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**Existence and uniqueness theorems
in the linear magnetohydrodynamics
with dissipative boundary conditions (**)**

1 - Introduction

This paper presents some *existence and uniqueness theorems* for a linear differential problem which characterizes the evolution of a magnetohydrodynamics system, without neglecting the current displacement and the separation of charges. Such a system is usually called a *plasma* [1], [3], [4], [7]. In particular, a condition on the boundary characterizing a large class of *dissipative boundary conditions* is also considered.

Furthermore, being the domain Ω unbounded [2], the existence and uniqueness are studied both in the case of a *finite energy* all over Ω and in the case of a *locally finite energy*. For this purpose, we will use a theorem on the domain dependence, which guarantees a finite speed of propagation and a theorem of continuous dependence on data.

The method, used in the following, refers to [5], [6], [8], [9], which are related to linear symmetric hyperbolic systems considered in this paper.

2 - Problem statement

Let Ω be a domain of the ordinary three dimensional space \mathbf{R}^3 , and $x = (x_1, x_2, x_3)$ an arbitrary point of \mathbf{R}^3 ; I is a finite interval of a timelike variable t . Vectors \mathbf{E} , \mathbf{H} , \mathbf{v} represent the *electric field*, the *magnetic field* and the

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speed of electrons respectively. The pressure of the electrons is characterized by a scalar function p .

This set of variables defines the electromagnetic state of a *magnetohydrodynamics system (plasma)* by means of the function $\mathcal{C} = [\mathbf{H}, \mathbf{E}, \mathbf{v}, p]$, which is related both to a supply of magnetic current \mathbf{K} , electric current \mathbf{J} , external body force \mathbf{F} , and to a supply of flow Φ .

Furthermore, on the plasma is applied a magnetic induction field denoted by \mathbf{B}_0 , while n_0 and v_0 will characterize the mean density and mean speed of electrons respectively.

For a given plasma, we will consider a differential problem with initial-boundary conditions, which satisfies the following system:

$$(1) \quad \mu_0 \frac{\partial \mathbf{H}}{\partial t} = -\nabla \times \mathbf{E} - \mathbf{K}$$

$$(2) \quad \varepsilon_0 \frac{\partial \mathbf{E}}{\partial t} = \nabla \times \mathbf{H} + en_0 \mathbf{v} - \mathbf{J} \quad Q = \Omega \times I$$

$$(3) \quad mn_0 \frac{\partial \mathbf{v}}{\partial t} = -\nabla p - en_0(\mathbf{E} + \mathbf{B}_0 \times \mathbf{v}) + \mathbf{F} \quad Q = \Omega \times I$$

$$(4) \quad \frac{1}{n_0 m v_0^2} \frac{\partial p}{\partial t} = -\nabla \cdot \mathbf{v} + \Phi$$

e being the charge of the electron whereas ε_0 and μ_0 represent the dielectric constant and the magnetic permeability, in vacuum, which are ruled by the condition $\frac{1}{\varepsilon_0 \mu_0} = c^2 > v_0^2$. The initial conditions are

$$(5) \quad \begin{aligned} \mathbf{H}(x, 0) &= \mathbf{h}(x) & \mathbf{E}(x, 0) &= \mathbf{e}(x) & \text{in } \Omega \\ \mathbf{v}(x, 0) &= \mathbf{v}(x) & p(x, 0) &= P(x) & \text{in } \Omega \end{aligned}$$

where the vectors $\mathbf{h}(x)$, $\mathbf{e}(x)$, $\mathbf{v}(x)$ and the scalar function $P(x)$ are assumed to be known in Ω , so that the initial state $\mathcal{C}^0 = [\mathbf{h}(x), \mathbf{e}(x), \mathbf{v}(x), P(x)]$ is given.

Furthermore the function \mathcal{C} satisfies a linear and homogeneous boundary conditions Γ such that the inequalities:

$$(6) \quad \mathbf{E} \times \mathbf{H} \cdot \mathbf{n} \leq 0 \quad \text{in } \partial\Omega \times I$$

$$(7) \quad p\mathbf{v} \cdot \mathbf{n} \leq 0 \quad \text{in } \partial\Omega \times I$$

are satisfied on the boundary $\partial\Omega$, being \mathbf{n} the outward normal vector. The field source is represented by the set of functions $\mathcal{F}(x, t) = [-\mathbf{K}, -\mathbf{J}, \mathbf{F}, \Phi]$.

Together with the differential system (1)-(4) we must consider, for the plasma state \mathcal{C} , the *rate of the energy* in the domain $\Omega \cap K$ at time t , i.e.

$$(8) \quad \frac{1}{2} \int_{\Omega \cap K} \frac{\partial}{\partial t} (\mu_0 \mathbf{H}^2 + \varepsilon_0 \mathbf{E}^2 + n_0 m \mathbf{v}^2 + \frac{p^2}{n_0 m v_0^2}) dx$$

for every cube $K \subset \mathbf{R}^3$.

If the integral (8) converges to a finite value for any measurable set K , then \mathcal{C} is called a *plasma state with finite energy*. If the integral (8) converges only for a bounded measurable set K then \mathcal{C} is called a *plasma state with locally finite energy*.

From a differential point of view, system (1)-(4) is an hyperbolic system of differential equations of the type:

$$(9) \quad E(x) \frac{\partial \mathbf{u}}{\partial t} = \mathcal{C}^i(x) \frac{\partial \mathbf{u}}{\partial x_i} + B(x) \mathbf{u} + \mathbf{f}(x, t) \quad i = 1, \dots, n$$

where the $m \times m$ matrices

$$E = (E_{\alpha\beta}) \quad \mathcal{C}^i = (A_{\alpha\beta}^i) \quad B = (B_{\alpha\beta}) \quad \alpha, \beta = 1, \dots, m$$

and the m vectors $\mathbf{u} = (u_\alpha)$, $\mathbf{f} = (f_\alpha)$ ($\alpha = 1, \dots, m$), are functions of $x = (x_\alpha)$ and t .

When E and \mathcal{C}^i , ($i = 1, \dots, n$) are symmetric matrices and E is definitely positive then system (9) is an hyperbolic symmetric system. Thus system (1)-(4) is an *hyperbolic symmetric system* like (9), assuming $n = 3$, $m = 10$ and with $\mathbf{u} = \mathcal{C}$, $\mathbf{f} = \mathcal{F}$.

The differential operator of the first order A acting on the class \mathcal{C}^1 of the continuous functions with continuous first derivative is defined as

$$A\mathcal{C} = \mathcal{C}^i \frac{\partial \mathcal{C}}{\partial x_i} \quad i = 1, 2, 3$$

and verifies the condition on the formal adjoint $A^* = A$.

3 - Preliminary definitions

In this section we give preliminary definitions and recall standard notations [8].

Let us consider only functions $u: D \rightarrow H$, defined in an arbitrary domain $D \subset \mathbf{R}^n$ and values into an Hilbert separable space H , where the scalar product is $(f, g)_H$ and the norm of a function is $\|f\|_H = (f, f)_H^{\frac{1}{2}}$. The space

$$L_2(D, H) = \{u \mid u \text{ is measurable in } H, \int_D \|u\|_H^2 dx < \infty\}$$

is an *Hilbert space*; moreover

$$L_2^{\text{loc}}(D, H) = \{u \mid u \in L_2(K \cap D, H) \quad \forall \text{ bounded measurable set } K \subset \mathbf{R}^n\}$$

is the *space of locally square-integrable functions*. Let $A: H \rightarrow H'$ be a linear differential operator with bounded and measurable coefficients and $A^+: H' \rightarrow H$ its formal adjoint. Let us now consider the linear subspaces:

$$L_2(A, D, H) = \{u \in L_2(D, H) \mid Au \in L_2(D, H')\}$$

$$L_2^{\text{loc}}(A, D, H) = \{u \in L_2^{\text{loc}}(D, H) \mid Au \in L_2^{\text{loc}}(D, H')\}$$

$$L_2^{\text{vox}}(A, D, H) = L_2(A, D, H) \cap \{u \mid u = 0 \text{ outside a bounded set}\}.$$

In particular, if $\mathbf{A}, \mathbf{B} \in L_2(\nabla \times, \Omega, \mathbf{R}^3)$ (where $\nabla \times$ is the *curl operator*) then

$$(10) \quad \int_{\Omega} (\mathbf{A} \cdot \nabla \times \mathbf{B} - \mathbf{B} \cdot \nabla \times \mathbf{A}) \, dx = \int_{\partial\Omega} \mathbf{A} \times \mathbf{B} \cdot \mathbf{n} \, ds$$

there follows that $\nabla \times$ is formally selfadjoint, i.e.

$$(11) \quad \int_{\Omega} \mathbf{A} \cdot \nabla \times \boldsymbol{\Psi} \, dx = \int_{\Omega} \nabla \times \mathbf{A} \cdot \boldsymbol{\Psi} \, dx$$

for any $\mathbf{A} \in L_2(\nabla \times, \Omega, \mathbf{R}^3)$ and $\boldsymbol{\Psi} \in C_1^0(\Omega, \mathbf{R}^3)$.

In fact, when \mathbf{A} fulfills the boundary condition $\mathbf{A} \times \mathbf{n} = 0$ in $\partial\Omega$ (inequality (6) is satisfied when $\mathbf{A} = \mathbf{E}$ or $\mathbf{A} = \mathbf{H}$), then equation (10) implies (11) for all \mathbf{B} . As consequence one is led to consider the space

$$L_2^0(\nabla \times, \Omega, \mathbf{R}^3) = \{\mathbf{A} \mid \mathbf{A} \in L_2(\nabla \times, \Omega, \mathbf{R}^3), \int_{\Omega} \mathbf{A} \cdot \nabla \times \mathbf{B} \, dx = \int_{\Omega} \mathbf{B} \cdot \nabla \times \mathbf{A} \, dx, \forall \mathbf{B} \in L_2(\nabla \times, \Omega, \mathbf{R}^3)\}$$

as a generalization of the class of fields satisfying the boundary condition

$$\mathbf{A} \times \mathbf{n} = 0 \quad \text{in } \partial\Omega.$$

The space $L_2^0(\nabla \times, \Omega, \mathbf{R}^3)$ is obviously a linear closed subspace of $L_2(\nabla \times, \Omega, \mathbf{R}^3)$. Analogously for the spaces $L_2(\nabla \cdot, \Omega, \mathbf{R}^3)$, $L_2(\nabla, \Omega, \mathbf{R}^3)$, according to the formula

$$(12) \quad \int_{\Omega} [(\nabla \cdot \mathbf{v})p + \nabla p \cdot \mathbf{v}] \, dx = \int_{\partial\Omega} p\mathbf{v} \cdot \mathbf{n} \, ds$$

the boundary condition $p\mathbf{v} \cdot \mathbf{n} = 0$ is satisfied in a generalized sense by the class of spaces:

$$L_2^0(\nabla \cdot, \Omega, \mathbf{R}^3) = \{\mathbf{v} \in L_2(\nabla \cdot, \Omega, \mathbf{R}^3) \mid \int_{\Omega} (\phi \nabla \cdot \mathbf{v} + \nabla \phi \cdot \mathbf{v}) \, dx = 0, \forall \phi \in L_2(\nabla, \Omega, \mathbf{R}^3)\}$$

$$L_2^0(\nabla, \Omega, \mathbf{R}^3) = \{p \in L_2(\nabla, \Omega, \mathbf{R}^3) \mid \int_{\Omega} (p \nabla \cdot \boldsymbol{\Phi} + \nabla p \cdot \boldsymbol{\Phi}) \, dx = 0, \forall \boldsymbol{\Phi} \in L_2(\nabla \cdot, \Omega, \mathbf{R}^3)\}.$$

In the following we shall assume $I = \{t \mid 0 < t < T\}$ and $H^1(I, H) = L_2(\frac{\partial}{\partial t}, I, H)$.

Let $\Gamma = \Gamma_1 \times \Gamma_2 \times \Gamma_3 \times \Gamma_4$ be a closed linear subspace of $L_2(A, \Omega, \mathbf{R}^{10})$, where $\Gamma_1, \Gamma_2, \Gamma_3, \Gamma_4$ are closed linear subspaces of $L_2(\nabla \times, \Omega, \mathbf{R}^3)$, $L_2(\nabla \times, \Omega, \mathbf{R}^3)$, $L_2(\nabla \cdot, \Omega, \mathbf{R}^3)$, $L_2(\nabla, \Omega, \mathbf{R})$. According to these definitions, we have

$$\mathbf{H} \in \Gamma_1 \quad \mathbf{E} \in \Gamma_2 \quad \mathbf{v} \in \Gamma_3 \quad p \in \Gamma_4.$$

The boundary conditions (6), (7) have to be generalized in order to include the dissipative condition derived from the identities (10), (11), i.e.:

$$\int_{\Omega} (\mathbf{E} \cdot \nabla \times \mathbf{H} - \mathbf{H} \cdot \nabla \times \mathbf{E}) dx \leq 0 \quad \forall \mathbf{H} \in \Gamma_1, \mathbf{E} \in \Gamma_2,$$

$$\int_{\Omega} (p \nabla \cdot \mathbf{v} + \mathbf{v} \cdot \nabla p) dx \leq 0 \quad \forall \mathbf{v} \in \Gamma_3, p \in \Gamma_4.$$

Together with the condition that Φ be locally dissipative we must have $\Phi \mathcal{C} \in \Gamma$ whenever $\Phi \in C_0^1(\mathbf{R}^3)$ and $\mathcal{C} \in \Gamma$. We assume also

$$L_2^0(\nabla \times, \Omega, \mathbf{R}^3) \subset \Gamma_1; L_2^0(\nabla \times, \Omega, \mathbf{R}^3) \subset \Gamma_2; L_2^0(\nabla \cdot, \Omega, \mathbf{R}^3) \subset \Gamma_3; L_2^0(\nabla, \Omega, \mathbf{R}) \subset \Gamma_4$$

so that any boundary condition Γ implies a further adjoint boundary condition Γ^* defined as

$$\Gamma^* = L_2(A, \Omega, \mathbf{R}^{10}) \cap \{ \mathcal{X} : \mathcal{X} \in L_2(A, \Omega, \mathbf{R}^{10}), \int_{\Omega} [\langle \mathcal{X}, A \mathcal{X} \rangle + \langle A \mathcal{X}, \mathcal{X} \rangle] dx = 0 \quad \forall \mathcal{C} \in \Gamma \}.$$

4 - Solutions with finite energy

In this section we define the solution with finite energy (FE-solution) and the solution with locally finite energy (LFE-solution) that satisfy given initial-boundary conditions.

Since we have that any linear homogeneous locally dissipative boundary condition Γ is made of a closed linear subspace of $L_2(A, \Omega, \mathbf{R}^{10})$, Γ is a *separable Hilbert space* with respect to the scalar product of $L_2(A, \Omega, \mathbf{R}^{10})$. The space

$$F = L_2(I, \Gamma) \cap H^1(I, L_2(A, \Omega, \mathbf{R}^{10}))$$

denotes the class of functions $\mathcal{C}(x, t)$ for which $E \frac{\partial \mathcal{C}}{\partial t}$, $A \mathcal{C}$, $B \mathcal{C}$ exist in $L_2(Q; \mathbf{R}^{10})$, and satisfies the boundary condition Γ , in the sense that $\mathcal{C}(t) \in \Gamma$ for almost all $t \in I$. Let us define the following sets:

$$\Gamma^{\text{vox}} = \Gamma \cap \{ \mathcal{C} \mid \mathcal{C} = 0 \text{ outside a bounded set } K \subset \Omega \}$$

$$\Gamma^{\text{loc}} = \Gamma_2^{\text{loc}}(A, \Omega) \cap \{ \mathcal{C} \mid \int_{\Omega} (\langle \mathcal{C}, A \mathcal{X} \rangle + \langle A \mathcal{C}, \mathcal{X} \rangle) dx = 0 \} \quad \forall \mathcal{X} \in (\Gamma^*)^{\text{vox}}$$

$$F^{\text{loc}} = \{ \mathcal{C} \mid \mathcal{C} \in L_2(I, L_2(A, K \cap \Omega)) \cap H^1(I, L_2(A, K \cap \Omega)) \}$$

for any bounded measurable set $K \subset \mathbf{R}^3$ and $\mathcal{C}(t) \in \Gamma^{\text{loc}} \forall t \in I$

$$F^* = L_2(I, \Gamma^*) \cap H^1((I, L_2 \Omega))$$

$$(F^*)^{\text{vox}} = F^* \cap \{\mathcal{X} : \mathcal{X} = 0 \text{ outside } K \times I \subset Q\}.$$

According to the above definitions we can distinguish the following kind of solutions:

Definition 1. \mathcal{C} is a *FE-solution* of system (1)-(4) with boundary condition Γ and $\mathcal{F}(x, t) \in L_2(Q; \mathbf{R}^{10})$, $\mathcal{C}^0(x) \in L_2(\Omega)$ given, iff $\mathcal{C}(t) \in F$ satisfies

$$(13) \quad \begin{aligned} E \frac{\partial \mathcal{C}}{\partial t} &= A\mathcal{C} + B\mathcal{C} + \mathcal{F} && \text{almost everywhere in } Q \\ \mathcal{C}(0) &= \mathcal{C}^0 && \text{almost everywhere in } \Omega. \end{aligned}$$

Definition 2. \mathcal{C} is a *weak FE-solution* of system (1)-(4) with boundary condition Γ and $\mathcal{F}(x, t) \in L_2(Q; \mathbf{R}^{10})$, $\mathcal{C}^0(x) \in L_2(\Omega)$ given, iff $\mathcal{C}(t) \in L_2(Q; \mathbf{R}^{10})$ satisfies

$$(14) \quad \begin{aligned} &\int_Q \langle E \frac{\partial \mathcal{X}}{\partial t} - A\mathcal{X} - B\mathcal{X}, \mathcal{C} \rangle + \langle \mathcal{F}, \mathcal{X} \rangle dx dt \\ &- \int_{\Omega} \langle E(x) \mathcal{X}(x, T), \mathcal{C}(x, T) \rangle dx + \int_{\Omega} \langle E(x) \mathcal{X}(x, 0), \mathcal{C}^0(x) \rangle dx = 0 \end{aligned}$$

for any $\mathcal{X} \in F^k$.

Definition 3. \mathcal{C} is a *LFE-solution* of system (1)-(4) with boundary condition Γ and $\mathcal{F}(x, t) \in L_2^{\text{loc}}(Q, \mathbf{R}^{10})$, $\mathcal{C}^0(x) \in L_2^{\text{loc}}(\Omega)$ given, iff $\mathcal{C}(t) \in F^{\text{loc}}$ satisfies system (13).

Definition 4. \mathcal{C} is a *weak LFE-solution* of system (1)-(4) with boundary condition Γ and $\mathcal{F}(x, t) \in L_2^{\text{loc}}(Q, \mathbf{R}^{10})$, $\mathcal{C}^0(x) \in L_2^{\text{loc}}(\Omega)$ given, iff $\mathcal{C}(t) \in L_2^{\text{loc}}(Q, \mathbf{R}^{10})$ satisfies system (14) $\forall \mathcal{X} \in (F^*)^{\text{vox}}$.

Following Wilcox methods [8], we give now some theorems on the existence, uniqueness and domain of dependence.

Theorem 1 (Existence and uniqueness). *Let Γ be a locally dissipative boundary condition for the system (1)-(4) and Γ^* its adjoint boundary condition, then in the initial-boundary problem of a plasma propagation there exists a unique:*

a. weak LFE-solution $\forall \mathcal{F} \in L_2^{\text{loc}}(Q, \mathbf{R}^{10})$ and $\mathcal{C}^0 \in L_2^{\text{loc}}(\Omega, \mathbf{R}^{10})$

b. weak FE-solution $\forall \mathcal{F} \in L_2(Q, \mathbf{R}^{10})$ and $\mathcal{C}^0 \in L_2(\Omega, \mathbf{R}^{10})$

c. FE-solution $\forall \mathcal{F} \in H^1(I, L_2(\Omega))$ and $\mathcal{C}^0 \in \Gamma$

d. LFE-solution $\forall \mathcal{F} \in H^1(I, L_2(\Omega))$ where K is any arbitrary bounded set $K \subset \mathbf{R}^3$ and $\mathcal{C}^0 \in \Gamma^{\text{loc}}$.

To prove this theorem, we have to show first the

Theorem 2 (Domain of dependence). *The LFE-solutions (as well as the FE-solutions) with locally dissipative boundary conditions fulfill the following inequality*

$$(15) \quad \int_{\Omega \cap S(x^0, a)} \langle E \mathcal{C}(T), \mathcal{C}(T) \rangle dx \\ \leq e^{kT} \left(\int_{\Omega \cap S(x^0, a+cT)} \langle E \mathcal{C}(0), \mathcal{C}(0) \rangle dx + 2 \int_{Q \cap C(x^0, a)} e^{-kt} \langle \bar{\mathcal{F}}, \mathcal{C} \rangle dx dt \right)$$

where: $S(x^0, a) = \{x: |x - x^0| \leq a\}$

$C(x^0, a) = \{(x, t): |x - x^0| \leq a + c(T - t), 0 \leq t \leq T\}$.

Proof. The proof makes use of an auxiliary function $\mu(x, t)$ defined as

$$\mu(x, t) = \mu_\delta(\tau) \quad \text{where} \quad \tau(x, t) = a - |x - x^0| + (T - t)$$

where the real-valued function μ_δ has the following properties:

$$(16) \quad \mu_\delta \in C^1(-\infty, \infty) \quad \mu'_\delta(\tau) = \frac{d\mu}{d\tau} \geq 0 \\ \mu_\delta(\tau) = 0 \quad \text{for } |\tau| \geq \delta \quad 0 \leq \mu_\delta(\tau) \leq 1 \quad \text{for } |\tau| \leq \delta.$$

According to this definition $\mu(x, t) \in C^1(Q \cap C(x^0, a))$ and $\mu(x, t) = 0$ outside $Q \cap C(x^0, a + \delta)$. We might also consider, without restrictions, that $k = 0$, so that from equation (14), assuming $\mathcal{X} = \mu \mathcal{C}$, we have

$$\int_{Q \cap C(x^0, a)} \left[\langle E \frac{\partial \mathcal{C}}{\partial t} - A \mathcal{C} - B \mathcal{C}, \mathcal{C} \rangle \mu + \langle \bar{\mathcal{F}}, \mathcal{C} \rangle \mu + \left\langle \left(E \frac{\partial \mu}{\partial t} - A \mu \right) \mathcal{C}, \mathcal{C} \right\rangle \right] dx dt \\ = \int_{Q \cap S(x^0, a)} \left[\langle E(x) \mathcal{C}(x, T), \mathcal{C}(x, T) \rangle - \langle E(x) \mathcal{C}(x, 0), \mathcal{C}^0(x) \rangle \right] \mu dx = 0.$$

and, according to (13) we have

$$\begin{aligned} & 2 \int_{Q \cap C(x^0, a)} \langle \mathcal{F}, \mathcal{C} \rangle \mu \, dx \, dt + \int_{Q \cap C(x^0, a)} \langle \langle E \frac{\partial \mu}{\partial t} - A\mu \rangle \mathcal{C}, \mathcal{C} \rangle \, dx \, dt \\ = & \int_{Q \cap S(x^0, a)} \langle E(x) \mathcal{C}(x, T), \mathcal{C}(x, T) \rangle \mu \, dx - \int_{Q \cap S(x^0, a)} \langle E(x) \mathcal{C}(x, 0), \mathcal{C}^0(x) \rangle \mu \, dx = 0. \end{aligned}$$

Taking into account the definition of μ and the value of its derivatives ($\frac{\partial \mu}{\partial t} = -\mu'_\delta(\tau)$, $|A\mu| = \mu'_\delta(\tau)|A\tau| = \mu'_\delta(\tau)$), it results

$$\begin{aligned} & 2 \int_{Q \cap C(x^0, a)} \langle \mathcal{F}, \mathcal{C} \rangle \mu \, dx \, dt = \int_{\Omega \cap S(x^0, a)} \langle \langle E(x) \mathcal{C}(x, T), \mathcal{C}(x, T) \rangle \mu \, dx \\ & - \int_{Q \cap S(x^0, a)} \langle E(x) \mathcal{C}(x, 0), \mathcal{C}^0(x) \rangle \mu \, dx - \int_{Q \cap C(x^0, a)} (\langle E \mathcal{C}, \mathcal{C} \rangle + \langle \mathcal{C}, \mathcal{C} \rangle) \mu'_\delta \, dx \, dt \end{aligned}$$

from where, as a consequence of the definition of μ and of the properties of the last integral (≥ 0), the inequality (15) easily follows.

We now give a result for the FE-solutions.

Theorem 3. *Let \mathcal{C} be a FE-solution with arbitrary given values $(\mathcal{F}, \mathcal{C}^0)$, then for any $t \in I$, we have*

$$(17) \quad \int_{\Omega} \langle E \mathcal{C}(T), \mathcal{C}(T) \rangle \, dx \leq C \left(\int_{\Omega} \langle E \mathcal{C}^0, \mathcal{C}^0 \rangle \, dx + \int_Q \langle E^{-1} \mathcal{F}, \mathcal{F} \rangle \, dx \, dt \right)$$

where E^{-1} denotes the inverse of matrix E .

Proof. Since Theorem 2 is still valid also for the FE-solutions we can use the inequality (15) and let $a \rightarrow \infty$, so we obtain

$$(18) \quad \int_{\Omega} \langle E \mathcal{C}(T), \mathcal{C}(T) \rangle \, dx \leq e^{kT} \left(\int_{\Omega} \langle E \mathcal{C}^0, \mathcal{C}^0 \rangle \, dx + 2 \int_Q e^{-kt} \langle \mathcal{F}, \mathcal{C} \rangle \, dx \, dt \right).$$

If we write $E(t) = \int_0^t \int_{\Omega} \langle E \mathcal{C}(T), \mathcal{C}(\tau) \rangle \, dx \, d\tau$, by using the Schwarz inequality and the expression of E

$$\begin{aligned} & 2 \int_Q e^{-kt} \langle \mathcal{F}, \mathcal{C} \rangle \, dx \, dt \leq 2 \left(\int_Q \langle E^{-1} \mathcal{F}, \mathcal{F} \rangle \, dx \, dt \right)^{\frac{1}{2}} \left(\int_Q \langle E \mathcal{C}, \mathcal{C} \rangle \, dx \, dt \right)^{\frac{1}{2}} \\ & \leq \int_Q \langle E^{-1} \mathcal{F}, \mathcal{F} \rangle \, dx \, dt + \int_Q \langle E \mathcal{C}, \mathcal{C} \rangle \, dx \, dt = \int_Q \langle E^{-1} \mathcal{F}, \mathcal{F} \rangle \, dx \, dt + E(t). \end{aligned}$$

From (18), exchanging T with any arbitrary $t \in I$, there follows

$$E'(t) \leq e^{kT} \left(E'(0) + E(t) + \int_0^T S(\tau) \, d\tau \right) \quad \text{where} \quad S(t) = \int_{\Omega} \langle E^{-1} \mathcal{F}, \mathcal{F} \rangle \, dx.$$

The above inequality can be written as

$$E'(t) - k_1 E(t) \leq k_2 = e^{k_1 T} (E'(0) + \int_0^T S(\tau) d\tau) \quad k_1 = e^{k_1 T}$$

then
$$\frac{d}{dt} (e^{-k_1 t} E(t)) \leq e^{-k_1 t} (E'(t) - k_1 E(t)) \leq e^{-k_1 t} k_2.$$

After integration, with respect to t , and since $E(0) = 0$, we have

$$e^{-k_1 t} E(t) \leq \frac{1 - e^{-k_1 t}}{k_1} k_2, \quad k_1 E(t) \leq (e^{k_1 t} - 1) k_2, \quad E'(t) \leq k_2 e^{k_1 t}$$

which is equivalent to the inequality (17).

5 - Proof of the uniqueness

In this section the proof of the uniqueness of Theorem 1 is given. Since the initial and boundary value problem is a linear differential system, we have only to show that if $\mathcal{F} = 0$ almost everywhere in Q and $\mathcal{X}^0 = 0$ almost everywhere in Ω , then $\mathcal{X} = 0$ almost everywhere in Q .

Theorem 4. *Given the data $\mathcal{F} = 0$ and $\mathcal{X}^0 = 0$, we derive:*

- a. *if \mathcal{X} is a LFE-solution, then $\mathcal{X} = 0$ almost everywhere in Q .*
- b. *if \mathcal{X} a weak LFE-solution, then $\mathcal{X} = 0$ almost everywhere in Q .*

Proof. The proof of a follows from Theorem 2, since the inequality (15) holds for any LFE-solution and arbitrary positive numbers a, T . In order to obtain the proof of the second part of Theorem 4 we give first

Theorem 5. *If \mathcal{X} is a weak LFE-solution of problem (1)-(4) with boundary condition Γ and given $(\mathcal{F}, \mathcal{X}^0)$ and if*

$$(19) \quad \mathcal{X}_1 = \int_0^t \mathcal{X}(\tau) d\tau \quad \mathcal{F}_1 = \int_0^t \mathcal{F}(\tau) d\tau$$

then \mathcal{X}_1 defines a LFE-solution with the same boundary conditions and values $(\mathcal{F}_1(t) + E\mathcal{X}^0, 0)$ given.

Proof. Assuming $\tau \in I$ and $\Psi = (\Psi_a) \in (\Gamma^*)^{\text{vox}}(\Omega, \mathbf{R}^{10})$ we define

$$(20) \quad \Phi(x, t) = \begin{cases} (\tau - t) \Psi(x) & 0 \leq t \leq \tau \\ 0 & \tau \leq t \leq T. \end{cases}$$

It easily follows that $\Phi \in (F^*)^{\text{vox}}$ and $\Phi(T) = 0$, so that Φ can be used in equation (14), which defines the weak LFE-solutions (see Definition 2), instead of \mathcal{X} . Thus we obtain

$$\int_0^\tau \int_\Omega (\langle -E\Psi - (\tau - t)(A\Psi - B\Psi), \mathcal{C} \rangle + (\tau - t)\langle \mathcal{F}, \Psi \rangle) dx d\tau + \tau \int_\Omega \langle E\Psi, \mathcal{C}^0 \rangle dx = 0.$$

Deriving the above with respect to τ we obtain

$$(21) \quad - \int_\Omega \langle E\Psi, \mathcal{C} \rangle dx - \int_0^\tau \int_\Omega (\langle A\Psi - B\Psi, \mathcal{C} \rangle + \langle \mathcal{F}, \Psi \rangle) dx dt + \int_\Omega \langle E\Psi, \mathcal{C}^0 \rangle dx = 0$$

for almost any $\tau \in I$. Taking into account Fubini's Theorem and equation

$$(22) \quad \frac{\partial \mathcal{C}_1(\tau)}{\partial t} = \mathcal{C}(\tau)$$

equation (21) can be rewritten as

$$(23) \quad \int_\Omega (\langle A\Psi, \mathcal{C}_1 \rangle - \langle E\Psi, \frac{\partial \mathcal{C}_1}{\partial t} \rangle - \langle B\Psi, \mathcal{C}_1 \rangle + \langle \mathcal{F}_1, \Psi \rangle - \langle E\Psi, \mathcal{C}^0 \rangle) dx = 0$$

for almost any $\tau \in I$ and $\forall \Psi \in (\Gamma^*)^{\text{vox}}$. Since $C_0^\infty(\Omega, \mathbf{R}^{10}) \subset (\Gamma^*)^{\text{vox}}$, equation (23) implies the existence of \mathcal{C}_1 for almost any $\tau \in I$ and also

$$(24) \quad E \frac{\partial \mathcal{C}_1}{\partial t} = A\mathcal{C}_1 + B\mathcal{C}_1 + \mathcal{F}_1 + E\mathcal{C}^0.$$

Thus \mathcal{C}_1 is a LFE-solution because it satisfies equation (24) and from

$$\mathcal{C} \in L_2^{\text{loc}}(Q, \mathbf{R}^{10}) \quad \mathcal{F} \in L_2^{\text{loc}}(Q, \mathbf{R}^{10})$$

according to equation (5) and Theorem 2, it results:

$$\mathcal{C}_1 \in F^{\text{loc}} \quad \mathcal{F}_1 \in F^{\text{loc}} \quad \mathcal{C}_1(x, 0) = 0.$$

Now we are able to show that Theorem 4 **b** is a direct consequence of Theorem 5. In fact, $(\mathcal{F}, \mathcal{C}^0) = (0, 0)$ implies $(\mathcal{F}_1 + E\mathcal{C}^0, 0) = (0, 0)$ and thus $\mathcal{C} = 0$ holds according to Theorem 4 **a**. This implies also $\mathcal{C} = \mathcal{C}_1 = 0$.

Let us now consider the linear operator $M: \mathcal{O}(A) \rightarrow \mathbf{F}$, defined as

$$M\mathcal{X} = (E \frac{\partial \mathcal{X}}{\partial t} - A\mathcal{X} - B\mathcal{X}, E\mathcal{X}(0))$$

where

$$\mathcal{O}(M) = \mathbf{F} \quad \mathbf{F} \in L_2(Q; \mathbf{R}^{10}).$$

Let $\mathcal{R}(M) \subset \mathbf{F}$ be the range of M ; as a consequence of the uniqueness we can prove

Theorem 6. *The set $\mathcal{R}(M)$ is dense in F .*

Proof. We have to prove that $\overline{\mathcal{R}(M)} = F$. Let us assume that $\overline{\mathcal{R}(M)} \neq F$, then there exists in F a non-zero element \mathcal{X} , \mathcal{X}^0 orthogonal to $\mathcal{R}(M)$, i.e.

$$\langle E \frac{\partial \mathcal{X}}{\partial t} - A\mathcal{X} - B\mathcal{X}, \mathcal{X} \rangle + \langle E\mathcal{X}(0), \mathcal{X}^0 \rangle = 0 \quad \forall \mathcal{X} \in \mathcal{D}(M).$$

Therefore
$$\int_0^T \langle \langle E \frac{\partial \mathcal{X}}{\partial t} - A\mathcal{X} - B\mathcal{X}, \mathcal{X} \rangle + \langle E\mathcal{X}(0), \mathcal{X}^0 \rangle \rangle dt = 0 \quad \forall \mathcal{X} \in \mathcal{D}(M).$$

According to (14), it results that \mathcal{X} is a weak FE-solution with $\mathcal{F} = 0$ and $\mathcal{X}(T) = 0$. In particular, by assuming $\mathcal{X} = 0$ we have $\mathcal{X}^0 = 0$ and, according to Theorem 4, it follows $\mathcal{X} = 0$ almost everywhere in $\mathcal{D}(M)$, thus $(\mathcal{X}, \mathcal{X}^0) = (0, 0)$, which is contrary to the hypothesis. Therefore $\overline{\mathcal{R}(M)} = F$.

6 - Proof of the existence

In this section we begin proving the existence of the weak FE-solutions, claimed by Theorem 1. From Theorems 4, 6 it follows that the space F^* is dense in the Hilbert space Γ . A sequence \mathcal{X}_n exists such that if $\mathcal{F}_n \in F^*$ and $\mathcal{F}_n \rightarrow \mathcal{F}$ in $L_2(Q; \mathbf{R}^{10})$, \mathcal{F}_n being the sequence

$$\mathcal{F}_n = E \frac{\partial \mathcal{X}_n}{\partial t} - A\mathcal{X}_n - B\mathcal{X}_n$$

then $\mathcal{X}_n \rightarrow \mathcal{X}$.

Applying Theorem 3 to the differences $\mathcal{F}_n - \mathcal{F}_m$, $\mathcal{X}_n - \mathcal{X}_m$ we have

$$(25) \quad \int_{\Omega} \langle E(\mathcal{X}_n(T) - \mathcal{X}_m(T)), (\mathcal{X}_n(T) - \mathcal{X}_m(T)) \rangle dx \\ \leq C \left(\int_{\Omega} \langle E(\mathcal{X}_n^0 - \mathcal{X}_m^0), (\mathcal{X}_n^0 - \mathcal{X}_m^0) \rangle dx + \int_Q \langle E^{-1}(\mathcal{F}_n - \mathcal{F}_m), (\mathcal{F}_n - \mathcal{F}_m) \rangle dx dt \right).$$

Therefore, according to the Riesz-Fischer Theorem, since \mathcal{X}_n is a Cauchy sequence there exists a function $\mathcal{X} \in L_2$ such that $\mathcal{X}(t) = \lim_{n \rightarrow \infty} \mathcal{X}_n(t)$ in $L_2(\Omega)$ exists, $\forall t \in I$.

Making $m \rightarrow \infty$ in equation (25), there follows the uniform convergence and then the existence of the limit in $L_2(Q; \mathbf{R}^{10})$ too, i.e.

$$\mathcal{X}(t) = \lim_{n \rightarrow \infty} \mathcal{X}_n(t) \quad \text{in } L_2(Q; \mathbf{R}^{10}).$$

Thus from definition (14), written for \mathcal{X}_n and \mathcal{F}_n , we just let $n \rightarrow \infty$, and using the convergence of the above limits we get the existence of the weak FE-solutions.

Now we prove the existence of FE-solutions. These solutions, when exist, are also weak FE-solutions, therefore a sufficient condition for the existence is

Theorem 7. *The weak FE-solutions $\tilde{\mathcal{C}}$ with data $(\mathcal{F}, \mathcal{C}^0)$ with $\mathcal{F} \in L_2^{\text{loc}}(Q)$ and $\mathcal{C}^0 \in L_2^{\text{loc}}(Q)$ are FE-solutions.*

Proof. If $\tilde{\mathcal{C}}$ is a weak FE-solution, then according to Theorem 5 and to the linearity of the operators E and B , the function \mathcal{C}_1 of equation (19) is a FE-solution with the same boundary conditions and data $[(A+B)(E^{-1}\mathcal{F}_1 + \mathcal{C}^0), 0]$. From equation (24), and taking into account equation (22) we have

$$(26) \quad E\tilde{\mathcal{C}} = A\mathcal{C}_1 + B\mathcal{C}_1 + (A+B)(E^{-1}\mathcal{F}_1 + \mathcal{C}^0)$$

and

$$(27) \quad E\tilde{\mathcal{C}} = A\mathcal{C}^* + B\mathcal{C}_1 + B(E^{-1}\mathcal{F}_1 + \mathcal{C}^0)$$

where

$$(28) \quad \mathcal{C}^* = \mathcal{C}_1 + E^{-1}\mathcal{F}_1 + \mathcal{C}^0$$

so that equation (27) becomes

$$(29) \quad E\tilde{\mathcal{C}} = A\mathcal{C}^* + B\mathcal{C}^* .$$

On the other hand, deriving \mathcal{C}^* with respect to t , we get

$$(30) \quad E \frac{\partial \mathcal{C}^*}{\partial t} = E\tilde{\mathcal{C}} + \mathcal{F}$$

and comparing with (29)

$$E \frac{\partial \mathcal{C}^*}{\partial t} = A\mathcal{C}^* + B\mathcal{C}^* + \mathcal{F} \quad \mathcal{C}^*(0) = \mathcal{C}^0$$

i.e. \mathcal{C}^* is a FE-solution, but for the uniqueness Theorem this solution must coincide with \mathcal{C} .

References

- [1] A. C. ERINGEN and G. A. MAUGIN, *Electrodynamics of Continua*, II, Springer, New York 1990.
- [2] M. FABRIZIO, *Su i teoremi di unicit  e di reciprocit  nella teoria di un plasma caldo*, Boll. Un. Mat. Ital. 4 (1969), 542-553.
- [3] A. A. GALEEV and R. N. SUDAN, *Plasma Physics* I, II, North-Holland, Amsterdam 1983, 1984.

- [4] J. D. JACKSON, *Classical Electrodynamics*, J. Wiley and Sons, New York 1962.
- [5] P. D. LAX and R. S. PHILLIPS, *Local boundary conditions for dissipative symmetric linear differential operators*, Comm. Pure Appl. Math. **13** (1960), 427-455.
- [6] R. S. PHILLIPS, *Dissipative operators and hyperbolic systems of partial differential equations*, Trans. Amer. Math. Soc. **90** (1959), 193-254.
- [7] P. SECCHI, *On the equations of ideal incompressible magneto-hydrodynamics*, Rend. Sem. Mat. Univ. Padova **90** (1993), 103-119.
- [8] C. H. WILCOX, *Initial-boundary value problems for linear hyperbolic partial differential equations of the second order*, Arch. Rational Mech. Anal. **10** (1962), 361-400.
- [9] C. H. WILCOX, *The domain of dependence inequality for symmetric hyperbolic system*, Bull. Amer. Math. Soc. **70** (1964), 149-154.

Sommario

Vengono stabiliti teoremi di esistenza, unicit  e dipendenza continua dai dati per un sistema lineare iperbolico, che descrive l'evoluzione di un plasma in un dominio non limitato con condizioni al contorno dissipative.
