1$.

For the second term, we argue as above: by Young inequality, we are reconducted to consider the L^1 norm of the function $|\rho'(\langle \xi \rangle t)| \langle \xi \rangle |t|$. Performing the same change of variable as before, we get

$$\langle \xi \rangle \int |\rho'(t\langle \xi \rangle)| |t| dt \leq C / \langle \xi \rangle,$$

and this concludes the proof of the proposition. \square

Let us now focus on the case $Y = LL$. Up to extend the symbols out of $[0, T]$ as explained above, without loss of generality we can focus just on the case $T = +\infty$.

Proposition 3.27. *Let $a \in LL(\Gamma_X^{m+\delta \log})$. Then the following properties hold true.*

p:LL-tilde

- (i) *One has $a - \tilde{a} \in L^\infty(\Gamma_X^{(m-1)+(\delta+1) \log})$ and $\partial_t \tilde{a} \in L^\infty(\Gamma_X^{m+(\delta+1) \log})$.*
- (ii) *The smoothed (in x) symbol $\sigma_{\tilde{a}}$, defined by formula (19), belongs to $LL(\Sigma_X^{m+\delta \log})$, and one has the ‘‘commutation’’ formula $\tilde{\sigma}_a = \sigma_{\tilde{a}}$. So, we adopt the notation $\tilde{\sigma}_a$ for it.*
- (iii) *The smoothing operator (in time and space) $\tilde{S} : a(t, x, \xi) \mapsto \tilde{\sigma}_a(t, x, \xi)$ maps continuously $LL(\Gamma_X^{m+\delta \log})$ into $LL(\Sigma_X^{m+\delta \log})$.*
- (iv) *The symbol $\tilde{\sigma}_a$ still satisfies the estimates of Lemma 3.19, uniformly in t . In addition, one has $\sigma_a - \tilde{\sigma}_a \in L^\infty(\Sigma_X^{(m-1)+(\delta+1) \log})$ and $\sigma_{\partial_t \tilde{a}} = \partial_t \tilde{\sigma}_a \in L^\infty(\Sigma_X^{m+(\delta+1) \log})$, with the estimates, for some constants C just depending on $|a|_{LL_t}$,*

$$\begin{aligned} |\sigma_a(t, x, \xi) - \tilde{\sigma}_a(t, x, \xi)| &\leq C (1 + |\xi|)^{m-1} \log^{\delta+1}(2 + |\xi|) \\ |\partial_t \tilde{\sigma}_a(t, x, \xi)| &\leq C (1 + |\xi|)^m \log^{\delta+1}(2 + |\xi|), \end{aligned}$$

and analogous formula for the higher order derivatives in ξ .

Proof. Assertion (i) comes from estimates in Lemma 3.25, while point (ii) is a direct consequence of formula (19) and definition of \tilde{a} (the integrals of the convolutions in time and space commute between themselves).

Next, we notice that $\tilde{\mathbb{S}}$ is a composition of the self-map $a \mapsto \tilde{a}$ of $LL(\Gamma_X^{m+\delta \log})$ with the map $\mathbb{S} : \tilde{a} \mapsto \tilde{\sigma}_a$ which goes from the previous space into $LL(\Sigma_X^{m+\delta \log})$: sentence (iii) follows.

Let us focus on assertion (iv): it basically relies on point (i) of the present proposition, and in particular on the fact that the convolution acts as an operator of order 0. Then, the first statement is immediate, the second one is implied by the linearity of operator \mathbb{S} together with Lemma 3.25. The formula for the time derivative derives directly from the definitions. Finally, for the estimates one has to use Lemma 3.25 again. \square

Now, given a symbol $a \in Y(\Gamma_X^{m+\delta \log})$, we define the paradifferential operator

$$(21) \quad \tilde{T}_a := T_{\tilde{a}} = \tilde{\sigma}_a(\cdot, D_x).$$

Notice that properties (ii), (iii) of Proposition 3.27 still hold true even if we replace LL by a generic Y , and in particular when $Y = L^\infty$. Therefore, fixing $X = LL(\mathbb{R}^N)$, from Theorem 3.20 we immediately get the following result.

Theorem 3.28. (i) If $a \in L^\infty(\Gamma_{LL}^{m+\delta \log})$, then the operator \tilde{T}_a is of order $m + \delta \log$.

(ii) For $a \in L^\infty(\Gamma_{LL}^{m+\delta \log})$ and $b \in L^\infty(\Gamma_{LL}^{n+\varrho \log})$, one has $\tilde{T}_a \circ \tilde{T}_b = T_{\tilde{a}\tilde{b}} + \tilde{R}_\circ$, where \tilde{R}_\circ is of order $(m + n - 1) + (\delta + \varrho + 1) \log$.

(iii) For $a \in L^\infty(\Gamma_{LL}^{m+\delta \log})$, denote by \tilde{T}_a^* the adjoint of \tilde{T}_a over L^2 . Then $\tilde{T}_a^* = \tilde{T}_a + \tilde{R}_*$, where \tilde{R}_* is of order $(m - 1) + (\delta + 1) \log$.

Remark 3.29. Thanks to point (iv) of Proposition 3.27, in point (ii) above we can substitute $T_{\tilde{a}\tilde{b}}$ by $\tilde{T}_{a\tilde{b}}$, up to another remainder which is still of order $(m + n - 1) + (\delta + \varrho + 1)$. But this would require regularity in time, that we do not want to use at this level.

In the same way, we can see that also the analogous of Theorem 3.21 holds true for the operator $\tilde{\mathcal{D}} := a - \tilde{T}_a$, whenever $a \in L^\infty(\Gamma_{\gamma+\varrho \log}^{m+\delta \log})$.

Finally, we take $Y = LL$ and $X = L^\infty$ and we exploit regularity in time. Still by use of Proposition 3.27, we get the following statement.

Theorem 3.30. Let $a \in LL(\Gamma_\infty^{m+\delta \log})$. Then the next properties are true:

- the operator $T_a - \tilde{T}_a$ is a remainder of order $(m - 1) + (\delta + 1) \log$;
- one has $[\partial_t, \tilde{T}_a] = T_{\partial_t \tilde{a}}$, and this is an operator of order $m + (\delta + 1) \log$.

4 Well-posedness on the whole \mathbb{R}^n

Thanks to the tools developed in the previous section, we are now able to prove Theorem 2.4, i.e. energy estimates for the global in space problem. The first part of the present section is devoted to this. In the estimates, we will keep track of the dependence of the different constants on the log-Lipschitz seminorms of the coefficients of the operator and of the symmetrizer.

In the final part (see Subsection 4.3), we will show how to derive Theorem 2.6 from the bounds of Theorem 2.4.

4.1 The energy

Let us start by defining the energy associated to our operator L . Roughly speaking, denoting by S a microlocal symmetrizer for \mathcal{A} , the leading idea is that the paradifferential operator T_S is an approximated symmetrizer for $iT_{\mathcal{A}}$, which represents the principal part of the operator L .

However, some ‘‘corrections’’ are needed. Indeed, on the one hand we need to smooth out the coefficients with respect to time in order to be able to perform energy estimates. On the other

hand, S is homogeneous of degree 0 in ξ , and in particular it can be singular in $\xi = 0$: therefore, we have to cut off the low frequencies, borrowing somehow an idea from paradifferential calculus with parameters (see e.g. [18], [21]).

So, let us proceed in the following way. First of all, given a symmetrizer $S(t, x, \xi)$ for our system, we smooth it out with respect to the time variable, according to the formula (20). Notice that the approximated symmetrizer S_ε still satisfies

$$0 < \lambda \text{Id} \leq S_\varepsilon(t, x, \xi) \leq \Lambda \text{Id}$$

for any $\varepsilon \in]0, 1]$ and all $(t, x, \xi) \in [0, T] \times \mathbb{R}_x^n \times (\mathbb{R}_\xi^n \setminus \{0\})$. In particular, the matrix symbol $S_\varepsilon^{1/2}$ is well-defined.

Furthermore, let us immediately set (accordingly with the original choice of [4])

$$\varepsilon = 1/|\xi| \quad \forall |\xi| \geq 1.$$

Notice that, since $|\xi| \geq 1$, it makes no special difference to take $|\xi|$ or $\langle \xi \rangle$. Hence, in what follows we will adopt the notations introduced in Paragraph 3.2.3.

Finally, let $\theta \in \mathcal{C}_0^\infty(\mathbb{R}^n)$ such that $0 \leq \theta \leq 1$, $\theta \equiv 1$ in the ball $B(0, 1)$ and $\theta \equiv 0$ for $|\xi| \geq 2$. For any $\mu > 0$, we set $\theta_\mu(\xi) := \theta(\mu^{-1}\xi)$, and we denote

$$\tilde{\Sigma}(t, x, \xi) := \tilde{S}^{1/2}(t, x, \xi) (1 - \theta_\mu(\xi)).$$

For all $(s, \alpha) \in \mathbb{R}^2$ fixed, we then define the quantities

$$E_{s, \alpha}[u] := \left\| \tilde{T}_\Sigma u \right\|_{H^{s+\alpha \log}}^2 + \|\theta_\mu(D_x)u\|_{H^{s+\alpha \log}}^2,$$

where \tilde{T}_Σ is the paradifferential operator associated to the just defined matrix symbol $\tilde{\Sigma}$, according to the construction explained in Subsection 3.2 (see in particular Paragraph 3.2.3).

We now establish positivity estimates: we show that this property does not depend on the considered Sobolev norm.

1:energy

Lemma 4.1. (i) *There exists a $C_0 > 0$ (just depending on the constant K_0 appearing in condition (6) for S) such that, for any $(s, \alpha) \in \mathbb{R}^2$ and any $u \in \mathcal{S}(\mathbb{R}^n)$,*

$$E_{s, \alpha}[u] \leq C_0 \|u\|_{H^{s+\alpha \log}}.$$

(ii) *There exists $\mu_0 \geq 2$ such that, for all $\mu \geq \mu_0$, the following property holds true: there exists a $C_\mu > 0$ for which, for all smooth $u \in \mathcal{S}(\mathbb{R}^n)$,*

$$\|u\|_{H^{s+\alpha \log}} \leq C_\mu E_{s, \alpha}[u].$$

The constant C_μ depends just on μ , on K_0 and on K_1 (recall conditions (6) and (7) for S); in particular, it is independent of (s, α) .

Proof. The first property is immediate, once noticing that $\tilde{\Sigma}$ is an operator of order $0 + 0 \log$, in the sense of Definition 3.15.

So, let us focus on the second inequality: according to the decomposition

$$u = (1 - \theta_\mu(D_x))u + \theta_\mu(D_x)u,$$

it is enough to prove it for the high frequency component $(1 - \theta_\mu(D_x))u$.

First of all, let us define $\psi_\mu(\xi) := (1 - \theta_\mu(\xi))(1 - \theta(\xi))$: by the properties of the support of θ , we easily infer that, for any $\mu \geq 2$, one has $\psi_\mu \equiv (1 - \theta_\mu)$. Therefore, denoting

$$\tilde{\Xi}(t, x, \xi) := \tilde{S}^{-1/2}(t, x, \xi) (1 - \theta(\xi)),$$

we deduce the equality $\tilde{\Xi} \tilde{\Sigma} = \text{Id}$ $\psi_\mu = \text{Id} (1 - \theta_\mu)$, which in turn gives, by symbolic calculus (recall Theorem 3.28),

$$\text{eq:pos_high} \quad (22) \quad (1 - \theta_\mu(D_x)) u = T_{\tilde{\Xi}\tilde{\Sigma}} u = T_{\tilde{\Xi}} T_{\tilde{\Sigma}} u + \tilde{R} u,$$

where the remainder \tilde{R} is given, at the principal order, by $\partial_\xi \tilde{\Xi} \partial_x \tilde{\Sigma}$. Using again the function ψ_μ defined above and introducing the matrix $\tilde{S}' = \tilde{S}^{1/2} (1 - \theta)$, it is easy to see that

$$\tilde{R} u = \tilde{R}' (1 - \theta_\mu(D_x)) u, \quad \text{with} \quad \tilde{R}' \sim \partial_\xi \tilde{\Xi} \partial_x \tilde{S}'.$$

Let us come back to (22): since $\tilde{\Xi}$ is of order $0 + 0 \log$ and \tilde{R}' is of order $-1 + \log$, we obtain

$$\|(1 - \theta_\mu(D_x)) u\|_{H^{s+\alpha \log}} \leq C_1 \|\tilde{T}_\Sigma u\|_{H^{s+\alpha \log}} + C_2 \|(1 - \theta_\mu(D_x)) u\|_{H^{(s-1)+(\alpha+1) \log}},$$

for some C_1 (depending on $|\tilde{S}|_{L^\infty,0}^{(0,0)}$) and C_2 (depending also on $|\tilde{S}|_{LL,0}^{(0,0)}$) large enough. Now, by spectral properties and Proposition 3.3, we have

$$\begin{aligned} \|(1 - \theta_\mu(D_x)) u\|_{H^{(s-1)+(\alpha+1) \log}}^2 &\leq \sum_{k \geq K_\mu} 2^{2(s-1)k} (1+k)^{2(\alpha+1)} \|\Delta_k (1 - \theta_\mu(D_x)) u\|_{L^2}^2 \\ &\leq 2^{-2K_\mu} (1+K_\mu)^2 \|(1 - \theta_\mu(D_x)) u\|_{H^{s+\alpha \log}}^2, \end{aligned}$$

where $K_\mu \sim \log_2 \mu$. Therefore, if μ is large enough, such that $C_2 2^{-K_\mu} (1+K_\mu) \leq 1/2$, we can absorb this term in the left-hand side of the previous inequality. This completes the proof of the statement (ii), and so also of the lemma. \square

Now, let $s \in]0, 1[$ be fixed. For $\beta > 0$, to be chosen in the course of the proof, we set

$$\text{eq:s} \quad (23) \quad s(t) = s - \beta t;$$

we then define the energies

$$\text{eq:E} \quad (24) \quad E(t) := E_{s(t),0}[u(t)] = \|\tilde{T}_\Sigma u(t)\|_{H^{s(t)}}^2 + \|\theta_\mu(D_x) u(t)\|_{H^{s(t)}}^2$$

$$\text{eq:E_log} \quad (25) \quad E_{\log}(t) := E_{s(t),1/2}[u(t)] = \|\tilde{T}_\Sigma u(t)\|_{H^{s(t)+(1/2) \log}}^2 + \|\theta_\mu(D_x) u(t)\|_{H^{s(t)+(1/2) \log}}^2,$$

where μ is fixed large enough, so that Lemma 4.1 holds true. In particular, we have

$$E(t) \sim \|u(t)\|_{H^{s(t)}}^2 \quad \text{and} \quad E_{\log}(t) \sim \|u(t)\|_{H^{s(t)+(1/2) \log}}^2.$$

r:energy **Remark 4.2.** Let us point out that the ‘‘true’’ energy associated to our operator is E . The second energy E_{\log} is introduced because of a logarithmic loss of derivatives in the estimates for E , due to the log-Lipschitz regularity of the coefficients and of the symmetrizer.

We conclude this part with an approximation result: we make a parilinearization of our operator. In this way, we create remainders, which nonetheless are regular enough.

l:L->T **Lemma 4.3.** *Let L be the operator defined by (3), and \tilde{L} be given by (5). Then*

$$\text{eq:paralin} \quad (26) \quad Lu = \partial_t u + iT_A u + T_B u + \mathcal{R}_L u,$$

where \mathcal{R}_L maps continuously $H^{s+\alpha \log}$ into $H^{s+(\alpha-1) \log}$ for all $0 < s < \gamma$ and all $\alpha \in \mathbb{R}$.

Analogously, we can write $\tilde{L}u = \partial_t u + iT_A u + T_B u + \mathcal{R}_{\tilde{L}} u$, where $\mathcal{R}_{\tilde{L}}$ is continuous from $H^{s+\alpha \log}$ to $H^{s+(\alpha-1) \log}$ for all $-\gamma < s < 0$ and all $\alpha \in \mathbb{R}$.

The operator norms of both \mathcal{R}_L and $\mathcal{R}_{\tilde{L}}$ depend only on the constants K_0, K_1 and K_2 appearing in conditions (6)-(7)-(8).

Proof. Let us focus on operator L first: decomposition (26) is easily obtained, if we set $\mathcal{R}_L = \mathcal{R}_A + \mathcal{R}_B$, where we have defined

$$\mathcal{R}_A u := \sum_{j=1}^n (A_j - T_{A_j}) \partial_j u \quad \text{and} \quad \mathcal{R}_B u := (B - T_B) u.$$

We observe that $\mathcal{R}_A(t, x, \xi) = \sum_j (A_j(t, x) - T_{A_j}(t, x)) \xi_j$ is a symbol in the class $L_T^\infty(\Gamma_{LL}^{1+0 \log})$: since $LL \hookrightarrow B_{\infty, \infty}^{1-\log}$, Theorem 3.21, point (iii), gives us

$$\|\mathcal{R}_A u\|_{H^{s+(\alpha-1)\log}} \leq C \left(\sup_{t \in [0, T]} \sup_{1 \leq j \leq n} \|A_j(t, \cdot)\|_{LL} \right) \|u\|_{H^{s+\alpha \log}}$$

for any $0 < s < 1$ and any $\alpha \in \mathbb{R}$. In particular, this is true for $0 < s < \gamma$.

On the other hand, $\mathcal{R}_B(t, x, \xi)$ belongs to $L_T^\infty(\Gamma_{\gamma+0 \log}^{0+0 \log})$ (actually, it does not even depend on ξ). Then, point (i) of Theorem 3.21 (with e.g. the particular choice $\sigma = s < \gamma$) implies

$$\|\mathcal{R}_B u\|_{H^{s+\alpha \log}} \leq C \left(\sup_{t \in [0, T]} \|B(t, \cdot)\|_{C^\gamma} \right) \|u\|_{H^{s+\alpha \log}}.$$

This inequality completes the proof of the statement for \mathcal{R}_L .

Let us now deal with \tilde{L} : we start by observing that

$$\begin{aligned} \tilde{L}u &= \partial_t u + \sum_{j=1}^n \partial_j \left((A_j - T_{A_j}) u \right) + \sum_{j=1}^n \partial_j (T_{A_j} u) + T_B u + \mathcal{R}_B u \\ &= \partial_t u + i T_A u + \sum_{j=1}^n \partial_j \left((A_j - T_{A_j}) u \right) + \sum_{j=1}^n [\partial_j, T_{A_j}] u + T_B u + \mathcal{R}_B u. \end{aligned}$$

Then, we just set $\mathcal{R}_{\tilde{L}} := \sum_{j=1}^n \partial_j \left((A_j - T_{A_j}) u \right) + \sum_{j=1}^n [\partial_j, T_{A_j}] u + \mathcal{R}_B u$. Let us consider each of its terms one by one.

\mathcal{R}_B is defined as before: this time, we apply item (iii) of Theorem 3.21 and we get, for any $-\gamma < s < 0$ and any $\alpha \in \mathbb{R}$,

$$\|\mathcal{R}_B u\|_{H^{(s+\gamma)+\alpha \log}} \leq C \left(\sup_{t \in [0, T]} \|B(t, \cdot)\|_{C^\gamma} \right) \|u\|_{H^{s+\alpha \log}}.$$

As for the commutator term, we notice that $[\partial_j, T_{A_j}] = \partial_j \sigma_{A_j}(t, x, D_x)$, where $\sigma_{A_j}(t, x, \xi)$ is the classical symbol associated to A_j via formula (19). Therefore, by Lemma 3.19 we deduce that $[\partial_j, T_{A_j}]$ is an operator of order $0 + \log$, and then, for any $(s, \alpha) \in \mathbb{R}^2$,

$$\|[\partial_j, T_{A_j}] u\|_{H^{s+(\alpha-1)\log}} \leq C \left(\sup_{t \in [0, T]} \sup_{1 \leq j \leq n} \|A_j(t, \cdot)\|_{LL} \right) \|u\|_{H^{s+\alpha \log}}.$$

Finally, for any j , the operator $A_j - T_{A_j}$ belongs to the class $L_T^\infty(\Gamma_{LL}^{0+0 \log})$: then again, point (iii) of Theorem 3.21 gives us

$$\begin{aligned} \left\| \partial_j \left((A_j - T_{A_j}) u \right) \right\|_{H^{s+(\alpha-1)\log}} &\leq C \left\| (A_j - T_{A_j}) u \right\|_{H^{(1+s)+(\alpha-1)\log}} \\ &\leq C \left(\sup_{t \in [0, T]} \sup_{1 \leq j \leq n} \|A_j(t, \cdot)\|_{LL} \right) \|u\|_{H^{s+\alpha \log}} \end{aligned}$$

for any $-1 < s < 0$, and in particular for $-\gamma < s < 0$.

This completes the proof of the assertion for $\mathcal{R}_{\tilde{L}}$, and then also of the lemma. \square

4.2 Energy estimates: proof of Theorem 2.4

ss:estimates

We are now ready to compute and estimate the time derivative of the energy. In a first moment, we aim at proving bounds for the parilinearized operator

$$T_L u := \partial_t u + i T_A u + T_B u.$$

Notice that $T_L u = Lu - \mathcal{R}_L u$ and also $T_L u = \tilde{L}u - \mathcal{R}_{\tilde{L}}u$.

In what follows, we will generically denote by C_{LL} a multiplicative constant which depends on the LL norms of the coefficients of $\mathcal{A}(t, x, \xi)$ and of the symmetrizer $S(t, x, \xi)$, i.e. quantities K_0 and K_1 in (6)-(7), but not on s neither on u . On the other hand, we will use the generic symbol C if the constant just depend on K_0 , i.e. the L^∞ bounds of the coefficients and of the symmetrizer, and on K_2 in (8), i.e. the Hölder norms of B . Finally, we will use the notation C_p to denote a constant which depends on the μ fixed for having the positivity estimates of Lemma 4.1.

1:est-T_L

Lemma 4.4. *Let $(A_j)_{1 \leq j \leq n}$ and B be as in the hypotheses of Theorem 2.4. For any $s \in \mathbb{R}$ and any $\beta > 0$, define $s(t)$ by formula (23). Then, for any smooth $u \in \mathcal{S}([0, T] \times \mathbb{R}^n)$ we have*

$$\begin{aligned} \frac{d}{dt} E(t) &\leq C_1 E(t) + (C_2 - C_3 \beta) E_{\log}(t) + \\ &\quad + \|\theta_\mu(D) T_L u\|_{H^{s(t)}} (E(t))^{1/2} + 2 \left| \operatorname{Re} \left(\Lambda^{s(t)}(D) \tilde{T}_\Sigma T_L u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right) \right|. \end{aligned}$$

The constant C_1 just depends on the μ fixed in the positivity estimates and on K_0 ; C_2 still depends on μ and K_0 , and also on K_1 ; finally, C_3 depends only on K_0 .

Proof. We start by noticing that, for any $v \in \mathcal{S}(\mathbb{R}^n)$, we have

$$\frac{d}{dt} \|v\|_{H^{s(t)}}^2 = s'(t) \int_{\mathbb{R}^n} \log(1 + |\xi|^2) (1 + |\xi|^2)^{s(t)} |\widehat{v}(\xi)|^2 d\xi \sim s'(t) \|v\|_{H^{s(t)+(1/2)\log}}^2.$$

For notation convenience, we set $\Lambda(D) := (1 - \Delta)^{1/2}$, i.e. $\Lambda(\xi) = \langle \xi \rangle = (1 + |\xi|^2)^{1/2}$.

Recalling definition (24) of the energy E , we get:

$$\begin{aligned} \frac{d}{dt} E(t) &= s'(t) \operatorname{Re} \left(\log(\Lambda^2(D)) \Lambda^{s(t)}(D) \tilde{T}_\Sigma u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2} + \\ &\quad + 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) T_{\partial_t \tilde{\Sigma}} u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2} + 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) \tilde{T}_\Sigma \partial_t u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2} + \\ &\quad + s'(t) \operatorname{Re} \left(\log(\Lambda^2(D)) \Lambda^{s(t)}(D) \theta_\mu(D) u, \Lambda^{s(t)}(D) \theta_\mu(D) u \right)_{L^2} + \\ &\quad + 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) \theta_\mu(D) \partial_t u, \Lambda^{s(t)}(D) \theta_\mu(D) u \right)_{L^2} = F_1 + F_2 + F_3 + F_4 + F_5. \end{aligned}$$

First of all, let us consider the terms with $s'(t)$: keeping in mind definitions (25) and (23), it is easy to see that, for some constant just depending on the C_0 appearing in Lemma 4.1,

$$\text{est:s'} \quad (27) \quad F_1 + F_4 \leq -C \beta E_{\log}(t).$$

On the other hand, Theorem 3.30 implies

$$\text{est:F_2} \quad (28) \quad |F_2| \leq \left\| T_{\partial_t \tilde{\Sigma}} u \right\|_{H^{s(t)-(1/2)\log}} \left\| \tilde{T}_\Sigma u \right\|_{H^{s(t)+(1/2)\log}} \leq C_{LL} \|u\|_{H^{s+(1/2)\log}} E_{\log}^{1/2} \leq C_{LL} C_p E_{\log}.$$

For both F_3 and F_5 , we have to use the equation for $T_L u$. We start by dealing with the low frequencies term:

$$F_5 = 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) \theta_\mu(D) \partial_t u, \Lambda^{s(t)}(D) \theta_\mu(D) u \right)_{L^2}$$

$$= 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) \theta_\mu(D) (T_L u - i T_{\mathcal{A}} u - T_B u), \Lambda^{s(t)}(D) \theta_\mu(D) u \right)_{L^2}.$$

First of all, Cauchy-Schwarz inequality immediately implies

$$\left| \operatorname{Re} \left(\Lambda^{s(t)}(D) \theta_\mu(D) T_L u, \Lambda^{s(t)}(D) \theta_\mu(D) u \right)_{L^2} \right| \leq \|\theta_\mu(D) T_L u\|_{H^{s(t)}} E^{1/2}.$$

As for the term with $T_{\mathcal{A}}$, we remark that, by spectral localization properties, one has

$$\theta_\mu(D) T_{i\mathcal{A}} u = \sum_{j=1}^n \sum_{k=0}^{k_\mu} S_{k-3}(A_j(t, x)) \Delta_k \partial_j u,$$

for some $k_\mu \sim \log_2 \mu$. Hence, Bernstein inequalities immediately imply

$$\left| 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) \theta_\mu(D) T_{i\mathcal{A}} u, \Lambda^{s(t)}(D) \theta_\mu(D) u \right)_{L^2} \right| \leq C \|u\|_{H^{s(t)}} E^{1/2} \leq C C_p E(t),$$

for a suitable constant C depending just on the L^∞ norms of the A_j 's. Exactly in the same way, we get an analogous estimate for the T_B term. In the end, putting these inequalities together, we deduce the control

$$\boxed{\text{est:F}_5} \quad (29) \quad |F_5| \leq C C_p E(t) + \|\theta_\mu(D) T_L u\|_{H^{s(t)}} (E(t))^{1/2} + C_p C_{LL} E_{\log}(t).$$

Let us consider now the term F_3 , which we rewrite as

$$F_3 = 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) \tilde{T}_\Sigma (T_L u - i T_{\mathcal{A}} u - T_B u), \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2}.$$

We leave the T_L term on one side: it will contribute to the last item appearing in our statement. In addition, the term with T_B is easy to control: since \tilde{T}_Σ and T_B are operators of order 0, we get

$$\left| \operatorname{Re} \left(\Lambda^{s(t)}(D) \tilde{T}_\Sigma T_B u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2} \right| \leq C \|u\|_{H^{s(t)}} E^{1/2} \leq C C_p E(t).$$

So, we have to focus just on the last term,

$$\tilde{F}_3 := 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) \tilde{T}_\Sigma T_{i\mathcal{A}} u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2} :$$

we are going to make a systematic use of symbolic calculus, taking advantage of the properties established in Theorems 3.20 and 3.28.

First of all, passing to the adjoints, we have the equality

$$\tilde{F}_3 = 2 \operatorname{Re} \left(\tilde{T}_\Sigma \Lambda^{2s(t)}(D) \tilde{T}_\Sigma T_{i\mathcal{A}} u, u \right)_{L^2} + 2 \operatorname{Re} \left(R_1 \Lambda^{2s(t)}(D) \tilde{T}_\Sigma T_{i\mathcal{A}} u, u \right)_{L^2},$$

where R_1 is of order $-1 + \log$. On the other hand, by Remark 3.14 we have $\tilde{T}_\Sigma \Lambda^{2s(t)}(D) = \Lambda^{2s(t)}(D) \tilde{T}_\Sigma + R_2$, where the operator R_2 is a remainder of order $2s(t) - 1 + \log$. Therefore, collecting the R_1 and R_2 terms into only one remainder $R(u, u)$, which can be estimated as

$$\boxed{\text{est:remainder}} \quad (30) \quad |R(u, u)| \leq C_{LL} \|u\|_{H^{s(t)+(1/2)\log}} E_{\log}^{1/2} \leq C_p C_{LL} E_{\log}(t),$$

we can write

$$\tilde{F}_3 = 2 \operatorname{Re} \left(\Lambda^{2s(t)}(D) \tilde{T}_\Sigma \tilde{T}_\Sigma T_{i\mathcal{A}} u, u \right)_{L^2} + R(u, u) = 2 \operatorname{Re} \left(\Lambda^{2s(t)}(D) T_{\tilde{\Sigma}^2} T_{i\mathcal{A}} u, u \right)_{L^2} + R(u, u).$$

Notice that, in the second step, we have included into $R(u, u)$ another rest, depending on a remainder R_3 which can be still bounded as in (30).

Now, we use the fact that, by definition, $\tilde{\Sigma}^2 = \tilde{S}(1 - \theta_\mu(\xi))^2$. Since, by Theorem 3.30, the difference operator $S - \tilde{S}$ contributes to the estimates as another remainder, in the sense of inequality (30), we arrive at the identity

$$\tilde{F}_3 = 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) T_{S(1-\theta_\mu)^2} T_{i\mathcal{A}} u, \Lambda^{s(t)}(D) u \right)_{L^2} + R(u, u).$$

Finally, by symbolic calculus again, up to adding another remainder R_4 of order $0 + \log$ (recall that \mathcal{A} is of order 1), we get

$$\tilde{F}_3 = 2 \operatorname{Re} \left(\Lambda^{s(t)}(D) T_{iS\mathcal{A}(1-\theta_\mu)^2} u, \Lambda^{s(t)}(D) u \right)_{L^2} + R(u, u).$$

At this point, we remark that $\operatorname{Re}(iS\mathcal{A}) = 0$, since S is a microlocal symmetrizer for \mathcal{A} : then, keeping in mind that $\operatorname{Re} P = (P + P^*)/2$, by symbolic calculus we deduce that

$$2 \operatorname{Re} (T_{iS\mathcal{A}(1-\theta_\mu)^2}) = R_5,$$

where also R_5 is an operator of order $0 + \log$. In the end, putting all these informations together, we find the estimate

$$\left| \tilde{F}_3 \right| \leq C_{LL} C_p E_{\log}(t),$$

which in turn gives us

$$\boxed{\text{est:F}_3} \quad (31) \quad |F_3| \leq C C_p E(t) + C_p C_{LL} E_{\log}(t) + 2 \left| \operatorname{Re} \left(\Lambda^{s(t)}(D) \tilde{T}_\Sigma T_L u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2} \right|.$$

Therefore, collecting inequalities (27), (28), (29) and (31) completes the proof of the energy estimates for the parilinearized operator T_L . \square

$\boxed{\text{r:s_values}}$ **Remark 4.5.** No restriction on s is needed at this level: its limitations derive just from product rules (see Proposition 3.9) and from the analysis of remainders (see Lemma 4.3 above).

We complete now the proof of Theorem 2.4 in the case of operator L . Operator \tilde{L} will be matter of Remark 4.6 below.

It remains us to deal with the term T_L in the estimates provided by Lemma 4.4. Recall that $T_L u = Lu - \mathcal{R}_L u$.

For the low frequencies term, it is an easy matter to see that

$$\begin{aligned} \|\theta_\mu(D) T_L u\|_{H^{s(t)}} &\leq \|Lu\|_{H^{s(t)}} + C_p \|\mathcal{R}_L u\|_{H^{s(t)-(1/2)\log}} \\ &\leq \|Lu\|_{H^{s(t)}} + C_p (C_{LL} + C) \|u\|_{H^{s(t)+(1/2)\log}} \end{aligned}$$

where C_p is, as usual, a constant which depends on the positivity estimates. The presence of the constant C in the last step is due to lower order terms, i.e. \mathcal{R}_B . Applying once more positivity estimates, we finally arrive to the bound

$$\boxed{\text{est:low-freq}} \quad (32) \quad \|\theta_\mu(D) T_L u\|_{H^{s(t)}} \leq \|Lu\|_{H^{s(t)}} + C_p (C_{LL} + C) E_{\log}(t).$$

We remark here that, up to let C_p depend also on s , the previous inequality holds true for any $s \in \mathbb{R}$: we have no need yet for using the condition $0 < s < \gamma$.

As for the high frequencies term, since \tilde{T}_Σ is an operator of order 0, we have the control

$$\begin{aligned} 2 \left| \operatorname{Re} \left(\Lambda^{s(t)}(D) \tilde{T}_\Sigma T_L u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2} \right| &\leq C \|Lu\|_{H^{s(t)}} E^{1/2} + C \|\mathcal{R}_L u\|_{H^{s(t)-(1/2)\log}} E_{\log}^{1/2} \\ &\leq C \|Lu\|_{H^{s(t)}} E^{1/2} + C (C_{LL} + C) \|u\|_{H^{s(t)+(1/2)\log}} E_{\log}^{1/2}, \end{aligned}$$

where we have used also Lemma 4.3, and hence the restriction on s . Therefore, by positivity estimates again we deduce

est:high-freq

$$(33) \quad 2 \left| \operatorname{Re} \left(\Lambda^{s(t)}(D) \tilde{T}_\Sigma T_L u, \Lambda^{s(t)}(D) \tilde{T}_\Sigma u \right)_{L^2} \right| \leq C \|Lu\|_{H^{s(t)}} E^{1/2} + C_p (C_{LL} + C) E_{\log}(t).$$

Now, putting (32) and (33) into the inequality given by Lemma 4.4, we find

$$\frac{d}{dt} E(t) \leq C_1 E(t) + C_3 \|Lu\|_{H^{s(t)}} (E(t))^{1/2} + (C_4 - \beta) E_{\log}(t),$$

where C_1 is the constant given by Lemma 4.4, $C_3 > 0$ depends on K_0 and K_1 , but just via positivity estimates, and C_4 depends on K_0 , K_1 and K_2 not only via Lemma 4.1, but also via Lemma 4.3.

Now, we chose $\beta > C_4$: then, setting $e(t) := (E(t))^{1/2}$, an application of Gronwall inequality leads us to the estimate

$$e(t) \leq M e^{Qt} \left(e(0) + \int_0^t \|Lu(\tau)\|_{H^{s(\tau)}} d\tau \right),$$

for two positive constants M and Q large enough. This completes the proof of the energy estimate stated in Theorem 2.4, for operator L .

r:adj

Remark 4.6. For operator \tilde{L} one can argue in a completely analogous way. The only difference is the presence of the remainder operator $\mathcal{R}_{\tilde{L}}$ instead of \mathcal{R}_L , as stated in Lemma 4.3. However, since the order of the two operators is the same, it is easy to see that the estimates do not change.

:precise-est

Remark 4.7. As already remarked in [11], we point out that, in fact, our proof gives a more accurate energy estimate: namely,

est:precise-LL

$$(34) \quad \sup_{t \in [0, T_*]} \|u(t)\|_{H^{s-\beta t}} + \left(\int_0^{T_*} \|u(\tau)\|_{H^{s-\beta\tau+(1/2)\log}}^2 d\tau \right)^{1/2} \leq \\ \leq C_1 e^{C_2 T} \left(\|u(0)\|_{H^s} + \int_0^{T_*} \|Lu(\tau)\|_{H^{s-\beta\tau}} d\tau \right)$$

for tempered distributions u as in the statement of Theorem 2.4. Moreover, if the last term is L^2 in time, one can replace it by $\int_0^{T_*} \|Lu(\tau)\|_{H^{s-\beta\tau-(1/2)\log}}^2 d\tau$.

r:t-LL

Remark 4.8. A careful inspection of our proof reveals that we just used the property $A_j \in L^\infty([0, T]; LL(\mathbb{R}^n; \mathcal{M}_m))$, namely only LL regularity in x is exploited for the A_j 's. On the contrary, for the symmetrizer S one needs log-Lipschitz continuity *both* in time and space variables.

Nonetheless, in general regularity of the symmetrizer is dictated by the regularity of the coefficients (keep in mind the discussion in Example 2.3). Furthermore, hypothesis (7) is invariant by change of coordinates, and then suitable for local analysis. This is why we required it.

4.3 Existence and uniqueness of solutions

ss:global_e

In this subsection, we prove Theorem 2.6, namely the existence and uniqueness of solutions to the global Cauchy problem (10). For this, we will exploit in a fundamental way the energy estimates of Theorem 2.4.

We will focus on the case of operator L , defined in (3). The same arguments hold over \tilde{L} , see its definition in (5), rather directly.

4.3.1 Regularity results, uniqueness

As done in [11] (see Definition 2.5 of that paper), for $(\sigma, \alpha) \in \mathbb{R}^2$ and $\beta > 0$ fixed, let us define the following spaces:

- $\mathcal{C}_{\sigma+\alpha \log, \beta}(T)$ is the set of functions v on $[0, T]$, with values in the space of tempered distributions, which verify the property

$$u \in \mathcal{C}([0, t]; H^{\sigma-\beta t+\alpha \log})$$

for all $t \in [0, T]$;

- $\mathcal{H}_{\sigma+\alpha \log, \beta}(T)$ is the set of functions w on $[0, T]$, with values in the space of tempered distributions, such that

$$\Lambda^{\sigma-\beta t}(D) \Pi^\alpha(D) w(t) \in L^2([0, T]; L^2(\mathbb{R}^n; \mathbb{R}^m)),$$

where we used the notations $\Lambda(D) = (1 - \Delta)^{1/2}$ and $\Pi(D) = \log(2 + |D|)$, introduced in the previous paragraphs;

- $\mathcal{L}_{\sigma+\alpha \log, \beta}(T)$ is the set of functions z on $[0, T]$, with values in the space of tempered distributions, with the property

$$\Lambda^{\sigma-\beta t}(D) \Pi^\alpha(D) z(t) \in L^1([0, T]; L^2(\mathbb{R}^n; \mathbb{R}^m)).$$

All these spaces are endowed with the natural norms induced by the conditions here above.

To begin with, we deal with *weak solutions* to the Cauchy problem (10). More precisely, fix $\gamma \in]0, 1[$ and $s \in]0, \gamma[$ as in the statement of Theorem 2.6, and take the positive constants β and T_* given by Theorem 2.4. For $u_0 \in H^s$ and $f \in \mathcal{L}_{s, \beta}(T_*)$, we are going to consider

$$u \in \mathcal{H}_{s, \beta}(T_*)$$

such that the equation $Lu = f$ is satisfied, together with the initial condition $u|_{t=0} = u_0$, in the sense of distributions.

We remark that, for such u , the product $A_j(t, x) \partial_j u$ is well-defined and belongs to the space $\mathcal{H}_{s-1, \beta}(T_*)$ for all j (recall Corollary 3.10), and $B(t, x)u$ belongs to $\mathcal{H}_{s, \beta}(T_*)$ (see Proposition 3.9). Hence, the equation $Lu = f$ makes sense in $\mathcal{D}'([0, T_*[\times \mathbb{R}^n; \mathbb{R}^m)$.

Let us show now that it makes sense also to impose the initial condition $u|_{t=0} = u_0$.

Lemma 4.9. *Let γ and s be as above. Let $u \in \mathcal{H}_{s, \beta}(T_*)$ verify the equation $Lu = f$ in \mathcal{D}' for some $f \in \mathcal{L}_{s, \beta}(T_*)$.*

Then one has $u \in \mathcal{C}_{s-(1/2), \beta}(T_)$. In particular, the trace $u|_{t=0}$ is well-defined in $H^{s-1/2}(\mathbb{R}^n)$, and the initial condition in (10) makes sense.*

Proof. The proof is somehow classical. Let us define $v(t) := u - \int_0^t f(\tau) d\tau$. Then, by hypotheses it easily follows that $\int_0^t f \in \mathcal{C}_{s, \beta}(T_*)$, and hence $v \in \mathcal{H}_{s, \beta}(T_*)$.

On the other hand, from the arguments exposed before, we deduce that $\partial_t v = \sum_j A_j \partial_j u + Bu$ belongs to the space $\mathcal{H}_{s-1, \beta}$. Then, an easy interpolation implies that $v \in \mathcal{C}_{s-(1/2), \beta}$, and therefore so does u . \square

As a result of this lemma and the previous considerations, we have clarified the sense to give to the Cauchy problem (10).

Next, let us state a “weak = strong” type result. Namely, we show that any weak solution u is in fact the limit of a suitable sequence of smooth approximate solutions, in the norm given by the left-hand side of inequality (9). In particular, this fact implies that u enjoys additional smoothness and it satisfies the energy estimates stated in Theorem 2.4.

th:w-s

Theorem 4.10. Fix $\gamma \in]0, 1[$ and $0 < s < \gamma$. Let $\beta > 0$ and T_* be the “loss parameter” and the existence time given by Theorem 2.4, together with the constants C_1 and C_2 .

Let $u_0 \in H^s(\mathbb{R}^n; \mathbb{R}^m)$, $f \in \mathcal{L}_{s,\beta}(T_*)$ and $u \in \mathcal{H}_{s,\beta}(T_*)$ be the corresponding weak solution to the Cauchy problem (10).

Then one has $u \in \mathcal{C}_{s,\beta}(T_*)$, and u satisfies the energy inequality (9).

Proof. For $\varepsilon \in]0, 1]$, let us introduce the smoothing operators $J_\varepsilon(D) := (1 - \varepsilon\Delta)^{-1/2}$, which are Fourier multipliers associated to the symbols $J_\varepsilon(\xi) = (1 + \varepsilon|\xi|^2)^{-1/2}$.

Recall that they are operators of order -1 , they are uniformly bounded from H^σ to $H^{\sigma+1}$ for any $\sigma \in \mathbb{R}$, and one has the strong convergence $J_\varepsilon(D)w \rightarrow w$ in H^σ , for $\varepsilon \rightarrow 0$.

For any $\varepsilon > 0$, we set $u_\varepsilon := J_\varepsilon(D)u$. Then, from the previous properties we deduce that the family $(u_\varepsilon)_\varepsilon$ is uniformly bounded in $\mathcal{H}_{s+1,\beta}(T_*) \hookrightarrow L^2([0, T_*]; H^1(\mathbb{R}^n; \mathbb{R}^m))$. Moreover, we have that $u_\varepsilon \rightarrow u$ in $\mathcal{H}_{s,\beta}(T_*)$ for $\varepsilon \rightarrow 0$.

In addition, if we apply operator $J_\varepsilon(D)$ to the equation $Lu = f$, we deduce that the u_ε 's solve the equation $Lu_\varepsilon = f_\varepsilon + g_\varepsilon$, with $f_\varepsilon = J_\varepsilon(D)f$ and

def:g_eps

$$(35) \quad g_\varepsilon := \sum_{j=1}^n [A_j(t, x), J_\varepsilon(D)] \partial_j u + [B(t, x), J_\varepsilon(D)] u.$$

By Lemma 4.6 of [11], it is easy to infer that $f_\varepsilon \rightarrow f$ in $\mathcal{L}_{s,\beta}(T_*)$ and $g_\varepsilon \rightarrow 0$ in $\mathcal{H}_{s,\beta}(T_*)$, in the limit for $\varepsilon \rightarrow 0$.

From this fact, Theorem 2.4 and Remark 4.7 together, we deduce the energy estimates

$$\begin{aligned} \sup_{t \in [0, T_*]} \|u_\varepsilon(t)\|_{H^{s-\beta t}} + \left(\int_0^{T_*} \|u_\varepsilon(\tau)\|_{H^{s-\beta\tau+(1/2)\log}}^2 d\tau \right)^{1/2} \leq \\ \leq C_1 e^{C_2 T} \left(\|u_\varepsilon(0)\|_{H^s} + \int_0^T (\|f_\varepsilon\|_{H^{s-\beta\tau}} + \|g_\varepsilon\|_{H^{s-\beta\tau}}) d\tau \right). \end{aligned}$$

By linearity of the equations, similar estimates are satisfied also by the difference $\delta u_{\varepsilon,\eta} := u_\varepsilon - u_\eta$, for all (say) $0 < \eta < \varepsilon \leq 1$, up to replace, in the right-hand side, $u_\varepsilon(0)$, f_ε and g_ε respectively by $\delta u_{\varepsilon,\eta}(0) = u_\varepsilon(0) - u_\eta(0)$, $\delta f_{\varepsilon,\eta} = f_\varepsilon - f_\eta$ and $\delta g_{\varepsilon,\eta} = g_\varepsilon - g_\eta$.

Since $\delta u_{\varepsilon,\eta}(0) \rightarrow 0$ in H^s , and so do $\delta f_{\varepsilon,\eta}$ and $\delta g_{\varepsilon,\eta}$ respectively in $\mathcal{L}_{s,\beta}(T_*)$ and in $\mathcal{H}_{s,\beta}(T_*)$, we gather that $(u_\varepsilon)_\varepsilon$ is a Cauchy sequence in $\mathcal{C}_{s,\beta}(T_*)$. Therefore, the limit $u \in \mathcal{H}_{s,\beta}(T_*)$ also belongs to $\mathcal{C}_{s,\beta}(T_*)$, and it verifies the energy estimate (9). \square

The previous theorem immediately implies uniqueness of weak solutions.

c:w_uniq

Corollary 4.11. Let γ and s be fixed as in the hypotheses of Theorem 4.10. Let $u \in \mathcal{H}_{s,\beta}(T_*)$ be a weak solution to the Cauchy problem (10), with initial datum $u_0 = 0$ and external force $f \equiv 0$.

Then $u \equiv 0$.

sss:weak

4.3.2 Existence of weak solutions

In order to complete the proof of Theorem 2.6, it remains us to prove existence of weak solutions.

w_existence

Proposition 4.12. Fix $\gamma \in]0, 1[$ and $0 < s < \gamma$. Let $\beta > 0$ and T_* be the “loss parameter” and the existence time given by Theorem 2.4.

Then, for all $u_0 \in H^s(\mathbb{R}^n; \mathbb{R}^m)$ and $f \in \mathcal{L}_{s,\beta}(T_*)$, there exists a weak solution $u \in \mathcal{H}_{s,\beta}(T_*)$ to the Cauchy problem (10).

Proof. For $\varepsilon \in]0, 1]$, let us introduce the smoothing operators $J_\varepsilon(D)$, of order -1 , as done in the previous proof. Denoting $A(t, x, D) = \sum_j A_j(t, x) \partial_j$, for any ε let us consider the linear system of m ODEs

$$\partial_t u_\varepsilon = -A(t, x, D) J_\varepsilon(D) u_\varepsilon - B(t, x) J_\varepsilon(D) u_\varepsilon + f,$$

with initial datum $(u_\varepsilon)|_{t=0} = u_0$.

For all ε fixed and all t , it is easy to see that the operators $A(t, x, D)J_\varepsilon(D)$ and $B(t, x)J_\varepsilon(D)$ are bounded in L^2 . Moreover, by our hypotheses we have $s > s - \beta T_* > 0$, so that we can solve the previous system in $\mathcal{C}([0, T_*]; L^2(\mathbb{R}^n; \mathbb{R}^m))$ by use of the Cauchy-Lipschitz Theorem, and find a solution u_ε . On the other hand, the hypothesis over f implies that, for any $t_0 \in [0, T_*]$ fixed, $f \in L^1([0, t_0]; H^{s-\beta t_0})$. Furthermore, by product rules (see Proposition 3.9 and Corollary 3.10), the operators $A(t, x, D)J_\varepsilon(D)$ and $B(t, x)J_\varepsilon(D)$ are self-maps of $H^{s-\beta t_0}$ into itself. Therefore, by Cauchy-Lipschitz Theorem again and uniqueness part, we infer that $u \in \mathcal{C}_{s,\beta}(T_*) \hookrightarrow \mathcal{H}_{s,\beta}(T_*)$.

Notice however that this is *not* enough to get uniform bounds on the family $(u_\varepsilon)_\varepsilon$, since we do not know if we have enough regularity in order to absorb, in energy estimates, the remainders which require an additional (1/2) log-regularity (see the computations in Subsection 4.2 above). Hence, we are going to argue in a slightly different way.

Let us define $w_\varepsilon := J_\varepsilon(D)u_\varepsilon$: by the previous argument, w_ε is in $L^2([0, T_*]; H^1(\mathbb{R}^n; \mathbb{R}^m))$ for all $\varepsilon \in]0, 1]$. Moreover, it satisfies

$$\begin{aligned} \partial_t w_\varepsilon &= -J_\varepsilon(D)A(t, x, D)w_\varepsilon - J_\varepsilon(D)B(t, x)w_\varepsilon + J_\varepsilon(D)f \\ &= -A(t, x, D)J_\varepsilon(D)w_\varepsilon - B(t, x)J_\varepsilon(D)w_\varepsilon + f_\varepsilon + h_\varepsilon, \end{aligned}$$

where $f_\varepsilon = J_\varepsilon(D)f$ as before, and h_ε is defined by the analogue of formula (35), but replacing u by w_ε itself. Notice that Lemma 4.6 of [11] implies the inequality

$$\|h_\varepsilon\|_{H^{s(t)}} \leq C \|w_\varepsilon\|_{H^{s(t)}},$$

where $s(t) = s - \beta t$ as above. Let us also remark that $(w_\varepsilon)|_{t=0} = \Lambda_\varepsilon(D)u_0$.

Then, we can apply energy estimates of Theorem 2.4 to w_ε . Indeed, $J_\varepsilon(\xi)$ being a scalar multiplier, $S(t, x, \xi)$ is still a microlocal symmetrizer for $A(t, x, \xi)J_\varepsilon(\xi)$. Moreover, Lemma 4.3 gives uniform bounds for the remainder operators $A(t, x, \xi)J_\varepsilon(\xi) - T_{A(t,x,\xi)J_\varepsilon(\xi)}$ and $B(t, x)J_\varepsilon(\xi) - T_{B(t,x)J_\varepsilon(\xi)}$ in suitable functional spaces. Finally, Lemma 4.4 provides with uniform bounds for the operators $\partial_t + T_{iA(t,x,\xi)J_\varepsilon(\xi)} + T_{B(t,x)J_\varepsilon(\xi)}$, since the symbols are uniformly bounded respectively in the classes $\Gamma_{LL}^{1+0\log}$ and $\Gamma_{\gamma+0\log}^{0+0\log}$.

So we find that $(w_\varepsilon)_\varepsilon$ is a bounded family in $\mathcal{C}_{s,\beta}(T_*) \hookrightarrow \mathcal{H}_{s,\beta}(T_*)$, and consequently, up to extraction of a subsequence, it weakly converges to some u in this space. On the other hand, by the equation for w_ε (recall also product rules of Proposition 3.9 and Corollary 3.10), we easily deduce that $(\partial_t w_\varepsilon)_\varepsilon$ is bounded in $\mathcal{C}_{s-1,\beta}(T_*)$, and hence the convergence holds true also in the weak-* topology of $H^1([0, T_*]; H^{s-\beta T_*-1}(\mathbb{R}^n; \mathbb{R}^m)) \hookrightarrow \mathcal{C}([0, T_*]; H^{s-\beta T_*-1})$. In particular, this implies that $(w_\varepsilon)|_{t=0} \rightharpoonup u|_{t=0}$ in the distributional sense.

Thanks to these properties, it is easy to pass to the limit in the weak formulation of the equations, obtaining thus that u solves the system $Lu = f$ in a weak sense, with initial datum u_0 . Finally, by uniform bounds we get $u \in \mathcal{H}_{s,\beta}(T_*)$, and this fact completes the proof of the existence of a weak solution. \square

5 The local Cauchy problem

s:local

We prove here local in space existence and uniqueness of solutions. First of all, let us show that it makes sense to consider the Cauchy problem for operator P on Σ : this is not clear *a priori*, due to the low regularity framework.

Here below we will use the notations introduced in Subsection 2.2. In particular, recall that we have set $\Omega_{\geq} := \Omega \cap \{\varphi \geq 0\}$ and $\Omega_{>} := \Omega \cap \{\varphi > 0\}$, where $\{\varphi = 0\}$ is a parametrization of the hypersurface Σ in Ω .

Recall that we have supposed that hypotheses from **(H-1)** to **(H-5)** (stated in Subsection 2.2) hold true. In particular, P has log-Lipschitz first order coefficients and a γ -Hölder continuous 0-th order coefficient, and it admits a family of full symmetrizers, which are smooth in the dual variable and log-Lipschitz in the z variable.

5.1 Giving sense to the local Cauchy problem

ss:sense

We start by noticing that, for smooth u and v , with $\text{supp } v$ compact in Ω_{\geq} , we can write the next identity only formally, due to the low regularity of the coefficients:

$$(Pu, v)_{L^2(\Omega_{>})} - (u, P^*v)_{L^2(\Omega_{>})} = (D_{\Sigma}u, N_{\nu, P_1}v)_{L^2(\Sigma)},$$

where $P^*(z, \zeta)v = -\sum_j \partial_{z_j}(A_j^*v) + B^*v$ is the adjoint operator of P , and we have defined

$$D_{\Sigma}u := u|_{\Sigma} \quad \text{and} \quad N_{\nu, P_1}v := \sum_{j=0}^n \nu^j D_{\Sigma}(A_j^*v).$$

As above, ν is the normal to the hypersurface Σ , which determines the integration form on Σ . Of course, the map $z \mapsto \nu(z) \neq 0$ is smooth on Σ .

In our framework the previous Green formula does not hold true a priori, due to the low regularity of the coefficients. The first step of the proof to Theorem 2.10 is to justify it for smooth enough functions. To begin with, let us study the regularity of the terms entering in the definition of Pu and P^*u : we have the following lemma.

1:reg_P

Lemma 5.1. (i) For $s \in]1-\gamma, 1+\gamma[$ and $u \in H_{loc}^s(\Omega_{\geq})$, all the terms entering in the definition of Pu are well-defined in $H_{loc}^{s-1}(\Omega_{\geq})$.

(ii) For $\sigma \in]1-\gamma, 1[$ and $v \in H_{loc}^{\sigma}(\Omega_{\geq})$, all the terms entering in the definition of P^*v are well-defined in $H_{loc}^{\sigma-1}(\Omega_{\geq})$.

(iii) For $s > 1/2$ and $u \in H_{loc}^s(\Omega_{\geq})$, the trace $D_{\Sigma}u$ is well-defined in $H_{loc}^{s-1/2}(\Sigma \cap \Omega)$.

(iv) For $\sigma \in]1/2, 1[$ and $v \in H_{loc}^{\sigma}(\Omega_{\geq})$, the normal trace $N_{\nu, P_1}v$ is well-defined in $H_{loc}^{s-1/2}(\Sigma \cap \Omega)$.

Proof. We repeatedly use product rules stated in Proposition 3.9 and Corollary 3.10 above.

In a first time, let us focus on P : for all $s \in]0, 2[$, $P_1(z, \partial_z)u$ belongs to $H_{loc}^{s-1}(\Omega_{\geq})$. Now, by the embedding $H^s \hookrightarrow H^{s-1}$, if $|s-1| < \gamma$ we have that $Bu \in H_{loc}^{s-1}(\Omega_{\geq})$. This completes the proof of point (i).

Concerning P^* , the argument is analogous: the principal part of P^*u belongs to $H_{loc}^{\sigma-1}(\Omega_{\geq})$ whenever $\sigma \in]0, 1[$. The B^*u term can be treated exactly as before. Also (ii) is proved.

Points (iii) and (iv) are straightforward. \square

Following the discussion in [11], we remark that $\mathcal{C}_0^{\infty}(\Omega_{>})$ is a dense subset of $H^{\sigma}(\Omega_{>})$ for $|\sigma| < 1/2$, and that, when $\sigma \in [0, 1/2[$ and $u \in H^{\sigma}(\Omega_{>})$, the pairing $(u, v)_{L^2(\Omega_{>})}$ for $v \in L^2$ extends to the duality $H^{\sigma} \times H^{-\sigma}$. From this and Lemma 5.1, we deduce the next statement, which tells us that Green formula makes sense for regular enough distributions.

1:duality

Lemma 5.2. Let us fix $1/2 < \gamma < 1$, and take $s \in]1/2, \gamma[$. For any $u \in H_{loc}^s(\Omega_{\geq})$ and $v \in H_{comp}^s(\Omega_{\geq})$, one has the equality

$$(36) \quad (Pu, v)_{H^{-\sigma} \times H^{\sigma}} - (u, P^*v)_{H^{\sigma} \times H^{-\sigma}} = (D_{\Sigma}u, N_{\nu, P_1}v)_{L^2(\Sigma)},$$

where $\sigma = 1 - s \in]1-\gamma, 1/2[$ (and in particular $0 \leq \sigma < 1/2$).

Proof. Since $s > 1/2$, we have $H^s \hookrightarrow H^{1-s}$. On the other hand, s belongs in particular to $]1-\gamma, 1[$, and then Lemma 5.1 applies.

To complete the proof, it is enough to remark that (see Lemma 1.3 of [11]) the Green's formula

$$(\partial_j u, v)_{H^{-\sigma} \times H^{\sigma}} = -(u, \partial_j v)_{H^{\sigma} \times H^{-\sigma}} + (\nu^j D_{\Sigma}u, D_{\Sigma}v)_{L^2(\Sigma)},$$

holds true for any $u \in H_{loc}^{1-\sigma}(\Omega_{\geq})$ and $v \in H_{comp}^{1-\sigma}(\Omega_{\geq})$, whenever $\sigma \in [0, 1/2[$. \square

As the final step, we want to justify Green formula for distributions which are in the domain of our operator. This is guaranteed by the next statement, which is the analogue to Proposition 1.4 of [11]. Thanks to this result, considering the Cauchy problem (CP) under the hypotheses of Theorem 2.10 makes sense.

Proposition 5.3. *Let us define the set $\mathcal{D}(P; H^s) := \{u \in H_{loc}^s(\Omega_{\geq}) \mid Pu \in H_{loc}^s(\Omega_{\geq})\}$.*

There exists a unique extension of the operator D_{Σ} to the set $\mathcal{D}(P) := \bigcup_{s>1-\gamma} \mathcal{D}(P; H^s)$, which acts continuously from $\mathcal{D}(P; H^s)$ into $H_{loc}^{s-1/2}(\Sigma \cap \Omega)$ for all $s \in]1-\gamma, \gamma[$.

The same property holds true for the operator N_{ν, P_1} .

Furthermore, for all $s_0 \in]1-\gamma, 1/2[$ such that $s_0 \leq s$, for all $v \in H_{comp}^{1-s_0}(\Omega_{\geq})$, one has the Green formula

$$(Pu, v)_{L^2} - (u, P^*v)_{H^{s_0} \times H^{-s_0}} = (D_{\Sigma}u, N_{\nu, P_1}v)_{H^{s-1/2} \times H^{1/2-s}}.$$

This proposition will be proved in the next subsection, and it will be derived by results on the global Cauchy problem. Then, in order to give sense to the Cauchy problem, the microlocally symmetrizability and hyperbolicity hypotheses are fundamental.

Let us notice here that, thanks to definitions, the properties for D_{Σ} easily pass also on N_{ν, P_1} . We observe also that all the terms entering in the last formula have sense. Indeed, $1/2 < 1-s_0 < \gamma$ and hence, by Lemma 5.1, we get $P^*v \in H_{comp}^{-s_0}(\Omega_{\geq})$, so that the pairing with u makes sense (recall that $s_0 \leq s$). Moreover one has $N_{\nu, P_1}v \in H_{comp}^{1/2-s_0}(\Omega_{\geq}) \hookrightarrow H_{comp}^{1/2-s}(\Omega_{\geq})$, and so also the last term in the equality is fine.

5.2 Invariance by change of variables, and regularity results

First of all, let us show here that our working hypotheses are invariant under smooth change of coordinates.

Indeed, let $z = \psi(y)$, for some ψ smooth, and denote $\tilde{f}(y) = f \circ \psi(y)$: derivatives change according to the rule $\nabla_z u \circ \psi = {}^t(\nabla_z \psi^{-1}) \circ \psi \cdot \nabla_y \tilde{u}$, namely

$$(\partial_{z_j} u^k) \circ \psi = \sum_{h=1}^m \partial_{z_j} (\psi^{-1})^h \circ \psi \partial_{y_h} \tilde{u}^k,$$

as well as covectors in the cotangent space, i.e. $\tilde{\zeta}(y) = ({}^t \nabla_z \psi^{-1} \cdot \zeta) \circ \psi(y)$. As a consequence, if we write the i -th component of our system $Pu = f$, that is to say

$$f^i = \sum_{j=0}^n \sum_{k=1}^m A_{j,ik}(z) \partial_{z_j} u^k(z) + \sum_{k=1}^m B_{ik}(z) u^k(z),$$

from the previous rules we obtain the expression in y coordinates:

$$\tilde{f}^i = \sum_{h=0}^n \sum_{k=1}^m M_{h,ik}(y) \partial_{y_j} \tilde{u}^k(y) + \sum_{k=1}^m \tilde{B}_{ik}(y) \tilde{u}^k(y),$$

where easy computations lead to the formula

$$M_{h,ik}(y) = \sum_{j=0}^n \tilde{A}_{j,ik}(y) \partial_{z_j} (\psi^{-1})^h(\psi(y)).$$

Namely, we have proved that $\tilde{P}u = \tilde{P}\tilde{u}$, where we have defined the operator

$$\tilde{P}v = \sum_{h=0}^n M_h \partial_{y_h} v + \tilde{B}v.$$

Notice that \tilde{P} has the same form as P , and the same regularity of its first and 0-th order coefficients. In addition, by Theorem 4.11 of [20] we deduce that it is still microlocally symmetrizable in the sense of Definition 2.9. We can check this property also by direct computations: as a matter of fact, in a very natural way, let us define

$$\tilde{\mathbf{S}}(y, \tilde{\zeta}) := \mathbf{S}(\psi(y), ({}^t\nabla_z \psi^{-1} \cdot \zeta) \circ \psi)$$

Of course, all the properties of \mathbf{S} hand over to $\tilde{\mathbf{S}}$; the only thing we have to check is that $\tilde{\mathbf{S}}\tilde{P}_1$ is self-adjoint, where we denote by \tilde{P}_1 the principal part of \tilde{P} . But this immediately follows from the definitions and the transformation rules:

$$\begin{aligned} (\tilde{\mathbf{S}}\tilde{P}_1)(y, \tilde{\zeta}) &= \tilde{\mathbf{S}}(y, \tilde{\zeta}) \sum_{h=0}^n \tilde{\zeta}^h M_h = \tilde{\mathbf{S}}(y, \tilde{\zeta}) \sum_{h=0}^n \sum_{j=0}^n \tilde{A}_j \partial_{z_j} (\psi^{-1})^h \tilde{\zeta}^h \\ &= \mathbf{S}(\chi(y), ({}^t\nabla_z \psi^{-1} \cdot \zeta) \circ \psi(y)) \sum_{j=0}^n A_j(\psi(y)) ({}^t\nabla_z \psi^{-1} \cdot \zeta) \circ \psi(y) \\ &= (\mathbf{S}P_1)(\chi(y), ({}^t\nabla_z \psi^{-1} \cdot \zeta)). \end{aligned}$$

Furthermore, as done in [11] (see Section 5), the Green formula (36) can be transported by ψ . As a byproduct of this discussion, we get that our statements, and in particular Proposition 5.3, are invariant by smooth change of variables.

Hence, we can reconduct the proof of Proposition 5.3 in the system of coordinates $z = (t, x) \in]-t_0, t_0[\times \omega$ and $\zeta = (\tau, \xi)$. In particular we will assume that $\Sigma = \{t = 0\}$, whose unit normal is dt , and that $z_0 = 0$. Moreover, by Proposition 4.12 of [20], up to shrink our domain we can suppose that the full symmetrizer is positive in the direction dt at any point (t, x) ; by assumption (13), we can also assume that $A_0(t, x)$ is invertible for any (t, x) .

In the end, we can recast our operator in the form $Lu = \partial_t u + \sum_j A_j(t, x) \partial_j u + B(t, x)u$, with existence of a microlocal symmetrizer (in the sense of Definition 2.1) $S(t, x, \xi) = \mathbf{S}(t, x, 1, \xi)$.

This having been established, we present a regularity result analogous to Proposition 4.9, which will be needed in the proof of Proposition 5.3.

1:regularity

Lemma 5.4. *Let γ and s as in the hypotheses of Theorem 2.10. Let $u \in H^s(]0, t_0[\times \omega)$ such that $Lu \in L^1([0, t_0]; H^s(\omega))$.*

Then one has $u \in \mathcal{C}([0, t_0]; H^{s-1/2}(\omega))$. In particular, the trace $u|_{t=0}$ is well-defined in $H^{s-1/2}(\omega)$, and the initial condition in (CP) makes sense.

Proof. The proof goes along the lines of Lemma 2.2 of [11], so let us just sketch it. First of all, we restrict our attention to the case $[0, t_0] \times \mathbb{R}^n$, the other one being obtained working with restrictions. In addition, we make use of the spaces $\mathcal{H}^{s,\sigma}([0, T_*] \times \mathbb{R}^n)$ of Hörmander (see e.g. Appendix B of [14]).

So, by hypothesis $u \in H^s([0, t_0] \times \mathbb{R}^n) = \mathcal{H}^{s,0}([0, t_0] \times \mathbb{R}^n) \hookrightarrow \mathcal{H}^{0,s}([0, t_0] \times \mathbb{R}^n)$. From this, we deduce that $L_1 u = \sum_j A_j \partial_j u \in \mathcal{H}^{0,s-1}([0, t_0] \times \mathbb{R}^n)$. On the other hand, the same is true also for the term Bu , since, by assumption, $|s-1| < \gamma$ (recall also Proposition 3.9).

Next, we notice that the hypothesis $Lu \in L^1([0, t_0]; H^s(\mathbb{R}^n))$ implies that $v(t) := \int_0^t Lu(\tau) d\tau$ belongs to $\mathcal{C}([0, t_0]; H^s(\mathbb{R}^n)) \hookrightarrow \mathcal{H}^{0,s}([0, t_0] \times \mathbb{R}^n)$.

From the previous properties, we get that $w := u - v \in \mathcal{H}^{0,s}([0, t_0] \times \mathbb{R}^n)$, while its time derivative $\partial_t w \in \mathcal{H}^{0,s-1}([0, t_0] \times \mathbb{R}^n)$, because $\partial_t w = \sum_j A_j \partial_j u + Bu$. Therefore, by properties of the Hörmander spaces we get $w \in \mathcal{H}^{1,s-1}([0, t_0] \times \mathbb{R}^n) \hookrightarrow \mathcal{C}([0, t_0]; H^{s-1/2}(\mathbb{R}^n))$ (see Theorem B.2.7 of [14]). As a conclusion, also $u = w + v$ belongs to $\mathcal{C}([0, t_0]; H^{s-1/2}(\mathbb{R}^n))$. \square

We are now ready to prove Proposition 5.3.

Proof of Proposition 5.3. We focus on uniqueness first. Notice that classical Green formula forces the definition of D_Σ to coincide with the usual trace operator for smooth functions. Let us argue by contradiction and suppose that there exist two extensions D_Σ^1 and D_Σ^2 , and a tempered distribution $u \in \mathcal{D}(P) = \bigcup_{s>1-\gamma} \mathcal{D}(P; H^s)$ such that $D_\Sigma^1 u = u_0^1$ and $D_\Sigma^2 u = u_0^2$, with $u_0^1 \neq u_0^2$. In the light of Lemma 5.2, we can suppose that $u \in \mathcal{D}(P; H^s)$, with $s \in]1 - \gamma, 1/2]$. Let us set $Pu = f \in H_{loc}^s(\Omega_\geq)$ and $\delta := \|u_0^1 - u_0^2\|_{H^{s-1/2}} > 0$.

Fix now $\varepsilon > 0$. By density of smooth functions, we can chose smooth g, v_0^1 and v_0^2 such that

$$\|f - g\|_{H^s(\Omega_\geq)} + \|u_0^1 - v_0^1\|_{H^{s-1/2}(\Sigma \cap \Omega)} + \|u_0^2 - v_0^2\|_{H^{s-1/2}(\Sigma \cap \Omega)} \leq \varepsilon.$$

Thanks to the previous discussion, we can work in local coordinates (t, x) , and hence suppose the following facts: first, that $\Sigma = \{t = 0\}$, and moreover that, in these coordinates, P takes the form of L as defined in (3). Therefore, we can apply Theorem 2.6 to the initial data v_0^1 and v_0^2 and external force g : we find two solutions v^1 and v^2 respectively, which belong to the space $\mathcal{C}([0, T_*]; H^s(\mathbb{R}^n; \mathbb{R}^m))$ and which solves the problems, for $i = 1, 2$,

$$Lv^i = g, \quad v^i|_{t=0} = v_0^i.$$

Notice that, in particular, we get $v^i \in \mathcal{H}^{0,s}([0, T_*] \times \mathbb{R}^n)$, and then (by the equation) $\partial_t v^i \in \mathcal{H}^{0,s-1}([0, T_*] \times \mathbb{R}^n)$, which implies $v^i \in \mathcal{H}^{1,s-1}([0, T_*] \times \mathbb{R}^n)$ (see also the proof of Lemma 5.4 here above). Since $s - 1 < 0$, from this property it is easy to deduce that $v^i \in \mathcal{H}^{s,0}([0, T_*] \times \mathbb{R}^n)$. Remark that, up to shrink Ω_\geq (i.e. take a smaller existence time T_*), we can also suppose that $v^i \in \mathcal{H}^{\sigma,0}([0, T_*] \times \mathbb{R}^n)$, with $\sigma > 1/2$. Furthermore, by linearity of L and Theorem 2.6 we infer the estimates

$$\|u - v^i\|_{H^s} \leq C(\varepsilon) \quad \implies \quad \|v^1 - v^2\|_{H^s} \leq 2C(\varepsilon).$$

But each v^i is smooth, i.e. it belongs $\mathcal{H}^{\sigma,0}([0, T_*] \times \mathbb{R}^n)$ with $\sigma > 1/2$, so that $D_\Sigma^1 v^i \equiv D_\Sigma^2 v^i = v_0^i$ for all $i = 1, 2$. Hence, by continuity of the trace operator on smooth functions we deduce

$$\|v_0^1 - v_0^2\|_{H^{s-1/2}} \leq \|v_0^1 - v_0^2\|_{H^{\sigma-1/2}} \leq C'(\varepsilon).$$

At this point, we write

$$\delta = \|u_0^1 - u_0^2\|_{H^{s-1/2}} \leq \|u_0^1 - v_0^1\|_{H^{s-1/2}} + \|v_0^1 - v_0^2\|_{H^{s-1/2}} + \|v_0^2 - u_0^2\|_{H^{s-1/2}} \leq 2\varepsilon + C'(\varepsilon) :$$

taking a $\varepsilon > 0$ small enough gives the contradiction, completing in this way the proof of the uniqueness part.

Let us consider the problem of existence of the trace operator onto Σ . Once again, we work in local coordinates $z = (t, x)$, with $z_0 = 0$ and $\Sigma = \{t = 0\}$, and the conormal given by $\nu(t, x) = \mu(x) dt$, for a suitable positive function μ . By Lemma 5.4, we can define the trace $u|_{t=0}$ as a distribution in $H^{s-1/2}(\Sigma)$: then, it is enough to check that, in this coordinates, Green formula in Proposition 5.3 makes sense with $D_\Sigma u = u|_{t=0}$. But from now on the arguments are analogous to the discussion in Section 5 of [11]: so we omit them.

The proposition is now completely proved. \square

5.3 Proof of local existence and uniqueness

Now, we can turn our attention to the statements about existence and uniqueness of solutions to the local Cauchy problem. Let us start with the proof of Theorem 2.10: it is analogous to the one of [11] for wave operators, so let us just sketch it.

Proof of Theorem 2.10. Up to a change of variables, we can suppose, as usual, that $z_0 = 0$, $\Sigma = \{t = 0\}$, with normal dt , and that the operator P , in these coordinates, assumes the form of $L = \partial_t + \sum_j A_j \partial_j + B$.

Consider a smooth function $\Phi : \mathbb{R}_y^{1+n} \rightarrow \Omega$ such that $\Phi(y) = y$ for all $y \in \Omega_1 \subset \Omega$, and $\Phi(y) = z_0$ if $|y|$ is large enough. Such a function can be built by taking Ω_1 to be a ball centered in z_0 and working along radial directions emanating from z_0 , for instance.

Changing the coefficients of L according to the rule $f^\sharp(y) := f(\Phi(y))$, we are led to consider a new operator L^\sharp , such that $L^\sharp \equiv L$ on Ω_1 , having coefficients with the same regularity as L and admitting a full symmetrizer $\mathbf{S}^\sharp(y, \zeta) = \mathbf{S}(\Phi(y), \zeta)$ which is positive in the directions $\nu(\Phi(y))$.

Fix now a $s \in]1 - \gamma, \gamma[$, and take another $s_1 \in]1 - \gamma, s[$. Let us apply Theorem 2.4 to operator L^\sharp : this provides us with a loss parameter β and with an existence time $T = (s - s_1)/\beta$. Define then $\Omega_0 := \Omega_1 \cap \{|t| < T\}$ and $\omega := \Omega_0 \cap \{t = 0\}$.

Let us take an initial datum $u_0 \in H^s(\omega)$: by definition (11), it can be seen as the restriction to ω of a $u_0^\sharp \in H^s(\mathbb{R}^n)$. In the same way, $f \in H^s(\Omega_0 \cap \{t > 0\})$ is the restriction of a $f^\sharp \in H^s(\mathbb{R}^{1+n} \cap \{t > 0\})$, and in particular $f^\sharp \in L^2(]0, T[; H^s(\mathbb{R}^n))$.

Then, by use of Theorem 2.4 we solve the Cauchy problem

$$L^\sharp u^\sharp = f^\sharp, \quad u^\sharp|_{t=0} = u_0^\sharp :$$

we find a solution u^\sharp on $[0, T] \times \mathbb{R}^n$ which, in particular, belongs to $L^2([0, T]; H^{s_1}(\mathbb{R}^n))$. But from this fact together with Proposition 3.9 and by use of the equation, we deduce that $\partial_t u^\sharp \in L^2(]0, T[; H^{s_1-1}(\mathbb{R}^n))$, which finally implies that $u^\sharp \in \mathcal{H}^{1, s_1-1}([0, T] \times \mathbb{R}^n) \hookrightarrow H^{s_1}([0, T] \times \mathbb{R}^n)$, the inclusion following from the fact that $s_1 - 1 < 0$.

Therefore, by restriction we infer the existence of a solution u to (CP) in Ω_0 . \square

This having been proved, we turn our attention to the question of local uniqueness. We start by establishing a result about propagation of zero across the surface $\{t = 0\}$.

1:0-prop

Lemma 5.5. *Let $s > 1 - \gamma$ and $\Omega =] - t_0, t_0[\times \omega$ for some $t_0 > 0$ and $\omega \subset \mathbb{R}^n$ neighborhood of 0. Suppose that $u \in H^s(\Omega \cap \{t > 0\})$ fulfills*

$$(37) \quad Lu = 0, \quad u|_{t=0} = 0.$$

Denote by u_e the extension of u by 0 on $\{t < 0\}$.

Then, there exists a neighborhood $\Omega_1 \subset \Omega$ of 0 such that $u_e \in H^s(\Omega_1)$ and $Lu_e = 0$ on Ω_1 .

Proof. Since $u \in H^s(\Omega \cap \{t > 0\})$, we have in particular that $u \in L^2([0, t_0]; H_{loc}^s(\omega))$, and therefore u_e belongs to $L^2([-t_0, t_0]; H_{loc}^s(\omega))$.

On the other hand, $\partial_t u \in L^2([0, t_0]; H_{loc}^{s-1}(\omega))$, with $u|_{t=0} = 0$. Hence, the (weak) derivative $\partial_t u_e$ is the extension of $\partial_t u$ by 0 for negative times (test it on functions of the form $\partial_t \varphi$). This implies $\partial_t u_e \in L^2([-t_0, t_0]; H_{loc}^{s-1}(\omega))$, and therefore $u_e \in H_{loc}^s(]-t_0, t_0[\times \omega)$ (using the same arguments as in the previous proof).

Finally, verifying that $Lu_e = 0$ is an easy matter. Indeed, testing the equations on smooth $\varphi \in \mathcal{D}(]-t_0, t_0[\times \omega)$, the integral in time reduces on intervals of the form $]0, t_1[$, because u_e vanishes for $t < 0$. But for $t > 0$, $u_e \equiv u$, and u is a weak solution of $Lu = 0$ with $u|_{t=0} = 0$: this entails that also the integral over $]0, t_1[$ is equal to 0. \square

From the previous result, we can deduce the local uniqueness of solutions, by mean of classical convexification arguments.

Proof of Theorem 2.12. Once again, we can fix local coordinates for which $\Omega =] - t_0, t_0[\times \omega$ and $\Sigma = \{t = 0\}$. Suppose that $u \in H^s(\Omega \cap \{t > 0\})$ satisfies (37). Denoting by u_e its extension by 0, provided by Lemma 5.5, we know that u_e is a H^s distribution on (say) $] - t_1, t_1[\times \omega_1$, for a suitably small $\omega_1 \subset \omega$, and in the same neighborhood $Lu_e = 0$.

We consider the change of variables $\psi : (t, x) \mapsto (\tilde{t}, \tilde{x})$ such that

$$\tilde{t} := t + |x|^2 \quad \text{and} \quad \tilde{x} := x.$$

Notice that this map sends $\{t < 0\}$ into $\{\tilde{t} < |\tilde{x}|^2\}$. Let us define $\tilde{u} = u_e \circ \psi$ and \tilde{L} the operator obtained by L under the transformation ψ .

Therefore, we have that \tilde{u} is defined in $\tilde{t} < \tilde{t}_1$, for a suitable $\tilde{t}_1 > 0$, and $\tilde{u} \equiv 0$ in $\{\tilde{t} < |\tilde{x}|^2\}$. Furthermore, up to take a smaller \tilde{t}_1 , we can suppose that \tilde{L} is defined on a neighborhood $\tilde{\Omega}$ of the origin which contains the closed lens $\tilde{\Theta} := \{|\tilde{x}|^2 \leq \tilde{t} \leq \tilde{t}_1\}$, and that $\tilde{L}\tilde{u} \equiv 0$ on $\tilde{\Omega}$.

Now, we repeat the construction explained in the proof of Theorem 2.10 above. Namely, we extend the coefficients of \tilde{L} and we obtain an operator L^\sharp , defined on the whole \mathbb{R}^{n+1} , which preserves the regularity of the coefficients and the microlocal symmetrizability assumption, and which coincides with \tilde{L} on a smaller neighborhood of $\tilde{\Theta}$.

We also extend \tilde{u} to u^\sharp . Then, on the set $] -\infty, \tilde{t}_1[\times \mathbb{R}^n$ we get

$$L^\sharp u^\sharp = 0, \quad u^\sharp \in H^s, \quad u^\sharp|_{\{\tilde{t} < |\tilde{x}|^2\}} = 0.$$

In particular, $u^\sharp|_{\tilde{t}=-\delta} = 0$ for $\delta > 0$ arbitrarily small. Therefore energy estimates of Theorem 2.4, applied to L^\sharp and the initial time $-\delta$ (for δ small enough), guarantee that $u^\sharp \equiv 0$ until a time \tilde{T}_0 . In particular, this is true for \tilde{u} , and coming back to coordinates (t, x) implies that $u \equiv 0$ in a neighborhood of the origin. \square

References

- B-C_1994 [1] H. Bahouri, J.-Y. Chemin: *Équations de transport relatives à des champs de vecteurs non-lipschitziens et mécanique des fluides*, Arch. Rational Mech. Anal., **127** (1994), n. 2, 159-181.
- B-C-D [2] H. Bahouri, J.-Y. Chemin, R. Danchin: “Fourier Analysis and Nonlinear Partial Differential Equations”. Grundlehren der Mathematischen Wissenschaften (Fundamental Principles of Mathematical Sciences), **343**, Springer, Heidelberg (2011).
- Bony [3] J.-M. Bony: *Calcul symbolique et propagation des singularités pour les équations aux dérivées partielles non linéaires*. Ann. Sci. École Norm. Sup. (4), **14** (1981), 209-246.
- C-DG-S [4] F. Colombini, E. De Giorgi, S. Spagnolo: *Sur les équations hyperboliques avec des coefficients qui ne dépendent que du temps*. Ann. Scuola Normale Sup. Pisa Cl. Scienze (4), **6** (1979), no. 3, 511-559.
- C-DS [5] F. Colombini, D. Del Santo: *A note on hyperbolic operators with log-Zygmund coefficients*, J. Math. Sci. Univ. Tokyo **16** (2009), no. 1, 95-111.
- C-DS-F-M_t1 [6] F. Colombini, D. Del Santo, F. Fanelli, G. Métivier: *Time-dependent loss of derivatives for hyperbolic operators with non-regular coefficients*. Comm. Partial Differential Equations, **38** (2013), n. 10, 1791-1817.
- C-DS-F-M_wp [7] F. Colombini, D. Del Santo, F. Fanelli, G. Métivier: *A well-posedness result for hyperbolic operators with Zygmund coefficients*. J. Math. Pures Appl. (9), **100** (2013), n. 4, 455-475.
- C-DS-F-M_Birk [8] F. Colombini, D. Del Santo, F. Fanelli, G. Métivier: *A note on complete hyperbolic operators with log-Zygmund coefficients*. In “Fourier Analysis”, Trends Math., Birkhäuser/Springer, Cham (2014), 47-72.
- DS-F-M_Z-sys [9] F. Colombini, D. Del Santo, F. Fanelli, G. Métivier: *The well-posedness issue in Sobolev spaces for hyperbolic systems with Zygmund-type coefficients*. Comm. Partial Differential Equations, **41** (2015), n. 11, 2082-2121.
- C-L [10] F. Colombini, N. Lerner: *Hyperbolic operators with non-Lipschitz coefficients*. Duke Math. J., **77** (1995), 657-698.
- C-M [11] F. Colombini, G. Métivier: *The Cauchy problem for wave equations with non-Lipschitz coefficients; application to continuation of solutions of some nonlinear wave equations*. Ann. Sci. École Norm. Sup. (4) **41** (2008), n. 2, 177-220.

- C-M_2015 [12] F. Colombini, G. Métivier: *Counterexamples to the well-posedness of the Cauchy problem for hyperbolic systems*. Anal. PDE, **8** (2015), n. 2, 499-511.
- D_2005 [13] R. Danchin: *Estimates in Besov spaces for transport and transport-diffusion equations with almost Lipschitz coefficients*, Rev. Mat. Iberoamericana, **21** (2005), n. 3, 863-888.
- Horm [14] L. Hörmander: "The analysis of linear partial differential operators III". Grundlehren der Mathematischen Wissenschaften (Fundamental Principles of Mathematical Sciences), **274**, Springer-Verlag, Berlin (1994).
- J-M-R_2005 [15] J.-L. Joly, G. Métivier, J. Rauch: *Hyperbolic domains of determinacy and Hamilton-Jacobi equations*. J. Hyperbolic Differ. Equ., **2** (2005), n. 3, 713-744.
- K [16] T. Kato: "Perturbation Theory for Linear Operators". Classics in Mathematics, Springer-Verlag, Berlin (1995).
- Iv-Pet [17] V. J. Ivrii, V. M. Petkov: *Necessary conditions for the correctness of the Cauchy problem for non-strictly hyperbolic equations*. Uspehi Math. Nauk, **29** (1974), n. 5(179), 3-70. English translation: Russian Math. Surveys, **29** (1974), n. 5, 1-70.
- M-1986 [18] G. Métivier: *Interactions de deux chocs pour un système de deux lois de conservation, en dimension deux d'espace*. Trans. Amer. Math. Soc., **296** (1986), 431-479.
- M-2008 [19] G. Métivier: "Para-differential calculus and applications to the Cauchy problem for nonlinear systems". Centro di Ricerca Matematica "Ennio De Giorgi" (CRM) Series, **5**, Edizioni della Normale, Pisa (2008).
- M-2014 [20] G. Métivier: *L^2 well-posed Cauchy problems and symmetrizability of first order systems*. J. Éc. Polytech. Math., **1** (2014), 39-70.
- M-Z [21] G. Métivier, K. Zumbrun: "Large viscous boundary layers for noncharacteristic nonlinear hyperbolic problems". Mem. Amer. Math. Soc., **175** (2005).
- Rauch_2005 [22] J. Rauch: *Precise finite speed with bare hands*. Methods Appl. Anal., **12** (2005), n. 3, 267-277.
- Rauch [23] J. Rauch: "Hyperbolic Partial Differential Equations and Geometric Optics". Graduate Studies in Mathematics, **133**, American Mathematical Society, Providence, RI (2012).
- Strang [24] G. Strang: *Necessary and insufficient conditions for well-posed Cauchy problems*. J. Differential Equations, **2** (1966), 107-114.
- Tar [25] S. Tarama: *Energy estimate for wave equations with coefficients in some Besov type class*. Electron. J. Differential Equations (2007), Paper No. 85, 12 pp. (electronic).
- T_1983 [26] H. Triebel: "Theory of function spaces". Monographs in Mathematics, **78**, Birkhäuser Verlag, Basel (1983).