Cauchy problem and quasi-stationary limit for the Maxwell-Landau-Lifschitz and Maxwell-Bloch equations

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Abstract. In this paper we continue the investigation of the Maxwell-Landau-Lifschitz and Maxwell-Bloch equations. In particular we extend some previous results about the Cauchy problem and the quasi-stationary limit to the case where the magnetic permeability and the electric permittivity are variable.

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1. Introduction

The models. This paper deals with two physical models which describe the propagation of electromagnetic waves, that is of the magnetic field H and of the electric field E, in some special medium which occupies an open subset Ω of $\mathbb{R}^3$, with magnetic permeability μ and electric permittivity ε. In both cases we denote by $\overline{f}$ the extension of a function f by 0 outside the set Ω. The time variable is $t \geqslant 0$, and the space variable is $x \in \mathbb{R}^3$.

The first model refers to **Maxwell-Landau-Lifschitz equations** (see [10] and [28] for Physics references). The magnetic field H and the electric field E satisfy the Maxwell equations in $\mathbb{R}^3$:

$$\begin{cases} \mu \partial_t H + \operatorname{curl} E = -\mu \partial_t \overline{M}, \\ \varepsilon \partial_t E - \operatorname{curl} H = 0, \\ \operatorname{div} \mu(H + \overline{M}) = 0, \\ \operatorname{div} \varepsilon E = 0, \end{cases} \tag{1.1}$$

where M stands for the magnetic moment in the ferromagnet Ω and takes values in the unit sphere of $\mathbb{R}^3$. It is solution to the Landau-Lifschitz equation:

$$\partial_t M = \gamma M \wedge H_T - \alpha M \wedge (M \wedge H_T) \quad \text{for } x \in \Omega, \tag{1.2}$$

where $\gamma \neq 0$ is the gyromagnetic constant, and $\alpha > 0$ is some damping coefficient. Neglecting the exchange phenomenon, the total magnetic field H_T is the sum

$$H_T = H + H_a(\overline{M}) + H_{\text{ext}}, \tag{1.3}$$

where the anisotropy field writes $H_a(\overline{M}) = \nabla_M \Phi(\overline{M})$, for some convex function Φ, and H_{ext} is some applied (exterior) magnetic field.

The second model refers to **Maxwell-Bloch equations** (see for example [7, 8, 16, 35, 38, 41]). In this setting Ω denotes some quantum medium with $N \in \mathbb{N}$ energy levels described by a Hermitian, non-negative, $N \times N$ density matrix ρ. Assuming the usual dipolar approximation, these quantum states change under the action of an electric field E by the quantum Liouville-Von Neumann (or Bloch) equation:

$$i \partial_t \rho = [\Lambda - E \cdot \Gamma, \rho] + i Q(\rho). \tag{1.4}$$

The $N \times N$ Hermitian symmetric matrix Λ, with entries in $\mathbb{C}$, represents the (electromagnetic field-) free Hamiltonian of the medium. The dipole moment operator Γ is a $N \times N$ Hermitian matrix, with entries in $\mathbb{C}^3$, and depends on the material considered. The (linear) relaxation term $Q(\rho)$ takes dissipative effects into account (see [5, 6, 32]). The polarization P of the matter is given by the constitutive law $P = \text{Tr}(\Gamma \rho)$ which influences back the electric field E. Again, the electromagnetic field satisfies the Maxwell equations in $\mathbb{R}^3$:

$$\begin{cases} \mu \partial_t H + \text{curl}\, E = 0, \\ \varepsilon \partial_t E - \text{curl}\, H = -\partial_t \overline{P}, \\ \text{div}(\varepsilon E + \overline{P}) = 0, \\ \text{div}\, \mu H = 0. \end{cases} \tag{1.5}$$

Cauchy problems. We first address the questions of global existence, uniqueness and stability for the Cauchy problem associated with these equations. The physically relevant solutions have finite energy: they satisfy the usual (L^2) energy estimates. Mathematically, this regularity leads to weak solutions and is usually not enough to ensure the desired uniqueness and stability properties (requiring for these hyperbolic semilinear systems in space dimension 3, in the general theory, H^s Sobolev regularity with $s > 3/2$).

However, in the case of the Maxwell-Landau-Lifschitz system, Joly, Métivier and Rauch [24] noticed that specific (algebraic) properties of the nonlinearities, as well as (geometric) properties of the differential operator involved, allowed to show the existence of global finite energy solutions (essentially, using compensated compactness arguments) enjoying stability properties. Furthermore, only a small amount of regularity (curl H and curl E in L^2) ensures uniqueness. This is achieved using dispersive properties of the system; namely, a limit Strichartz estimate controlling the $L_t^2 L_x^\infty$ norm of (a limited frequency part of) the fields H and E. These results were obtained for equations posed in the whole space ($\Omega = \mathbb{R}^3$) and for constant coefficients ε and μ.

In practice, the various coefficients of the system may not be constant. Typically, the magnetic permeability and electric permittivity may depend on the space variable x and have jumps across the boundary of the domain Ω.

Adapting the above mentioned compensated compactness argument, Jochmann established in [23] the existence and weak stability of global finite energy solutions for the Maxwell-Landau-Lifschitz system, considering any domain $\Omega \subset \mathbb{R}^3$, and variable, possibly discontinuous coefficients ($\varepsilon, \mu \in L^\infty(\mathbb{R}^3)$). In the (space) 2-dimensional case, we refer to the work of Haddar [21].

Concerning the Maxwell-Bloch system, the first author noticed that it shares with the Maxwell-Landau-Lifschitz some of its structural properties. This author thus showed in [18] results on existence and uniqueness of global finite energy solutions, similar to the ones of Joly, Métivier and Rauch, but for some general class of systems including the two models above. Again, these results where obtained for equations posed in the whole space and for constant coefficients ε and μ.

Here, we continue this study, again for a general class of systems including the Maxwell-Landau-Lifschitz equations and the Maxwell-Bloch equations, so as to enlight the similarities and differences between these two models. Adapting Jochmann's method, we show the existence and stability of global finite energy solutions, for a given domain $\Omega \subset \mathbb{R}^3$, and L^∞ coefficients. Then, for smooth coefficients, constant out of some compact set, we prove a limit Strichartz estimate analogous to the one obtained by Joly, Métivier and Rauch in the constant coefficient case. This allows us to show propagation of regularity and uniqueness when initially, curl H and curl E belong to $L^2(\mathbb{R}^3)$. As a corollary of a result of Saint-Raymond [40], we also infer generic uniqueness of the global finite energy solutions.

Quasi-stationary limits.   Next, we turn to the problem of the so-called quasi-stationary limit. Physically, this regime appears when the domain Ω is small compared to the wavelength. Mathematically, it amounts to some long-time asymptotics (replacing in the equations ∂_t by $\eta\partial_t$, for some small parameter η) with weak non-linearities (also scaled so as to have an amplitude of size η).

Jochmann showed in [23] the weak convergence of the corresponding solutions to the Maxwell-Landau-Lifschitz system towards the solutions of some reduced system driven by the magnetization, using the weak stability property. Starynkevitch extended this result, proving strong and global-in-time convergence in the constant coefficient case in [42], thanks to local energy estimates performed on the explicit fundamental solution of the associated wave equation. He also obtained the same result in the case of smooth coefficients, constant out of some compact set, in [43], thanks to dispersive estimates obtained from resolvent estimates on elliptic operators.

Here, we apply the same methods to our general systems to get weak and strong convergence in the quasi-stationary limit. For the latter however, some time integrability assumption is needed to conclude, which is satisfied by the Maxwell-Landau-Lifschitz system (since $\partial_t M \in L^2((0, \infty) \times \Omega)$), but we do not know if the Maxwell-Bloch system enjoys such a property.

Remark 1.1. Taking exchange energy into account, one should add to the total magnetic field in (1.3) a term $-K\Delta M$ with $K > 0$. The resulting system is then parabolic. We refer to [1, 2, 12–14] and [45] for works on the corresponding (weak or strong) Cauchy problem, long-time asymptotics and quasi-stationary limits.

2. Main results

Let us stress that we do not assume that Ω is bounded, for the moment. To deal with both the Maxwell-Landau-Lifschitz system (1.1)-(1.2) and the Maxwell-Bloch system (1.4)-(1.5), we put these two models above into a single class of systems consisting in the coupling of the Maxwell equations (with the fields H and E as unknowns) with some ODE (corresponding to a third unknown variable). The resulting sytem is symmetrizable hyperbolic, with semilinear nonlinearity, and some structure assumptions are made, such as affine dependence of the nonlinearity with respect to the electromagnetic field, and *a priori* pointwise estimates on the third unknown variable. One of the key points in our study is that the electromagnetic fields decompose into an "irrotational" part, which is directly related to this third unknown, and a "divergence free" part, which solves some wave equation.

2.1. An abstract setting

On any finite-dimensional vector space $\mathbb{R}^N$, we denote by $u \cdot u'$ the usual scalar product between vectors u and u', and by $|\cdot|$ the associated norm. For all $r > 0$, B_r denotes the (closed) ball centered at 0, with radius r.

We consider two scalar functions $\kappa_1(x)$ and $\kappa_2(x)$, which are uniformly positive:

$$\text{for } i = 1, 2, \quad \kappa_i \in L^\infty(\mathbb{R}^3), \quad \text{and} \quad \exists c > 0, \kappa_i \geqslant c. \tag{2.1}$$

We denote by H_{curl} the space of functions f in $L^2(\mathbb{R}^3, \mathbb{R}^3)$ with curl f in $L^2(\mathbb{R}^3, \mathbb{R}^3)$. We consider the operator B defined by

$$B(u_1, u_2) = (\kappa_1^{-1}\,\text{curl}\,u_2, -\kappa_2^{-1}\,\text{curl}\,u_1) \text{ for } u := (u_1, u_2) \in D(B) := H_{\text{curl}} \times H_{\text{curl}}.$$

This is a skew self-adjoint operator on the Hilbert space $L^2(\mathbb{R}^3, \mathbb{R}^6)$ endowed with the scalar product

$$\langle (u_1, u_2), (u_1', u_2') \rangle_{\kappa_1, \kappa_2} := \int_{\mathbb{R}^3} (\kappa_1 u_1 \cdot u_1' + \kappa_2 u_2 \cdot u_2') dx.$$

We denote by $P(u_1, u_2) := (P_1 u_1, P_2 u_2)$ the orthogonal projector on $(\ker B)^\perp$ with respect to the weighted scalar product above, so that for $i = 1, 2$,

$$\begin{aligned}
\text{ran } P_i &= \{u_i \in L^2(\mathbb{R}^3, \mathbb{R}^3) \mid \text{div}(\kappa_i u_i) = 0\}, \\
\text{ran } (Id - P_i) &= \{u_i \in L^2(\mathbb{R}^3, \mathbb{R}^3) \mid \text{curl}(u_i) = 0\}.
\end{aligned} \tag{2.2}$$

We consider a function $F : \mathbb{R}^3 \times \mathbb{R}^d \times \mathbb{R}^6 \to \mathbb{R}^d$, where $d \in \mathbb{N}$, affine in its third variable, and written

$$F(x, v, u) = F_0(x, v) + F_1(x, v)u. \tag{2.3}$$

For each $j = 0, 1$, F_j is measurable with respect to x and continuously differentiable with respect to v. Furthermore,

$$\text{for } j = 0, 1, \quad \text{for almost all } x \in \mathbb{R}^3, \ F_j(x, 0) = 0,$$
$$\text{and} \quad \forall R > 0, \ \text{for almost all } x \in \mathbb{R}^3, \ \forall v \in B_R, \tag{2.4}$$
$$|F_j(x, v)| + |\partial_v F_j(x, v)| \leqslant C_F(R).$$

Finally, we assume that there exists $K \geqslant 0$ such that:

$$\text{for almost all } x \in \mathbb{R}^3, \forall (v, u) \in \mathbb{R}^d \times \mathbb{R}^6, \quad F(x, v, u) \cdot v \leqslant K|v|^2. \tag{2.5}$$

Remark 2.1. The constant K above may sometimes be taken equal to zero. In this case, Estimate (i) in Theorem 2.6 is improved, since v does not undergo any growth. This is the case for the Maxwell-Landau-Lifschitz model, as well as for the Maxwell-Bloch model, when only transverse relaxation is taken into account ($Q(\rho) = -\gamma \rho_{od}$, for some $\gamma \geqslant 0$, and with ρ_{od} the off-diagonal part of ρ).

We also consider a function $l = (l_1, l_2) \in (L^\infty(\mathbb{R}^3, L(\mathbb{R}^d, \mathbb{R}^3)))^2$, where $L(\mathbb{R}^d, \mathbb{R}^3)$ denotes the space of linear functions from $\mathbb{R}^d$ to $\mathbb{R}^3$. We introduce the following shorthand notation: for any $x \in \mathbb{R}^3$, $(\kappa^{-1} \cdot l)(x)$ is the mapping from $\mathbb{R}^d$ to $\mathbb{R}^6$, such that

for almost all $x \in \mathbb{R}^3$, $\forall v \in \mathbb{R}^d$, $(\kappa^{-1} \cdot l)(x)v := (\kappa_1(x)^{-1} l_1(x)v, \kappa_2(x)^{-1} l_2(x)v)$.

Then, for any $U := (u, v)$ in

$$\mathbf{L}^2 := L^2(\mathbb{R}^3, \mathbb{R}^6) \times L^2(\Omega, \mathbb{R}^d),$$

the conditions

$$\text{div}(\kappa_1 u_1 - l_1 \overline{v}) = 0, \quad \text{div}(\kappa_2 u_2 - l_2 \overline{v}) = 0,$$

may be equivalently written

$$(Id - P)(u - (\kappa^{-1} \cdot l)\overline{v}) = 0. \tag{2.6}$$

We look for $U \in C([0, \infty), \mathbf{L}^2)$, solution to

$$(\partial_t + B)u = (\kappa^{-1} \cdot l)F(x, \overline{v}, u) \quad \text{for } x \in \mathbb{R}^3, \tag{2.7}$$
$$\partial_t v = F(x, v, u) \quad \text{for } x \in \Omega, \tag{2.8}$$

and (2.6). We recall that for any function f defined on Ω, we denote by $\overline{f}$ the extension of f (to $\mathbb{R}^3$) by 0 outside Ω. We also impose the condition

$$v \in L^\infty_{\text{loc}}((0, \infty), L^\infty(\Omega, \mathbb{R}^d)). \tag{2.9}$$

Equations (2.6)-(2.8) are understood in the distributional sense, noticing that (2.9) gives sense to the nonlinear term, since the function $F(x, v, u)$ is affine in u.

Remark 2.2. Equations (2.6)-(2.8) reduce to the Maxwell-Landau-Lifschitz system (1.1)-(1.2) when $u_1 = H$, $u_2 = E$, $v = M$ (with $d = 3$), $\kappa_1 = \mu$, $\kappa_2 = \varepsilon$, $l_1 = -\mu$, $l_2 = 0$, $F(x, v, u) = \gamma v \wedge (u_1 + H_a(v) + H_{\text{ext}}) - \alpha v \wedge (v \wedge (u_1 + H_a(v) + H_{\text{ext}}))$ and to the Maxwell Bloch system (1.4)-(1.5) when $u_1 = H$, $u_2 = E$, $v = \rho$ (with $d = N^2$), $\kappa_1 = \mu$, $\kappa_2 = \varepsilon$, $l_1 = 0$, $l_2 = \text{Tr}(\Gamma \cdot)$, $F(x, v, u) = -i[\Lambda - u_2 \cdot \Gamma, v] + Q(v)$. The exterior magnetic field above is usually depending on time. We did not consider such time-dependent coefficients in our study, since it would have made notations more intricate; up to some integrability assumptions, this extension is straightforward.

Definition 2.3. We call $U = (u, v) \in C([0, \infty), \mathbf{L}^2)$ a global finite energy solution to (2.6)-(2.8) if (2.9) holds true and U is a solution to (2.6)-(2.8) in the distributional sense.

Remark 2.4. Equation (2.6) has to be seen as a (linear) constraint, which propagates from $t = 0$ for solutions to (2.7)-(2.8):

$$\partial_t (Id - P)(u - (\kappa^{-1} \cdot l)\overline{v}) = 0. \tag{2.10}$$

Indeed, by definition of the projector P, we have $(Id - P)B = 0$, so that we get (2.10) when applying $(Id - P)$ to (2.7), using (2.8) (which extends to all $x \in \mathbb{R}^3$ since $F(x, 0, u) \equiv 0$) and commuting the derivative ∂_t with $(Id - P)$ and $\kappa^{-1} \cdot l$.

We therefore have to consider initial data U_{init} satisfying (2.6), and for such constrained initial data, the solutions to (2.7)-(2.8) also satisfy (2.6) as long as they exist. We shall write $U_{\text{init}} := (u_{\text{init}}, v_{\text{init}})$ with $u_{\text{init}} := (u_{\text{init},1}, u_{\text{init},2})$.

Definition 2.5. Let L_{div} be the set of functions $U := (u, v) \in L^2(\mathbb{R}^3, \mathbb{R}^6) \times (L^2(\Omega, \mathbb{R}^d) \cap L^\infty(\Omega, \mathbb{R}^d))$ satisfying (2.6).

2.2. Cauchy problems

Our first result states the existence of global finite energy solutions to (2.6)-(2.8).

Theorem 2.6. *Assume* (2.1) *and* (2.3)-(2.5). *For any U_{init} in L_{div}, there exists $U := (u, v) \in C([0, \infty), \mathbf{L}^2)$, global finite energy solution to* (2.6)-(2.8) *with U_{init} as initial data. Moreover, for all $T > 0$, there is $C = C(T, F, l, \|v_{\text{init}}\|_{L^\infty})$ such that*

(i) *for almost all $x \in \mathbb{R}^3$, for all $t \geqslant 0$, $|v(t, x)| \leqslant |v_{\text{init}}(x)|e^{Kt}$ (with K from* (2.5));
(ii) *for all $t \in [0, T]$, $\|(u, v)(t)\|_{\mathbf{L}^2} \leqslant C\|U_{\text{init}}\|_{\mathbf{L}^2}$;*
(iii) *$v \in W^{1,\infty}_{\text{loc}}((0, \infty), L^2(\Omega, \mathbb{R}^d))$, and for almost every $t \in [0, T]$, $\|\partial_t v(t)\|_{L^2(\Omega)} \leqslant C\|U_{\text{init}}\|_{\mathbf{L}^2}$.*

Finally, if $\mathcal{U}_{\text{init}}$ is a bounded set of L_{div} which is compact in $\mathbf{L}^2$, then for all $T > 0$, the set $\mathcal{U}$ of the above solutions with Cauchy data in $\mathcal{U}_{\text{init}}$ is compact in $C([0, T], \mathbf{L}^2)$.

To establish this first result, we follow the strategy of Jochmann in [23], which is itself an improvement of the method by Joly, Métivier and Rauch in [24]. This is the classical regularization method, in which (global-in-time) approximate solutions $U^n = (u^n, v^n)$ are built first (Section 3.2); the delicate step consists of course in passing to the limit $n \to \infty$ in the regularization (Section 3.3). Pointwise bounds are available for v^n, which imply L^p bounds for $(Id - P)u^n = (Id - P)(\kappa^{-1} \cdot l)\overline{v^n}$, for finite p. The main argument relies on compensated compactness, applied to Pu^n (Lemma 3.4).

As a byproduct of the proof of Theorem 2.6, we also have the following version of stability, where we assume strong convergence only for the v part of the initial data. It will be useful below (*cf.* proof of Theorem 2.10) when considering the weak quasi-stationary limit.

Proposition 2.7. *Let* $(U^n)_{n \in \mathbb{N}}$ *be a sequence in* $L^\infty_{\text{loc}}((0, \infty), L_{\text{div}})$, *bounded in* $L^\infty_{\text{loc}}((0, \infty), \mathbf{L}^2)$, *with* $(v^n)_{n \in \mathbb{N}}$ *bounded in* $W^{1,\infty}_{\text{loc}}((0, \infty), L^2(\Omega)) \cap L^\infty_{\text{loc}}((0, \infty),$ $L^\infty(\Omega))$ *satisfying* (2.8), $v^n|_{t=0} \to v_{\text{init}}$ *in* $L^2(\Omega)$ *and* $Bu^n = \partial_t D^n$ *with* $(D^n)_n$ *bounded in* $L^\infty_{\text{loc}}((0, \infty), L^2(\mathbb{R}^3))$. *Then, up to a subsequence,* v^n *converges to* v *in* $L^\infty_{\text{loc}}((0, \infty), L^p(\Omega))$ *for any* $p \geqslant 2$ *and in* $L^\infty_{\text{loc}}((0, \infty), L^\infty(\Omega))$ *weak* $*$, u^n *converges to* u *in* $L^\infty_{\text{loc}}((0, \infty), L^2(\mathbb{R}^3))$ *weak* $*$, *and* $U := (u, v)$ *satisfies* (2.6), (2.8), *as well as* $v|_{t=0} = v_{\text{init}}$.

Let us now turn our atttention to smoother solutions. We need to assume more smoothness on the coefficients ε and μ. Of course, when considering in some physical situation a domain Ω with boundaries, the coefficients ε and μ experiment discontinuity jumps. Since we do not know how to tackle this physical case, we assume from now on that

$$\Omega \text{ is bounded,} \quad \text{and with } \mathcal{K} = \overline{\Omega}, \quad \kappa_i - 1 \in \mathcal{C}^\infty_\mathcal{K}(\mathbb{R}^3), \ i = 1, 2, \qquad (2.11)$$

where $\mathcal{C}^\infty_\mathcal{K}(\mathbb{R}^3)$ is the space of $\mathcal{C}^\infty$ functions on $\mathbb{R}^3$, with support contained in $\mathcal{K}$. In order to get a uniqueness result, we only need to ensure that the "divergence free" part Pu of the fields has the H^1 regularity. To this end, once a finite energy solution is given, we make use of the linear system solved by Pu, with coefficients depending on the rest of the solution. As in [24], we proceed in two steps: we begin with the propagation of H^μ regularity, for $\mu \in (0, 1)$, using Strichartz estimates (Proposition 4.6). Applying this result with $\mu = 1/2$ provides enough integrability for the coefficients of the above mentioned linear equation to ensure propagation of H^1 regularity. This implies that u is "almost" L^∞, a natural condition to prove uniqueness of the solution. Technically, a L^∞ approximation of u is built thanks to a limit Strichartz estimate for low frequencies (Proposition 4.7). We also need a decoupling assumption, which was introduced in [18], and is satisfied by the Maxwell-Landau-Lifschitz system as well as by the Maxwell-Bloch system.

Theorem 2.8. *In addition to the assumptions of Theorem 2.6, assume* (2.11). *Let* $\mu \in]0, 1]$, *and* $U_{\text{init}} \in L_{\text{div}}$ *with* $\text{curl}\, u_{\text{init},i} \in H^{\mu-1}(\mathbb{R}^3)$, *for* $i = 1, 2$. *Then, the following holds true*:

(a) *Any solution U to* (2.6)-(2.8) *with U_{init} as initial data given by Theorem* 2.6 *satisfies* $\text{curl}\, u_i \in C([0, \infty], H^{\mu-1}(\mathbb{R}^3))$, *for* $i = 1, 2$.
(b) *If $\mu = 1$, assuming moreover that*

$$
\begin{aligned}
&\text{there exists } j \in \{1, 2\} \text{ such that } l_{3-j}F = 0 \\
&\text{and such that } F \text{ depends only on } (x, v, u_j),
\end{aligned}
\tag{2.12}
$$

there exists only one solution to (2.6)-(2.8) *with U_{init} as initial data as in Theorem* 2.6.

Theorem 2.8 asserts that the uniqueness property holds for initial data U_{init} in L_{div} with $\text{curl}\, u_{\text{init},i} \in L^2(\mathbb{R}^3)$, $i = 1, 2$, which are dense in L_{div} for the topology of $\mathbf{L}^2$. The following theorem says that the uniqueness property even holds generically for the following topologies. Let τ_s and τ_w denote respectively the strong and weak topologies of $L^2(\mathbb{R}^3, \mathbb{R}^6)$ and let $\tilde{\tau}_s$ denote the strong topology of $L^2(\Omega, \mathbb{R}^d)$. We consider the product topology τ_{ss} (respectively τ_{ws}) on $\mathbf{L}^2$ obtained from τ_s (respectively τ_w) and $\tilde{\tau}_s$.

Theorem 2.9. *Under the assumptions of Theorem 2.8 (b), for any $C_{\text{init}} > 0$, there exists a G_δ dense set $\widetilde{L_{\text{div}}}$ in the set $\{U_{\text{init}} \in L_{\text{div}} \mid \|v_{\text{init}}\|_{L^\infty(\Omega)} \leqslant C_{\text{init}}\}$ for the topology τ_{ss} and τ_{ws}, such that for any $U_{\text{init}} \in \widetilde{L_{\text{div}}}$, there exists only one solution to* (2.6)-(2.8) *with U_{init} as initial data, with the same properties as in Theorem* 2.6.

Let us stress that we cannot expect that the problem (2.6)-(2.8) admits smoother solutions than the ones given by Theorem 2.8 since, by definition, $\overline{v}$ is discontinuous across the boundary $\partial\Omega$. However it follows from the general theory of discontinuous solutions of hyperbolic semilinear systems [34, 39] that the problem (2.6)-(2.8) admits piecewise regular solutions discontinuous accross $\partial\Omega$ (let us also refer to the appendix of [44]). Yet the general theory only guaranties local-in-time solutions. We do not know if in the particular case of the problem (2.6)-(2.8) global-in-time solutions can be obtained.

2.3. Quasi-stationary limits

As described in the Introduction, the quasi-stationary regime consists in the limit $\eta \to 0^+$ for (2.6)-(2.8), where ∂_t is replaced with $\eta\partial_t$, and F is replaced with ηF. Equations (2.6) and (2.8) are invariant under this rescaling, whereas (2.7) becomes

$$
(\eta\partial_t + B)u = \eta(\kappa^{-1} \cdot l)F(x, \overline{v}, u), \quad \text{for } x \in \mathbb{R}^3.
\tag{2.13}
$$

For this semi-classical version of (2.8), it is still true that the constraint (2.6) is propagated from the initial data. Formally, in the limit $\eta \to 0^+$, v still satisfies

(2.8), whereas u satifies (2.6) and $Bu = 0$. But for $U = (u, v) \in C([0, \infty), \mathbf{L}^2)$, these last two conditions are equivalent to the fact that for all $t \geqslant 0$, $u(t)$ is directly determined by $v(t)$, and more precisely:

$$u = (Id - P)u = (Id - P)(\kappa^{-1} \cdot l)\overline{v}, \tag{2.14}$$

with P from (2.2). Then, (2.8) becomes

$$\partial_t v = F(x, v, ((Id - P)(\kappa^{-1} \cdot l)\overline{v})_{|\Omega}). \tag{2.15}$$

Using the stability result given by Proposition 2.7, we have a first result of convergence towards the quasi-stationary limit, weakly for u and locally in time for v:

Theorem 2.10. *Assume (2.1) and (2.3)-(2.5). For any U_{init} in L_{div}, for any $\eta \in (0, 1)$, let $U^\eta := (u^\eta, v^\eta)$ be a global finite energy solution to (2.6), (2.8) and (2.13) with U_{init} as initial data. Then, up to a subsequence, v^η converges in $L^\infty_{\text{loc}}((0, \infty), L^p(\Omega))$ for all $p \geqslant 2$ and in $L^\infty_{\text{loc}}((0, \infty), L^\infty(\Omega))$ weak $*$ towards a solution v to (2.15), with v_{init} as initial data; Pu^η converges to 0 in $L^\infty_{\text{loc}}((0, \infty), L^2(\Omega))$ weak $*$, and $(Id - P)u^\eta = (Id - P)(\kappa^{-1} \cdot l)\overline{v^\eta}$ converges in $L^\infty_{\text{loc}}((0, \infty), L^p(\mathbb{R}^3))$ for all $p \geqslant 2$ towards u, given by (2.14).*

Convergence of the whole sequence U^η is ensured as soon as the Cauchy problem associated with the limiting equation (2.15) has a unique solution. This is given by the following proposition, which extends [43, Theorem 3.1] by Starynkevitch.

Proposition 2.11. *Assume (2.11), and let $v_{\text{init}} \in L^\infty(\Omega)$. Then, there is a unique $v \in C([0, \infty), L^2(\Omega)) \cap L^\infty_{\text{loc}}((0, \infty), L^\infty(\Omega))$ solution to (2.15), with v_{init} as initial data.*

We also prove strong and global-in-time convergence for u, assuming (2.11) again, as well as integrability in time for $\|\partial_t v\|^2_{L^2(\Omega)}$ and non-trapping for some wave operator:

Theorem 2.12. *Under the assumptions of Theorem 2.10, assume moreover (2.11), the non-trapping hypothesis (6.6) and that $\partial_t v^\eta$ is bounded (with respect to η) in $L^2((0, \infty) \times \Omega)$. Then, Pu^η goes to zero in $L^2((0, \infty), L^2_{\text{loc}}(\mathbb{R}^3))$.*

Remark 2.13. In the case of the Maxwell-Landau-Lifschitz system (1.1)-(1.2), $\partial_t M$ actually belongs to $L^2((0, \infty) \times \Omega)$. Define the energy $\mathcal{E}(t)$ as

$$\mathcal{E}(t) = \frac{1}{2} \int_{\mathbb{R}^3} (\varepsilon |E|^2 + \mu |H|^2)\, dx + \int_\Omega \mu \left(\Phi(M) + \frac{1}{2}|H_{\text{ext}} - M|^2 \right) dx.$$

Differentiating formally this expression with respect to time, we see that the integral of $H \cdot \text{curl}\, E - E \cdot \text{curl}\, H = \text{div}(E \wedge H)$ vanishes, as well as $M \cdot \partial_t M$ (since $|M|$ is constant). Using the orthogonality relations of the nonlinearity, we get

$$\partial_t M \cdot H_T = \alpha |M \wedge H_T|^2 \quad \text{and} \quad |\partial_t M|^2 = (\alpha^2 + \gamma^2)|M \wedge H_T|^2,$$

so that estimate (ii) in Theorem 2.6 is improved to

$$\mathcal{E}(t) + \frac{\alpha}{\alpha^2 + \gamma^2} \int_0^t \|\sqrt{\mu}\partial_t M(t')\|_{L^2(\Omega)}^2 dt' - \frac{1}{2}\|\sqrt{\mu}H_{\text{ext}}(t)\|_{L^2(\Omega)}^2$$
$$+ \int_0^t \left(\int_\Omega \mu M \cdot \partial_t H_{\text{ext}} dx \right) dt' = cst,$$

and the same is true with the quasi-stationary scaling. Assuming for example that $H_{\text{ext}} \in L_t^\infty L_x^2$ and $\partial_t H_{\text{ext}} \in L_{t,x}^1$, we deduce that $\mathcal{E}$ is bounded, and $\partial_t M$ belongs to $L^2((0, \infty) \times \Omega)$.

In the case of the Maxwell-Bloch system, we do not know if such an estimate is available for $\partial_t \rho$.

3. Existence of global finite energy solutions: proof of Theorem 2.6

3.1. Technical interlude 1

3.1.1. *Intersections and sums of Banach spaces*

We recall some useful properties of the intersection and the sum of Banach spaces. Consider two Banach spaces X_1 and X_2 that are subsets of a Hausdorff topological vector space X. Then

$$X_1 \cap X_2 := \{f \in X \mid f \in X_1, \ f \in X_2\}$$
$$(\text{respectively } X_1 + X_2 := \{f \in X \mid E(f) \neq \emptyset\},$$
$$\text{where } E(f) := \{(f_1, f_2) \in X_1 \times X_2 \mid f_1 + f_2 = f\})$$

is a Banach space endowed with the norm

$$\|f\|_{X_1 \cap X_2} := \|f\|_{X_1} + \|f\|_{X_2}$$
$$(\text{respectively } \|f\|_{X_1 + X_2} := \inf\{\|f_1\|_{X_1} + \|f_2\|_{X_2} \mid (f_1, f_2) \in E(f)\}).$$

If furthermore $X_1 \cap X_2$ is a dense subset of both X_1 and X_2, then $(X_1 \cap X_2)' = X_1' + X_2'$ and $(X_1 + X_2)' = X_1' \cap X_2'$ (*cf.* Bergh and Löfström [4, Lemma 2.3.1 and Theorem 2.7.1]).

3.1.2. *Mollifiers*

We shall use the following symmetric operators $R^n : L^2(\mathbb{R}^3) \to L^2(\mathbb{R}^3)$, defined by

$$(R^n f)(x) := \int_{\mathbb{R}^3} f(y) w^n(x - y) dy \quad \text{for } x \in \mathbb{R}^3, \tag{3.1}$$

where $w^n \in C_0^\infty(\mathbb{R}^3)$ is a mollifier with supp $w^n \subset B(0, 1/(1+n))$ and $\int_{\mathbb{R}^3} w^n = 1$. These operators have the following well-known properties: there exists $C > 0$ such that for all $f \in L^2(\mathbb{R}^3)$, $r > 1$ and $n \in \mathbb{N}$,

$$\|f - R^n f\|_{L^2(\mathbb{R}^3)} \to 0, \qquad \|R^n f\|_{L^2(\mathbb{R}^3)} \leqslant C\|f\|_{L^2(\mathbb{R}^3)}, \qquad (3.2)$$

$$\|R^n f\|_{L^2(B_r)} \leqslant C\|f\|_{L^2(B_{r+1})}, \text{ and } \|R^n f\|_{L^2(\mathbb{R}^3 \setminus B_r)} \leqslant C\|f\|_{L^2(\mathbb{R}^3 \setminus B_{r-1})}. \quad (3.3)$$

Moreover for all $n \in \mathbb{N}$, there exists $C_n > 0$ such that for all $f \in L^2(\mathbb{R}^3)$,

$$\|R^n f\|_{L^\infty(\mathbb{R}^3)} \leqslant C_n \|f\|_{L^2(\mathbb{R}^3)}. \qquad (3.4)$$

3.2. Approximate solutions

The following lemma claims the existence of global solutions to some regularized problem.

Lemma 3.1. *For all $n \in \mathbb{N}$, there exists $U^n := (u^n, v^n) \in C([0, \infty), \mathbf{L}^2)$, with*

$$v^n \in C([0, \infty), L^\infty(\Omega, \mathbb{R}^d)) \cap C^1([0, \infty), L^2(\Omega, \mathbb{R}^d)), \qquad (3.5)$$

solution to the regularized problem:

$$(\partial_t + B)u^n = (\kappa^{-1} \cdot l)F^n \quad \text{for } x \in \mathbb{R}^3, \qquad (3.6)$$

$$\partial_t v^n = F^n \qquad \text{for } x \in \Omega, \qquad (3.7)$$

where

$$F^n(t, x) := F(x, \overline{v^n}(t, x), R^n u^n(t, x)), \qquad (3.8)$$

with U_{init} as initial data. Moreover, for all $n \in \mathbb{N}$,

(a) *For almost all $x \in \mathbb{R}^3$, for all $t \geqslant 0$, $|v^n(t, x)| \leqslant |v_{\text{init}}(x)|e^{Kt}$ (with K from (2.5)).*

(b) *For all $T > 0$, there is $C = C(T, F, l, \kappa, \|v_{\text{init}}\|_{L^\infty})$ such that, for all $t \in [0, T]$, $\|(u^n, v^n)(t)\|_{\mathbf{L}^2} + \|\partial_t v^n(t)\|_{L^2(\Omega)} \leqslant C\|U_{\text{init}}\|_{\mathbf{L}^2}$.*

Proof. The local-in-time solution is constructed *via* a usual fixed point argument for the mapping $\mathcal{A}^n : C([0, T], \mathbf{L}^2) \to C([0, T], \mathbf{L}^2)$,

$$\mathcal{A}^n(u, v)(t, \cdot)$$
$$= \left(\exp(-tB)u_{\text{init}} + \int_0^t \exp((t-t')B)(\kappa^{-1} \cdot l)\mathcal{F}^n(t', \cdot)dt', \ v_{\text{init}} + \int_0^t \mathcal{F}^n(t', \cdot)dt'\right),$$

where

$$\mathcal{F}^n(t, \cdot) := F(\cdot, \overline{v}(t, \cdot), R^n u(t, \cdot)).$$

For $T > 0$ small enough, $\mathcal{A}^n$ is shown to be a contraction mapping thanks to properties (2.3)-(2.5) of F, (3.4), and because B is a skew self-adjoint operator in the Hilbert space $L^2(\mathbb{R}^3, \mathbb{R}^6)$ endowed with the scalar product $\langle \cdot, \cdot \rangle_{\kappa_1, \kappa_2}$.

Global existence is given by the *a priori* bounds (a) and (b). The first one follows directly from (2.5) and Gronwall's lemma. In the same way, taking the L^2 norm of $(u^n, v^n)(t) = \mathcal{A}^n(u^n, v^n)(t)$, one gets

$$\|(u^n, v^n)(t)\|_{\mathbf{L}^2} \leqslant \|U_{\text{init}}\|_{\mathbf{L}^2} + \int_0^t (1 + \|\kappa^{-1} \cdot l\|_{L^\infty}) \|F^n(t', \cdot)\|_{L^2} dt'.$$

One may add to this inequality the one obtained from (3.7),

$$\|\partial_t v^n(t, \cdot)\|_{L^2} \leqslant \|F^n(t, \cdot)\|_{L^2}.$$

From (2.3), (2.4), we have

$$\|F^n(t, \cdot)\|_{L^2} \leqslant C_F(\|v_{\text{init}}\|_{L^\infty} e^{Kt}) \|(u^n, v^n)(t)\|_{\mathbf{L}^2},$$

so that Gronwall's lemma concludes. $\qquad\qquad\qquad\qquad\qquad\qquad\qquad\square$

3.3. Passing to the limit $n \to \infty$

Let us stress that Estimate (3.4) is not uniform with respect to n. However, we have:

Proposition 3.2. *For all $T > 0$, there is a subsequence of $(U^n)_{n \in \mathbb{N}}$ given by Lemma 3.1 that strongly converges in $C([0, T], \mathbf{L}^2)$ to $U := (u, v) \in C([0, \infty), \mathbf{L}^2)$, global finite energy solution to (2.6)-(2.8) with U_{init} as initial data, and satisfying the estimates* (i), (ii), (iii) *of Theorem 2.6.*

Proof. First we infer from the bounds (a)-(b) in Lemma 3.1 that there exists a subsequence, still denoted (u^n, v^n), such that u^n (respectively F^n) tends to u (respectively to F_{lim}) in $L^\infty((0, T), \mathbf{L}^2)$ weak $*$ (respectively $L^\infty((0, T), L^2(\mathbb{R}^3))$ weak $*$) and v^n tends to v in $W^{1,\infty}((0, T), L^2(\Omega))$ weak $*$ and in $L^\infty((0, T), L^\infty(\Omega))$ weak $*$. This is enough to ensure that (u, v) satisfies (2.6). Moreover, Fatou's lemma yields that u and v satisfy (i)-(ii) of Theorem 2.6 for almost every t in $(0, T)$.

Since the function F is not linear, these weak limits do not suffice to pass to the limit in Equation (3.7). The strategy is to carefully study the nonlinear term F^n to prove that the solutions U^n of the regularized problems (3.6)-(3.7) actually converge (strongly) in L^2. The key step consists in proving the strong convergence of v^n.

It shall be useful several times to keep in mind that, thanks to the growth conditions (2.4) on F and to the pointwise bound Lemma 3.1, (a) of the v^n, there holds, for all $n, m \in \mathbb{N}$, for all $(t, x) \in [0, T] \times \mathbb{R}^3$,

$$|F_i^n(t, x)| \leqslant C_F(e^{KT} \|v_{\text{init}}\|_{L^\infty}), \tag{3.9}$$

$$|F_i^n(t, x) - F_i^m(t, x)| \leqslant C_F(e^{KT} \|v_{\text{init}}\|_{L^\infty}) |\overline{v^n}(t, x) - \overline{v^m}(t, x)|, \tag{3.10}$$

where $F_i^n = F_i(x, \overline{v^n})$, $i = 0, 1$.

Strong convergence of v^n. We perform energy estimates on (3.6)-(3.7). Since u may be unbounded, we introduce a weight function, which precisely depends on u. More exactly, we choose a positive function $\rho_0(x)$ in $L^\infty(\mathbb{R}^3) \cap L^2(\mathbb{R}^3)$ and define

$$\rho(t,x) := \rho_0(x)e^{-L\int_0^t |u(s,x)|ds}, \tag{3.11}$$

with $L \geqslant C_F(e^{KT}\|v_{\text{init}}\|_{L^\infty})$. First, using (3.7) we get

$$\frac{1}{2}\frac{d}{dt}\left(\|\rho(v^n - v^m)\|_{L^2(\Omega)}^2\right)(t)$$
$$= \int_\Omega \rho^2(v^n - v^m) \cdot (F^n - F^m)dx - L\int_\Omega \rho^2|u||v^n - v^m|^2dx.$$

Next, decompose $F^n - F^m$ according to (2.3) to get

$$\frac{1}{2}\frac{d}{dt}\left(\|\rho(v^n - v^m)\|_{L^2(\Omega)}^2\right)(t) = \int_\Omega \rho^2(v^n - v^m) \cdot (F_0^n - F_0^m)dx$$
$$+ \int_\Omega \rho^2(v^n - v^m) \cdot (F_1^n R^n u^n - F_1^m R^m u^m)dx - L\int_\Omega \rho^2|u||v^n - v^m|^2dx. \tag{3.12}$$

The first term in the right hand side of (3.12) can be estimated by $C\|\rho(v^n - v^m)\|_{L^2(\Omega)}^2(t)$ thanks to (3.10). Now, decompose $F_1^n R^n u^n - F_1^m R^m u^m$ into

$$F_1^n R^n u^n - F_1^m R^m u^m = F_1^n\left(R^n u^n - u\right) - F_1^m\left(R^m u^m - u\right) + \left(F_1^n - F_1^m\right)u.$$

The terms produced by the third parenthesis are estimated thanks to (3.10), and absorbed by the last term in (3.12), so that

$$\frac{1}{2}\frac{d}{dt}\left(\|\rho(v^n - v^m)\|_{L^2(\Omega)}^2\right)(t) \leqslant C\|\rho(v^n - v^m)\|_{L^2(\Omega)}^2(t)$$
$$+ \int_\Omega \rho^2(v^n - v^m) \cdot F_1^n(R^n u^n - u)dx$$
$$+ \int_\Omega \rho^2(v^m - v^n) \cdot F_1^m(R^m u^m - u)dx.$$

Then, decompose $R^n u^n$ and $R^m u^m$ according to the orthogonal projector P to get

$$\frac{1}{2}\frac{d}{dt}\left(\|\rho(v^n - v^m)\|_{L^2(\Omega)}^2\right)(t) \leqslant C\|\rho(v^n - v^m)\|_{L^2(\Omega)}^2(t)$$
$$+ \sum_{j=1}^3 h_{j,m,n}(t) + h_{j,n,m}(t), \tag{3.13}$$

where

$$h_{1,m,n}(t) := \int_{\Omega} \rho^2(v^n - v^m) \cdot F_1^n R^n P(u^n - u)dx,$$

$$h_{2,m,n}(t) := \int_{\Omega} \rho^2(v^n - v^m) \cdot F_1^n R^n(Id - P)(u^n - u)dx,$$

$$h_{3,m,n}(t) := \int_{\Omega} \rho^2(v^n - v^m) \cdot F_1^n(R^n u - u)dx.$$

The following lemma deals with the term $h_{1,m,n}(t)$.

Lemma 3.3. *There holds*

$$\forall \delta > 0, \ \exists N_\delta \in \mathbb{N}, \ \forall n \geqslant N_\delta, \ \forall m \in \mathbb{N}, \qquad \left| \int_0^T h_{1,m,n}(t)dt \right| \leqslant 2\delta. \qquad (3.14)$$

Proof. First notice that

$$h_{1,m,n}(t) = \int_{\mathbb{R}^3} R^n \left(\rho^2(\overline{v^n} - \overline{v^m}) \cdot F_1^n \right) P(u^n - u)dx.$$

We first handle the case where x is outside of a large ball. Using the Cauchy-Schwarz inequality, the second property of R^n in (3.3), the uniform bound in $L^\infty([0, T], L^2(\mathbb{R}^3))$ for v^n given in Lemma 3.1, (b) and the bound (3.9) for F_1^n, we get that

$$\int_0^T \int_{\mathbb{R}^3 \setminus B_r} \left| R^n \left(\rho^2(\overline{v^n} - \overline{v^m})F_1^n \right) \cdot P(u^n - u) \right| dxdt$$

$$\leqslant C \int_0^T \| \rho(t)^2 \|_{L^2(\mathbb{R}^3 \setminus B_{r-1})} dt. \qquad (3.15)$$

By definition of ρ there exists $r > 0$ so that this integral is less than δ.

It remains to tackle the case where $x \in B_r$. We use the following compactness lemma:

Lemma 3.4 (Jochmann[23, Lemma 3.4]). *Let $(G^n)_{n \in \mathbb{N}}$ and $(K^n)_{n \in \mathbb{N}}$ be bounded sequences in $L^\infty([0, T), L^2(\mathbb{R}^3, \mathbb{R}^6))$, with K^n converging to 0 in $L^\infty([0, T), L^2(\mathbb{R}^3, \mathbb{R}^6))$ weak $*$. Suppose that $(G^n)_{n \in \mathbb{N}}$ is equicontinuous from $[0, T]$ to $L^2(\mathbb{R}^3, \mathbb{R}^6)$ and that $BK^n = \partial_t C^n$ with $(C^n)_{n \in \mathbb{N}}$ bounded in $L^\infty([0, T), L^2(\mathbb{R}^3, \mathbb{R}^6))$. Then for all $r > 0$,*

$$\sup_{p \in \mathbb{N}} \left| \int_0^T \int_{B_r} G^p(t) \cdot PK^n(t)dxdt \right| \underset{n \to \infty}{\longrightarrow} 0. \qquad (3.16)$$

Let us denote $G^{k,l} = R^k \left(\rho^2 (\overline{v^k} - \overline{v^l}) \cdot F_1^k \right)$ and $K^n = u^n - u$. Thanks to (3.9), (3.2) and to Lemma 3.1, (b), we get that $(G^{k,l})_{k,l \in \mathbb{N}}$ and $(K^n)_{n \in \mathbb{N}}$ are bounded in $L^\infty([0,T), L^2(\mathbb{R}^3))$. Moreover K^n tends to zero in $L^\infty([0,T), L^2(\mathbb{R}^3))$ weak $*$, by definition of u. Let us denote $\mathcal{F}^n = \int_0^t (\kappa^{-1} \cdot l) F^n dt'$ and $\mathcal{F} = \int_0^t (\kappa^{-1} \cdot l) F_{\lim} dt'$. From (3.6) we infer that

$$BK^n = \partial_t C^n, \quad \text{with} \quad C^n := \mathcal{F}^n - u^n - (\mathcal{F} - u).$$

The sequence $(C^n)_{n \in \mathbb{N}}$ is bounded in $L^\infty([0,T), L^2(\mathbb{R}^3))$. In the same way, equicontinuity is obtained from the bounds on $\partial_t v^n = F^n$. We therefore apply the lemma observing that, for all $m, n \in \mathbb{N}$,

$$\left| \int_0^T \int_{B_r} G^{m,n}(t) \cdot PK^n(t) dx dt \right| \leqslant \sup_{k,l \in \mathbb{N}} \left| \int_0^T \int_{B_r} G^{k,l}(t) \cdot PK^n(t) dx dt \right|.$$

Lemma 3.4 therefore ensures that there is $N_{r,\delta} \in \mathbb{N}$ such that, for $n \geqslant N_{r,\delta}$ and for all $m \in \mathbb{N}$,

$$\left| \int_0^T \int_{B_r} R^n \left(\rho^2 (\overline{v^n} - \overline{v^m}) F_1^n \right) \cdot P(u^n - u) dx dt \right| \leqslant \delta,$$

and Lemma 3.3 is proved. $\qquad\square$

We now deal with the term $h_{2,m,n}(t)$.

Lemma 3.5. *There holds*

$$\forall \delta > 0, \ \exists N_\delta \in \mathbb{N}, \ \forall n \geqslant N_\delta, \ \forall m \in \mathbb{N},$$
$$\left| \int_0^T h_{2,m,n}(t) dt \right| \leqslant \delta + C \| \rho(\overline{v^n} - \overline{v^m}) \|_{L^2_{t,x}} \| \rho(\overline{v^n} - \overline{v}) \|_{L^2_{t,x}}. \tag{3.17}$$

Proof. The (u^n, v^n) satisfy (2.6) and so does their weak limit (u, v). Thus

$$h_{2,m,n}(t) = -\int_\Omega \rho^2(v^n - v^m) \cdot F_1^n R^n (Id - P)(\kappa^{-1} \cdot l)(\overline{v^n} - \overline{v}) dx$$
$$= -\int_{\mathbb{R}^3} \rho^2(\overline{v^n} - \overline{v^m}) \cdot F_1^n R^n (Id - P)(\kappa^{-1} \cdot l)(\overline{v^n} - \overline{v}) dx.$$

Then we decompose

$$\int_0^T |h_{2,m,n}(t)| dt \leqslant \int_0^T \left| \int_{\mathbb{R}^3} (\overline{v^n} - \overline{v^m}) \cdot F_1^n \left[\rho^2 R^n (Id - P)(\kappa^{-1} \cdot l)(\overline{v^n} - \overline{v}) \right. \right.$$
$$\left. \left. - \rho R^n (Id - P)\rho(\kappa^{-1} \cdot l)(\overline{v^n} - \overline{v}) \right] dx \right| dt \tag{3.18}$$
$$+ \int_0^T \left| \int_{\mathbb{R}^3} \rho(\overline{v^n} - \overline{v^m}) \cdot F_1^n R^n (Id - P)\rho(\kappa^{-1} \cdot l)(\overline{v^n} - \overline{v}) dx \right| dt.$$

The second integral in the right hand side of (3.18) is estimated thanks to Hölder's inequality:

$$\int_0^T \left| \int_{\mathbb{R}^3} \rho(\overline{v^n} - \overline{v^m}) \cdot F_1^n R^n (Id - P) \rho(\kappa^{-1} \cdot l)(\overline{v^n} - \overline{v}) dx \right| dt$$
$$\leqslant C \|\rho(\overline{v^n} - \overline{v^m})\|_{L_{t,x}^2} \|\rho(\overline{v^n} - \overline{v})\|_{L_{t,x}^2},$$

where C depends only on T, F, l, $\|v_{\text{init}}\|_{L^\infty}$. To deal with the first integral in the right-hand side of (3.18), we use the following commutation lemma.

Lemma 3.6 (Jochmann [23, Lemma 3.5]). *Let ρ belong to $L^2((0,T), L^2(\mathbb{R}^3)) \cap L^\infty((0,T), L^\infty(\mathbb{R}^3))$, and let $(M^n)_{n \in \mathbb{N}}$ be a bounded sequence in $W^{1,\infty}((0,T), L^2(\mathbb{R}^3)) \cap L^\infty((0,T), L^\infty(\mathbb{R}^3))$ which converges to 0 in $L^\infty((0,T), L^2(\mathbb{R}^3))$ weak $*$. Then*

$$\int_0^T \|\rho(t)^2 R^n (Id-P)M^n - \rho(t) R^n (Id-P)\rho(t)M^n\|_{L^1(\mathbb{R}^3)+L^2(\mathbb{R}^3)} dt \xrightarrow[n \to \infty]{} 0, \quad (3.19)$$

and $\int_0^T \|\rho(t)^2 (Id-P)M^n - \rho(t)(Id-P)\rho(t)M^n\|_{L^1(\mathbb{R}^3)+L^2(\mathbb{R}^3)} dt \xrightarrow[n \to \infty]{} 0. \quad (3.20)$

We apply Lemma 3.6, (3.19) with $M^n = (\kappa^{-1} \cdot l)(\overline{v^n} - \overline{v})$: the first integral in the right-hand side of (3.18) is estimated by

$$C(T, F, \|v_{\text{init}}\|_{L^\infty}) \int_0^T \|\overline{v^n} - \overline{v^m}\|_{L^\infty \cap L^2} \|\rho^2 R^n (Id-P)M^n - \rho R^n (Id-P)\rho M^n\|_{L^1 + L^2} dt,$$

and thus goes to zero as n goes to infinity, uniformly with respect to m. Hence, we get (3.17). □

For all $\delta > 0$, we also bound $h_{3,m,n}$ by δ for all $n \geqslant N_\delta$ and all $m \in \mathbb{N}$ thanks to (3.2). Finally, summing up with (3.14) and (3.17), we have from (3.13):

$$\forall \delta > 0, \exists N_\delta \in \mathbb{N}, \forall n, m \geqslant N_\delta, \forall t \in [0, T],$$
$$\|\rho(v^n - v^m)\|_{L^2(\Omega)}^2 (t) \qquad\qquad\qquad (3.21)$$
$$\leqslant C \left(\delta + \|\rho(v^n - v^m)\|_{L^2((0,T) \times \Omega)}^2 + \|\rho(v^n - v)\|_{L^2((0,T) \times \Omega)}^2 \right).$$

Use Gronwall's Lemma, then let m go to ∞, and use Gronwall's Lemma again to deduce:

$$\forall \delta > 0, \exists N_\delta \in \mathbb{N}, \forall n \geqslant N_\delta, \forall t \in [0, T], \qquad \|\rho(v^n - v)\|_{L^2(\Omega)}^2 (t) \leqslant C\delta,$$

which implies that v^n converges towards v strongly in $L^2((0, T) \times \Omega, \rho(t,x)^2 dt dx)$. Up to a subsequence, convergence then holds almost everywhere, for the measure $\rho^2 dt dx$, or $dt dx$, since ρ is positive almost everywhere in $(0, T) \times \Omega$. Thanks to the

pointwise estimates (a) from Lemma 3.1, dominated convergence thus ensures that v^n converges towards v strongly in $L^2((0, T) \times \Omega, dt dx)$. Then, equicontinuity of $\{v\} \cup \{v^n\}_{n \in \mathbb{N}}$ in $C([0, T], L^2(\Omega))$ implies (by Ascoli's Theorem) the strong convergence of v^n in $C([0, T], L^2(\Omega))$. This, together with the uniform bounds on $\{v\} \cup \{v^n\}_{n \in \mathbb{N}}$ and with the weak convergence of u^n, is enough to pass to the limit in (3.6), (3.7) to get (2.7), (2.8).

Strong convergence of u^n. Since u^n and u satisfy (3.6) and (2.7) respectively, their difference is solution to a hyperbolic equation with source term in $L^1((0, T), L^2(\mathbb{R}^3))$,

$$(\partial_t + B)(u^n - u) = (\kappa^{-1} \cdot l)(F^n - F(x, \overline{v}, u)).$$

The standard energy estimate then gives

$$
\begin{aligned}
\|u^n - u\|_{L^2}(t) &\leqslant C \int_0^t \|F^n - F(x, \overline{v}, u)\|_{L^2}(t') dt' \\
&\leqslant C \int_0^t \left(\|F_0^n - F_0(x, \overline{v})\|_{L^2}(t') + \|F_1^n R^n u^n - F_1(x, \overline{v})u\|_{L^2}(t') \right) dt'.
\end{aligned}
\tag{3.22}
$$

Thanks to the growth conditions (2.4) of F and to the pointwise bound Lemma 3.1, (a) of the v^n, there holds, for all $n, m \in \mathbb{N}$ and $(t, x) \in [0, T] \times \mathbb{R}^3$,

$$|F_i^n(t, x) - F_i(x, \overline{v}(t, x))| \leqslant C_F(e^{KT} \|v_{\text{init}}\|_{L^\infty})|\overline{v^n}(t, x) - \overline{v}(t, x)|. \tag{3.23}$$

In particular this yields that for any $t' \in [0, T]$, $\|F_0^n - F_0(x, \overline{v})\|_{L^2}(t')$ goes to zero as n goes to infinity.

Furthermore,

$$|F_1^n R^n u^n - F_1(x, \overline{v})u| \leqslant |F_1^n||R^n(u^n - u)| + |F_1^n||(R^n - Id)u| + |F_1^n - F_1(x, \overline{v})||u|.$$

Thanks to the L^∞ bounds on v^n (*cf.* Lemma 3.1, (a)) and on F_1^n (*cf.* (3.9)), and to the property (3.2) of the operator R^n, the first term in the right hand side above is bounded by $C_F(e^{KT} \|v_{\text{init}}\|_{L^\infty})|u^n - u|$. In the same way, the second term goes to zero in L^2 as n goes to infinity. Finally, up to a subsequence, the third term tends to zero almost everywhere, and is bounded by $C(F, T, \|v_{\text{init}}\|_{L^\infty})|u|$. By dominated convergence, it thus goes to zero in L^2. Finally, we get from (3.22):

$$\|u^n - u\|_{L^2}(t) \leqslant C(F, T, \|v_{\text{init}}\|_{L^\infty}) \int_0^t \|u^n - u\|_{L^2}(t') dt' + o(1),$$

and Gronwall's Lemma shows that u^n converges to u in $C([0, T], L^2)$. $\qquad \square$

Thanks to a diagonal extraction process (using times $T \in \mathbb{N}^\star$), Proposition 3.2 produces $U \in C([0, \infty), L^2)$, solution to (2.6)-(2.8) with U_{init} as initial data. Estimates (i) and (ii) are then straightforward. To prove Theorem 2.6, there remains to show its last statement: the stability property. To this end, consider a

sequence $(U_{\text{init}}^n)_{n\in\mathbb{N}}$, bounded in L_{div}, and converging to U_{init} in L^2. It gener-
ates a (sub)sequence of solutions $(U^n)_{n\in\mathbb{N}}$, with, from the bounds (i), (ii) in Theo-
rem 2.6, u^n converging to u in $L^\infty_{\text{loc}}((0,\infty), L^2(\mathbb{R}^3))$ weak $*$ and v^n converging to
v in $W^{1,\infty}_{\text{loc}}((0,\infty), L^2(\Omega))$ weakly $*$ and in $L^\infty_{\text{loc}}((0,\infty), L^\infty(\Omega))$ weak $*$. Then,
define the weight ρ from (3.11) and estimate $v^n - v^m$ as in (3.12), with $R^n u^n$ and
$R^m u^m$ replaced with u^n and u^m, respectively. This leads to the analogue to (3.13),
with no $h_{3,m,n}$ and $h_{3,n,m}$ terms, and no R^n in $h_{1,m,n}$ and $h_{2,m,n}$. Apply Lemma 3.4
and Lemma 3.6, (3.20) (instead of (3.19)), to get strong $C([0,T], L^2)$ convergence
of v^n towards v (in (3.21), the term $\|v_{\text{init}}^n - v_{\text{init}}^m\|_{L^2}$ goes to zero, and contributes
to δ). Strong convergence of the fields u^n is then obtained as above, with an initial
term $\|u_{\text{init}}^n - u_{\text{init}}\|_{L^2}$ going to zero added to the right hand side of (3.22).

The same process proves Proposition 2.7.

4. Propagation of smoothness and uniqueness: proof of Theorem 2.8

It is worth noting that, under the smoothness assumption on ε and μ in (2.11),
$u \in L^2(\mathbb{R}^3, \mathbb{R}^6)$ with $Pu \in H^\mu(\mathbb{R}^3, \mathbb{R}^6)$ iff $u \in L^2(\mathbb{R}^3, \mathbb{R}^6)$ with $\operatorname{curl} u_i \in$
$H^{\mu-1}(\mathbb{R}^3, \mathbb{R}^6)$ for $i = 1, 2$.

We thus split the proof of Theorem 2.8 in several steps. In Section 4.1, we
isolate a Cauchy problem for the projection Pu of u. This allows some dispersive
estimates that we etablish in Section 4.3, while in Section 4.2, Littlewood-Paley
decompositions are introduced. We consider first the case where μ is in $(0,1)$,
then we prove the part (a) of Theorem 2.8 in the case $\mu = 1$, which concerns the
propagation of smoothness, and finally the part (b), which concerns uniqueness.

Remark 4.1. Let us mention that in the proof of the propagation of H^1 regularity
given in [18], the step "$\mu \in (0,1)$" is missing, and the resulting estimates (collected
here in Lemma 4.18) are claimed without proof.

4.1. Preliminaries

Lemma 4.2. *For any solution $U := (u,v)$ to (2.6)-(2.8) with $U_{\text{init}} := (u_{\text{init}}, v_{\text{init}}) \in$
L_{div} as initial data given by Theorem 2.6, the part $\mathbf{u} := Pu$ solves for $x \in \mathbb{R}^3$,*

$$(\partial_t + B)\mathbf{u} = P(A\mathbf{u}) + Pg, \tag{4.1}$$

$$\mathbf{u}|_{t=0} = Pu_{\text{init}}, \tag{4.2}$$

where

$$A(t,x) := (\kappa^{-1} \cdot l) F_1(x, \overline{v}), \tag{4.3}$$

$$g(t,x) := (\kappa^{-1} \cdot l) F(x, \overline{v}, (Id - P)(\kappa^{-1} \cdot l)\overline{v}), \tag{4.4}$$

$$= (\kappa^{-1} \cdot l) F_0(x, \overline{v}) + (\kappa^{-1} \cdot l) F_1(x, \overline{v})(Id - P)(\kappa^{-1} \cdot l)\overline{v}. \tag{4.5}$$

Proof. First, apply the projector P to the system (2.7), observing that P commutes with both ∂_t and B. Then, split F according to (2.3), split u into $u = \mathbf{u} + (Id - P)u$ and finally use the constraint (2.6). $\quad\square$

The projectors P_i, $i = 1, 2$, defined on $L^2(\mathbb{R}^3, \mathbb{R}^3)$, extend to $L^p(\mathbb{R}^3, \mathbb{R}^3)$ (this result extends the classical one by Calderón and Zygmund [11] on singular integrals, in the spirit of the extension by Judovič [46]):

Lemma 4.3 (Starynkevitch [43, Lemma 3.13]). *Under assumption* (2.11)*, the projectors* P_i, $i = 1, 2$*, extend to* $L^p(\mathbb{R}^3, \mathbb{R}^3)$ *and for all* $p_0 > 1$*, there exists* $C > 0$ *such that for all* $p \in [p_0, \infty)$*, their norm from* $L^p(\mathbb{R}^3, \mathbb{R}^3)$ *into itself are less than* Cp*.*

We deduce estimates for the right-hand side of (4.1):

Lemma 4.4. *As in Theorem* 2.8*, assume* (2.3)-(2.5) *and* (2.11)*. Let* $U_{\text{init}} := (u_{\text{init}}, v_{\text{init}}) \in L_{\text{div}}$*, and let* $U := (u, v)$ *be any solution to* (2.6)-(2.8) *with* $U_{\text{init}} := (u_{\text{init}}, v_{\text{init}}) \in L_{\text{div}}$ *as initial data given by Theorem* 2.6*. The following holds true for A and g given by* (4.3)-(4.5):

$$A \in L^\infty_{\text{loc}}((0, \infty), L^\infty(\mathbb{R}^3)), \tag{4.6}$$

$$A \in C([0, \infty), L^2(\mathbb{R}^3)), \tag{4.7}$$

$$\partial_t A \in L^\infty_{\text{loc}}((0, \infty), L^2(\mathbb{R}^3)), \tag{4.8}$$

$$g \in \cap_{1 \leqslant p < \infty} C([0, \infty), L^p(\mathbb{R}^3)), \tag{4.9}$$

$$\partial_t g \in \cap_{1 \leqslant q \leqslant 2} L^\infty_{\text{loc}}((0, \infty), L^p(\mathbb{R}^3)). \tag{4.10}$$

Proof. For all $t, t' \geqslant 0$, there holds

$$\|A(t)\|_{L^\infty(\mathbb{R}^3)} \leqslant \|\kappa^{-1} \cdot l\|_{L^\infty(\mathbb{R}^3)} C_F(\|v_{\text{init}}\|_{L^\infty} e^{Kt}), \tag{4.11}$$

$$\|A(t) - A(t')\|_{L^2(\mathbb{R}^3)} \leqslant \|\kappa^{-1} \cdot l\|_{L^\infty(\mathbb{R}^3)} C_F(\|v_{\text{init}}\|_{L^\infty} e^{Kt}) \|v(t) - v(t')\|_{L^2(\Omega)}, \tag{4.12}$$

$$\|\partial_t A(t)\|_{L^2(\mathbb{R}^3)} \leqslant \|\kappa^{-1} \cdot l\|_{L^\infty(\mathbb{R}^3)} C_F(\|v_{\text{init}}\|_{L^\infty} e^{Kt}) \|\partial_t v(t)\|_{L^2(\Omega)}, \tag{4.13}$$

what yields estimates (4.6)-(4.8).

Since $v \in C([0, \infty), L^2(\Omega)) \cap L^\infty_{\text{loc}}((0, \infty), L^\infty(\Omega))$, we have $v \in \cap_{p \geqslant 1} C([0, \infty), L^p(\Omega))$ – using the boundedness of Ω for $p < 2$, and by interpolation for $p > 2$. Lemma 4.3 then yields (4.9). Next, using that $|F_i(x, v)| \leqslant C_F(\|v_{\text{init}}\|_{L^\infty} e^{Kt})|v|$ for $i = 0, 1$, we infer from (2.8) that $\partial_t v \in \cap_{1 \leqslant q \leqslant 2} L^\infty_{\text{loc}}((0, \infty), L^p(\mathbb{R}^3))$. Since

$$\partial_t g(t, x) = (\kappa^{-1} \cdot l)\{\partial_v F_0(x, \overline{v}) \cdot \partial_t \overline{v} + F_1(x, \overline{v})(Id - P)(\kappa^{-1} \cdot l)$$
$$\cdot \partial_t \overline{v} + (\partial_v F_1(x, \overline{v}) \cdot \partial_t \overline{v}) \cdot (Id - P)(\kappa^{-1} \cdot l)\overline{v}\},$$

thanks to Lemma 4.3, we finally get (4.10). $\quad\square$

Also, for a given v, the "fields part" u is in fact uniquely determined:

Lemma 4.5. *Let $A \in L_{\mathrm{loc}}^{\infty}((0, \infty), L^{\infty}(\mathbb{R}^3))$ and $g \in L_{\mathrm{loc}}^{1}((0, \infty), L^2(\mathbb{R}^3))$. For any $u_{\mathrm{init}} \in L^2(\mathbb{R}^3)$, there exists only one solution $\mathbf{u} \in C([0, \infty), L^2(\mathbb{R}^3))$ to (4.1)-(4.2) with $P u_{\mathrm{init}}$ as initial data. Furthermore, it satisfies $\mathbf{u} = P\mathbf{u}$.*

Proof. Existence is given by Lemma 4.2. To prove uniqueness, consider two solutions $\mathbf{u}_1$ and $\mathbf{u}_2$ in $C([0, \infty), L^2(\mathbb{R}^3))$ to (4.1)-(4.2), and $T > 0$. Then

$$\forall t \geqslant 0, \quad (\mathbf{u}_1 - \mathbf{u}_2)(t) = \int_0^t e^{i(t-s)B} P A (\mathbf{u}_1 - \mathbf{u}_2)(s) ds,$$

so that, using (4.6), for $t \in [0, T]$,

$$\|(\mathbf{u}_1 - \mathbf{u}_2)(t)\|_{L^2(\mathbb{R}^3)} \leqslant C(\kappa, T) \int_0^t \|(\mathbf{u}_1 - \mathbf{u}_2)(s)\|_{L^2(\mathbb{R}^3)} ds.$$

Hence, by Gronwall's Lemma, $\mathbf{u}_1 = \mathbf{u}_2$ on $[0, T]$, for any $T > 0$. Thus, there is only one solution $\mathbf{u}$. Finally, in the same way, $\mathbf{u} - P\mathbf{u}$ simply satisfies:

$$\forall t \geqslant 0, \quad (\partial_t + B)(\mathbf{u} - P\mathbf{u})(t) = 0,$$

so that $\mathbf{u} = P\mathbf{u}$. $\qquad\qquad\qquad\qquad\qquad\qquad\qquad\qquad\qquad\qquad\quad\square$

4.2. Technical interlude 2: Fourier analysis

We recall the existence of a smooth dyadic partition of unity: there exist two radial bump functions χ and ϕ valued in the interval $[0, 1]$, supported respectively in the ball $B(0, 4/3) := \{|\xi| < 4/3\}$ and in the annulus $C(3/4, 8/3) := \{3/4 < |\xi| < 8/3\}$, such that

$$\forall \xi \in \mathbb{R}^3, \quad \chi(\xi) + \sum_{j \geqslant 0} \phi(2^{-j}\xi) = 1, \quad \forall \xi \in \mathbb{R}^3 \setminus \{0\}, \quad \sum_{j \in \mathbb{Z}} \phi(2^{-j}\xi) = 1,$$

$$|j - j'| \geqslant 2 \Rightarrow \mathrm{supp}\, \phi(2^{-j} \cdot) \cap \mathrm{supp}\, \phi(2^{-j'} \cdot) = \emptyset,$$

$$j \geqslant 1 \Rightarrow \mathrm{supp}\, \chi(2^{-j} \cdot) \cap \mathrm{supp}\, \phi(2^{-j} \cdot) = \emptyset.$$

The Fourier transform $\mathcal{F}$ is defined on the space of integrable functions $f \in L^1(\mathbb{R}^3)$ by $(\mathcal{F}f)(\xi) := \int_{\mathbb{R}^3} e^{-2i\pi x \cdot \xi} f(x) dx$, and extended to an automorphism of the space $\mathcal{S}'(\mathbb{R}^3)$ of tempered distributions, which is the dual of the Schwartz space $\mathcal{S}(\mathbb{R}^3)$ of rapidly decreasing functions.

The so-called dyadic blocks Δ_j correspond to the Fourier multipliers $\Delta_j := \phi(2^{-j}D)$, that is

$$\Delta_j u(x) := 2^{3j} \int_{\mathbb{R}^3} \tilde{h}(2^j y) u(x - y) dy \quad \text{for } j \geqslant 0, \quad \text{where } \tilde{h} := \mathcal{F}^{-1}\phi.$$

We also introduce $S_0 := \chi(D)$, that is

$$S_0 u(x) := \int_{\mathbb{R}^3} h(y)u(x-y)dy, \quad \text{where } h := \mathcal{F}^{-1}\chi.$$

We will use the inhomogeneous Littlewood-Paley decomposition $Id = S_{-1} + \sum_{j \in \mathbb{N}} \Delta_j$, which holds in the space of tempered distributions $\mathcal{S}'(\mathbb{R}^3)$, and the homogeneous Littlewood-Paley decomposition $Id = \sum_{j \in \mathbb{Z}} \Delta_j$, which holds in $\mathcal{S}'_h(\mathbb{R}^3)$, the space of tempered distributions u such that $\lim_{j \to -\infty} \|\sum_{k \leqslant j} \Delta_k u\|_{L^\infty(\mathbb{R}^3)} = 0$.

We now recall the definition of the inhomogeneous (respectively homogeneous) Besov spaces $B_{p,q}^\lambda$ (respectively $\dot{B}_{p,q}^\lambda$) on $\mathbb{R}^3$ which are, for $\lambda \in \mathbb{R}$ (the smoothness index), $p, q \in [1, +\infty]$ (respectively the integral-exponent and the sum-exponent), the spaces of tempered distributions u in $\mathcal{S}'(\mathbb{R}^3)$ (respectively $\mathcal{S}'_h(\mathbb{R}^3)$) such that

$$\|f\|_{B_{p,q}^\lambda(\mathbb{R}^3)} := \|S_0 f\|_{L^p(\mathbb{R}^3)} + \|(2^{j\lambda}\|\Delta_j f\|_{L^p(\mathbb{R}^3)})_j\|_{l^q(\mathbb{N})}$$

(respectively $\|f\|_{\dot{B}_{p,q}^\lambda(\mathbb{R}^3)} := \|(2^{j\lambda}\|\Delta_j f\|_{L^p(\mathbb{R}^3)})_j\|_{l^q(\mathbb{Z})}$)

is finite. These Banach spaces do not depend on the choice of the dyadic partition above (*cf.* for instance the book [3]).

4.3. Dispersion

Propagation of smoothness or singularities for solutions to hyperbolic Cauchy problems, such as

$$Lu := (\partial_t + B)u = f, \text{ with } u|_{t=0} = u_{\text{init}}, \tag{4.14}$$

obeys the laws of geometrical optics. Let us refer here to the survey [19] by Gårding for an introduction to the subject. The characteristic variety of the operator L is defined as

$$\text{Char}(L) := \{(t, x, \tau, \xi) \in \mathbb{R} \times \mathbb{R}^3 \times \mathbb{C} \times (\mathbb{R}^3 \setminus \{0\}) \mid \det L(t, x, \tau, \xi) = 0\},$$

where $L(t, x, \tau, \xi)$ denotes the (principal) symbol of the operator L, which is the 6×6 matrix

$$L(t, x, \tau, \xi) \equiv L(x, \tau, \xi) := \tau Id + B(x, \xi),$$

$$\text{where } B(x, \xi) := \begin{bmatrix} 0 & \kappa_1(x)^{-1}\xi \wedge \cdot \\ \kappa_2(x)^{-1}\xi \wedge \cdot & 0 \end{bmatrix}.$$

For all $(x, \xi) \in \mathbb{R}^3 \times (\mathbb{R}^3 \setminus \{0\})$, the matrix $B(x, \xi)$ admits three eigenvalues: $\lambda_\pm(x, \xi) := \pm(\kappa_1\kappa_2(x))^{-1/2}|\xi|$ and 0, each one with multiplicity 2. The eigenspace associated with the eigenvalue 0 is precisely $\text{ran}(Id - P)$. We introduce

$$P_\pm(x, \xi) := \frac{1}{2\pi i} \oint_{|z - \lambda_\pm(x,\xi)| = r} L(x, -z, \xi)^{-1}dz,$$

where r is chosen small enough for $\lambda_\pm(x, \xi)$ being the only eigenvalue inside the circle of integration. The matrix $P_\pm(x, \xi)$ is the spectral projection associated with the eigenvalue $\lambda_\pm(x, \xi)$, that is the projection onto the kernel of $L\big(x, -\lambda_\pm(x, \xi), \xi\big)$ along its range. These spectral projections are homogeneous of degree 0 with respect to ξ, and the associated pseudo-differential operators $P_\pm$ satisfy $P_+ + P_- = Id - P$. In addition, $Id - P$, P_+ and P_- are orthogonal projectors (in the weighted L^2 space introduced in Section 2.1) commuting with B, and acting on Besov spaces.

Now the point is that considering the Cauchy problem (4.14) for solutions u satisfying $(Id - P)u = 0$, we select the branch of the characteristic variety which are curved, what generates dispersion. We shall need the following indices $p_1, r_1, q_1, s_1, \mu, \sigma$ and ρ:

$$p_1 \in [2, \infty) \quad \text{and} \quad 1/r_1 + 1/p_1 = 1/2; \tag{4.15}$$

$$q_1 \in (1, 2] \quad \text{and} \quad 1/s_1 + 1/q_1 = 3/2; \tag{4.16}$$

$$\mu \in \mathbb{R}, \quad \sigma := \mu - 1 + 2/p_1 \quad \text{and} \quad \rho := \mu - 1 + 2/q_1. \tag{4.17}$$

Proposition 4.6. *Let $p_1, r_1, q_1, s_1, \mu, \sigma$ and ρ be given by (4.15)-(4.17). Under assumption (2.1), there is a non-decreasing function $C : (0, \infty) \to (0, \infty)$ such that, for any $T > 0$, for any initial data u_{init} in $\dot{H}^\mu$ such that $(Id - P)u_{\text{init}} = 0$, and for any source term f in $L^{s_1}((0, T), \dot{B}^\rho_{q_1,2}(\mathbb{R}^3))$ such that $(Id - P)f = 0$, any (weak) solution u to the Cauchy problem (4.14) belongs to $L^{r_1}((0, T), \dot{B}^\sigma_{p_1,2}(\mathbb{R}^3))$ and satisfies $u = Pu$, as well as*

$$\|Pu\|_{L^{r_1}\big((0,T), \dot{B}^\sigma_{p_1,2}(\mathbb{R}^3)\big)} \leqslant C(T) \left(\|Pu_{\text{init}}\|_{\dot{H}^\mu} + \|Pf\|_{L^{s_1}\big((0,T), \dot{B}^\rho_{q_1,2}(\mathbb{R}^3)\big)} \right). \tag{4.18}$$

In the case $p_1 = 2$ (hence $r_1 = \infty$, $\sigma = \mu$), the function $u = Pu$ is even in $C([0, T], \dot{H}^\mu(\mathbb{R}^3))$.

The result of Proposition 4.6 is false for $r_1 = 2$, $p_1 = \infty$, $s_1 = 1$, $q_1 = 2$ and $\mu = \rho = 1$, $\sigma = 0$. However, it is true when truncating frequencies. We use, for $\lambda > 0$, the low frequency cut-off operator S^λ, which is the Fourier multiplier with symbol $\chi_\lambda := \chi(\cdot/\lambda)$, where the cut-off function $\chi \in C^\infty_c(\mathbb{R}^d, [0, 1])$ takes value 1 when $|\xi| \leqslant 1/2$, and 0 when $|\xi| \geqslant 1$. Then, we have:

Proposition 4.7. *Under assumption (2.1), there is a non-decreasing function $C : (0, \infty) \to (0, \infty)$ such that, for all $\lambda, T > 0$ and for any $u \in C([0, T], H^1(\mathbb{R}^3))$ solution to (4.14),*

$$\begin{aligned} &\|S^\lambda Pu\|_{L^2((0,T), L^\infty(\mathbb{R}^3))} \\ &\leqslant C(T)\sqrt{\ln(1 + \lambda T)} \left(\|\partial_x Pu_{\text{init}}\|_{L^2(\mathbb{R}^3)} + \|\partial_x Pf\|_{L^1((0,T), L^2(\mathbb{R}^3))} \right). \end{aligned} \tag{4.19}$$

Estimate (4.19) is proved in the case where the operator P is $P := -\Delta$ in [24] Proposition 6.3. Even in this case Estimate (4.19) without the cut-off S^λ is false [26, 30]. Proposition 4.7 extends [24]'s result to (smooth) variable coefficients.

Let us mention that one can deduce from these results some similar estimates for the $\mathbb{R}^N$-valued solutions u of wave equations of the form:

$$(\partial_t^2 + A)u = f \text{ in } \mathbb{R}^3, \quad \partial_t^\nu u|_{t=0} = u_\nu \quad \text{for } \nu = 0, 1,$$

where $A := -a(x)\Delta + \sum_{j=1}^n B_j(x)\partial_j + C(x)$, when $a - 1 \in \mathcal{C}_K^\infty(\mathbb{R}^3)$, $a(x) \geqslant c_0 > 0$ and the B_j and C are in $\mathcal{C}_K^\infty(\mathbb{R}^3, \mathcal{M}_{N \times N}(\mathbb{R}))$. These equations stand for the propagation of waves in an inhomogeneous isotropic compact medium $K \subset \mathbb{R}^3$ surrounded by vacuum.

To prove Proposition 4.6 and Proposition 4.7 we use the Lax method, that is an explicit representation of the solution which allows to take advantage of oscillations *via* the method of stationary phase, for each dyadic block, to get a pointwise dispersive estimate. The final step relies on the $TT^\star$ argument and the summation over the dyadic blocks.

This kind of strategy is now very classical. In particular, the $TT^\star$ argument is due to Ginibre and Velo [20], and the use of bilinear interpolation was initiated by Keel and Tao [25]. We refer here to the book [3] by Bahouri, Chemin and Danchin for a larger overview of its use and of its consequences. However we did not find Proposition 4.7 in the literature so that we now detail a little bit its proof.

Proof of Proposition 4.6. Let us first remark that it is sufficient to prove Estimate (4.18) for smooth data (by the usual regularization process) and locally in time. More precisely, it suffices to prove that there is a constant $C > 0$ and $T_1 > 0$ such that, for all $\lambda > 0$ and for all $u \in \mathcal{C}([0, T_1], H^1(\mathbb{R}^3))$ solution to (4.14),

$$\|Pu\|_{L^{r_1}\left((0,T_1), \dot{B}_{p_1,2}^\sigma(\mathbb{R}^3)\right)} \leqslant C\left(\|Pu_{\text{init}}\|_{\dot{H}^\mu} + \|Pf\|_{L^{s_1}\left((0,T_1), \dot{B}_{q_1,2}^\rho(\mathbb{R}^3)\right)}\right). \quad (4.20)$$

Indeed, apply several times Estimate (4.20) on time intervals of the form $(kT_1, (k+1)T_1)$, with k an integer ranging from zero to the integral part K of T/T_1 (plus the interval (KT_1, T)): on the right hand side, $\|Pu(kT_1)\|_{\dot{H}^\mu}$ is estimated by $C\left(\|Pu((k-1)T_1)\|_{\dot{H}^\mu} + \|Pf\|_{L^{s_1}\left(((k-1)T_1,kT_1), \dot{B}_{q_1,2}^\rho(\mathbb{R}^3)\right)}\right)$. Summing up gives (4.18).

Now, consider the operator $S(t) := e^{-tB}P$. It admits a parametrix, and thus is given by a sum $S(t) = I_+(t) + I_-(t)$ of operators on the dispersive eigenspaces, which are Fourier Integral Operators: for any smooth function $u(x)$,

$$(I_\pm(t)u)(x) = \int_{\mathbb{R}^3 \times \mathbb{R}^3} e^{i\left(\Psi(t,x,\xi) - 2\pi y \cdot \xi\right)} a(t, x, \xi)u(y)d\xi dy,$$

and we drop the subscript $\pm$ in the sequel. The phase $\Psi(t, x, \xi)$ is real, positively homogeneous of degree one in ξ, $\mathcal{C}^\infty$ for $\xi \neq 0$, and satisfies the eikonal equation:

$$\partial_t \Psi(t, x, \xi) = \pm(\kappa_1 \kappa_2(x))^{-1/2} |\Psi'_x(t, x, \xi)|, \quad (4.21)$$

with

$$\Psi|_{t=0}(x,\xi) = 2\pi x \cdot \xi. \tag{4.22}$$

The amplitude a is in Hörmander's class S^0, and admits an asymptotic expansion whose successive orders satisfy a sequence of linear hyperbolic equations. Such a method was initiated by Lax in his pioneering paper [29]. Because of the caustic phenomenon, the lifespan of smooth solutions to Equation (4.21) is limited. However, since Ψ is homogeneous in ξ and the set $K \times S^2$ is compact, there exists $T_1 > 0$ such that the solution to (4.21)-(4.22) remains smooth on $[-T_1, T_1]$. Let us mention here that Ludwig [33] succeeded in extending Lax' analysis into a global-in-time result. The arguments have been refined thanks to Hörmander's theory of Fourier Integral Operators [17, 22]. But we use the parametrices (and the solution operators $I(t)$) only locally in time and we refer for their construction to the work of Chazarain [15], Nirenberg and Treves [37] and [36], Kumano-go [27] and Brenner [9]. In particular we refer to the last one for the following precious informations about the phase ([9, Lemma 2.1]): there exist $c, C > 0$ such that

(i) $c|\xi| \leqslant |\Psi'_x| \leqslant C|\xi|$ on $[-T_1, T_1] \times \mathbb{R}^3 \times (\mathbb{R}^3 \setminus \{0\})$;
(ii) $cId \leqslant \pm\Psi''_{x\xi} \leqslant CId$, $\Psi''_{x\xi}$ being real symmetric, on $[-T_1, T_1] \times \mathbb{R}^3 \times (\mathbb{R}^3 \setminus \{0\})$;
(iii) $\Psi''_{\xi\xi}$ is semi-definite with rank 2 for $|\xi| \neq 0$, $t \neq 0$; and for $|\xi| = 1$, $x \in K$, there is a constant $c_0 > 0$ such that the moduli of the non-zero eigenvalues of $\Psi''_{\xi\xi}$ are bounded from below by $c_0|t|$;
(iv) for $x \notin K$, the above results are consequence of the exact formula: $\Psi(t, x, \xi) = 2\pi x \cdot \xi \pm 2\pi(\kappa_1\kappa_2(x))^{-1/2}t|\xi|$.

Note that these results imply that the kernel of $I(t)$ is a Lagrangian distribution.
We use the $TT^\star$ method for the frequency localized operators

$$T_j(t) := \Delta_j I(t), \ j \in \mathbb{Z}.$$

The composed operator $T_j(t)T_j(t')^\star$ is then

$$T_j(t)T_j(t')^\star = \Delta_j I(t - t')\Delta_j.$$

Here is the $TT^\star$ result.

Lemma 4.8. *There exist $0 < c < C$ such that for all $j \in \mathbb{Z}$, $u \in L^2(\mathbb{R}^3)$, $p, r \in [1, \infty]$ and $f \in L^{r'}((0, T), L^{p'}(\mathbb{R}^3))$,*

$$\|T_j(t)u\|_{L^r((0,T),L^p(\mathbb{R}^3))} \leqslant \|u\|_{L^2(\mathbb{R}^3)} \sup_{g \in \mathcal{B}^j_{r,p}} (b_j(g,g))^{\frac{1}{2}}, \tag{4.23}$$

and $$\left\| \int_0^T T_j(t)T_j(t')^\star f(t')dt' \right\|_{L^r((0,T),L^p(\mathbb{R}^3))} = \sup_{g \in \mathcal{B}^j_{r,p}} |b_j(f,g)|, \tag{4.24}$$

with

$$b_j(f, g) := \int_{(0,T)\times(0,T)} \langle T_j(t)T_j(t')^\star f(t'), g(t)\rangle_{L^2(\mathbb{R}^3)} dt dt',$$

$\mathcal{B}_{r,p}^j$ *the space of functions* $g(t, x)$ *in* $\mathcal{B}_{r,p}$ *whose Fourier transform is supported in* $c2^j \leqslant |\xi| \leqslant C2^j$, *and* $\mathcal{B}_{r,p}$ *the space of smooth functions* $g(t, x)$ *satisfying* $\|g\|_{L^{r'}\left((0,T), L^{p'}(\mathbb{R}^3)\right)} \leqslant 1$.

Proof. Begin with

$$\|T_j(t)u\|_{L^r((0,T),L^p(\mathbb{R}^3))} = \sup_{g\in\mathcal{B}_{r,p}} \left| \int_{(0,T)\times\mathbb{R}^3} T_j(t)u(x) \cdot g(t, x) dt dx \right|.$$

Then, use the Plancherel identity plus the properties of the support of the Fourier transform of $T_j(t)u$ to get

$$\|T_j(t)u\|_{L^r((0,T),L^p(\mathbb{R}^3))} = \sup_{g\in\mathcal{B}_{r,p}^j} \left| \int_{(0,T)\times\mathbb{R}^3} T_j(t)u \cdot g dt dx \right|.$$

Using Fubini's principle, transposition and Cauchy-Schwarz inequality, observe that

$$\|T_j(t)u\|_{L^2((0,T),L^\infty(\mathbb{R}^3))} \leqslant \|u\|_{L^2(\mathbb{R}^3)} \sup_{g\in\mathcal{B}_{r,p}^j} \left\| \int_0^T T_j(t)^\star g dt \right\|_{L^2(\mathbb{R}^3)}. \quad (4.25)$$

Using again Fubini's principle and transposition, get

$$\left\| \int_0^T T_j(t)^\star g(t) dt \right\|_{L^2(\mathbb{R}^3)}^2 = b_j(g, g). \quad (4.26)$$

This leads to (4.23). Moreover,

$$\left\| \int_0^T T_j(t)T_j(t')^\star f(t') dt' \right\|_{L^r((0,T),L^p(\mathbb{R}^3))}$$
$$= \sup_{g\in\mathcal{B}_{r,p}^j} \left| \int_{(0,T)\times\mathbb{R}^3} \left(\int_0^T T_j(t)T_j(t')^\star f(t') dt' \right) \cdot g(t) dt dx \right|$$
$$= \sup_{g\in\mathcal{B}_{r,p}^j} |b_j(f, g)|. \qquad \square$$

Now, we need to estimate $b_j(f, g)$. We start with a pointwise estimate.

Lemma 4.9. *There exists $C > 0$ such that for all $j \in \mathbb{Z}$, $(t, t') \in (0, T) \times (0, T)$ and $u \in L^1(\mathbb{R}^3)$,*

$$\|T_j(t)(T_j)(t')^\star u\|_{L^\infty(\mathbb{R}^3)} \leqslant C 2^{3j}(1 + 2^j|t - t'|)^{-1}\|u\|_{L^1(\mathbb{R}^3)}.$$

Proof. Since Δ_j is a bounded operator on $L^p(\mathbb{R}^3)$ for all $p \in [1, \infty]$, it is sufficient to prove that for all $t \in (-T, T)$ and $u \in L^1(\mathbb{R}^3)$,

$$\|I(t)\Delta_j u\|_{L^\infty(\mathbb{R}^3)} \leqslant C 2^{3j}(1 + 2^j|t|)^{-1}\|u\|_{L^1(\mathbb{R}^3)}. \tag{4.27}$$

Writing down $I(t)\Delta_j u$, we get, for all $x \in \mathbb{R}^3$:

$$
\begin{aligned}
(I(t)\Delta_j u)(x) &= \int_{\mathbb{R}^3 \times \mathbb{R}^3} e^{i(\Psi(t,x,\xi) - 2\pi y \cdot \xi)} a(t, x, \xi)\varphi(2^{-j}D)u(y)d\xi dy \\
&= \int_{\mathbb{R}^3} e^{i\Psi(t,x,\xi)} a(t, x, \xi)\varphi(2^{-j}\xi)\hat{u}(\xi)d\xi \\
&= 2^{3j} \int_{\mathbb{R}^3} e^{i 2^j \Psi(t,x,\eta)} a(t, x, 2^j\eta)\varphi(\eta)\hat{u}(2^j\eta)d\eta,
\end{aligned}
$$

so that

$$|(I(t)\Delta_j u)(x)| \leqslant 2^{3j} \sup_{y \in \mathbb{R}^3} \left| \int_{\mathbb{R}^3} e^{i 2^j(\Psi(t,x,\eta) - y \cdot \eta)} a(t, x, 2^j\eta)\varphi(\eta)d\eta \right| \|u\|_{L^1(\mathbb{R}^3)}.$$

To get (4.27), simply apply the following lemma of stationary phase.

Lemma 4.10 (Littman [31]). *Let $\Psi(\xi)$ be a real function $\mathcal{C}^\infty$ such that the rank of its Hessian matrix $\Psi''_{\xi\xi}$ is at least ρ and let $v(\xi)$ be a function supported in a ring. Then there exists $M \in \mathbb{N}$ and $C > 0$ (which depends only on a finite number of derivatives of Ψ, of a lower bound of the maximum of the absolute values of the minors of order ρ of $\Psi''_{\xi\xi}$, on supp v) such that, for all $\Lambda \in \mathbb{R}$,*

$$\|\mathcal{F}^{-1}(e^{i\Lambda\Psi}v)\|_{L^\infty(\mathbb{R}^3))} \leqslant C(1 + |\Lambda|)^{-\frac{\rho}{2}} \cdot \sum_{|\alpha| \leqslant M} \|D^\alpha v\|_{L^1(\mathbb{R}^3)}.$$

To use this lemma, distinguish between short times, for which the eikonal equation (4.21) implies that the phase Ψ admits the expansion

$$\Psi(t, x, \xi) = 2\pi\left(x \cdot \xi \pm t(\kappa_1\kappa_2(x))^{-1/2}|\xi|\right) + O(t^2),$$

(thus $\rho = 2$, $\Lambda = 2^j(\kappa_1\kappa_2(x))^{-1/2}t$ is suitable), and subsequent times, for which Estimate (4.3) on the phase gives $\rho = 2$ (and $\Lambda = 2^j$). This yields

$$\sup_{y \in \mathbb{R}^3} \left| \int_{\mathbb{R}^3} e^{i 2^j(\Psi(t,x,\eta) - y \cdot \eta)} a(t, x, 2^j\eta)\varphi(\eta)d\eta \right| \leqslant C(1 + 2^j t)^{-1}. \tag{4.28}$$

$\square$

Since $T_j(t)T_j(t')^\star$ is also bounded on $L^2(\mathbb{R}^3)$, from the above result and the Riesz-Thorin Interpolation Theorem, we infer the following.

Lemma 4.11. *There exists $C > 0$ such that for all $j \in \mathbb{Z}$, $(t, t') \in (0, T) \times (0, T)$, $p \in [2, \infty]$ and $u \in L^{p'}(\mathbb{R}^3)$,*

$$\|T_j(t)T_j(t')^\star u\|_{L^p(\mathbb{R}^3)} \leqslant C 2^{3j(1-2/p)}(1+2^j|t-t'|)^{-(1-2/p)}\|u\|_{L^{p'}(\mathbb{R}^3)}. \quad (4.29)$$

Lemma 4.12. *For any p_1, $p_2 \in [2, \infty)$, define r_1 and r_2 by $1/r_1 + 1/p_1 = 1/2$ and $1/r_2 + 1/p_2 = 1/2$. Then, there exists $C > 0$ such that for all $j \in \mathbb{Z}$, $f \in L^{r'_1}((0, T), L^{p'_1}(\mathbb{R}^3))$ and $g \in L^{r'_2}((0, T), L^{p'_2}(\mathbb{R}^3))$,*

$$|b_j(f, g)| \leqslant C 2^{2j(1/r_1+1/r_2)} \|f\|_{L^{r'_1}\left((0,T),L^{p'_1}(\mathbb{R}^3)\right)} \|g\|_{L^{r'_2}\left((0,T),L^{p'_2}(\mathbb{R}^3)\right)}. \quad (4.30)$$

Proof. Using (4.29), we apply Hölder's inequality to get

$$|b_j(f, g)| \leqslant C 2^{6j/r_1} \left\| \left((1 + 2^j|\cdot|)^{-2/r_1} \star \|f\|_{L^{p'_1}(\mathbb{R}^3)} \right)(t') \, \|g(t')\|_{L^{p'_1}(\mathbb{R}^3)} \right\|_{L^1(0,T)}.$$

Using again Hölder's inequality, we get

$$|b_j(f, g)|$$
$$\leqslant C 2^{6j/r_1} \left\| \left((1 + 2^j|\cdot|)^{-2/r_1} \star \|f\|_{L^{p'_1}(\mathbb{R}^3)} \right)(t') \right\|_{L^{r_1}(0,T)} \|g\|_{L^{r'_1}\left((0,T),L^{p'_1}(\mathbb{R}^3)\right)}.$$

Thus, the Hardy-Littlewood-Sobolev inequality implies

$$|b_j(f, g)| \leqslant C 2^{4j/r_1} \|f\|_{L^{r'_1}\left((0,T),L^{p'_1}(\mathbb{R}^3)\right)} \|g\|_{L^{r'_1}\left((0,T),L^{p'_1}(\mathbb{R}^3)\right)}. \quad (4.31)$$

Moreover, (4.26) yields

$$\left\| \int_0^T T_j(t)^\star f(t) dt \right\|_{L^2(\mathbb{R}^3)} = b_j(f, f)^{1/2} \leqslant C 2^{2j/r_1} \|f\|_{L^{r'_1}\left((0,T),L^{p'_1}(\mathbb{R}^3)\right)}. \quad (4.32)$$

Now, we observe that

$$b_j(f, g) = \int_0^T \left\langle \int_0^T T_j(t')^\star f(t') dt', T_j(t)^\star g(t) \right\rangle_{L^2(\mathbb{R}^3)} dt.$$

We use the Cauchy-Schwarz and the uniform boundedness of the $T_j(t)^\star$, for t in $(0, T)$, in $L^2(\mathbb{R}^3)$ to get

$$|b_j(f, g)| \leqslant b_j(f, f)^{1/2} \|g\|_{L^1((0,T),L^2(\mathbb{R}^3))}. \quad (4.33)$$

Using now Equation (4.32) yields

$$|b_j(f, g)| \leqslant C 2^{2j/r_1} \|f\|_{L^{r'_1}\left((0,T),L^{p'_1}(\mathbb{R}^3)\right)} \|g\|_{L^1((0,T),L^2(\mathbb{R}^3))}. \quad (4.34)$$

Interpolating between (4.31) and (4.34), we find (4.30). $\qquad \square$

To complete the proof of Proposition 4.6, we need to sum over all frequencies. Apply Duhamel's principle to get

$$Pu(t) = e^{-tB}Pu_{\text{init}} + \int_0^t e^{(t'-t)B}Pf(t')dt'. \qquad (4.35)$$

The first term in $\Delta_j Pu(t)$ is thus

$$\Delta_j I(t)u_{\text{init}} = \sum_{k\in\mathbb{Z}} \Delta_j I(t)\Delta_k u_{\text{init}}.$$

An important fact in order to estimate $\Delta_j I(t)u_{\text{init}}$ is that the quasi-orthogonality of the dyadic blocks is not destroyed by the parametrix $I(t)$. Actually, according to [9, Proposition 1.1 and Lemma 1.4], there exists $L \in \mathbb{N}^*$ such that for all $j \in \mathbb{Z}$,

$$\Delta_j I(t) = \sum_{|j-k|\leqslant L} \Delta_j I(t)\Delta_k + R_j(t),$$

where $R(t) := \sum_{j\in\mathbb{Z}} R_j(t)$ satisfies

$$\|R(t)v\|_{L^{r_1}((0,T),\dot{B}^{\sigma}_{p_1,2}(\mathbb{R}^3))} \leqslant C\|v\|_{\dot{H}^{\mu}}. \qquad (4.36)$$

This shows that there is $C > 0$ such that

$$\|\Delta_j I(t)u_{\text{init}}\|_{L^{p_1}(\mathbb{R}^3)} \leqslant C\|\Delta_j I(t)\Delta_j u_{\text{init}}\|_{L^{p_1}(\mathbb{R}^3)} + \|R_j(t)u_{\text{init}}\|_{L^{p_1}(\mathbb{R}^3)},$$

so that

$$\|I(t)u_{\text{init}}\|_{L^{r_1}((0,T),\dot{B}^{\sigma}_{p_1,2}(\mathbb{R}^3))}$$

$$\leqslant C \left\| \|(2^{j\sigma}\|\Delta_j I(t)\Delta_j u_{\text{init}}\|_{L^{p_1}(\mathbb{R}^3)})_j\|_{l^2(\mathbb{Z})} \right\|_{L^{r_1}(0,T)}$$

$$+ \left\| \|(2^{j\sigma}\|R_j(t)u_{\text{init}}\|_{L^{p_1}(\mathbb{R}^3)})_j\|_{l^2(\mathbb{Z})} \right\|_{L^{r_1}(0,T)}$$

$$\leqslant C\|(2^{j\sigma}\|\Delta_j I(t)\Delta_j u_{\text{init}}\|_{L^{r_1}((0,T),L^{p_1}(\mathbb{R}^3))})_j\|_{l^2(\mathbb{Z})}$$

$$+ \|(2^{j\sigma}\|R_j(t)u_{\text{init}}\|_{L^{r_1}((0,T),L^{p_1}(\mathbb{R}^3))})_j\|_{l^2(\mathbb{Z})} \qquad \text{by Minkowski's inequality}$$

$$\leqslant C\|(2^{j(\sigma+2/r_1)}\|\Delta_j u_{\text{init}}\|_{L^2(\mathbb{R}^3)})_j\|_{l^2(\mathbb{Z})}$$

$$+ \|(2^{j\sigma}\|R_j(t)u_{\text{init}}\|_{L^{r_1}((0,T),L^{p_1}(\mathbb{R}^3))})_j\|_{l^2(\mathbb{Z})} \qquad \text{by (4.23) and (4.30)}$$

$$\leqslant C\|u_{\text{init}}\|_{\dot{H}^{\mu}} \text{ for some new constant } C, \text{ using } \mu = \sigma + 2/r_1 \text{ and (4.36).}$$

The estimate for $\int_0^t e^{(t'-t)B}Pf(t')dt'$ follows the same lines, using (4.24) instead of (4.23).

To conclude the proof of Proposition 4.6, consider the case $p_1 = 2$. As (4.18) holds, it is easy to see that u is continuous with respect to time: using a smooth approximation $f_k \in L^{s_1}((0,T), \dot{H}^{\mu})$ of f, we get $Lu_k = f_k$, and since $s_1 \geqslant 1$, the usual energy estimates show that $(u_k)_k$ is a Cauchy sequence in $C([0,T], \dot{H}^{\mu})$. By (4.18), as f_k tends to f in $L^{s_1}((0,T), \dot{B}^{\rho}_{q_1,2})$, u_k tends to u in $C([0,T], \dot{H}^{\mu})$. $\square$

Proof of Proposition 4.7. Now, in order to prove Proposition 4.7, it suffices to come back to the proof of Lemma 4.12, taking $p_1 = \infty$ (and $r_1 = 2$). Using the standard Young inequality instead of the Hardy-Littlewood-Sobolev inequality leads to

$$|b_j(f, g)| \leqslant C 2^{2j} \ln(1 + 2^j T) \|f\|_{L^2((0,T),L^1(\mathbb{R}^3))} \|g\|_{L^2((0,T),L^1(\mathbb{R}^3))}. \quad (4.37)$$

In particular for any g in $L^2((0, T), L^1(\mathbb{R}^3))$, this yields

$$|b_j(g, g)|^{1/2} \leqslant C 2^j \sqrt{\ln(1 + 2^j T)} \|g\|_{L^2((0,T),L^1(\mathbb{R}^3))}.$$

Then we apply Lemma 4.8 with $(p, r) = (2, \infty)$ and use (4.33), where we commute f and g, to estimate the right hand side of (4.24). We find that for any $u \in L^2(\mathbb{R}^3)$, for any $f \in L^1((0, T), L^2(\mathbb{R}^3))$,

$$\|T_j(t)u\|_{L^2((0,T),L^\infty(\mathbb{R}^3))} \leqslant C 2^j \sqrt{\ln(1+2^j T)} \|u\|_{L^2(\mathbb{R}^3)},$$

$$\left\| \int_0^T T_j(t) T_j(t')^\star f(t') dt' \right\|_{L^2((0,T),L^\infty(\mathbb{R}^3))} \leqslant C 2^j \sqrt{\ln(1+2^j T)} \|f\|_{L^1((0,T),L^2(\mathbb{R}^3))}.$$

Then we proceed as in the proof of Proposition 4.6 to sum over the dyadic blocks whose frequencies are below the cut-off parameter λ, with an extra factor $\sqrt{\ln(1 + \lambda T)}$ (and (4.36) holds for $r_1 = 2$, $p_1 = \infty$). $\square$

4.4. Propagation of smoothness: proof of Theorem 2.8, case where $\mu \in (0, 1)$

In the case where $\mu \in (0, 1)$, Theorem 2.8 is a consequence of the following proposition:

Proposition 4.13. *Suppose that A and g satisfy the estimates* (4.6)-(4.10), *that* $\mu \in (0, 1)$ *and* $u_{\text{init}} \in L^2(\mathbb{R}^3, \mathbb{R}^6)$ *with* $P u_{\text{init}} \in H^\mu(\mathbb{R}^3, \mathbb{R}^6)$. *Then the unique solution* $\mathbf{u} \in C([0, \infty), L^2(\mathbb{R}^3))$ *of* (4.1)-(4.2) *given by Lemma* 4.5 *belongs to* $C([0, \infty), H^\mu(\mathbb{R}^3))$.

In order to prove Proposition 4.13 we introduce, for $T > 0$, the space

$$Y^\mu(T) := C([0,T], H^\mu(\mathbb{R}^3)) \cap C^1([0,T], H^{\mu-1}(\mathbb{R}^3)) \cap L^r((0,T), B^0_{p,2}(\mathbb{R}^3)), \quad (4.38)$$

where $p := 2/(1 - \mu)$ and $r := 2/\mu$. These indices belong to $(2, \infty)$. We shall also use the space $Z^\mu(T) := Z^\mu_1(T) + Z^\mu_2(T)$, where

$$Z^\mu_1(T) := L^1((0, T), H^\mu(\mathbb{R}^3)) \cap C([0, T], L^2(\mathbb{R}^3)),$$

$$Z^\mu_2(T) := \{ f \in C([0, T], L^2(\mathbb{R}^3)) \cap L^r((0, T), B^{-1}_{p,2}(\mathbb{R}^3)) \mid$$

$$\partial_t f \in L^1((0, T), H^{\mu-1}(\mathbb{R}^3)) + L^s((0, T), B^0_{q,2}(\mathbb{R}^3)) \},$$

where $q = 2/(2 - \mu)$ and $s = 2/(1 + \mu)$. These indices belong to $(1, 2)$.

Lemma 4.14. *Suppose that $\mu, T \in (0,1)$, and $f \in Z^\mu(T)$. Then there exists $C > 0$ such that for any $\mathbf{u}_{\text{init}}$ in $H^\mu(\mathbb{R}^3, \mathbb{R}^6)$ such that $P\mathbf{u}_{\text{init}} = \mathbf{u}_{\text{init}}$, the unique corresponding solution $\mathbf{u} \in C([0,\infty), L^2(\mathbb{R}^3))$ to $L\mathbf{u} = f$, with $\mathbf{u}|_{t=0} = \mathbf{u}_{\text{init}}$, belongs to $Y^\mu(T)$ and*

$$\|\mathbf{u}\|_{Y^\mu(T)} \leqslant C(\|\mathbf{u}_{\text{init}}\|_{H^\mu(\mathbb{R}^3)} + \|f\|_{Z^\mu(T)}). \tag{4.39}$$

Proof. By definition f splits into $f = f_1 + f_2$, with $f_i \in Z_i^\mu(T)$, and $\mathbf{u}$ into $\mathbf{u} = u_0 + u_1 + u_2$, where the u_i solve the following hyperbolic Cauchy problems:

$$\begin{aligned} Lu_0 &= S_0 f && \text{with} && u_0|_{t=0} = S_0\mathbf{u}_{\text{init}}, \\ Lu_1 &= (Id - S_0)f_1 && \text{with} && u_1|_{t=0} = (Id - S_0)\mathbf{u}_{\text{init}}, \\ Lu_2 &= (Id - S_0)f_2 && \text{with} && u_2|_{t=0} = 0. \end{aligned}$$

We already have by energy estimates that the u_i are in $C([0,\infty), L^2(\mathbb{R}^3))$. Then one gets the estimates of the u_i in $C^1([0,T], H^{\mu-1}(\mathbb{R}^3))$ by using the equations. To get the estimates in $Y^\mu(T)$ it therefore only remains to get the estimates in $L^r((0,T), B_{p,2}^0(\mathbb{R}^3))$.

For u_0 this just follows the energy estimate thanks to Bernstein lemma.

In order to estimate u_1 we apply Proposition 4.6 with $p_1 = p, r_1 = r, q_1 = 2$, $s_1 = 1, \sigma = 0, \rho = \mu$. This gives the estimate of u_1 in $L^r((0,T), B_{p,2}^0(\mathbb{R}^3))$ by the right-hand side of (4.39).

In order to estimate u_2 we observe that $\partial_t u_2$ satisfies $L\partial_t u_2 = (Id - S_0)\partial_t f_2$, with $\partial_t u_2|_{t=0} = (Id - S_0)f_2|_{t=0}$. By assumption there exists $g_a \in L^1((0,T), H^{\mu-1}(\mathbb{R}^3))$ and $g_b \in L^s((0,T), B_{q,2}^0(\mathbb{R}^3))$ such that $\partial_t f_2 = g_a + g_b$. We split accordingly $\partial_t u_2$ into $\partial_t u_2 = u_a + u_b$, where u_a solves $Lu_a = g_a$, with $u_a|_{t=0} = (Id - S_0)f_2|_{t=0}$, and u_b solves $Lu_b = g_b$, with $u_b|_{t=0} = 0$. In order to estimate u_a (respectively u_b) we apply Proposition 4.6 with $p_1 = p, r_1 = r, q_1 = 2$, $s_1 = 1, \sigma = -1, \rho = \mu - 1$ (respectively with $p_1 = p, r_1 = r, q_1 = q$, $s_1 = s, \sigma = -1, \rho = 0$ and $\mu - 1$ instead of μ). This yields the estimate of $\partial_t u_2$ in $L^r((0,T), B_{p,2}^{-1}(\mathbb{R}^3))$ by $\|f_2\|_{Z_2^\mu(T)}$. As a consequence $Bu_2 = -\partial_t u_2 + (Id - S_0)f_2 \in L^r((0,T), B_{p,2}^{-1}(\mathbb{R}^3))$. Since $Pu_2 = u_2$ this entails that $u_2 \in L^r((0,T), B_{p,2}^0(\mathbb{R}^3))$. $\qquad\square$

The proof of Proposition 4.13 follows by induction of the following lemma:

Lemma 4.15. *Suppose that A and g satisfy the estimates (4.6)-(4.10), and $\mu \in (0,1)$. Then there exists $T_1 > 0$ and $C > 0$ such that for any $u_{\text{init}} \in L^2(\mathbb{R}^3, \mathbb{R}^6)$ with $Pu_{\text{init}} \in H^\mu(\mathbb{R}^3, \mathbb{R}^6)$, the unique corresponding solution $\mathbf{u} \in C([0,\infty), L^2(\mathbb{R}^3))$ to (4.1)-(4.2) given by Lemma 4.5 belongs to $Y^\mu(T_1)$ and*

$$\|\mathbf{u}\|_{Y^\mu(T_1)} \leqslant C(\|Pu_{\text{init}}\|_{H^\mu(\mathbb{R}^3)} + \|g\|_{Z^\mu(T_1)}). \tag{4.40}$$

Proof. In order to prove Lemma 4.15 we recall the following estimate:

Lemma 4.16 (Joly-Métivier-Rauch [24, Lemma 5.3]). *There is a constant C, which depends only on $\|A\|_{L^\infty((0,T),L^\infty(\mathbb{R}^3))}$, $\|\partial_t A\|_{L^\infty((0,T),L^2(\mathbb{R}^3))}$ and μ, such that for all $T \in (0, 1)$ and $u \in Y^\mu(T)$, Au belongs to $Z^\mu(T)$, and*

$$\|Au\|_{Z^\mu(T)} \leqslant CT^{\mu/2}\|u\|_{Y^\mu(T)} + C\|u\|_{L^\infty((0,T),L^2(\mathbb{R}^3))}. \tag{4.41}$$

Let us warn the reader that there is a small misprint in the right hand side in [24, Lemma 5.3], which is corrected above. We also have:

Lemma 4.17. *If g satisfies (4.9)-(4.10), then for any $T > 0$, g belongs to $Z^\mu(T)$.*

Proof. Since $1 < q < 2 < p < \infty$ there hold continuous embeddings $L^q \subset B^0_{q,2}$ and $L^p \subset B^0_{p,2} \subset B^{-1}_{p,2}$, so that, using (4.9)-(4.10), we get, for any $T > 0$, that g is in $Z^\mu_2(T) \subset Z^\mu(T)$. □

According to Lemma 4.16 and Lemma 4.17 there is a constant C such that for all $T \in (0, 1)$, if $\mathbf{u} \in Y^\mu(T)$ then $f := P(A\mathbf{u}) + Pg \in Z^\mu(T)$ and

$$\|f\|_{Z^\mu(T)} \leqslant CT^{\mu/2}\|\mathbf{u}\|_{Y^\mu(T)} + C\|\mathbf{u}\|_{L^\infty((0,T),L^2(\mathbb{R}^3))}. \tag{4.42}$$

Therefore applying Lemma 4.14 and choosing T_1 small enough we get Lemma 4.15. □

4.5. Propagation of smoothness: proof of Theorem 2.8 (a), case where $\mu = 1$

We now consider U_{init} in L_{div} with $\operatorname{curl} u_{\text{init},i} \in L^2(\mathbb{R}^3)$, for $i = 1, 2$, and we consider U solution to (2.6)-(2.8) with U_{init} as initial data given by Theorem 2.6. The idea is to estimate $B(Pu) = P(Au) + Pg - \partial_t Pu$ from (4.1).

Lemma 4.18. *Define A and g by (4.3)-(4.5). The following holds true*:

$$\partial_t A \in L^\infty_{\text{loc}}((0, \infty), L^3(\mathbb{R}^3)), \tag{4.43}$$

$$\partial_t g \in L^\infty_{\text{loc}}((0, \infty), L^2(\mathbb{R}^3)). \tag{4.44}$$

Proof. We apply Theorem 2.8 in the case $\mu = 1/2$. This yields that u belongs to $C([0, \infty), H^{1/2}(\mathbb{R}^3))$ and thus also to $L^\infty_{\text{loc}}((0, \infty), L^3(\mathbb{R}^3))$. We then infer that $\partial_t v \in L^\infty_{\text{loc}}((0, \infty), L^3(\mathbb{R}^3))$ and then we get the estimates (4.43) and (4.43). □

Lemma 4.19. *Define A and g by (4.3)-(4.5). If $\mathbf{u} \in L^\infty_{\text{loc}}((0, \infty), H^1(\mathbb{R}^3)) \cap W^{1,\infty}_{\text{loc}}((0, \infty), L^2(\mathbb{R}^3))$ then $f := P(A\mathbf{u} + g)$ satisfies $\partial_t f \in L^\infty_{\text{loc}}((0, \infty), L^2(\mathbb{R}^3))$, and for any $T > 0$, there exists $C > 0$, which only depends on T, A and g, such that*

$$\|\partial_t f\|_{L^\infty((0,T),L^2(\mathbb{R}^3))} \leqslant C(\|\mathbf{u}\|_{L^\infty((0,T),H^1(\mathbb{R}^3))} + \|\partial_t \mathbf{u}\|_{L^\infty((0,T),L^2(\mathbb{R}^3))} + 1). \tag{4.45}$$

Proof. We have $\partial_t f = P(\partial_t A\mathbf{u} + A\partial_t\mathbf{u} + \partial_t g)$, so that using Hölder's inequality, the continuous embedding $L^6(\mathbb{R}^3) \subset H^1(\mathbb{R}^3)$ and the estimates (4.43)-(4.44), we get (4.45). □

Now we observe that $\partial_t\mathbf{u}$ solves for $x \in \mathbb{R}^3$,

$$(\partial_t + B)\partial_t\mathbf{u} = \partial_t(P(A\mathbf{u}) + Pg), \tag{4.46}$$

$$\partial_t\mathbf{u}|_{t=0} = -Bu_{\text{init}} + P(A|_{t=0}\mathbf{u}_{\text{init}}) + Pg|_{t=0}. \tag{4.47}$$

This provides an estimate of $\partial_t\mathbf{u}$ in $L^\infty_{\text{loc}}((0,\infty), L^2(\mathbb{R}^3))$, and using (4.1)-(4.2), of $B\mathbf{u}$ in $L^\infty_{\text{loc}}((0,\infty), L^2(\mathbb{R}^3))$, hence of $\mathbf{u}$ in $L^\infty_{\text{loc}}((0,\infty), H^1(\mathbb{R}^3))$.

4.6. Uniqueness: proof of Theorem 2.8 (b)

Let us recall that in this section we assume that there exists $j \in \{1,2\}$ such that $l_{3-j}F = 0$ and such that F depends only on (x, v, u_j). Let $U_{\text{init}} \in L_{\text{div}}$ with $\operatorname{curl} u_{\text{init},i} \in L^2(\mathbb{R}^3)$, for $i = 1, 2$. Let U and U' be two solutions to (2.6)-(2.8), given by Theorem 2.6, both with U_{init} as initial data. The difference $\delta U := U' - U = (\delta u, \delta v)$ between $U = (u, v)$ and $U' = (u', v')$ is solution to the following hyperbolic system

$$M(\delta U) := ((\partial_t + B)\delta u, \partial_t \delta v) = ((\kappa^{-1} \cdot l)\delta F, \delta F),$$
$$\text{where} \quad \delta F = F(x, v', u'_j) - F(x, v, u_j). \tag{4.48}$$

Thanks to (2.3) we have

$$\delta F = F_0(x, v', u'_j) - F_0(x, v, u_j) + (F_1(x, v') - F_1(x, v)) \cdot u_j$$
$$+ F_1(x, v') \cdot (u_j - u'_j). \tag{4.49}$$

The first and last terms in (4.49) are easily estimated in $L^2(\mathbb{R}^3)$ by $C_F(\|v_0\|_{L^\infty(\mathbb{R}^3)}e^{Kt})\|\delta U\|_{L^2(\mathbb{R}^3)}$. To deal with the second one we construct a L^∞ approximation of the field u_j (analogous to the ones of [24, Lemma 6.2], and [21, Lemma 2.7]).

Lemma 4.20. *There is a non-decreasing function $C : (0,\infty) \to (0,\infty)$ and for all $T > 0$, there exists $(u^\lambda_{\parallel,j})_{\lambda \geqslant e} \subset L^\infty((0,T) \times \mathbb{R}^3)$ such that for all $\lambda \geqslant e$,*

$$\|u^\lambda_{\parallel,j}(t)\|_{L^\infty((0,T)\times\mathbb{R}^3)} \leqslant C(T)\ln\lambda,$$
$$\text{and} \quad \|((Id - P_j)u_j - u^\lambda_{\parallel,j})(t)\|_{L^2(\mathbb{R}^3)} \leqslant C(T)/\lambda, \tag{4.50}$$
$$\text{for all} \quad t \in [0, T],$$

$$\|S^\lambda P_j u_j\|_{L^2((0,T),L^\infty(\mathbb{R}^3))} \leqslant C(T)\sqrt{\ln\lambda},$$
$$\text{and} \quad \|((Id - S^\lambda)P_j u_j(t)\|_{L^2(\mathbb{R}^3)} \leqslant C(T)/\lambda, \tag{4.51}$$
$$\text{for all} \quad t \in [0, T].$$

Let us admit for a while Lemma 4.20 in order to finish the proof of Theorem 2.8. Fix $T > 0$ and consider $t \in (0, T)$. In the second term on the right-hand side of (4.49), decompose $u_j = u_{\parallel, j}^{\lambda} + \left((Id - P_j)u_j - u_{\parallel, j}^{\lambda} \right) + S^{\lambda} P_j u_j + (Id - S^{\lambda}) P_j u_j$. We infer from Lemma 4.20 that

$$\|\delta F\|_{L^2(\mathbb{R}^3)} \leqslant C_F \left(\|v_{\text{init}}\|_{L^{\infty}(\mathbb{R}^3)} e^{KT} \right) \tag{4.52}$$

$$\cdot \left(\|\delta U\|_{L^2(\mathbb{R}^3)} + \left(C(T) \ln \lambda + \|S^{\lambda} P_j u_j\|_{L^{\infty}(\mathbb{R}^3)} \right) \|\delta v\|_{L^2(\mathbb{R}^3)} + 2 \frac{C(T)}{\lambda} \|\delta v\|_{L^{\infty}(\mathbb{R}^3)} \right).$$

The energy estimate, together with Gronwall's Lemma, gives

$$\|\delta U(t)\|_{L^2(\mathbb{R}^3)} \leqslant 2 C_F C(T) \frac{t}{\lambda} \exp \left(C_F \int_0^t (1 + C(T) \ln \lambda + \|S^{\lambda} P_j u_j(t')\|_{L^{\infty}(\mathbb{R}^3)}) dt' \right).$$

Since $\int_0^t (1 + C(T) \ln \lambda + \|S^{\lambda} P_j u_j(t')\|_{L^{\infty}(\mathbb{R}^3)}) dt' \leqslant \overline{C}(T) \ln \lambda$ with $\overline{C}(T) \xrightarrow[T \to 0]{} 0$, we choose T_0 small enough (in order to have $\overline{C}(T_0) < 1$), and let λ go to infinity. This shows that $\delta U(t)$ vanishes on $[0, T_0]$. Repeat this procedure on intervals of size T_0 to get the desired conclusion.

Proof of Lemma 4.20. We define $u_{\parallel, j}^{\lambda}$ by setting $u_{\parallel, j}^{\lambda}(t, x) := (Id - P_j)u_j(t, x)$ if $|(Id - P_j)u_j(t, x)| \leqslant C \ln \lambda$, and $u_{\parallel, j}^{\lambda}(t, x) := 0$ otherwise, where the constant C is chosen below (independently of $(t, x) \in [0, T] \times \mathbb{R}^3$). Therefore, for $p \in [2, \infty)$,

$$\|((Id - P_j)u_j - u_{\parallel, j}^{\lambda})(t)\|_{L^2}^2$$
$$= \int_{|(Id - P_j)u_j| \geqslant C \ln \lambda} |(Id - P_j)u_j(t)|^2 dx \leqslant (C \ln \lambda)^{2-p} \|(Id - P_j)u_j(t)\|_{L^p}^p. \tag{4.53}$$

Now, according to Lemma 4.3, the projection $Id - P$ acts continously in any L^p with a norm less than $C_0 p$. Furthermore, we have

$$\|v_{\text{init}}\|_{L^p}^p \leqslant \|v_{\text{init}}\|_{L^{\infty}}^p |\Omega|,$$

so that $(\|v_{\text{init}}\|_{L^p})_{1 \leqslant p \leqslant \infty}$ is bounded. Thus, using Equation (2.6) and the bound from Theorem 2.6 (iii), we infer from (4.53) that

$$\|((Id - P_j)u_j - u_{\parallel, j}^{\lambda})(t)\|_{L^2}^2 \leqslant \frac{(C_0 \, p \, e^{KT} \, \|v_{\text{init}}\|_{L^p})^p}{(C \ln \lambda)^{p-2}}$$

$$= (C \ln \lambda)^2 \lambda^{2 \ln \left(2 \frac{C_0}{C} e^{KT} \sup_{1 \leqslant q \leqslant \infty} \|v_{\text{init}}\|_{L^q} \right)},$$

choosing $p = 2 \ln \lambda$. With C big enough, we obtain (4.50) ($C(T) = 2 C_0 e^{KT+1} \sup_{1 \leqslant q \leqslant \infty} \|v_{\text{init}}\|_{L^q}$ is suitable).

We are now concerned with the first inequality in (4.51). Coming back to (4.1)-(4.2), the idea is to use the Strichartz estimate (4.19). However, since we are not able to bound $\partial_x f$ in $L^\infty_{\text{loc}}((0, \infty), L^2(\mathbb{R}^3))$, we cannot apply (4.19) directly. To overcome this difficulty we introduce some potential vectors. Since $Pu_{\text{init}} \in H^1(\mathbb{R}^3)$ (respectively since $f = (f_1, f_2) \in W^{1,\infty}_{\text{loc}}((0, \infty), L^2(\mathbb{R}^3))$ (*cf.* (4.42) and (4.45)) and $Pf = f$), for $i = 1, 2$, there exists $\phi_{\text{init}} := (\phi_{\text{init},1}, \phi_{\text{init},2}) \in H^2(\mathbb{R}^3)$ (respectively $\psi := (\psi_1, \psi_2) \in W^{1,\infty}_{\text{loc}}((0, \infty), H^1(\mathbb{R}^3)))$ such that

$$\text{div}(\kappa_{3-i}\phi_{\text{init},i}) = 0, \quad \text{curl}(\phi_{\text{init},i}) = \kappa_i u_{\text{init},i}, \tag{4.54}$$

$$\text{(respectively } \text{div}(\kappa_{3-i}\psi_i) = 0, \quad \text{curl}(\psi_i) = \kappa_i f_i). \tag{4.55}$$

We consider the operator $\check{B}$ defined by

$$\check{B}(\phi_1, \phi_2) = (\kappa_2^{-1} \operatorname{curl} \phi_2, -\kappa_1^{-1} \operatorname{curl} \phi_1)$$
$$\text{for } \phi := (\phi_1, \phi_2) \in D(\check{B}) := D(B) = H_{\text{curl}} \times H_{\text{curl}}.$$

The operator $\check{B}$ is simply deduced from B by switching κ_1 and κ_2. It therefore shares the same properties and estimates. In addition it satisfies the identity:

$$\begin{bmatrix} \kappa_1^{-1} \operatorname{curl} & 0 \\ 0 & \kappa_2^{-1} \operatorname{curl} \end{bmatrix} \check{B} = B \begin{bmatrix} \kappa_1^{-1} \operatorname{curl} & 0 \\ 0 & \kappa_2^{-1} \operatorname{curl} \end{bmatrix}. \tag{4.56}$$

Let $\phi := (\phi_1, \phi_2)$ be the solution (for $x \in \mathbb{R}^3$) of

$$(\partial_t + \check{B})\phi = \psi, \quad \text{with } \phi|_{t=0} = \phi_{\text{init}}. \tag{4.57}$$

Using the identity (4.56) and Lemma 4.5 we obtain

$$\text{curl}(\phi_i) = \kappa_i u_i, \quad \text{for } i = 1, 2. \tag{4.58}$$

Now observe that $\partial_t \phi$ verifies $(\partial_t + \check{B})\partial_t \phi = \partial_t \psi$, with $\partial_t \phi|_{t=0} = \psi|_{t=0} - \check{B}\phi_{\text{init}}$. Applying the Strichartz estimate (4.19) we obtain

$$\begin{aligned} &\|S^\lambda \partial_t \phi\|_{L^2((0,T), L^\infty(\mathbb{R}^3))} \\ &\leqslant C(T)\sqrt{\ln(1 + \lambda T)}\Big(\|\partial_x \partial_t \phi|_{t=0}\|_{L^2(\mathbb{R}^3)} + \|\partial_t \partial_x \psi\|_{L^1((0,T), L^2(\mathbb{R}^3))}\Big). \end{aligned} \tag{4.59}$$

From the definitions of ϕ and ψ and from the estimates (4.42) and (4.45) of the previous sections we get :

$$\|\partial_x \partial_t \phi|_{t=0}\|_{L^2(\mathbb{R}^3)} + \|\partial_t \partial_x \psi\|_{L^1((0,T), L^2(\mathbb{R}^3))} \leqslant C(T)\|Pu_{\text{init}}\|_{H^1(\mathbb{R}^3)}. \tag{4.60}$$

Using now (4.57), observing that f_{3-j} and therefore ψ_{3-j} vanish because of Assumption 2.12, and using (4.58), we obtain the first inequality in (4.51). The second one follows by applying Bernstein lemma. $\square$

5. Generic uniqueness: proof of Theorem 2.9

We apply the following general result of generic uniqueness for evolution equations by Saint-Raymond.

Theorem 5.1 (Saint-Raymond [40, Theorem 1]). *Let $\mathcal{E}_{\text{init}}$ be a topological space and $\mathcal{E}$ a metric space. Let (S) be an evolution equation admitting a solution in $\mathcal{E}$ for any initial data in $\mathcal{E}_{\text{init}}$. Consider the following hypotheses.*

(H1) *For any initial data $U_{\text{init}} \in \mathcal{E}_{\text{init}}$, for any $(U_{\text{init}}^{\varepsilon})_{\varepsilon}$ tending to U_{init} in $\mathcal{E}_{\text{init}}$, for any $(U^{\varepsilon})_{\varepsilon}$ in $\mathcal{E}$ respective solutions to (S) with $U_{\text{init}}^{\varepsilon}$ as initial data,*

   (i) *there exists a limit point of $(U^{\varepsilon})_{\varepsilon}$ in $\mathcal{E}$;*
   (ii) *any limit point of $(U^{\varepsilon})_{\varepsilon}$ in $\mathcal{E}$ is solution to (S) with U_{init} as initial data.*

(H2) *There exists $\mathcal{D}$, dense subset of $\mathcal{E}$, such that for any U_{init} in $\mathcal{D}$, there exists only one solution to (S) in $\mathcal{E}$ with U_{init} as initial data.*

Under these two hypotheses, there exists a G_{δ} dense $\widetilde{\mathcal{E}_{\text{init}}}$ of $\mathcal{E}_{\text{init}}$ such that for any $U_{\text{init}} \in \widetilde{\mathcal{E}_{\text{init}}}$, there exists only one solution to (S) in $\mathcal{E}$ with U_{init} as initial data.

Recall that we denote by τ_s and τ_w respectively the strong and weak topologies of $L^2(\mathbb{R}^3, \mathbb{R}^6)$, and by $\tilde{\tau}_s$ the strong topology of $L^2(\Omega, \mathbb{R}^d)$. We consider the product topology τ_{ss} (respectively τ_{ws}) on $\mathbf{L}^2$ obtained from τ_s (respectively τ_w) and $\tilde{\tau}_s$.

For any $C_{\text{init}} > 0$, consider

$$\mathcal{E}_{\text{init}} := \{U_{\text{init}} \in L_{\text{div}} \mid \|v_{\text{init}}\|_{L^{\infty}(\Omega)} \leqslant C_{\text{init}}\},$$

endowed with the topology τ_{ss} (respectively τ_{ws}) inherited from $\mathbf{L}^2$, and

$$\mathcal{E} := \{U \in C([0, \infty), \mathbf{L}^2), \text{ satisfying } (2.9)$$
$$\text{and the estimates (i), (ii), (iii) of Theorem 2.6}\},$$

endowed with the strong topology (respectively the weak $*$ topology relative to τ_{ws}) of $C([0, \infty), \mathbf{L}^2)$. Hypothesis (H1) is a direct consequence of the stability property stated in Theorem 2.6 (respectively Proposition 2.7). Now, set

$$\mathcal{D} := \{U_{\text{init}} \in L_{\text{div}} \text{ with } \operatorname{curl} u_{\text{init},i} \in L^2(\mathbb{R}^3), \text{ for } i = 1, 2\},$$

which is dense in $\mathcal{E}_{\text{init}}$ for the topology τ_{ss} inherited from $\mathbf{L}^2$. Moreover Theorem 2.8 yields that Hypothesis (H2) is satisfied. We can therefore apply Theorem 5.1, what proves Theorem 2.9.

6. Quasi-stationary limits: proof of Theorem 2.10, Proposition 2.11 and Theorem 2.12

Proof of Theorem 2.10. We first observe that the bounds (i), (ii) given by Theorem 2.6 for U^η are uniform in $\eta \in (0, 1)$. Therefore, up to a subsequence, U^η converges to $U := (u, v)$ in $W^{1,\infty}_{loc}((0, \infty), L^2(\Omega))$ weak $*$ and v^η converges to v in $L^\infty_{loc}((0, \infty), L^\infty(\Omega))$ weak $*$. In addition, there holds $Bu^\eta = \partial_t D^\eta$, with $D^\eta := -\eta(u^\eta - (\kappa^{-1} \cdot l)\overline{v^\eta})$. Passing to the limit already yields that U satisfies the (linear) equations (2.6) and $Bu = 0$. Using Proposition 2.7, we also get that U satisfies (2.8), which means that v solves (2.15). $\square$

Proof of Proposition 2.11. The proof is very similar to the uniqueness proof in Theorem 2.8 (b): it relies on some L^∞ approximation of $(Id - P)(\kappa^{-1} \cdot l)v$. Consider $v_1, v_2 \in C([0, \infty), L^2(\Omega)) \cap L^\infty_{loc}((0, \infty), L^\infty(\Omega))$, solutions to (2.15) with the same initial data v_{init}, and define $\delta v := v_1 - v_2$. Fix $T > 0$. From the properties (2.3)-(2.5) of F, we get the pointwise estimate

$$\partial_t(|\delta v|^2) \leqslant C_F\Big((|(Id-P)(\kappa^{-1} \cdot l)v_1|+1)|\delta v|^2 + |(Id-P)(\kappa^{-1} \cdot l)\delta v|\,|\delta v|\Big) \quad (6.1)$$
on $\quad [0, T] \times \Omega$,

for some constant $C_F = C_F\left(\|v_{init}e^{KT}\|_{L^\infty(\Omega)}\right)$. Now, defining for $M > 0$

$$w^M_\| := \mathbf{1}_{|(Id-P)(\kappa^{-1}\cdot l)v_1|\leqslant M}(Id - P)(\kappa^{-1} \cdot l)v_1,$$

we get from Lemma 4.20 that there is $C(T) > 0$ such that

$$\forall M \geqslant 1, \quad \|w^M_\|\|_{L^\infty((0,T)\times\Omega)} \leqslant M,$$
$$\|(Id - P)(\kappa^{-1} \cdot l)v_1 - w^M_\|\|_{L^\infty((0,T),L^2(\Omega))} \leqslant C(T)e^{-M/C(T)}.$$

Integrating (6.1) over Ω, using the Cauchy-Schwarz inequality and increasing the constant C (which is still independent of M), we obtain

$$\partial_t\left(\|\delta v\|^2_{L^2(\Omega)}\right) \leqslant C_F\left((M+1)\|\delta v\|^2_{L^2(\Omega)} + C(T)e^{-M/C(T)}\right).$$

Then, Gronwall's lemma yields

$$\forall t \in [0, T], \quad \|\delta v(t)\|^2_{L^2(\Omega)} \leqslant \frac{C(T)}{M+1}e^{C_F(M+1)T-M/C(T)}.$$

Now, choose T so small that $C_F MT - M/C(T) < 0$ (which is possible, since C is a non-decreasing function of T), and let M go to infinity. This shows that δv vanishes on $[0, T]$. Repeating the argument on successive time intervals yields $v_1 = v_2$. $\square$

Proof of Theorem 2.12. For each $\eta \in (0, 1)$, consider a solution U^η (given by Theorem 2.6) to (2.6), (2.8) and (2.13). Convergence of v^η (and $(Id - P)u^\eta$) is obtained as in the proof of Theorem 2.10 above. Now, drop the index η for simplicity. Then, symmetrizing the system by the change of dependent variables

$$\widetilde{u}_i = \kappa_i^{1/2} u_i,$$

we get in the distributional sense:

$$\text{for } i = 1, 2, \quad \eta \partial_t \widetilde{u}_i + (-1)^{3-i} R_{3-i} \widetilde{u}_{3-i} = \eta \kappa_i^{-1/2} \cdot l_i \partial_t \bar{v},$$

and therefore, applying ∂_t and combining,

$$\text{for } i = 1, 2, \quad \eta^2 \partial_t^2 \widetilde{u}_i + R_i^\star R_i \widetilde{u}_i = (-1)^i \eta R_{3-i} \kappa_{3-i}^{-1/2} l_{3-i} \partial_t \bar{v} + \eta^2 \kappa_i^{-1/2} l_i \partial_t^2 \bar{v}, \quad (6.2)$$

where we have set

$$i = 1, 2, \quad R_i := \kappa_{3-i}^{-1/2} \operatorname{curl} \kappa_i^{-1/2} (= R_{3-i}^\star, \text{ for the duality in } L^2(\mathbb{R}^3, dx)).$$

System (6.2) shall be understood as a system of wave equations for the "divergence free" parts $\pi_i \widetilde{u}_i$, when

$$\text{for } i = 1, 2, \quad \pi_i := \kappa_i^{1/2} P_i \kappa_i^{-1/2}.$$

Then, $\pi = (\pi_1, \pi_2)$ is an orthogonal projector in the space $L^2(\mathbb{R}^3, dx) \times L^2(\mathbb{R}^3, dx)$. Furthermore, from the description of ran P_i and ran $(1 - P_i)$ in (2.2), we deduce that

$$\text{for } i = 1, 2, \quad R_i \pi_i = R_i \quad (\text{and} \quad \pi_i R_i^\star = R_i^\star \quad \text{by transposition}).$$

Thus, we have finally:

$$\text{for } i = 1, 2, \quad \eta^2 \partial_t^2 \pi_i \widetilde{u}_i - Q_i \pi_i \widetilde{u}_i = (-1)^i \eta R_{3-i} \kappa_{3-i}^{-1/2} l_{3-i} \partial_t \bar{v} + \eta^2 \pi_i \kappa_i^{-1/2} l_i \partial_t^2 \bar{v}, \quad (6.3)$$

with

$$\text{for } i = 1, 2, \quad Q_i := -R_i^\star R_i + \kappa_i^{1/2} \nabla \left(\kappa_i^{-2} \kappa_{3-i}^{-1} \operatorname{div}(\kappa_i^{1/2} \cdot) \right).$$

From [43, Lemma 3.10], we know that for $i = 1, 2$, the differential second-order operator $(-Q_i)$ is a self-adjoint, positive, and elliptic. Thus, with $v \in L^\infty_{\text{loc}}((0, \infty),$ $L^\infty(\Omega)) \cap W^{1,\infty}_{\text{loc}}((0, \infty), L^2(\Omega))$ given, and for given initial data,

$$\pi_i \widetilde{u}_{i \,|_{t=0}} = \pi_i \kappa_i^{1/2} u_{\text{init}, i} \text{ and } \eta (\partial_t \pi_i \widetilde{u}_i)_{|_{t=0}}$$
$$= \eta \pi_i \kappa_i^{-1/2} l_i F(\overline{v_{\text{init}}}, u_{\text{init}}) + (-1)^i R_{3-i} \kappa_{3-i}^{1/2} u_{\text{init}, 3-i},$$

the solution $(\pi_1 \widetilde{u}_1, \pi_2 \widetilde{u}_2)$ to the linear wave equation system (6.3) is uniquely determined. We recover it *via* vector potentials: defining

$$\text{for } i = 1, 2, \ \eta^2 \partial_t^2 \phi_i - Q_i \phi_i = \eta \pi_i \kappa_i^{-1/2} l_i \partial_t \bar{v}, \ \phi_{i|_{t=0}} = 0, \ \eta \partial_t \phi_{i|_{t=0}} = \pi_i \widetilde{u}_{i|_{t=0}}, \quad (6.4)$$

we have

$$\text{for } i = 1, 2, \quad \pi_i \widetilde{u}_i = \eta \partial_t \phi_i + (-1)^i R_{3-i} \phi_{3-i}. \tag{6.5}$$

The problem (6.4) also determines uniquely the vector potentials ϕ_i. Since $\pi_i \phi_i$ also satisfies the problem (6.4) we infer that $\pi_i \phi_i = \phi_i$. Furthermore, [43, Lemma 3.10] ensures that for $i = 1, 2$, Q_i does not admit 0 as a resonance. One then needs to assume the following:

$$\text{for } i = 1, 2, \ Q_i \text{ is non-trapping.} \tag{6.6}$$

This is enough to apply

Theorem 6.1 (Starynkevitch [43, Theorem 3.2]). *Let Q be a non-trapping, $(L^2$-)self-adjoint, negative, and elliptic differential second-order operator, for which 0 is not a resonance. Let $s > 1/2$, $\gamma \in (-3/2, 1/2)$ and $R > 0$. Then, there exists $C \geqslant 0$ such that: for all $(u_0, u_1) \in \dot{H}_Q^{\gamma+1}(\mathbb{R}^3) \times \dot{H}_Q^{\gamma}(\mathbb{R}^3)$, and f such that $\langle x \rangle^s (-Q)^{\gamma/2} f \in L^2((0, \infty) \times \mathbb{R}^3)$, the solution u to*

$$\partial_t^2 u - Qu = f \text{ on } (0, \infty) \times \mathbb{R}^3, \quad \text{with } u_{|_{t=0}} = u_0, \quad \partial_t u_{|_{t=0}} = u_1,$$

satisfies

$$\|(u, \partial_t u)\|_{L^2((0,\infty), \dot{H}_Q^{\gamma+1}(B_R) \times \dot{H}_Q^{\gamma}(B_R))}$$
$$\leqslant C \left(\|u_0\|_{\dot{H}_Q^{\gamma+1}(\mathbb{R}^3)} + \|u_1\|_{\dot{H}_Q^{\gamma}(\mathbb{R}^3)} + \|\langle x \rangle^s (-Q)^{\gamma/2} f\|_{L^2((0,\infty) \times \mathbb{R}^3)} \right).$$

For all $\mu \in \mathbb{R}$, the space $\dot{H}_Q^{\mu}(\mathbb{R}^3)$ is defined by the norm

$$\|v\|_{\dot{H}_Q^{\mu}(\mathbb{R}^3)} = \|(-Q)^{\mu/2} v\|_{L^2(\mathbb{R}^3)}.$$

We apply the result above to $\phi_i(\eta t, x)$, whith $\gamma = 0$ and $s = 1$. This leads to:

$$\forall R > 0, \exists C_R > 0,$$
$$\|(\phi_i, \eta \partial_t \phi_i)\|_{L^2((0,\infty), \dot{H}^1(B_R) \times L^2(B_R))} \leqslant C_R \big(\eta^{1/2} \|\pi_i \kappa_i^{1/2} u_{\text{init},i}\|_{L^2(\mathbb{R}^3)}$$
$$+ \eta \|\langle x \rangle \pi_i \kappa_i^{-1/2} l_i \partial_t \bar{v}\|_{L^2((0,\infty), L^2(\mathbb{R}^3))} \big).$$

The right-hand side is controlled thanks to

Lemma 6.2 (Starynkevitch [43, Lemma 3.11]). *For all $R > 0$, there exists $C_R > 0$ such that, if $m \in L^2(\mathbb{R}^3)$ and $\mathrm{supp}(m) \subset B_R$, then*

$$\text{for } i = 1, 2, \quad |\pi_i m(x)| \leqslant C_R \langle x \rangle^{-3} \|m\|_{L^2(\mathbb{R}^3)} \quad \text{for a.e. } x \in \mathbb{R}^3.$$

Since $\pi_i \phi_i = \phi_i$, by the usual $TT^\star$ argument, $\|R_i \phi_i\|_{L^2(B_R)} \leqslant \|\phi_i\|_{\dot{H}^1(B_R)}$, and we deduce from (6.5):

$$\text{for } i = 1, 2, \quad \pi_i \widetilde{u}_i = O(\sqrt{\eta}) \quad \text{in } L^2((0, \infty), L^2_{\text{loc}}(\mathbb{R}^3)),$$

which yields the convergence of Pu^η to zero in $L^2((0, \infty), L^2_{\text{loc}}(\mathbb{R}^3))$. □

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