ANALYSIS OF THE ENERGY DECAY OF A DEGENERATED THERMOELASTICITY SYSTEM

AMEL ATALLAH-BARAKET AND CLOTILDE FERMANIAN KAMMERER

Abstract: In this paper, we study a system of thermoelasticity with a degenerated second order operator in the Heat equation. We analyze the evolution of the energy density of a family of solutions. We consider two cases: when the set of points where the ellipticity of the Heat operator fails is included in a hypersurface and when it is an open set. In the first case and under special assumptions, we prove that the evolution of the energy density is the one of a damped wave equation: propagation along the rays of geometric optic and damping according to a microlocal process. In the second case, we show that the energy density propagates along rays which are distortions of the rays of geometric optic.

1. Introduction

We consider $\Omega$ an open subset of $\mathbb{R}^d$ and the following system of thermoelasticity

\begin{equation}
\begin{aligned}
\partial_t^2 u - \Delta u + \nabla \cdot (\gamma(x) \theta) &= 0, \quad (t,x) \in \mathbb{R}^+ \times \Omega, \\
\partial_t \theta - \nabla \cdot (B(x) \nabla \theta) + \gamma(x) \nabla \partial_t u &= 0, \quad (t,x) \in \mathbb{R}^+ \times \Omega, \\
u|_{t=0} = u_0 \quad \text{and} \quad \partial_t u|_{t=0} = u_1 & \quad \text{in } \Omega, \\
\theta|_{t=0} = \theta_0 \quad \text{in } \Omega, \\
u|_{\partial \Omega} = 0, \quad \partial_t u|_{\partial \Omega} = 0 \quad \text{and} \quad \theta|_{\partial \Omega} = 0.
\end{aligned}
\end{equation}

where $u$ and $\theta$ are scalar real-valued functions. The matrix-valued function $x \mapsto B(x)$ and the vector-valued function $x \mapsto \gamma(x)$ are supposed to be defined on $\mathbb{R}^d$ (and thus on $\Omega$) and to depend smoothly on the variable $x \in \mathbb{R}^d$. The matrix $B(x)$ is assumed to be symmetric, non-negative: there exists $C_2 > 0$ such that

\begin{equation}
\forall (x, \xi) \in \mathbb{R}^d, \quad 0 \leq B(x) \xi \cdot \xi \leq C_2 |\xi|^2.
\end{equation}

Note that we may have $\det B(x) = 0$.

System (1.1) has an energy

\begin{equation}
E(u, \theta, t) := \int_\Omega |\partial_t u(t,x)|^2 dx + \int_\Omega |\nabla u(t,x)|^2 dx + \int_\Omega \theta^2(t,x) dx
\end{equation}

which decreases in time according to

\begin{equation}
E(u, \theta, t) - E(u, \theta, 0) = -2 \int_0^t \int_\Omega B(x) \nabla \theta(s,x) \cdot \nabla \theta(s,x) ds dx \leq 0.
\end{equation}

Equation (1.3) gives a priori estimates for initial data $u_0 \in H^1_0(\Omega), u_1 \in L^2(\Omega)$ and $\theta_0 \in L^2(\Omega)$ and yields by classical arguments the existence of a unique solution $(u, \theta) \in C^0(\mathbb{R}^+, H^1_0(\Omega)) \cap C^1(\mathbb{R}^+, L^2(\Omega)) \times C^0(\mathbb{R}^+, L^2(\Omega))$. We are interested in characterizing the way the energy decay: our aim is to describe the weak limits of the energy densities

\begin{equation}
e_n(t,x) = |\partial_t u_n(t,x)|^2 + |\nabla x u_n(t,x)|^2 + (\theta_n(t,x))^2, \quad n \in \mathbb{N},
\end{equation}
associated with families of solutions \((u_n, \theta_n)_{n \in \mathbb{N}}\) of \((1.1)\) corresponding to families of initial data \((u_{0,n},n \in \mathbb{N})\) in one hand and \((u_{1,n},n \in \mathbb{N})\) in the other hand, uniformly bounded in \(H^1(\Omega)\) and \(L^2(\Omega)\) respectively. Without loss of generality, we suppose that \((u_{1,n}^2)_{n \in \mathbb{N}}\) and \((\theta_{1,n})_{n \in \mathbb{N}}\) goes to zero weakly in \(L^2(\Omega)\) and that \((u_{0,n})_{n \in \mathbb{N}}\) goes to zero weakly in \(H^1(\Omega)\).

Main results on the subject are devoted to the situation where the matrix \(B\) is positive. It is known since the work of Dafermos [6] that the energy decays and the description of this decay has been the subject of several contributions. In particular, in [18], G. Lebeau and E. Zuazua [18] have proved that the decay of solutions of \((1.1)\) is not uniform. They crucially use a result of Henry, Perissinoto and Lopes [15] which show that the semigroup associated to \((1.1)\) is equal up to a compact operator to the semi-group of a system consisting of a damped wave equation coupled with the equation on the temperature \(\theta\):

\[
\begin{aligned}
\partial_t^2 u - \Delta u + \Gamma u &= 0, \quad (t,x) \in \mathbb{R}^+ \times \Omega \\
\partial_t \theta - \nabla \cdot (B(x) \nabla \theta) + \gamma(x) \nabla \partial_t u &= 0.
\end{aligned}
\]

(1.4)

The operator \(\Gamma\) is a damping operator given by

\[
\Gamma = G^* Q^{-1} G, \quad G = \gamma(x) \cdot \nabla \quad \text{and} \quad Q = \nabla \cdot (B(x) \nabla \cdot).
\]

The behaviour of the energy density \(\|u_n(t,x)\|^2\) associated with families of solutions of this damped wave equation (for initial data \((u_{0,n},n \in \mathbb{N})\) and \((u_{1,n},n \in \mathbb{N})\) uniformly bounded in \(H^1(\Omega)\) and \(L^2(\Omega)\) respectively) can be studied in the same manner than in the paper of G. Lebeau and N. Burq (see [16], [17] and [4]). In these papers, the damping is somehow different than the one of \((1.4)\), however their method applies. One obtains that the energy propagates along the rays of geometric optic associated with the wave operator \(\partial_t^2 - \Delta\) with a damping depending simultaneously on the position and the speed of the trajectory. These trajectories are trajectories of the phase space \((x(t),\omega(t))\) and the damping is microlocal: it is given by the principal symbol of the operator \(\Gamma\), the function \((x,\omega) \mapsto \frac{(\gamma(x) \omega)^2}{B(x,\omega)}\).

The method of [4] is based on the use of microlocal defect measures in the spirit of the article [13] by P. Gérard and E. Leichtnam who are at the origin of that method. Then, similar works have been achieved for the Lamé system in [4], for the equations of magnetoelasticity in [7] (see also [8]) and for the equation of viscoelastic waves by the authors [2].

Of course, the strategy that we have just described, fails if the kernel of \(B\) does not reduce to zero and we are interested in this situation: we assume that the set

\[
\Lambda = \{(x,\omega) \in \Omega \times \mathbb{S}^{d-1}, \quad B(x,\omega) = 0\}
\]

is not empty. Then, the operator \(\nabla \cdot (B(x) \nabla \cdot)\) is no longer elliptic and something else has to be done. The method we use to treat the coupling between the temperature \(\theta\) and the amplitude \(u\) is mainly inspired by the analysis of semiclassical systems performed in [12]. Of course, we recover the result sketched in the previous paragraph when \(\Lambda = \emptyset\) and we are also able to extend it to situations where \(\Lambda \neq \emptyset\) provided a weak degeneracy assumption stated below (see Assumption 2.2). This assumption consists first in a geometric assumption: the projection of \(\Lambda\) on \(\Omega\) is included in a hypersurface \(\Sigma\) and for \((x,\omega) \in \Lambda\), the vector \(\omega\) is transverse to \(\Sigma\) at
the point \( x \). Then, Assumption 2.2 contains a compatibility relation between the vector \( \gamma(x) \) and the matrix \( B(x) \): \( \gamma(x) \in \text{Ran} B(x) \). With these assumptions, we are able to prove that the energy density is still damped along the rays of geometric optic even though they pass through \( \Sigma \). At the opposite, if \( B(x) = 0 \) in an open subset \( \Omega \) of \( \Omega \), then the damping disappears and we have transport of the energy along rays which are distortion of the rays observed before. Precise statements of our results are given in Section 2 and the organization of the paper is discussed at the end of this section.

Notations. We will say that a sequence \( (u_n)_{n \in \mathbb{N}} \) is u.b. in the functional space \( F \) if the sequence \( (u_n)_{n \in \mathbb{N}} \) is a uniformly bounded family of \( F \). We denote by \( |X|_{C^{d+2}} \) the hermitian norm of \( X = (X_1, \cdots, X_{d+2}) \in \mathbb{C}^{d+2} : 
\end{align*}

\begin{align*}
|X|^2_{C^{d+2}} = X_1^2 + \cdots + X_{d+2}^2.
\end{align*}

Similarly, we will use the notation \( (X|Y) \) for the hermitian scalar product of \( \mathbb{C}^{d+2} : 
\end{align*}

\begin{align*}
\forall X = (X_1, \cdots, X_{d+2}) \in \mathbb{C}^{d+2}, \forall Y = (Y_1, \cdots, Y_{d+2}) \in \mathbb{C}^{d+2}, 
\end{align*}

\begin{align*}
(X|Y)_{C^{d+2}} = X_1 \bar{Y}_1 + \cdots + X_{d+2} \bar{Y}_{d+2}.
\end{align*}

2. Main results

In this section, we present our results which crucially rely on the use of microlocal defect measures that we define in the first subsection: the evolution of the energy density is a corollary of the analysis of the behavior of microlocal defect measures associated with the sequence \( (u_n, \theta_n)_{n \in \mathbb{N}} \). The second subsection is devoted to the analysis of properties of the thermoelasticity operator that are important for our purpose. Then, in the third subsection, we mainly consider the situation where the determinant of \( B \) vanishes on points of \( \Omega \) which are simultaneously included in a hypersurface and in a compact subset of \( \Omega \) (so that \( B \) is non negative in a neighborhood of \( \partial \Omega \)). Finally, in the fourth subsection, we discuss what happens if \( B \) vanishes in an open subset of \( \Omega \).

2.1. Microlocal defect measures. Microlocal defect measures allow to treat quadratic quantities like the energy density by taking into account microlocal effects. They describe up to a subsequence the limit of quantities of the form \( (\mathcal{A}^\infty(\Omega \times \mathbb{R}^d), \mathcal{L}^\infty(\Omega \times \mathbb{R}^d)) \) such that \( a = a(x, \xi) \) is characterized by its symbol \( a(x, \xi) \) which is a smooth function taken, for example, in the space \( \mathcal{A}_m^m \) of symbols of order \( m \); this set contains the functions \( a = a(x, \xi) \) of \( \mathcal{L}^\infty(\Omega \times \mathbb{R}^d) \) such that \( a \) is compactly supported in \( \Omega \) as a function of \( x \) and satisfies

\begin{align*}
\forall \alpha, \beta \in \mathbb{N}^d, \exists C_{\alpha, \beta} > 0, \forall (x, \xi) \in \Omega \times \mathbb{R}^d, \left| \partial_\alpha^2 \partial_\xi^\beta a(x, \xi) \right| \leq C_{\alpha, \beta}(1 + |\xi|)^{m-|\beta|}.
\end{align*}

Then, the operator \( a(x, D) \), defined in the Weyl quantization by

\begin{align*}
\forall f \in \mathcal{L}^2(\mathbb{R}^d), \quad a(x, D)f(x) = (2\pi)^{-d} \int a \left( \frac{x + y}{2}, \xi \right) f(y) e^{i\xi (x-y)} dy d\xi
\end{align*}

maps \( H^s(\Omega) \) into \( H^{s-m}(\Omega) \) (see [1]). We will also assume that for \( a \in \mathcal{A}_m^m \), there exists a function \( a_\infty(x, \xi) \) homogeneous of degree 0 such that for

\begin{align*}
\forall \omega \in \mathbb{S}^{d-1}, \lim_{R \to \infty} R^{-m} a(x, R\omega) = a_\infty(x, \omega).
\end{align*}
The symbols of $\mathcal{A}^m$ are called interior symbols because they are compactly supported inside $\Omega$. Such symbols are of no help for studying the behaviour of $(f_n)_{n \in \mathbb{N}}$ close to $\partial \Omega$: one then use tangential symbols which are defined in the Appendix. It is easy to convince oneself that the limits of quadratic quantities $(a(x, D)f_n, f_n)$ (for $(f_n)_{n \in \mathbb{N}}$ u.b. in $L^2(\Omega)$ and $a \in \mathcal{A}^0$) only depends on the function $a_{\infty}$. Then, following [12] and [20], it is possible to prove that these limits are characterized by a positive Radon matrix-valued measure $\mu$ on $\Omega \times S^{d-1}$ such that, up to the extraction of a subsequence,

$$(a(x, D)f_n, f_n) \xrightarrow{n \to \infty} (a_{\infty}, \mu).$$

Such a measure $\mu$ is called a microlocal defect measure of the family $(f_n)_{n \in \mathbb{N}}$.

Let us come back to the sequences $(\partial_t u_n)_{n \in \mathbb{N}}$, $(\nabla_x u_n)_{n \in \mathbb{N}}$ and $(\theta_n)_{n \in \mathbb{N}}$ that we need to study simultaneously. In that purpose, we set

$$U_n = \begin{pmatrix} \partial_t u_n \\ \nabla_x u_n \\ \theta_n \end{pmatrix}$$

and the energy density $e_n$ writes

$$e_n(t, x) = |U_n(t, x)|^2_{C^{d+2}}.$$

We shall consider quadratic quantities $(a(x, D)U_n(t), U_n(t))$ where the symbol

$$a(x, \xi) = (a_{i,j}(x, D))_{i,j}$$

is a $(d+2) \times (d+2)$ matrix of interior symbols of $\mathcal{A}^m$. Then, a microlocal defect measure of $(U_n(t))_{n \in \mathbb{N}}$ will be matrix-valued and will describe up to the extraction of a subsequence, the limits of the quantities $(a(x, D)U_n(t), U_n(t))$. Of course, the $t$-dependence of these measures is an issue by itself. Therefore, since $U_n \in L^2_{loc}(\mathbb{R}, L^2(\Omega))$, we test $(a(x, D)U_n(t), U_n(t))$ against smooth compactly supported functions of the variable $t$ and consider the limits of

$$I(a, \theta) := \int \theta(t)(a(x, D)U_n(t), U_n(t))dt.$$

A microlocal defect measure $M$ of $(U_n(t))_{n \in \mathbb{N}}$ is a positive Radon measure on the set $\mathbb{R}^+ \times \Omega \times S^{d-1}$ such that, up to a subsequence, we have: $\forall \theta \in C^\infty_0(\mathbb{R})$

$$I(a, \theta) \xrightarrow{n \to +\infty} \int_{\mathbb{R} \times \Omega \times S^{d-1}} \theta(t) \text{tr}(a_{\infty}(x, \omega)M(dt, dx, d\omega)).$$

In particular, $\forall \theta \in C^\infty_0(\mathbb{R})$, $\forall \chi \in C^\infty_0(\Omega)$,

$$\int \theta(t)\chi(x)e_n(t, x)dxd\omega \xrightarrow{n \to +\infty} \int_{\mathbb{R} \times \Omega \times S^{d-1}} \theta(t)\chi(x)\text{tr}M(dt, dx, d\omega).$$

The matrix-valued measure $M(t, x, \omega)$ is positive in the sense that its diagonal components $m_{i,i}$ are positive Radon measure and its off-diagonal components $m_{i,j}$ are absolutely continuous with respect to $m_{i,i}$ and $m_{i,j}$. Using the special form of the components of the vector $U_n(t)$, one can write

$$M(t, x, \omega) = \begin{pmatrix} m_1(t, x, \omega) & m_{0,1}(t, x, \omega) & m_{0,2}(t, x, \omega) \\ m_{0,1}(t, x, \omega) & m_0(t, x, \omega) & m_{1,2}(t, x, \omega) \\ m_{0,2}(t, x, \omega) & m_{1,2}(t, x, \omega) & m_0(t, x, \omega) \end{pmatrix}$$

where
• \( m_1(t, x, \omega), \nu_0(t, x, \omega) \) are the microlocal defect measures of \((\partial_t u_n)_{n \in \mathbb{N}}\) and of \((\theta_n)_{n \in \mathbb{N}}\) respectively,
• \( m_0(t, x, \omega) \omega_j \) is the joint measure of \((\partial_x u_n)_{n \in \mathbb{N}}\) and \((\partial_x u_n)_{n \in \mathbb{N}}\),
• \( m_{0,1}(t, x, \omega) \omega_j \) is the joint measure of \((\partial_t u_n)_{n \in \mathbb{N}}\) and \((\partial_x u_n)_{n \in \mathbb{N}}\),
• \( m_{0,2}(t, x, \omega) \) is the joint measure of \((\theta_n)_{n \in \mathbb{N}}\) and \((\partial_t u_n)_{n \in \mathbb{N}}\),
• \( m_{1,2}(t, x, \omega) \omega_j \) is the joint measure of \((\theta_n)_{n \in \mathbb{N}}\) and \((\partial_x u_n)_{n \in \mathbb{N}}\).

By “joint measure” of two u.b. families of \( L^2(\Omega) \), \((f_n)_{n \in \mathbb{N}}\) and \((g_n)_{n \in \mathbb{N}}\), we mean a measure which describes the limit up to extraction of a subsequence of quantities \((a(x,D)f_n, g_n)_{n \in \mathbb{N}}\) for symbols \( a \) of order 0. In the following, we assume that \((U_n(0))_{n \in \mathbb{N}}\) has only one microlocal defect measure \( M(0, x, \omega) \).

2.2. Analysis of the thermoelasticity operator. Let us now come back to our system \([1.1]\) that we rewrite as

\[
i \partial_t U_n = P(x, D) U_n
\]

where \( U_n \) is defined in \([2.2]\) and

\[
P(x, D) = \begin{pmatrix}
0 & i \nabla & -i \nabla \cdot (\gamma) \\
i \nabla & 0 & 0 \\
-i \gamma \cdot \nabla & 0 & i \nabla \cdot (B(x) \nabla \cdot \cdot)
\end{pmatrix}.
\]

We first study the eigenspaces of the matrix \( P(x, \xi) \) which is not self adjoint if \( B \neq 0 \). However, we have the following proposition:

**Proposition 2.1.** There exists \( R_0 > 0 \) such that for \( |\xi| \geq R_0 \), the matrix \( P(x, \xi) \) diagonalize with smooth eigenvalues and smooth eigenprojectors. Besides, the kernel of \( P(x, \xi) \) is of dimension greater or equal to \( d - 1 \).

Denote by \( \lambda_- \), \( \lambda_0 \) and \( \lambda_+ \) the three eigenvalues of \( P(x, \xi) \) which are not identically 0, and by \( \Pi^- \), \( \Pi^0 \) and \( \Pi^+ \) the smooth associated eigenprojectors. Then, for all \( R > R_0 \) and for \( \chi \in C^\infty(\mathbb{R}^d) \) such that \( \chi(\xi) = 0 \) for \( |\xi| \leq 1/2 \), \( \chi(\xi) = 1 \) for \( |\xi| \geq 1 \) and \( 0 \leq \chi \leq 1 \), we have in \( D'(\Omega) \),

\[
(2.3) \quad \left| U_n(t, x) \right|_{C^{d+2}}^2 = \sum_{k \in \{0, -, +\}} \left| \Pi^k(x, D) \chi \left( \frac{D}{R} \right) U_n(t, x) \right|_{C^{d+2}}^2 + o(1).
\]

This proposition is a consequence of Propositions \([3.1]\) and \([3.2]\) below where we study the asymptotics of \( P(x, \xi) \) for large \( \xi \) (which differ whether \( \left( x, \frac{\xi}{|\xi|} \right) \in \Lambda \) or not). Each eigenprojector \( \Pi^k(x, \xi) \) \((k = 0, -, +)\) of \( P(x, \xi) \) characterizes a mode and for each mode, we will analyze the microlocal defect measure of the component \((\Pi^k(x, D) U_n)_{n \in \mathbb{N}}\).

2.3. Propagation and damping for weakly degenerated \((B, \gamma)\). Let us precise first our assumptions.

**Assumption 2.2.** We say that the pair \((B, \gamma)\) is weakly degenerated if \( B \) and \( \gamma \) satisfy the following conditions.

1. There exists a hypersurface \( \Sigma \) of \( \mathbb{R}^d \) such that \{det \( B(x) = 0 \}\} \subset \Sigma \) and for all \((x, \omega) \in \Lambda \), then \( \omega \) is transverse to \( \Sigma \) in \( x \).
2. There exists \( \tilde{\gamma} \in C^\infty(\Omega, \mathbb{R}^d) \) such that \( \gamma(x) = B(x) \tilde{\gamma}(x) \).
Example 2.3. Suppose $d = 2$, $\Omega = \mathbb{R}$, $\gamma(x) = c_2$, $B(x) = \begin{pmatrix} b(x_1) & 0 \\ 0 & 1 \end{pmatrix}$. Suppose that $b(y) = 0$ if and only if $y = 0$. Then,

$$\Lambda = \{(0, y), (\pm 1, 0)\}, y \in \mathbb{R},$$

and one can check that $(B, \gamma)$ is weakly degenerated.

The sequences

$$v_{n,0}^\pm = \frac{1}{\sqrt{2}} \left( u_{n,1} \pm \frac{D}{|D|} \cdot \nabla_x u_{n,0} \right)$$

have only one microlocal defect measure $\mu_0^\pm$ given by

$$\mu_0^\pm (x, \omega) = m_1(0, x, \omega) + m_0(0, x, \omega) \pm \Re (m_0(0, x, \omega)).$$

Theorem 2.4. Suppose that $(B, \gamma)$ is weakly degenerated, that $\text{Supp}(\mu_0^\pm) \subset \Lambda^c$ and that there exists $\tau_0 \in \mathbb{R}^{++}$ such that for $(x, \omega) \in \text{Supp}(\mu_0^\pm)$ and for $t \in [0, \tau_0]$, $(x \pm t \omega, \omega) \in \Omega$. Then, for $\chi \in C_0^\infty([0, \tau_0])$ and $\phi \in C_0^\infty(\{\det B(x) \neq 0\}),$

$$\lim_{n \to +\infty} \int_{\Omega \times \mathbb{R}^d} \chi(t) \phi(t) e_n(t, x) dt dx = \int_{\Omega \times \mathbb{R}^d} \chi(t) \phi(t) d\mu_0^+(x, \omega) dt + \int_{\Omega \times \mathbb{R}^d} \chi(t) \phi(t) d\mu_0^-(x, \omega) dt$$

where for all $a \in C_0^\infty(\Omega \times S^{d-1}),$

$$(a, \mu_0^\pm) = \int_{\Omega \times S^{d-1}} a(x \pm t \omega, \omega) \exp \left[ - \int_0^t \frac{(\gamma(x \pm \sigma \omega) \cdot \omega)^2}{B(x \pm \sigma \omega) \omega \cdot \omega} d\sigma \right] d\mu_0^\pm (x, \omega).$$

Remark 2.5. 1) Note that even though the support of $\phi$ does not intersect the set $\{\det B(x) = 0\}$, the trajectories $x \pm t \omega$ which reach the support of $\phi$ for $t \in [0, \tau_0]$ may pass through it.

2) Besides, because of (1) in Assumption 2.2, a trajectory $x \pm t \omega$ crosses $\Sigma$ at a finite number of times $0 < t_1 < \cdots < t_N \leq \tau_0$. Assuming $B(x + t_j \omega) \omega = 0$, by (2) of Assumption 2.2, for $j \in \{1, \cdots, N\}$, there exists $c_j \in \mathbb{R}$ such that

$$\frac{(\gamma(x \pm \sigma \omega) \cdot \sigma \omega)^2}{\omega \cdot B(x \pm \sigma \omega) \omega} \sim c_j (s - t_j) \quad \text{as} \quad s \sim t_j.$$

This implies that the integral in (2.4) is well defined.

This result is proved in Section 4.2. The measures $\mu_0^\pm$ contain the part of the energy corresponding to the projection of $U_n$ on the $\pm$-mode. There is no contribution of the 0-mode (which corresponds to the temperature) because of the smoothing effect of the Heat equation. Note finally that the damping in (2.4) can be 0 if for all times $t \in [0, \tau_0]$, we have $\gamma(x \pm t \omega) \cdot \omega = 0$. We refer to section 4.2 for a discussion of what happens when (2) fails in Assumption 2.2: all the energy may be damped in finite time (see Remark 1.3).

Finally, in the Appendix, we shortly discuss what happens close to the boundary under the following assumptions:

Assumption 2.6. (1) The rays of geometric optics have no contact of infinite order with the tangent to $\partial \Omega$.

(2) There exists a compact $K$ such that $\{\det B(x) = 0\} \subset K \subset \Omega$. 

depend on the cut-off function $\chi$ and suppose now that $B$ where $W$ the function arguments which are exactly the same ones.

in Appendix A how to reduce to the analysis of [4], without writing in details the energy along the generalized bicharacteristics as defined in [19]. We shortly explain of a family of solutions of a wave equation. One can prove the propagation of the papers [3], [4] or [16] for the analysis at the boundary of microlocal measures modifications. We set

Before stating the result, let us introduce a notation. Define

Theorem 2.7. Suppose that $\mu_0$ holds, then one can use the methods developed in [13] and the papers [3], [4] or [16] for the analysis at the boundary of microlocal measures a family of solutions of a wave equation. One can prove the propagation of the energy along the generalized bicharacteristics as defined in [19]. We shortly explain in Appendix A how to reduce to the analysis of [4], without writing in details the arguments which are exactly the same ones.

2.4. Distorted propagation in an open set included in $\{B(x) = 0\}$. We suppose now that $B(x) = 0$ an open subset of $\Omega$. Then, the symbol $P(x, \xi)$ is self-adjoint on $\Omega \times \mathbb{R}^d$ and the method of [13] can be adapted with straightforward modifications. We set

For all $x \in \partial \Omega$, $\gamma(x) \in T_x(\partial \Omega)$

If Assumption 2.6 holds, then one can use the methods developed in [13] and the papers [3], [4] or [16] for the analysis at the boundary of microlocal measures modifications. We set

we get

and $\tilde{\nu}_0$ the microlocal defect measures of the sequences

and we consider $\tilde{\nu}_0^{\pm}$ and $\tilde{\nu}_0$ the microlocal defect measures of the sequences

we get

Before stating the result, let us introduce a notation. Define

the Hamiltonian vector field associated with $c$ and $H_c^\infty$ the vector field induced by $H_c$ on $S^*\Omega$.

Theorem 2.7. Suppose that $\text{Supp}(\tilde{\nu}_0) \subset \tilde{\Omega}$, $\text{Supp}(\tilde{\mu}_0^\pm) \subset \tilde{\Omega}$ and that there exists $\tau_0 \in \mathbb{R}^+$ such that for all $x \in \tilde{\Omega}$ and $\omega \in S^{d-1}$, the projection on $\tilde{\Omega}$ of the integral curve of $H_c^\infty$ issued from $(x, \omega)$ stays in $\tilde{\Omega}$ on the time interval $[0, \tau_0]$. Then, for $\chi \in C_0^\infty(\mathbb{R})$ and $\phi \in C_0^\infty(\tilde{\Omega})$,

we get

the limit

$$
\lim_{n \to +\infty} \int_{\mathbb{R}^+ \times \Omega} \chi(t) \phi(x) e_n(t, x) dt dx = \int_{\mathbb{R}^+ \times \Omega \times S^{d-1}} \chi(t) \phi(x) d\tilde{\mu}_0^+(x, \omega) dt
$$

and

$$
\lim_{n \to +\infty} \int_{\mathbb{R}^+ \times \Omega} \chi(t) \phi(x) e_n(t, x) dt dx = \int_{\mathbb{R}^+ \times \Omega \times S^{d-1}} \chi(t) \phi(x) d\tilde{\nu}_0(x, \omega) dt
$$
where \( \tilde{\mu}_t^\pm \) satisfy the transport equations \( \partial_t \tilde{\mu}_t^\pm + H^\infty(x,\omega)\tilde{\mu}_t^\pm = 0 \) with initial data \( \tilde{\mu}_0^\pm \).

**Remark 2.8.** Let us call distorted bicharacteristic curves the trajectories of \( S^*\Omega \) associated with \( H^\infty : \) the transport equation for \( \tilde{\mu}_t^\pm \) implies that the energy propagates along these trajectories.

Note that inside \( \tilde{\Omega} \), one cannot separate the contribution to the energy density of \((u_n)\) and of \((\theta_n)\). The measures \( \tilde{\mu}_t^\pm \) depend of the value at time \( t = 0 \) of both quantities \((u_n)\) and \((\theta_n)\).

These results call for further works: it would be interesting to know what happens at the boundary of \( \tilde{\Omega} \) and how transitions occur between the two regimes. It would be also interesting to know whether this result in \( \{ B(x) = 0 \} \) extend to the set \( \Lambda \) provided \( \Lambda \) is invariant by the distorted rays. The following example shows such a situation.

**Example 2.9.** Suppose \( d = 2 \), \( \Omega = \mathbb{R} \), \( \gamma(x) = e_1 \), \( B(x) = \begin{pmatrix} b(x_2) & 0 \\ 0 & 1 \end{pmatrix} \). Suppose that \( b(y) = 0 \) if and only if \( y = 0 \). Then, \( \Lambda = \{(x,\omega), \exists y \in \mathbb{R}, x = (y,0), \omega = (\pm 1,0)\} \) is invariant by the distorted bicharacteristic curves issued from points of \( \Lambda \):

\[
(x^+_s, \xi^+_s) = ((y \pm s\sqrt{2},0), (\pm 1,0)), \quad s \in \mathbb{R}.
\]

Note that, in that case, the distorted trajectories issued from points of \( \Lambda \) coincide with the usual ones; however, they are described with different speed. Note also that (2) of Assumptions 2.2 is not satisfied here.

2.5. **Organization of the paper.** The main part of the article consists in the analysis of the m.d.m.s associated with the families \((\Pi^k(x,D)U_n)\) (for \( k \in \{0,+,−\} \)) where the functions \( \Pi^k \) are defined in Proposition 2.1. We begin in Section 3 by studying the symbol \( P(x,\xi) \), which allows to prove Proposition 2.1. Then, in Section 4 we prove Theorems 2.3 and 2.7; they rely on the analysis of the propagation of the m.d.m.s associated with the sequences \((\Pi^k(x,D)U_n)\) which is the object of Section 5. Finally, the Appendix A is devoted to a discussion of the reflection of the measures on the boundary.

3. **Analysis of the symbol of** \( P(x,D) \)

In this section we analyze the properties of the matrix \( P(x,\xi) \). The main interest of Weyl quantization is that the symbol of a self-adjoint operator is real-valued. We denote by \( \sigma(A) \) the symbol of an operator \( A \) and we have in particular

\[
\sigma(\nabla) = i\xi,
\]

\[
(3.1) \quad \sigma(\gamma(x) \cdot \nabla) = i\gamma(x) \cdot \xi - \frac{1}{2} \nabla \cdot \gamma(x),
\]

\[
(3.2) \quad \sigma(\nabla \cdot (B(x)\nabla)) = -B(x)\xi \cdot \xi + b_0(x),
\]

where

\[
(3.3) \quad b_0(x) = -\frac{1}{4} \sum_{1 \leq j,k \leq d} \partial^2_{x_j,x_k} B_{jk}(x).
\]
Observe that, if \( d = 1 \), the function \( b_0 \) has a sign on \( \Lambda \). Indeed, if \( d = 1 \), the points \((x, \omega) \in \Lambda \) correspond to values \( x \) which are minima of \( B(x) \) and in this case \( b_0 \leq 0 \). However, in higher dimension, the function \( b_0(x) \) can be positive or negative indifferently as shows the following example: choose \( d = 2 \) and \( B(x) \) such that \( \Lambda = \{(0, y), (0, \pm 1), y \in \mathbb{R} \} \) with
\[
B(x) = \begin{pmatrix}
x_1^4 & 0 \\
0 & 1 + x_2^3
\end{pmatrix} + O(|x|^5)
\]
close to \((0, 0)\). Then we have
\[
b_0(x) = -\frac{1}{4} (12x_1^2 + 6x_2) + O(|x|^3).
\]
Therefore, \( b_0(0, y) = -\frac{3}{2} y + O(|y|^3) \) and \( b_0 \) changes of sign on \( \Lambda \). Note also that if \( B(x) = 0 \) in an open subset \( \tilde{\Omega} \) of \( \Omega \), then \( b_0 = 0 \) in \( \tilde{\Omega} \).

The eigenvalues of the matrix \( P(x, \xi) \) satisfies the following proposition:

**Proposition 3.1.** The matrix \( P(x, \xi) \) has a kernel of dimension greater or equal to \( d - 1 \). Moreover, \( P(x, \xi) \) have three smooth eigenvalues \( \lambda_-, \lambda_0 \) and \( \lambda_+ \) with:

1) If \( (x, \frac{\xi}{|\xi|}) \notin \Lambda \), then \( \dim [\text{Ker} \ P(x, \xi)] = d - 1 \) and
\[
\begin{align*}
\lambda_0(x, \xi) &= -i \xi \cdot B(x) \xi + ib_0(x) + i \frac{(\gamma(x) \cdot \xi)^2}{\xi \cdot B(x) \xi} + O(|\xi|^{-1}), \\
\lambda_{\pm}(x, \xi) &= \pm \beta(x, \xi) \pm i \alpha(x, \xi),
\end{align*}
\]
with
\[
\beta(x, \xi) = |\xi| + O(1), \quad \alpha(x, \xi) = -\frac{1}{2} \frac{(\gamma(x) \cdot \xi)^2}{\xi \cdot B(x) \xi} + O(|\xi|^{-1}).
\]

2) If \( (x, \frac{\xi}{|\xi|}) \in \Lambda \), then, if \( b_0(x) = 0 \), \( \dim [\text{Ker} \ P(x, \xi)] = d \) and if \( b_0(x) \neq 0 \), \( \dim [\text{Ker} \ P(x, \xi)] = d - 1 \). Moreover,
\[
\begin{align*}
\lambda_0(x, \xi) &= ib_0(x) \frac{|\xi|^2}{c(x, \xi)^2} + O(|\xi|^{-1}), \\
\lambda_{\pm}(x, \xi) &= \pm \beta(x, \xi) \pm i \alpha(x, \xi)
\end{align*}
\]
with
\[
\beta(x, \xi) = c(x, \xi) + O(1), \quad \alpha(x, \xi) = \frac{1}{2} \frac{b_0(x) (\gamma(x) \cdot \xi)^2}{c(x, \xi)^2} + O(|\xi|^{-1}),
\]
and \( c(x, \xi) = \sqrt{(\gamma(x) \cdot \xi)^2 + |\xi|^2} \).

The modes \( \pm \) (corresponding to the eigenvalues \( \lambda_{\pm} \)) give the wave feature of the equation. The speed of propagation is characterized by the function \( \beta \) and the function \( \alpha \) corresponds to the damping. Note that in the second case, the speed of propagation are distorted by comparison with the initial wave operator \( \partial_t^2 - \Delta \). Outside \( \Lambda \), the eigenvalue \( \lambda_0 \) encounters of the Heat aspect.

**Proof:** We write \( P(x, \xi) = iQ(x, \xi) \) and for simplicity we work with \( Q(x, \xi) \). For \( p, q \in \mathbb{N}^* \), we denote by \( 0_{p,q} \) the \( p \times q \) matrix with all coefficients equal to 0. We have
\[
Q(x, \xi) = \begin{pmatrix}
0 & i\xi & -\bar{k}(x, \xi) \\
\bar{i} \xi & 0_{d,d} & 0_{d,1} \\
k(x, \xi) & 0_{1,d} & b(x, \xi)
\end{pmatrix},
\]
where, in view of (3.1) and (3.2), \( b = b_2 + b_0 \) with \( b_2 = -B(x)\xi \cdot \xi \) real-valued and \( k = ik_1 + k_0 \) with
\[ k_1 = -\gamma(x) \cdot \xi, \quad k_0 = \frac{1}{2} \nabla \cdot \gamma(x). \]
The vector \((x, Y, y) \in \mathbb{R} \times \mathbb{R}^d \times \mathbb{R}\) is an eigenvector of \( Q \) for the eigenvalue \( \nu \) if and only if
\[
\begin{align*}
(1) & \quad \nu x - i\xi \cdot Y + \nabla y = 0, \\
(2) & \quad ix\xi - \nu Y = 0, \\
(3) & \quad kx + (b - \nu)y = 0.
\end{align*}
\]
Let us suppose first that \( \nu = 0 \), then for \( \xi \neq 0 \), equation (2) gives \( x = 0 \). For \( b \neq 0 \), equation (3) then implies \( y = 0 \) and (1) gives \( \xi \cdot Y = 0 \). Therefore, for large \( \xi \)
the kernel of \( Q(x, \xi) \) is of dimension at least \( d - 1 \) and contains the smooth vectors
\[ V_j(\xi) = (0, e_j(\xi), 0) \]
where \((e_j(\xi))_{1 \leq j \leq d}\) is a smooth basis of the hyperplane orthogonal to \( \xi \) (such a basis exists for \( \xi \neq 0 \)).

Let us now suppose \( \nu 
eq 0 \). Equation (2) implies that \( Y \) is colinear to \( \xi \) and \( x = -i\nu r \) where \( Y = r\xi \). Equations (1) and (3) become
\[
\begin{align*}
\begin{cases}
\nabla y - i(\nu^2 + |\xi|^2)r &= 0, \\
(b - \nu)y - ivkr &= 0.
\end{cases}
\end{align*}
\]
Therefore, the non zero eigenvalues are the roots of the real-valued polynomial
\[ f(X) = -X^3 + X^2b - X(|\xi|^2 + |k|^2) + b|\xi|^2. \]
It is easy to see that for large \( \xi \), \( f \) has one real-valued root \( \nu_0 \) (with \( \nu_0 \neq 0 \)
for \( b_0 \neq 0 \)) and two conjugated complex-valued roots \( \nu_+ \) and \( \nu_- \). These three roots of the polynomial \( f \) are simple and thus smooth; they give three smooth
eigenvalues of \( Q \). Consider three associated eigenvectors \( V_0, V_+ \) and \( V_- \), they are independent from the vectors \( V_j(\xi) \) defined above. Therefore, we are left with a basis of eigenvectors: the matrix \( Q(x, \xi) \) diagonalize.

Let us now study more precisely the asymptotics of the eigenvalues.
1) Suppose \((x, \xi) \notin \Lambda \), \( \xi \) large and denote by \( X_1 < X_2 < 0 \) the two negative roots
of \( f'(X) \). We have \( X_1 = \frac{2}{b_0} b_0 + O(|\xi|) \) and \( X_2 = O(|\xi|) \). Then, using \( f'(X_2) = 0 \), we obtain
\[ f(X_2) = b(X_2^2 + |\xi|^2) < 0 \]
and deduce that \( f \) has only one real-valued root \( \nu_0 \) with \( \nu_0 < X_1 \). We set \( \nu_0 = \phi b \) with \( \phi \geq \frac{2}{3} + o(1) \) and
\[(3.8) \quad 0 = f(\nu_0) = b^3 \phi^2 (1 - \phi) - \phi b(|k|^2 + |\xi|^2) + b|\xi|^2. \]
Necessarily, \( \phi = 1 + r \) with \( r = -k_0^2/b_0^2 + O(|\xi|^{-3}) \), whence \( \nu_0 = b - \frac{k_0^2}{b_0^2} + O(|\xi|^{-1}) \)
and \((3.4)\).
Let \( \nu_\pm \) be the two other (non real-valued) roots, we set \( \nu_\pm = \alpha \mp i\beta \). We observe
\[(3.9) \quad \begin{cases}
\nu_+ + \nu_- = 2\alpha = b - \nu_0, \\
\nu_0(\nu_+ + \nu_-) + \nu_+ \nu_- = 2\alpha \nu_0 + \alpha^2 + \beta^2 = |k|^2 + |\xi|^2,
\end{cases}
\]
whence \( \alpha = \frac{k_0^2}{b_0^2} + O(|\xi|^{-1}) \) and \( \beta^2 = |\xi|^2 + O(|\xi|) \). This implies \((35)\).
2) Suppose now \((x, \xi) \in \Lambda\), then \(b(x, \xi) = b_0(x)\). The polynomial function \(f'(X)\) has no real-valued root and \(f(X)\) has only one real-valued root \(\nu_0\). Since \(f(0)f(b_0) < 0\), the function \(\nu_0(x, \xi) = O(1)\) as \(\xi\) grows and

\[
\nu_0(|\xi|^2 + |k|^2) = b_0|\xi|^2 + O(1),
\]

whence (3.6). Besides, denoting as before by \(\nu_\pm = \alpha \mp i\beta\) the two other roots, (3.9) gives

\[
\alpha = \frac{b_0}{2} \frac{k_1^2}{k_1^2 + |\xi|^2} + O(|\xi|^{-1}) \quad \text{and} \quad \beta^2 = k_1^2 + |\xi|^2 + O(1),
\]

whence (3.7). ♦

Let us now describe the eigenspaces of \(P(x, \xi)\) for large \(\xi\).

**Proposition 3.2.** For \(\xi\) large enough, there exists a smooth basis of \(\mathcal{C}^{2d+2}\),

\[
(V_0, V_+, V_-, V_j), \quad 1 \leq j \leq d - 1 \quad \text{such that}
\]

\[
P(x, \xi)V_k = \lambda_k V_k \quad \text{for} \quad k \in \{0, +, -\},
\]

\[
P(x, \xi)V_j = 0 \quad \text{for} \quad 1 \leq j \leq d - 1,
\]

with the following expansions

1) If \(\left(x, \frac{\xi}{|\xi|}\right) \notin \Lambda\),

\[
V_j = (0, e_j(\xi), 0), \quad e_j(\xi) \cdot \xi = 0, \quad 1 \leq j \leq d - 1,
\]

(3.10)

\[
V_\pm = \frac{i}{\sqrt{2}} \left( \mp 1, \mp \frac{\xi}{|\xi|}, 0 \right) + O(|\xi|^{-1}),
\]

(3.11)

\[
V_0 = (0, 0, 1) + O(|\xi|^{-2}).
\]

2) If \(\left(x, \frac{\xi}{|\xi|}\right) \in \Lambda\),

\[
V_j = (0, e_j(\xi), 0), \quad e_j(\xi) \cdot \xi = 0, \quad 1 \leq j \leq d - 1,
\]

(3.13)

\[
V_\pm = \frac{1}{\sqrt{2}} \left( -1, \pm \frac{\xi}{c(x, \xi)}, \mp \frac{\gamma(x, \xi)}{c(x, \xi)} \right) + O(|\xi|^{-1}),
\]

(3.14)

\[
V_0 = \left(0, \frac{\gamma(x, \xi)}{c(x, \xi)}, \frac{|\xi|}{c(x, \xi)} \right) + O(|\xi|^{-1}).
\]

(3.15)

Note that there exits smooth eigenvectors but their asymptotics are discontinuous; similarly, even though they are not orthogonal projectors, their asymptotics give a decomposition of \(\mathcal{C}^{2d+2}\) on orthogonal subspaces. Moreover, when \(b_0 = 0\), there is an eigenvalue crossing between the eigenspace for \(\lambda_0\) which merges into the kernel of \(P\). However, there still exists a smooth basis of eigenvectors and we will take advantage of this fact in the following sections. Besides, we have the following remark.

**Remark 3.3.** Assuming (2) of Assumption 2.2, we have for \((x, \omega) \in \Lambda\),

\[
\gamma(x) \cdot \omega = \gamma(x) \cdot (B(x)\omega) = 0.
\]
Therefore, if \( (x, \frac{\xi}{|\xi|}) \in \Lambda \), \( c(x, \xi) = |\xi| \) and

\[
\begin{align*}
V_j &= (0, e_j(\xi), 0), \quad e_j(\xi) \cdot \xi = 0, \quad 1 \leq j \leq d - 1, \\
V_\pm &= \frac{1}{\sqrt{2}} \left( -1, \pm \frac{\xi}{|\xi|}, 0 \right) + O(|\xi|^{-1}), \\
V_0 &= (0, 0, 1) + O(|\xi|^{-1}).
\end{align*}
\]

Let us now prove Proposition 3.2.

**Proof:** In the proof of Proposition 3.1 we have already built the vectors \( V_j \). Besides, we have seen that the eigenvectors of \( Q(x, \xi) \) associated with \( \nu_0, \nu_+ \) and \( \nu_- \) are of the form \((x, r\xi, y)\) with \( x = -ivr \) and

\[
\begin{align*}
\{ (b - v) & y - ivkr &= 0, \\
\frac{1}{ky} - i(\nu^2 + |\xi|^2)r &= 0.
\end{align*}
\]

Let us consider first the \( \pm \)-modes. We have \( \nu_\pm = O(|\xi|) \). Therefore, \( b - \nu_\pm \neq 0 \) for large \( \xi \) independently of the fact that \((x, \xi/|\xi|)\) \( \in \Lambda \) or not. The vectors

\[
V_\pm(x, \xi) = \tilde{r}(\nu_\pm(b - \nu), i(\nu_\pm - b)\xi, k\nu_\pm),
\]

with \( \tilde{r} = (|\nu_\pm - b|^2(|\nu_\pm|^2 + |\xi|^2) + k^2|\nu_\pm|^2)^{-1/2} \) are smooth non-zero eigenvectors associated with \( \nu_\pm \). In view of the asymptotics of Proposition 3.1 we obtain asymptotics for \( V_\pm \):

1. In \( \Lambda^c \), \( \nu_\pm = \mp i|\xi| + O(1), \) \( k = -i\gamma(x) \cdot \xi + O(1) \) and \( b = -B(x)\xi \cdot \xi + O(1) \). Therefore

\[
\tilde{r} = \frac{1}{\sqrt{2}} \left( |\xi|B(x)\xi \cdot \xi \right)^{-1} \left( 1 + O(|\xi|^{-1}) \right),
\]

\[
V_\pm(x, \xi) = \tilde{r} \left[ (\mp i|\xi|(B(x)\xi \cdot \xi), i(B(x)\xi \cdot \xi)\xi, \mp|\xi|\gamma(x) \cdot \xi + O(|\xi|^2) \right],
\]

whence 3.11.

2. In \( \Lambda \), \( \nu_\pm = \mp ic(x, \xi) + O(1), \) \( b = O(1) \) and we still have \( k = -i\gamma(x) \cdot \xi + O(1) \). Therefore

\[
\tilde{r} = 1/(\sqrt{2}c(x, \xi)^2)(1 + O(1)),
\]

\[
V_\pm(x, \xi) = \tilde{r} \left[ (-c(x, \xi)^2, \pm c(x, \xi)\xi, \mp c(x, \xi)\gamma(x) \cdot \xi + O(|\xi|) \right],
\]

whence 3.11.

Let us consider now the 0-mode. We have \( \nu_0 = O(1) \) in \( \Lambda \) and \( \nu_0 + O(1) \in \mathbb{R} \) in \( \Lambda^c \). Therefore, \( \nu_0^2 + |\xi|^2 \neq 0 \) for large \( \xi \) and the vector

\[
V_0(x, \xi) = r_0 (\nu_0 \overline{k}, -i\overline{k} \xi, \nu_0^2 + |\xi|^2)
\]

with

\[
r_0 = \left( |\nu_0|^2|k|^2 + |\xi|^2|k|^2 + (\nu_0^2 + |\xi|^2)^2 \right)^{-1/2}
\]

is a smooth eigenvector associated with the eigenvalue \( \nu_0 \). We are now left with a smooth basis of eigenvectors. Let us now study the asymptotics of this vector.

1. In \( \Lambda^c \), \( \nu_0 = -B(x)\xi \cdot \xi + O(1) \), therefore

\[
r_0 = (B(x)\xi \cdot \xi)^{-2} (1 + O(|\xi|^{-1})) \quad \text{and} \quad V_0(x, \xi) = (0, 0, 1) + O(|\xi|^{-2}).
\]
In $\Lambda$, $\nu_0 = O(1)$, whence
\[
    r_0 = 1/(c(x, \xi) |\xi|)(1 + O(|\xi|^{-1}),
\]
\[
    V_0(x, \xi) = r_0 \left[ (0, (\gamma(x) \cdot \xi) \xi, |\xi|^2) + O(|\xi|) \right],
\]
which gives (3.15).

Before concluding this section, we point out that these asymptotics imply Proposition 2.1. Indeed, we have obtained the existence of $R_0$ such that for $|\xi| > R_0$, $P(x, \xi)$ diagonalize with smooth eigenprojectors given by the Gram matrix of the basis of eigenvectors of Proposition 3.2. Asymptotically, these projectors are orthogonal and we have
\[
    \Pi^k(x, \xi) = V_k(x, \xi) \otimes V_k(x, \xi) + o(1).
\]
Since $U_n(t, x)$ goes weakly to 0 in $L^2(\Omega)$, we have
\[
    |U_n(t, x)|^2_{C^{s+2}} = \left| \chi \left( \frac{D}{R} \right) U_n(t, x) \right|^2 + o(1) \text{ in } D'(\Omega).
\]
Besides for $R > R_0$, we observe that $\Pi^k(x, D) \chi(D/R)U_n(t, x)$ is well-defined and we have in $D'(\Omega)$,
\[
    \left| \Pi^k(x, D) \chi \left( \frac{D}{R} \right) U_n(t, x) \right|^2_{C^{s+2}} = \left| \left( V_k(x, D) \chi \left( \frac{D}{R} \right) U_n(t, x) \right) \right|_{C^{s+2}}^2 + o(1)
\]
We can now use the asymptotics of Proposition 3.2 which, combined with the weak convergence to 0 of $U_n(t, x)$ gives
\[
    \sum_{k \in \{0, -, +\}} \left| \left( V_k(x, D) \chi \left( \frac{D}{R} \right) U_n(t, x) \right) \right|_{C^{s+2}}^2
\]
\[
    = (\partial_t u_n(t, x))^2 + |\nabla u_n(t, x)|^2 + (\theta_n(t, x))^2 + o(1) \text{ in } D'(\Omega),
\]
whence Proposition 2.1.

4. PROOF OF THE MAIN RESULTS

The proofs of Theorem 2.4 and Theorem 2.7 are inspired by the method developed in [14] for analyzing semi-classical measures associated with solutions of a system of p.d.e’s. The proof of Theorem 2.7 is a direct adaptation of the results of [14] in the microlocal defect measures setting while the proof of Theorem 2.4 requires non-trivial adaptations due to the fact that $P(x, \xi)$ is not self-adjoint and that one of its eigenvalue is a symbol of order 2. Therefore, we focus on the proof of Theorem 2.4 and we let to the reader the simple adaptation of these arguments to prove Theorem 2.7. Theorem 2.4 relies on Propositions 4.1 and 4.2 stated in Subsection 4.1, then the proof of Theorem 2.4 is performed in Subsection 4.2.

4.1. Preliminaries. We state technical results (Propositions 4.1 and 4.2) that we will use in the next subsection. The proof of Proposition 4.11 is done at the end of this section while the one of Proposition 4.12 is postponed to Section 5.

The first result describes the evolution of the temperature $\theta_n$:

**Proposition 4.1.** Let $\nu$ be a microlocal defect measure of the sequence $(\theta_n)_{n \in \mathbb{N}}$ of $L^2_{\text{loc}}(\mathbb{R}, L^2(\mathbb{R}^d))$, then, as a measure on $\mathbb{R}^+ \times \Lambda^c$, $\nu = 0$. 

Our second result concerns the contribution to the energy density of the sequences $(\Pi^\pm(x,D)\chi(D/R)u_n)_{n \in \mathbb{N}}$ for $R > R_0$ and $\chi$ as in Proposition 4.1. We set

$$U^\pm_{n,R}(t,x) = \Pi^\pm(x,D)\chi\left(\frac{D}{R}\right)U_n.$$  

The sequences $(U^\pm_{n,R})_{n \in \mathbb{N}}$ are uniformly bounded in $L^2(\Omega)$ and their microlocal defect measures are matrix-valued measures independent of $R > 0$. Besides, by the definition of $U^\pm_{n,R}$, these measures are of the form $\mu^\pm(t,x,\omega)\Pi^\pm(x,\omega)$ where the measures $\mu^\pm(t,x,\omega)$ can be understood as the traces of $\mu^\pm(t,x,\omega)\Pi^\pm(x,\omega)$. We prove the following result.

**Proposition 4.2.** Assume condition (2) of Assumption 2.2 is satisfied in $\Omega_1 \subset \Omega$ and let $T > 0$. There exists a subsequence $n_k$ and a continuous map $t \mapsto \mu^\pm(t)$ from $[0,T]$ into the set of positive Radon measures on $\Omega \times S^{d-1}$ such that for all $t \in [0,T]$ and for all scalar symbol $a \in \mathcal{A}_0^0$, we have

$$\left(\begin{array}{l}
(a(x,D_x)U^\pm_{n_k,R}(t) | U^\pm_{n_k,R}(t)) 
\end{array}\right) \xrightarrow{k \to +\infty} \int a_{\infty}(x,\omega)d\mu^\pm(t,x,\omega).$$

Moreover, in $\mathcal{D}'(\{t \geq 0\} \times \Omega_1 \times S^{d-1})$, we have

$$\partial_t \mu^\pm \pm H_\beta(\mu^\pm) - 2\alpha_{\infty}\mu^\pm = \nu^\pm$$

where $\nu^\pm$ is a measure supported on $\Lambda$ absolutely continuous with respect to $\mu^\pm$ and where for all $a \in \mathcal{A}_0^0$,

$$\int a_{\infty}(x,\omega)d(H_\beta(\mu^\pm)) = -\int (H_\beta a)_{\infty}d\mu^\pm = \int \langle \{a,\beta\}\rangle_{\infty}d\mu^\pm.$$  

Proposition 4.2 is proved in Section 5 below.

Let us now prove Proposition 4.1.

**Proof:** It is enough to prove that if $q$ is a symbol of order 0 such that $q(x,\omega) \in C^\infty(\Lambda')$, then $q(x,D)\theta_n$ goes to 0 in $L^2_{\theta_n}(\mathbb{R}, L^2(\Omega))$. We observe that we only need to consider large values of $\xi$. Indeed, if $\chi \in C^\infty(\mathbb{R}^d)$, $\chi(\xi) = 0$ for $|\xi| \leq 1$ and $\chi(\xi) = 1$ for $|\xi| \geq 2$ with $0 \leq \chi \leq 1$, we have

$$q(x,D)\theta_n = q(x,D)\chi\left(\frac{D}{R}\right)\theta_n + o(1)$$

in $L^2(\Omega)$ because $\theta_n$ goes weakly to 0 and the operator $1 - \chi(\frac{D}{R})$ is compact. We write

$$q(x,\xi)\chi\left(\frac{\xi}{R}\right) = \frac{1}{R}Q_R(x,\xi) \cdot \sqrt{B(x)}\xi$$

where $Q_R$ is the vector-valued symbol of order $-1$:

$$Q_R(x,\xi) = R\chi\left(\frac{\xi}{R}\right)\frac{q(x,\xi)}{B(x)\xi} \sqrt{B(x)}\xi = \bar{\chi}\left(\frac{\xi}{R}\right)\frac{|\xi|q(x,\xi)}{B(x)\xi} \sqrt{B(x)}\xi$$

with $\bar{\chi}(u) = |u|^{-1}\chi(u)$. Note that $Q_R$ is smooth since $q = 0$ in a neighborhood of $\Lambda$ and because $\bar{\chi} \in C^\infty_0(\mathbb{R}^d)$. Note also that $Q_R$ satisfies symbols estimates...
uniformly with respect to $R$. Therefore, there exists a constant $C > 0$ such that for all $f \in L^2(\Omega)$,
$$
\|Q_R(x, D)f\|_{L^2(\Omega)} \leq C \|f\|_{L^2(\Omega)}.
$$
We write
$$
q(x, D)\theta_n = R^{-1}Q_R(x, D)\sqrt{B(x)}\nabla x \theta_n + o(1)
$$
in $L^2(\Omega)$ where the $o(1)$ is uniform in $R > 1$ when $n$ goes to $+\infty$. As a consequence, if $f \in L^2(\Omega)$ and $\psi \in C_0^\infty(\mathbb{R})$, we have
$$
\left|\int \psi(t)(q(x, D)\theta_n(t) \mid f) \, dt\right|
$$
for some constant $C > 0$ and where the $o(1)$ is uniform in $R > 1$ as $n$ goes to $+\infty$. We observe that the energy equality (13) gives that the family $\left(\sqrt{B(x)}\nabla x \theta_n(t)\right)$ is u. b. in $L^2([0, T] \times \Omega, \mathbb{R}^d)$ for all $T > 0$. Therefore, letting $n$ and $R$ go to $\infty$, we obtain the result. 

4.2. Proof of Theorem 2.4
We suppose that $(B, \gamma)$ is weakly degenerated, that is that the conditions (1), (2) and (3) of Assumption 2.2 are satisfied. Let $\chi, \phi$ be as in Theorem 2.4. We observe that the energy density $e_n(t, x)$ writes
$$
\frac{d}{dt}e_n(t, x) = \left|\partial_t u_n - \frac{D}{|D|} \cdot \nabla u_n\right|^2 + \frac{1}{2} \left|\partial_t u_n + \frac{D}{|D|} \cdot \nabla u_n\right|^2.
$$
Therefore, it remains to analyze the limit of each of these terms for $\phi$ supported outside $\{\det B(x) = 0\}$.

First, we observe that by Proposition 4.1 if $\phi \in C_0^\infty(\{\det B(x) \neq 0\})$,
$$
\lim_{n \to +\infty} \int_\mathbb{R} \int_\Omega \chi(t)\phi(x)|\theta_n(t, x)|^2 \, dt \, dx = \int_\mathbb{R} \int_\Omega \chi(t)\phi(x) \, d\nu(t, x, \omega) = 0.
$$
Therefore, the weak limit of the energy density express only in terms of the sequences $\partial_t u_n \pm \frac{D}{|D|} \cdot \nabla u_n$, of which microlocal defect measures are $\mu_\pm(t, x, \omega)$ by Remark 3.3. By Proposition 4.2 the measures $\mu_\pm$ satisfy
$$
\partial_t \mu_\pm \pm \omega \nabla x \mu_\pm = -1 \Lambda(\gamma(x) \cdot \omega)^2 \frac{B(x)}{\omega \cdot \omega} \mu_\pm + \nu_\pm
$$
where we have used $\alpha_\infty(x, \omega)1_\Lambda = \frac{1}{2} b_0(x) \frac{(\gamma(x) \cdot \omega)^2}{c(x, \omega)^2} = 0$ by Remark 3.3.

Let us now prove that the fact that $\omega$ is transverse to the hypersurface $\Sigma$, implies that $\mu_\pm 1_\Sigma = 0$. Let $f(x) = 0$ be a local equation of $\Sigma$ in a subset $\Omega_2$ of $\Omega$ and let $a$ be a symbol supported in $\Omega_2 \times \mathbb{R}^d$ and such that $a \geq 0$ and $\omega \cdot \nabla f(x) > 0$ for $(x, \xi) \in \text{Supp } a, \xi \neq 0$. We choose a function $\theta \in C_0(\mathbb{R})$ such that $\theta(0) = 1$ and we use the test symbol
$$
b_\delta(x, \xi) = \delta a(x, \xi) \theta \left(\frac{f(x)}{\delta}\right)
$$
where $\delta > 0$. The transport equations for $\mu_\pm$ imply

$$
\int a_\infty(x, \omega) \omega \cdot \nabla f(x) \partial^t \left( \frac{f(x)}{\delta} \right) d\mu_\pm(x, \omega) = O(\delta).
$$

By letting $\delta$ go to 0, we obtain $\int_{f(x)=0} a_\infty(x, \omega) \omega \cdot \nabla f(x) d\mu_\pm(x, \omega) = 0$, whence $\mu_\pm \mathbf{1}_\Sigma = 0$ on the support of $a$. By this way (inspired from [9]), we finally obtain $\mu_\pm \mathbf{1}_\Lambda = 0$ since $\Lambda \subset \Sigma \times \mathbb{R}^d$.

Besides, since the measure $\nu_\pm$ is supported on $\Lambda$ and absolutely continuous with respect to $\mu_\pm$, we deduce $\nu_\pm = 0$.

Finally, we observe that the function $F(x, \omega) = \frac{(\gamma(x) \omega)^2}{B(x)}$ extends continuously to the set $\Omega \times S^{d-1}$ with $F(x, \omega) = 0$ on $\Lambda$ (since $\gamma(x) \omega = \gamma(x) \cdot (B(x) \omega) = O(|B(x)\omega|)$ by (2) of Assumptions [2,2]). Therefore, we can write

$$
\partial_t \mu_\pm \pm \omega \nabla_x \mu_\pm = -\frac{(\gamma(x) \cdot \omega)^2}{B(x)} \omega \mu_\pm.
$$

As a conclusion, we obtain Theorem [2,4]. Indeed, take $a \in C_0^\infty(\Omega \times S^{d-1})$, $0 \leq s \leq t \leq \tau_0$ and set

$$
a_t(s, x, \omega) = a(x \pm (t-s)\omega, \omega) \exp \left[ \int_0^{t-s} F(x \pm \sigma \omega, \omega) d\sigma \right],
$$

then $a_t(t, x, \omega) = a(x, \omega)$

$$
\frac{d}{ds} \langle a_t(s), \mu_{\pm}(s) \rangle = 0
$$

whence formula [2,4].

Let us conclude this section by a remark.

**Remark 4.3.** If (2) of Assumption [2,2] fails, then all the energy may be damped in finite time.

**Proof:** Suppose $\mu^-_\rho = 0$ and $\mu^+_0 = a_0 \delta(x - x_0) \otimes \delta(\omega - \omega_0)$ for $(x_0, \omega_0) \notin \Lambda$ and suppose that there exists $\tau_0$ such that

$$
\forall t \in [0, \tau_0], \ (x_0 + t \omega_0, \omega_0) \notin \Lambda \ \text{and} \ (x_0 + t_0 \omega_0, \omega_0) \in \Lambda.
$$

Set $y_0 = x_0 + t_0 \omega_0$ and suppose that $dB(y_0)[\omega_0, \omega_0, \omega_0] \neq 0$ and $\gamma(y_0) \cdot \omega \neq 0$ then

$$
\exp \left[ -\int_0^{t_1} \frac{(\gamma(x_0 + s \omega_0) \cdot \omega_0)^2}{B(x_0 + s \omega_0)\omega_0 \cdot \omega_0} ds \right] \to 0
$$

and all the energy is damped between times $0$ and $t_0$.

5. **Propagation of Microlocal defect measures**

In this section, we prove Proposition [4,2]. We consider the + mode; the − mode can be treated in the same way. We proceed in three steps:

- First, we analyze the time derivative of $\left( a(x, D) U^+_{n,R}(t) | U^+_{m,R}(t) \right)$ for scalar valued symbols $a \in \mathcal{A}_0^1$ and prove that there exists a symbol $T \in \mathcal{A}_1$ such
that, uniformly in $R$ as $n$ goes to $+\infty$,

$$
(5.1) \quad \frac{d}{dt} \left( a(x, D)U_{n,R}^+ | U_{n,R}^+ \right) = -\{(a, \beta)(x, D)U_{n,R}^+ | U_{n,R}^+ \} \\\quad + 2((\alpha \alpha)(x, D)U_{n,R}^+ | U_{n,R}^+) \\\quad + \left( T(x, D)\chi \left( \frac{D}{R} \right) U_n | \chi \left( \frac{D}{R} \right) U_n \right) + o(1).
$$

- We calculate precisely the symbol $T(x, \xi)$ and show that $T \in A^0_0$. Therefore, the quantity $$
(5.1) \quad \frac{d}{dt} \left( a(x, D)U_{n,R}^+ | U_{n,R}^+ \right) = -\{(a, \beta)(x, D)U_{n,R}^+ | U_{n,R}^+ \} \\\quad + 2((\alpha \alpha)(x, D)U_{n,R}^+ | U_{n,R}^+) \\\quad + \left( T(x, D)\chi \left( \frac{D}{R} \right) U_n | \chi \left( \frac{D}{R} \right) U_n \right) + o(1).
$$

- Finally, we prove that for all $\psi \in C^\infty_0(R)$,

$$
\lim_{n \to +\infty} \int \psi(t) \left( T(x, D)\chi \left( \frac{D}{R} \right) U_n(t) \right) + \chi \left( \frac{D}{R} \right) U_n(t) \ dt \quad = \quad \int \psi(t) a_\infty(x, \omega) d\nu_\pm(x, \omega) dt
$$

where $\nu_\pm$ is a measure supported on $\Lambda$ and absolutely continuous with respect to $\mu_\pm$.

At the end of these three steps, we have obtained Proposition 4.2. We now detail each of these steps.

5.1. 1st. step: proof of (5.1). The family of functions $U_{n,R}^+$ is solution of the equation

$$
(5.3) \quad i\partial_t U_{n,R}^+(t, x) = \Pi^+(x, D)\chi \left( \frac{D}{R} \right) P(x, D)U_n.
$$

Since $\nabla \chi$ is compactly supported, we have

$$
\chi \left( \frac{D}{R} \right) P(x, D) = P(x, D)\chi \left( \frac{D}{R} \right) + K_R(x, D),
$$

where $K_R$ is a compact operator. Moreover we have

$$
\Pi^+(x, D)P(x, D) = \lambda_+(x, D)\Pi^+(x, D) + R_1(x, D),
$$

where the symbol $R_1$ will be precisely calculated in the next subsection. So (5.3) becomes

$$
(5.4) \quad i\partial_t U_{n,R}^+(t, x) = (\lambda_+(x, D)\Pi^+(x, D) + R_1(x, D)) \chi \left( \frac{D}{R} \right) U_n + \Pi^+(x, D)K_R(x, D)U_n
$$

$$
(5.5) \quad F_{n,R}^+(x, D) = R_1(x, D)\chi \left( \frac{D}{R} \right) U_n + \Pi^+(x, D)K_R(x, D)U_n.
$$
For a real-valued symbol \( a \in \mathcal{A}^0_r \), \[ (5.6) \] implies \( \frac{d}{dt} \langle a(x, D)U_{n,R}^+|U_{n,R}^+ \rangle = I_1 + I_2 \), with

\[
I_1 = \frac{1}{i} \langle (a(x, D)\lambda_+(x, D)U_{n,R}^+|U_{n,R}^+ \rangle - \frac{1}{i} \langle (a(x, D)U_{n,R}^+|\lambda_+(x, D)U_{n,R}^+ \rangle,
\]

\[
I_2 = \frac{1}{i} \langle (a(x, D)F_{n,R}^+(x, D)|U_{n,R}^+ \rangle - \frac{1}{i} \langle (a(x, D)U_{n,R}^+|F_{n,R}^+ \rangle.
\]

The term \( I_1 \) will give the transport by the vector field \( H_\beta \) and the damping by \( \alpha \). The term \( I_2 \) is a rest term and its main contribution will be described by a symbol \( T \).

Let us study \( I_1 \).

\[
I_1 = \frac{1}{i} \langle (a(x, D)\lambda_+(x, D) - (\lambda_+(x, D))^*a(x, D))U_{n,R}^+|U_{n,R}^+ \rangle.
\]

We recall that \( \lambda_+(x, D) = \beta(x, D) + i\alpha(x, D) \), so, since we use the Weyl quantification for the symbols we have \( (\lambda_+(x, D))^* = \beta(x, D) - i\alpha(x, D) \) and

\[
I_1 = \frac{1}{i} \langle (a(x, D), \beta(x, D)|U_{n,R}^+|U_{n,R}^+ \rangle + 2\langle (a\alpha)(x, D)U_{n,R}^+|U_{n,R}^+ \rangle
\]

\[ + \langle r_{-1}(x, D)U_{n,R}^+|U_{n,R}^+ \rangle \]

\[ = -\langle \{a, \beta\}(x, D)U_{n,R}^+|U_{n,R}^+ \rangle + 2\langle (a\alpha)(x, D)U_{n,R}^+|U_{n,R}^+ \rangle + \langle r_{-1}(x, D)U_{n,R}^+|U_{n,R}^+ \rangle,
\]

where \( r_{-1}(x, \xi) \) will denote from now on a generic symbol in \( \mathcal{A}_{\epsilon}^i \). Using \[ 5.5 \], we then obtain if \( n, R \) tend to infinity

\[
(5.6) \quad I_1 \to -\int \{a, \beta\}_\infty(x, \omega)d\mu_+(x, \omega) + 2 \int a_\infty(x, \omega)\alpha_\infty(x, \omega)d\mu_+(x, \omega).
\]

Note that, for these terms, we do not need to integrate in time to get the convergence.

Let us now study \( I_2 \). We set \( U_{n,R} = \chi \begin{pmatrix} D \\ R \end{pmatrix} U_n \) and by \[ 5.5 \] we obtain \( I_2 = I_{2,1} + I_{2,2} \) with

\[
I_{2,1}(t) = \frac{1}{i} \langle (a(x, D)\Pi^+(x, D)K_RU_{n}(t))|\Pi^+(x, D)U_{n,R}(t) \rangle
\]

\[ - \frac{1}{i} \langle (a(x, D)\Pi^+(x, D)U_{n,R}(t))|\Pi^+(x, D)K_RU_n(t) \rangle \]

\[
I_{2,2}(t) = \frac{1}{i} \langle (a(x, D)R_1(x, D)U_{n,R}(t))|\Pi^+(x, D)U_{n,R}(t) \rangle
\]

\[ - \frac{1}{i} \langle (a(x, D)\Pi^+(x, D)U_{n,R}(t))|R_1(x, D)U_{n,R}(t) \rangle \]

Since \( (U_n)_{n \in \mathbb{N}} \) goes weakly to 0 as \( n \) goes to \( +\infty \) and \( K_R \) is a compact operator, the sequence \( (K_RU_n)_{n \in \mathbb{N}} \) goes strongly to 0. Therefore

\[
(5.7) \quad \forall t \in \mathbb{R}^+, \quad I_{2,1}(t) \to 0.
\]

Besides, \( I_{2,2}(t) = (T(x, D)U_{n,R} | U_{n,R}) \) with

\[
T(x, D) := \frac{1}{i} \langle \Pi^+(x, D)^*a(x, D)R_1(x, D) - R_1(x, D)^*a(x, D)\Pi^+(x, D) \rangle.
\]
5.2. 2nd. step: the symbol of the rest term. In the following, rest terms in the symbol class $\mathcal{A}_i^{-k}$ for $k \in \mathbb{N}$ will be denoted by $r_{-k}$. We prove the following Lemma.

**Lemma 5.1.** Suppose $\Omega_1$ is like in Proposition 4.2. Then, in $\Omega_1$, $T \in \mathcal{A}_i^0$ and
\[
T(x, \xi) = T_1(x, \xi) + T_2(x, \xi) + r_{-1}(x, \xi)
\]
with
\[
(5.8) \quad T_1 = -\frac{a}{2} \left( \Pi^+ \{ \Pi^+, \beta \} + \Pi^+ \{ \Pi^+, \beta \} \Pi^+ - \Pi^- \{ \Pi^+, \beta \} \Pi^- \right)
\]
\[
(5.9) \quad T_2 = \frac{a}{2i} (\Pi^+ \{ \Pi^+, b \} K - K \{ \Pi^+, b \} \Pi^+)
\]
Besides, $T_1 \Pi_0$ and $\Pi_0 T_1$ are symbols of $\mathcal{A}_i^{-1}$.

Note that $T_2 \in \mathcal{A}_i^0$ because $\Pi^+ K, K \Pi^+ \in \mathcal{A}_i^{-1}$ by Remark 3.2.

**Proof**
We first calculate $R_1$, recall that
\[
R_1(x, D) = \Pi^+ (x, D) P(x, D) - \lambda_+ (x, D) \Pi^+ (x, D).
\]
We write $P(x, \xi) = P_1(x, \xi) + ib(x, \xi) K$, with
\[
P_1(x, \xi) = \begin{pmatrix} 0 & -\xi & \gamma(x) \cdot \xi \\ -\xi & 0_{d,d} & 0_{d,1} \\ \gamma(x) \cdot \xi & 0_{1,d} & 0 \end{pmatrix} \quad \text{and} \quad K = \begin{pmatrix} 0_{d+1,d+1} & 0_{d+1,1} \\ 0_{1,d+1} & 1 \end{pmatrix}.
\]
Since $b(x, \xi)$ is of order 2, $R_1$ is a priori of order 1 and in view of $\Pi^+ P = \lambda_+ \Pi^+$, we have
\[
R_1 = \frac{1}{2} \left( \{ \Pi^+, P \} - \{ \lambda_+, \Pi^+ \} \right) + QK + r_{-1}
\]
where the matrix-valued symbol $Q$ is in $\mathcal{A}_i^0$ and is the sum of second derivatives of $\Pi^+$ multiplied by second derivatives of $b$. More precisely and after ordering the terms of higher degrees, we write
\[
R_1 = \frac{1}{2} \left( \{ \Pi^+, P \} - \lambda_+, \Pi^+ \right) + QK + r_{-1},
\]
where we have used that $\lambda_+ = \beta + i\alpha$ and $\{ \alpha, \Pi^+ \} \in \mathcal{A}_i^{-1}$. However, since $\gamma(x) \in \text{Ran} \ B(x)$, then $\gamma(x) \cdot \xi = 0$ for $\left( x, \frac{\xi}{|\xi|} \right) \in \Lambda$ and, by Remark 3.2
\[
(5.10) \quad K \Pi^+, \Pi^+ K \in \mathcal{A}_i^{-1}.
\]
Derivations of these relations imply $\{ \Pi^+, b \} K \in \mathcal{A}_i^0$ and $QK \in \mathcal{A}_i^{-1}$. Finally, we find
\[
R_1 = \frac{1}{2} \left( \{ \Pi^+, b \} K - \frac{i}{2} \left( \{ \Pi^+, \beta \} + \{ \Pi^+, P_1 \} \right) + r_{-1} \in \mathcal{A}_i^0,
\]
whence $T \in \mathcal{A}_i^0$ with
\[
T = \frac{1}{i} \left( \Pi^+ a R_1 - R_1^* a \Pi^+ \right) + r_{-1}
\]
\[
= \frac{a}{2i} \left( \Pi^+ \{ \Pi^+, b \} K - K \{ \Pi^+, b \} \Pi^+ \right)
\]
\[
- \frac{a}{2} \left( \{ \Pi^+, \beta \} + \{ \Pi^+, \beta \} \Pi^+ + \Pi^+ \{ \Pi^+, P_1 \} - \{ P_1, \Pi^+ \} \Pi^+ \right) + r_{-1},
\]
where we have used \((\Pi^+)^* = \Pi^+ + r_{-1}\) and
\[
R_1^i = \frac{1}{2} K \{ \Pi^+, b \} + \frac{i}{2} \left( \{ \Pi^+, \beta \} - \{ P_1, \Pi^+ \} \right) + r_{-1}.
\]
We now transform the expression of the principal symbol of \(T\). We write
\[
P_1 = P - ibK = \lambda_+ \Pi^+ + \lambda_- \Pi^- + \lambda_0 \Pi^0 - ibK.
\]
By Proposition 3.1, 3.2 and Remark 3.3, we notice that
- in \(\Lambda^c\), \(\Pi^0 = K + O(|\xi|^{-2})\) and \(\lambda_0 = ib + O(1)\),
- in \(\Lambda\), \(\Pi^0 = K + O(|\xi|^{-1})\), \(\lambda_0 = O(1)\) and \(b = O(1)\).
Therefore,
\[
P_1 = \lambda_+ \Pi^+ + \lambda_- \Pi^- + r_0.
\]
In view of (5.9), we deduce \(T = T_2 + \widetilde{T}\) with
\[
\widetilde{T} = -\frac{a}{2} \left( \Pi^+ \{ \Pi^+, \beta \} + \{ \Pi^+, \beta \} \Pi^+ + \sum_{\ell \in \{+,-\}} (\Pi^+ \{ \Pi^+, \beta \} + \{ \Pi^+, \beta \} \Pi^+) \right) + r_{-1}
\]
\[
= -\frac{a}{2} \left( \Pi^+ \{ \Pi^+, \beta \} + \{ \Pi^+, \beta \} \Pi^+ + \sum_{\ell \in \{+,-\}} (\Pi^+ \{ \Pi^+, \beta \} + \{ \Pi^+, \beta \} \Pi^+) \right) + r_{-1}
\]
Observing that \(\{ \Pi^+, \Pi^+ \}, \{ \Pi^+, \Pi^\pm \} \in \mathcal{A}_i^{-1}\) by Remark 3.3 and equation (5.8), we obtain
\[
\widetilde{T} = -\frac{a}{2} \left( \Pi^+ \{ \Pi^+, \beta \} + \{ \Pi^+, \beta \} \Pi^+ + \sum_{\ell \in \{+,-\}} (\Pi^+ \{ \Pi^+, \beta \} + \{ \Pi^+, \beta \} \Pi^+) \right) + r_{-1}
\]
\[
= T_1 + r_{-1}
\]
where we have used \(\Pi^+ \{ \beta, \Pi^+ \} \Pi^+ = 0\) (which comes from \((\Pi^+)^2 = \Pi^+\) whence \(\{ \beta, \Pi^+ \} = \Pi^+ \{ \beta, \Pi^+ \} + \{ \beta, \Pi^+ \} \Pi^+\) and, multiplying by \(\Pi^+\) on both sides, we obtain \(\Pi^+ \{ \beta, \Pi^+ \} \Pi^+ = 0\)). Notice that \(\{ \beta, \Pi^+ \} \Pi^0 \in \mathcal{A}_i^{-2}\) and \(\Pi^0 \{ \beta, \Pi^+ \} \in \mathcal{A}_i^{-2}\), whence \(\Pi^0 T, T \Pi^0 \in \mathcal{A}_i^{-1}\).

5.3. 3rd. step: passing to the limit in the rest term. We use the following Lemma.

**Lemma 5.2.** Consider \(A\) a smooth symbol of order 0 supported in \(\Omega_1\) and such that \(A = \Pi^+ A \Pi^\ell\) or \(A = \Pi^\ell A \Pi^+\) where \(\Pi^\ell = \Pi^-\) or \(\Pi^\ell = \Pi^+\), \(1 \leq j \leq d - 1\). Then if \(\psi \in C_0^\infty (\mathbb{R})\),
\[
\int \psi(t) (A(x, D)U_n(t), U_n(t)) dt \underset{n \to +\infty}{\to} 0.
\]
Since the symbol $T_1$ is the sum of terms of the form $\Pi^+\{\Pi^+, \beta\}\Pi^\ell$ or $\Pi^\ell\{\Pi^+, \beta\}\Pi^+$ with $\ell = - \ell = j$, $1 \leq j \leq d - 1$, we can apply the Lemma and we obtain

$$\int \psi(t) \left( T(x, D) \chi \left( \frac{D}{R} \right) U_n(t) \right) \left| \chi \left( \frac{D}{R} \right) U_n(t) \right| dt = \int \psi(t) \left( T_2(x, D) \chi \left( \frac{D}{R} \right) U_n(t) \right) \left| \chi \left( \frac{D}{R} \right) U_n(t) \right| dt + o(1)$$

as $n$ goes to $+\infty$. In view of (5.9) we have

$$\limsup_{n \to +\infty} \int \psi(t) \left( T_2(x, D) \chi \left( \frac{D}{R} \right) U_n(t) \right) \left| \chi \left( \frac{D}{R} \right) U_n(t) \right| dt = \int \psi(t) \left( T_2(x, D) \chi \left( \frac{D}{R} \right) U_n(t) \right) \left| \chi \left( \frac{D}{R} \right) U_n(t) \right| dt$$

where $\tilde{\nu}_+^+$ is the joint measure of $KU_n = \theta_n$ and of $U_n^+$. The measure $\nu_+^+$ is absolutely continuous with respect to $\mu_+$ and $\nu$. By Proposition 5.1, $\tilde{\nu}_+^+$ is supported on $\Lambda$ and we obtain (5.2) with $\nu_+^+(T_2 \alpha) \tilde{\nu}_+^+$ which is absolutely continuous with respect to $\mu_+$ and supported on $\Lambda$. We now focus on the proof of Lemma 5.2.

**Proof**

Let $\Pi^\ell$ be one of the projectors $\Pi^-$ or $\Pi^\ell$, $1 \leq j \leq d - 1$ and let us consider a term of the form $\Pi^+ A\Pi^\ell$; the proof is similar for the other terms. We have

$$[\Pi^+ A\Pi^\ell, P] = (\lambda_\ell - \lambda_+)\Pi^+ A\Pi^\ell.$$

Let us denote by $C^\ell$ the symbol of $A_\ell^{-1}$

$$C^\ell = (\lambda_\ell - \lambda_+)\Pi^+ A\Pi^\ell.$$

We have $C^\ell \in A_\ell^{-1}$ and, by Remark 3.3, we have $C^\ell K, KC^\ell \in A_\ell^{-2}$. We write

$$\int \psi(t) \left( (\Pi^+ A\Pi^\ell)(x, D)U_n(x, D) \right)_{L^2(\mathbf{R}^d)} dt = \int \psi(t) \left( (C^\ell, P)(x, D)U_n(x, D) \right)_{L^2(\mathbf{R}^d)} dt.$$

Besides,

$$[C^\ell, P](x, D) = [C^\ell(x, D), -i\partial_t + P(x, D)] - \frac{1}{2i} \left( \{C^\ell, P\} - \{P, C^\ell\} \right)(x, D) + r_1(x, D).$$

Using $P = P_1 + i b K$ and $\{C^\ell, P_1\}, \{P_1, C^\ell\} \in A_\ell^{-1}$, we obtain,

$$\{C^\ell, P\} - \{P, C^\ell\} = i\{C^\ell, b\}K - iK\{b, C^\ell\} + r_1.$$

Therefore, we have

$$[C^\ell, P](x, D) = [C^\ell(x, D), -i\partial_t + P(x, D)] - \frac{1}{2} \left( \{C^\ell, b\}K - K\{b, C^\ell\} \right)(x, D) + r_1(x, D).$$

where we have used $\{C^\ell, b\}K = \{C^\ell K, b\} \in A_\ell^{-1}$ and $K\{C^\ell, b\} = \{K C^\ell, b\} \in A_\ell^{-1}$ (as a consequence of $KC^\ell, C^\ell K \in A_\ell^{-2}$). In view of $P^s(x, D) = P(x, D) -
We work in space-time variables and set \( L = \mathbb{R}_t \times \Omega \). We denote by \((z, \zeta)\) the points of \( T^*L \): \( z = (t, x) \) and \( \zeta = (\tau, \xi) \). Then, the Melrose-Sjöstrand compressed bundle to \( L \) is given by

\[ T^*_bL = (T^*L \setminus \{0\}) \cup (T^*\partial L \setminus \{0\}). \]

Quotienting by the action of \( \mathbb{R}^+ \) through homotheties, one obtains the normal compressed bundle to \( \Omega \)

\[ S^*_bL = T^*_bL / \mathbb{R}^+. \]

The projection

\[ \pi : (T^*\mathbb{R}^{d+1})_L \to T^*_bL \]
induces a topology on \( T^*_bL \). On \( T^*_L \), we denote by \( p_0 \) the symbol of the wave operator and by \( \Sigma_0 \) the projection on \( T^*_bL \) of its characteristic set
\[
\rho(z, \zeta) = \tau^2 - |\xi|^2, \quad \Sigma_0 = \pi \{ (z, \zeta), \quad p_0(z, \zeta) = 0 \}.
\]
Locally, near a point of \( \partial L \), we use normal geodesic coordinates
\[
(y, \eta) = (y_0, y_1, \ldots, y_d; \eta_0, \eta_1, \ldots, \eta_d) \in \mathbb{R}^{2d+2}
\]
in an open set \( V \) of \( T^*_L \) such that \( L \cap V = \{ y_d > 0 \} \cap V \). Therefore, the wave operator \( \partial^2_t - \Delta \) writes \( -\partial^2_{y_d} - R(y, D_y) \) where the symbol \( R(y, \eta') \) is a homogeneous polynomial of degree 2 in \( \eta' \); we denote by \( r(y, \eta') \) its principal symbol. We can now distinguish between different sorts of points of \( T^*_L \setminus \{ 0 \} \): those who are not in \( \Sigma_0 \) and those who are in \( \Sigma_0 \) depending whether \( R(y, \eta') < 0 \) or not. In the case where \( \rho \in \Sigma_0 \), there exists at most two points of \( \{ \tau^2 = |\xi|^2 \} = \{ \eta_d^2 = R(y, \eta') \} \) which are in \( \pi^{-1}(\{ \rho \}) \), they correspond to the two roots of the equation \( \eta_d^2 = R(y, \eta') \).

Consider \( \rho \in T^*_L \setminus \{ 0 \} \).

- If \( \rho \notin \Sigma_0 \), one says that \( \rho \) is elliptic.
- If Card \( (\pi^{-1}(\{ \rho \}) \cap \{ \tau^2 = |\xi|^2 \}) = 2 \), \( \rho \) is said to be hyperbolic.
- If Card \( (\pi^{-1}(\{ \rho \}) \cap \{ \tau^2 = |\xi|^2 \}) = 1 \), \( \rho \) is said to be glancing.

We denote by \( \mathcal{H} \) (resp. \( \mathcal{G} \)) the hyperbolic (resp. glancing) points of \( \partial L \). We say that \( \rho \in \mathcal{G} \) is

- non strictly glancing if \( \partial_{y_d} r(\rho) \geq 0 \),
- strictly glancing if \( \partial_{y_d} r(\rho) < 0 \),
- diffractive if \( \partial_{y_d} r(\rho) > 0 \),
- glancing of order \( k \) if

\[
H^j_{r(y', 0, \eta')}(\partial_{y_d} r|_{y_d=0})(\rho) = 0, \quad 0 \leq j < k - 2 \quad \text{and} \quad H^{k-2}_{r(y', 0, \eta')}(\partial_{y_d} r|_{y_d=0})(\rho) \neq 0.
\]

We denote by \( \mathcal{G}_k \) (resp. \( \mathcal{G}_d \)) the set of points which are glancing of order \( k \) (resp. diffractive). The assumption that \( \Omega \) has no contact of infinite order with its tangents \((1)\) in Assumptions 2.6 consists in assuming
\[
(A.2) \quad \mathcal{G} = \bigcup_{k \in \mathbb{N}} \mathcal{G}_k.
\]

The generalized bicharacteristic are defined by taking the Hamiltonian trajectory of \( p_0 \) inside \( \Omega \) and by specifying how the connection is made on the boundary. The only problematic points are the glancing ones where the trajectory arrives tangentially to the boundary. Indeed, for \( \rho \in \mathcal{G} \), we have \( y_d = \eta_d = r(y, \eta') = 0 \). Recall that the geodesic trajectories are generated by the Hamiltonian flow \( H_{p_0} \) which writes in coordinates \((y, \eta)\)
\[
H_{p_0} = (-\nabla_y r(y, \eta'), 2\eta_d, \nabla_y r(y', \eta'), \partial_{y_d} r(y, \eta')).
\]
If \( \rho \in \mathcal{G}_d, \partial_{y_d} r(y, \eta') > 0 \), therefore \( \eta_d \) and thus \( y_d \) increase on the trajectory passing in \( \rho \) and this trajectory remains in \( \Omega \). On the contrary, if \( \rho \notin \mathcal{G}_d \), the coordinates on \( \partial \eta_d \) is nonpositive and the trajectory will leave \( \Omega \): if \( \eta_d \) decreases, it becomes non positive and so for \( y_d \). To overcome this difficulty, one uses the vector
\[
H(\rho) = H_{p_0}(\rho) - \partial_{y_d} r(y, \eta') \partial_{\eta_d}.
\]
This vector has a coordinate on \( \partial_{\eta_d} \) which is 0. One then defines the generalized bicharacteristic as follows (see [19] or [3]).

**Definition 1.** A generalized bicharacteristic is a continuous map \( \gamma : \mathbb{R} \to T^*_bL \) such that there exists a set \( I \) of isolated points with
\begin{itemize}
  \item \( \gamma(s) \in T^*L \cup \mathcal{G} \) for \( s \notin I \) and \( \gamma(s) \in \mathcal{H} \) for \( s \in I \).
  \item for \( s \notin I \), \( \gamma \) is differentiable with
  \[
  \begin{cases}
    \gamma'(s) = H_{p_0}(\gamma(s)) & \text{if } \gamma(s) \in T^*L \cup \mathcal{G}_d, \\
    \gamma'(s) = \tilde{H}(\rho) & \text{if } \gamma(s) \in \mathcal{G} \setminus \mathcal{G}_d.
  \end{cases}
  \]
\end{itemize}

It is proved in [19] that these definitions are intrinsic and that if assumption \( (A.2) \) is satisfied, then for \( \rho_0 \in T^*_p L \cap \Sigma_0 \), there exists a unique generalized bicaracteristic such that \( \gamma(0) = \rho_0 \).

\textbf{A.2. Propagation near the boundary.} Working in space-time variables, one first defines tangential symbols by use of the system of local normal geodesic coordinates: in an open set \( \mathcal{O} \) where we have such coordinates \((y, \eta)\), the function \( a(y, \eta) \in C^\infty(\overline{L} \times \mathbb{R}^d) \) is said to be a symbol of \( \mathcal{A}_w^m \) if \( a \) is compactly supported in \( \mathcal{O} \) in the variable \( y \) and satisfies: \( \forall \alpha, \beta \in \mathbb{N}, \exists C_{\alpha, \beta} > 0, \)
\[
  (A.3) \quad \forall (y, \eta) \in \overline{L} \times \mathbb{R}^{d+1}, \quad \left| \partial_y^a \partial_\eta^b (a(y, \eta)) \right| \leq C_{\alpha, \beta}(1 + |\eta'|)^{m-|\beta|}
\]
This definition implies that the sets \( \mathcal{A}_w^m \) depend on the choice of the open sets \( (\mathcal{O}_j)_{j \in J} \). Then, one defines
\[
  \mathcal{A}^m = \{ a \in C^\infty(\overline{\Omega} \times \mathbb{R}^d), \quad \exists a^1 \in \mathcal{A}_w^m, \quad \exists a^b \in \mathcal{A}_w^m, \quad a = a^b + a^i \},
\]
and one associates with our sequence \( u_n(t, x) \) its \( H^1 \) space-time microlocal defect measure \( \mu \) by: up to extraction of a subsequence, for all \( q = q_0 + q_1 \in \mathcal{A}^2 \),
\[
  (q(z, D_z)u_n, u_n) \underset{n \to +\infty}{\longrightarrow} \int_{\mathbb{T} \times S^d(2)} q_{\infty}(z, \zeta) d\mu(z, \zeta)
\]
where \( S^d(2) \) is the sphere of radius 2 of \( \mathbb{R}^{d+1} \) and for \( (z, \zeta) \in S^d(2), \)
\[
  q_{\infty}(z, \zeta) = \lim_{R \to +\infty} R^{-2} q(z, R\zeta).
\]

Note that on \( S^*\Omega \), the measure \( \mu \) is the usual microlocal defect measure (or H-measure) as introduced in [12] and [20]. The link between this measure \( \mu \) and the measures \( \mu^\pm \) of Theorem 2.4 is described in the following Proposition.

\textbf{Proposition A.1.} In \( \mathbb{R} \times \Omega_1 \times S^d(2) \), we have
\[
  (A.4) \quad \mu(t, x, \tau, \xi) = \frac{1}{2} \delta(\tau + 1) \otimes \mu^+(x, \xi) \otimes dt + \frac{1}{2} \delta(\tau - 1) \otimes \mu^-(x, \xi) \otimes dt.
\]

This proposition is proved in section \( A.3 \) below. The next step consists in analyzing \( \mu \) on the boundary. We recall that if \( \rho \in T^*\partial L \) is an hyperbolic or a glancing point, then there exists two points \( \rho^\pm \) such that \( \pi(\rho^\pm) = \rho \) and \( \rho^\pm \in \{ \tau^2 = |\xi|^2 \} \).

These two points are equal if \( \rho \in \mathcal{G} \). When \( \rho \in \mathcal{H} \), they differ by their \( \xi \)-component that we will denote by \( \xi^\pm \). Besides, as in [4], because of the equation satisfied by \( (u_n)_{n \in \mathbb{N}} \), the sequence of the normal derivatives \( (\partial_N u_n)_{n \in \mathbb{N}} \) is a uniformly bounded family of \( L^2_{\text{loc}}(\mathbb{R}, L^2(\partial \Omega)) \) and we denote by \( \lambda \) its microlocal defect measure (we suppose that we have extracted a subsequence so that \( \lambda \) is uniquely determined). Following [4], one aims at calculating
\[
  \ell := H_{p_0} \mu - 4\tau \alpha_{\infty} \mu.
\]
We already know that if $q$ is an interior symbol, $(q, \ell) = 0$. For analyzing the action of $\ell$ on tangential symbols, one computes the quantity

$$I_n(t) = \left\{ \left[ \partial_{y_\delta}^2 + R(y, D_{y'})q_b(y, D_{y'}) \right] u_n(t), u_n(t) \right\}$$

for $q_b$ a tangential symbol of order 1.

The computation of $I_n$ and the passage $n \to +\infty$ arises terms on the boundary which are the same than in [4], namely a distribution $\ell_1 + \ell_2$ where $\ell_1$ is a measure absolutely continuous with respect to $\mu$ and where, denoting by $N(x)$ the exterior normal vector to $\partial \Omega$, $\ell_2$ satisfies

$$\ell_2 = \frac{\delta(\xi - \xi^+) - \delta(\xi - \xi^-)}{(\xi_+ - \xi_-) \cdot N(x)} \times \mathbf{1}_{H \cup G}.$$ 

The reader can refer to [3] (pages 14-15) and [4] where the computations are carefully carried out.

On the other hand, if one uses the equation to transform $I_n$, there appears a new term which was not in the preceding articles. We are going to discuss why these terms are harmful when one has (3) of Assumptions 2.6. This term involves the quantity $\nabla \cdot (\xi)$. The fact that $\xi$ is tangent to the boundary implies that

$$\tilde{\gamma}(y', 0) = (\tilde{\gamma}', y \tilde{\gamma}_d(0), 0);$$

therefore, we can write $\tilde{\gamma} = (\tilde{\gamma}', y_d \tilde{\gamma}_d)$. Then, the worst term to estimate – which is the one which involves $\partial_{y_d}$ derivatives of $\theta_n$ – writes

$$\tilde{J}_{n, \delta}(t) = \left( \chi \left( \frac{y_d}{\delta} \right) q_b(y, D_{y'}) u_n, \tilde{\gamma}(y) \cdot \nabla \theta_n \right);$$


and we observe

$$\left| \tilde{J}_{n, \delta}(t) \right| = \left| \int_0^{+\infty} \chi \left( \frac{y_d}{\delta} \right) q_b(y, D_{y'}) y_d \tilde{\gamma}_d(y) \nabla \theta_n(t), u_n(t) \right| \leq \delta \left\| u_n(t) \right\|_{L^2(\Omega)} \left\| \nabla \theta_n(t) \right\|_{L^2(\Omega)}.$$ 

Therefore, for any $\Theta \in C_0^\infty(R)$, there exists a constant $C$

$$\limsup_{n \to +\infty} \int \Theta(t) J_{n, \delta}(t) dt \leq C \delta.$$ 

Finally, letting $\delta$ go to 0, we obtain that this term has no contribution on the boundary. As a conclusion, we obtain $\ell = \ell_1 + \ell_2$ and we can conclude like in [4].
A.3. The link between the measure $\mu$ and the measures $\mu^\pm$. In this section, we prove Proposition A.1.

Note first that the second assertion is a simple consequence of the first one. Suppose that we have A.4 where, by Proposition A.2, the measures $\mu^\pm$ satisfy
$$\partial_t \mu^\pm + \xi \cdot \nabla_x \mu^\pm = 2\alpha \mu^\pm, \quad \text{in } \Omega_1 \times S^{d-1}.$$ We obtain
$$H_{\partial_\nu} \mu = 2(\tau \partial_\nu \mu - \xi \cdot \nabla_x \mu)$$
$$= (-\partial_t \mu^+ - \xi \cdot \nabla_x \mu^+) \otimes \delta(\tau + 1) \otimes dt$$
$$+ (\partial_t \mu^- - \xi \cdot \nabla_x \mu^-) \otimes \delta(\tau - 1) \otimes dt$$
$$= -2\alpha \mu^+(x, \xi) \otimes \delta(\tau + 1) \otimes dt + 2\alpha \mu^-(x, \xi) \otimes \delta(\tau - 1) \otimes dt$$
$$= 4\alpha \mu.$$

Let us now prove A.4. We take $q \in A_0^0$ and apply $q(z, D_z)$ to the first equation of (1.3). We get
$$0 = (q(z, D_z)(\partial_t^2 u_n - \Delta u_n + \nabla \cdot (\gamma(x) \theta_n)), u_n).$$
Using that $(\theta_n)_{n \in \mathbb{N}}$ is u.b. in $L^2_{loc}(\mathbb{R}, H^1(\Omega))$ and that $u_n$ goes to 0 weakly in $H^1(\Omega)$, we obtain passing to the limit
$$\int_{\mathbb{R} \times \Omega_1 \times S^d} q(z, \zeta)(|\xi|^2 - \tau^2) d\mu(z, \zeta) = 0.$$ Therefore, on the support of $\mu$, we have $\tau^2 = |\xi|^2$. Since moreover
$$|\xi|^2 = \tau^2 + |\gamma|^2 = 2 \text{ on } S^d,$$
we obtain that $\tau^2 = |\xi|^2 = 1$ on the support of $\mu$, whence the existence of two measures $\tilde{\mu}^\pm$ on $\mathbb{R} \times \Omega_1 \times S^{d-1}$ such that
$$\mu(t, x, \tau, \xi) = \delta(\tau + 1) \otimes \tilde{\mu}^+(t, x, \xi) + \delta(\tau - 1) \otimes \tilde{\mu}^-(t, x, \xi) \quad \text{in } \mathcal{D}'(\mathbb{R} \times \Omega_1 \times S^d).$$

Let us now link the measure $\tilde{\mu}^\pm$ with the measures $\mu^\pm$ of Section 4.2 inside $\Omega$. Let us consider as in Section 4.2
$$f_{\pm, n}(t, x) = \frac{1}{\sqrt{2}}(i\partial_t u_n(t, x) \pm |D_x|u_n(t, x)).$$
These families are uniformly bounded in $L^2(\Omega)$ for all $t \in \mathbb{R}$ and their microlocal defect measures are the measures $\mu^\pm$. Besides, by the definition of $\mu$, for $q \in A_0^0$, we have
$$\int_{\mathbb{R} \times \Omega \times S^d} \mu \chi(t, x, \tau, \xi) d\mu^\pm(t, x, \xi) = \int_{\mathbb{R} \times \Omega \times S^d} \chi(t, x, \tau, \xi) d\mu^\pm(t, x, \xi).$$
where we have used (A.5) for the last equality. Let us choose $q(z, \zeta) = \chi(t) b(x, \xi)$ with $b(x, \xi)$ symbol of order 0 compactly supported in $\Omega$, we obtain in the one hand
$$\int_{\mathbb{R} \times \Omega \times S^d} \mu \chi(t, x, \tau, \xi) d\mu^\pm(t, x, \xi) = \int_{\mathbb{R} \times \Omega \times S^d} \chi(t) b(x, \xi) d\mu^\pm(t, x, \xi).$$
In the other hand, we have for all $t \in \mathbb{R}$,

$$
(b(x, D)f_{\pm,n}(t), f_{\pm,n}) \rightarrow_{n \rightarrow +\infty} \int_{\Omega \times S^{d-1}} q(x, \omega)d\mu_{1}^{\pm}(x, \omega).
$$

Whence $\tilde{\mu}_{\pm}(t, x, \omega) = \frac{1}{2}\mu_{1}^{\pm}(x, \omega) \otimes dt$.

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(A. Atallah) Dpt. de Mathématiques, Faculté des Sciences de Tunis, Université de Tunis El Manar, 1060, Tunis, Tunisie.
E-mail address: amel.atallah@fst.rnu.tn

(C. Fermanian) LAMA UMR CNRS 8050, Université Paris Est - Créteil, 61, avenue du Général de Gaulle, 94010 Créteil Cedex, France.
E-mail address: Clotilde.Fermanian@u-pec.fr