# 55. Initial Boundary Value Problem for the Equations of Ideal Magneto-Hydro-Dynamics with Perfectly Conducting Wall Condition 

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1. In this paper we consider the initial boundary value problem for the equations of ideal MHD that describe the motion of an ideal plasma filling an open subset of $\boldsymbol{R}^{3}$, surrounded by a rigid and perfectly conducting wall. (See [1].) Our problem is to solve
(1) ${ }_{b}$
(1) $\mathrm{c}_{\mathrm{c}} \quad \partial_{t} H-\nabla \times(u \times H)=0 \quad$ in $[0, T] \times \Omega$, (1) ${ }_{\mathrm{d}}$
$\left(\partial_{t}+(u \cdot \nabla)\right) S=0$
(1)
$\nabla \cdot H=0$
(2)
(3)
$\left.(p, u, H, S)\right|_{t=0}=\left(p_{0}, u_{0}, H_{0}, S_{0}\right) \equiv U_{0} \quad$ in $\Omega$,
$u \cdot n=0, H \cdot n=0 \quad$ on $[0, T] \times \Gamma$.

Here $\Omega$ is a bounded or unbounded domain in $R^{3}$ with a smooth and compact boundary $\Gamma$, or a half space $\boldsymbol{R}_{+}^{3}$; the pressure $p$, the velocity $u=\left(u^{1}, u^{2}, u^{3}\right)$, the magnetic field $H=\left(H^{1}, H^{2}, H^{3}\right)$, and the entropy $S$ are the unknown functions of $t$ and $x$; the density $\rho$ is determined by the equation of state $\rho=\rho(p, S) ; \rho>0$ and $\rho_{p}=\partial \rho / \partial p>0$ for $p>0$; the magnetic permeability $\mu$ is assumed to be constant; we write $\partial_{t}=\partial / \partial t, \partial_{i}=\partial / \partial x_{i}, \nabla=\left(\partial / \partial x_{1}, \partial / \partial x_{2}, \partial / \partial x_{3}\right)$ and use the conventional notations in vector analysis; $n=n(x)=\left(n_{1}, n_{2}, n_{3}\right)$ denotes the unit outward normal at $x \in \Gamma$.
2. We set $U={ }^{t}(p, u, H, S)$ and rewrite the system (1) $)_{a-d}$ in the symmetric form

$$
\begin{equation*}
A_{0}(U) \partial_{t} U+\sum_{i=1}^{3} A_{i}(U) \partial_{i} U=0 . \tag{4}
\end{equation*}
$$

In order to solve the problem by iteration, we consider the linearization of (4) around an arbitrary function $U^{\prime}={ }^{t}\left(p^{\prime}, u^{\prime}, H^{\prime}, S^{\prime}\right)$ near the initial data, satisfying $u^{\prime} \cdot n=0$ and $H^{\prime} \cdot n=0$ on $\Gamma$. The linearized equation forms a symmetric hyperbolic system with singular boundary matrix. In fact, the boundary matrix has constant rank 2 on $\Gamma$. We define $X^{m}(T, \Omega)$ to be the space of functions $U(t, x)$ taking values in $R^{8}$ and satisfying the following property: Let $\beta \geq 0$ be an integer and let $\Lambda_{1}, \cdots, \Lambda_{\beta}$ be an arbitrary $\beta$-tuple of smooth and bounded vector fields tangential to $\Gamma$, namely, let $\left\langle\Lambda_{i}(x), n(x)\right\rangle$ $=0$ for $x \in \Gamma, i=1, \cdots, \beta$. Then $\partial_{t}^{\alpha} \Lambda_{1} \cdots \Lambda_{\beta} \delta_{n}^{k} U(t, x) \in L^{\infty}\left(0, T ; L^{2}(\Omega)\right)$ for $\alpha+\beta \leq m-2 k, k=0,1, \cdots,[m / 2]$. Here $\partial_{n}$ denotes the partial differentia-

[^0]tion in the direction normal to $\Gamma$.
Our main results are the following two theorems.
Theorem 1. Let $\Omega$ be a bounded domain in $\boldsymbol{R}^{3}$ with smooth boundary
$\Gamma$. Let $m \geq 8$ be an integer. Suppose that $U_{0} \in H^{m}(\Omega)$ and that $U_{0}$ satisfies the following conditions
(5)
\[

$$
\begin{aligned}
& \nabla \cdot H_{0}=0, p_{0}>0 \quad \text { in } \Omega, H_{0} \cdot n=0 \quad \text { on } \Gamma, \\
& \partial_{t}^{k} u(0) \cdot n=0, k=0,1, \cdots, m-1, \quad \text { on } \Gamma .
\end{aligned}
$$
\]

Then there exists a constant $T_{0}>0$ such that the problem (1) $)_{\text {a-e }}$ (2) (3) has a unique solution $U \in X^{m}\left(T_{0}, \Omega\right)$.

Theorem 2. Let $\Omega$ be an unbounded domain in $\boldsymbol{R}^{3}$ with smooth and compact boundary $\Gamma$ or a half space $\boldsymbol{R}_{+}^{3}$. Let $m \geq 8$ be an integer. Suppose that $U_{0}-^{t}(c, 0) \in H^{m}(\Omega)$ for some constant $c>0$ and that $U_{0}$ satisfies the conditions given in Theorem 1. Then there exists a constant $T_{1}>0$ such that the problem (1) $)_{\mathrm{a}-\mathrm{e}}(2)(3)$ has a unique solution $U$ satisfying $U-^{t}(c, 0)$ $\in X^{m}\left(T_{1}, \Omega\right)$.

Remark 1. Let $U_{0}$ satisfy $\nabla \cdot H_{0}=0$ in $\Omega$ and $H_{0} \cdot n=0, \partial_{t}^{k} u(0) \cdot n=0$, $k=0,1, \cdots, m-1$, on $\Gamma$. Then the solution of (1) ${ }_{\mathrm{a}-\mathrm{d}}(2)$ satisfying $u \cdot n=0$ on $[0, T] \times \Gamma$ automatically satisfies $\nabla \cdot H=0$ in $[0, T] \times \Omega, H \cdot n=0$ on $[0, T]$ $\times \Gamma$ and $\partial_{t}^{k} H(0) \cdot n=0, k=0,1, \cdots, m-1$, on $\Gamma$. This means that we may $\operatorname{regard} \nabla \cdot H=0$ and $H \cdot n=0$ as the restrictions on the initial data $U_{0}$.

Remark 2. The characteristic boundary value problem was studied in [2]-[7]. Our approach is close to that of [6] and [7]. But some further considerations are needed for our problem.
3. Let $\Omega$ be a half space $R_{+}^{3}=\left\{x \mid x_{1}>0\right\}$. The general case can be reduced to this case by localization and flattening of the boundary. We introduce the new unknown function $V={ }^{t}(q-c, u, H, S)$ in place of $U$ $={ }^{t}(p, u, H, S)$, where $q=p+(1 / 2)|H|^{2}$ is the magnetic pressure, and rewrite the equations (1) $)_{\mathrm{a}-\mathrm{d}}$ in the form

$$
\begin{align*}
& \left(\begin{array}{cccc}
\alpha & 0 & -\alpha H & 0 \\
0 & \rho I_{3} & 0 & 0 \\
-\alpha^{t} H & 0 & I_{3}+\alpha H \otimes H & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \partial_{t} V  \tag{7}\\
& \quad+\left(\begin{array}{cccc}
\alpha(u \cdot \nabla) & \nabla & -\alpha H(u \cdot \nabla) & 0 \\
{ }^{t} \nabla & (u \cdot \nabla) I_{3} & -(H \cdot \nabla) I_{3} & 0 \\
-\alpha^{t} H(u \cdot \nabla) & -(H \cdot \nabla) I_{3} & \left(I_{3}+\alpha H \otimes H\right)(u \cdot \nabla) & 0 \\
0 & 0 & 0 & (u \cdot \nabla)
\end{array}\right) V \\
& \\
& \equiv A_{0}(V) \partial_{t} V+\sum_{i=1}^{3} A_{i}(V) \partial_{i} V=0
\end{align*}
$$

Here we set $\alpha=\rho_{q} / \rho$ and $H \otimes H=\left(H^{i} H^{j} \mid i \rightarrow 1,2,3, j \downarrow 1,2,3\right)$. Note that $\rho=\rho(q, H)>0, \rho_{q}>0$ for $q-(1 / 2)|H|^{2}>0$. We write

$$
A_{i}(V)=\left(\begin{array}{ll}
P_{i}(V) & Q_{i}(V)  \tag{8}\\
t Q_{i}(V) & R_{i}(V)
\end{array}\right) \quad i=0,1,2,3
$$

where $P_{i}(V), Q_{i}(V)$, and $R_{i}(V)$ are $2 \times 2,2 \times 6$ and $6 \times 6$ matrices, respectively. We write also $v=\left(q-c, u^{1}\right), w=\left(u^{2}, u^{3}, H^{1}, H^{2}, H^{3}, S\right)$. Hence,
$V={ }^{t}(v, w)$. Notice that

$$
\left.P_{1}(V)\right|_{x_{1}=0}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right),\left.\quad Q_{1}(V)\right|_{x_{1}=0}=0,\left.\quad R_{1}(V)\right|_{x_{1}=0}=0
$$

if $\left.u^{1}\right|_{x_{1}=0}=\left.H^{1}\right|_{x_{1}=0}=0$. For a function $f(t, x)$ valued in $\boldsymbol{R}^{d}, d=6$, 8 , we set
(9) $\quad\|f(t)\|_{m}^{2}=\sum_{k=0}^{[m / 2]} \sum_{|\theta| \leq m-2 k}\left|\partial_{*}^{e} \partial_{1}^{k} f(t)\right|_{0}^{2},\|f\|_{m, T}=\underset{t \in[0, T]}{\operatorname{esss} \sup }\|f(t)\|_{m}$,
where $\ell=\left(\alpha, \beta_{1}, \beta_{2}, \beta_{3}\right)$ and $\partial_{*}^{e}=\partial_{t}^{\alpha}\left(\phi\left(x_{1}\right) \partial_{1}\right)^{\beta_{1}} \partial_{2}{ }^{\beta_{2}} \partial_{3}{ }^{\beta_{3}}$. The weight $\phi\left(x_{1}\right)$ is a smooth and positive function such that $\phi\left(x_{1}\right)=x_{1}$ for $x_{1}$ small enough and $\phi\left(x_{1}\right)=1$ for $x_{1} \geq 1$, and $|\cdot|_{0}$ denotes $L^{2}\left(\boldsymbol{R}_{+}^{3}\right)$-norm. Then $X^{m}\left(T, \boldsymbol{R}_{+}^{3}\right)$ consists of all functions $f(t, x)$ for which $\|f\|_{m, T}<\infty$. This is a Banach space with $\|\cdot\|_{m, T}$ taken as the norm. Now we study the linearized problem.

$$
\begin{align*}
& A_{0}\left(V^{\prime}\right) \partial_{t} V+\sum_{i=1}^{3} A_{i}\left(V^{\prime}\right) \partial_{i} V=0 \text { in }[0, T] \times \boldsymbol{R}_{+}^{3},  \tag{10}\\
& \left.V\right|_{t=0}=\left(p_{0}+(1 / 2)\left|H_{0}\right|^{2}-c, u_{0}, H_{0}, S_{0}\right) \equiv V_{0} \text { in } \boldsymbol{R}_{+}^{3},  \tag{10}\\
& u^{1}=0 \text { on }[0, T] \times \partial \boldsymbol{R}_{+}^{3} .
\end{align*}
$$

Let $\kappa, M_{m-1}$, and $M_{m}$ be positive constants and let $X^{m}\left(T, R_{+}^{3} ; \kappa, M_{m-1}, M_{m}\right)$ be the set of functions $V^{\prime}$ satisfying the following conditions

$$
\begin{cases}V^{\prime} \in X^{m}\left(T, \boldsymbol{R}_{+}^{3}\right), \partial_{t}^{k} V^{\prime}(0) \in H^{m-k}\left(\boldsymbol{R}_{+}^{3}\right) & \text { for } k=0,1, \cdots, m-1  \tag{11}\\ u^{\prime \prime}=H^{\prime 1}=0 & \text { on }[0, T] \times \partial \boldsymbol{R}_{+}^{3} \\ q^{\prime}-(1 / 2)\left|H^{\prime}\right|^{2} \geq \kappa & \text { for }(t, x) \in[0, T] \times \boldsymbol{R}_{+}^{3}\end{cases}
$$

Then we have
Proposition 3. Let $m \geq 6$ and let $V^{\prime} \in X^{m}\left(T, \boldsymbol{R}_{+}^{3} ; \kappa, M_{m-1}, M_{m}\right)$. Then, (i) the null space of the boundary condition $(10)_{3}$ is the maximally non-negative subspace of the boundary matrix $-A_{1}\left(V^{\prime}\right)$ for $(t, x) \in[0, T] \times \partial R_{+}^{3}$, (ii) any smooth solution of $(10)_{1-3}$ satisfies $H^{1}=0$ on $[0, T] \times \partial \boldsymbol{R}_{+}^{3}$ if $H_{0}^{1}=0$ on $\partial \boldsymbol{R}_{+}^{3}$.

To get a counterpart of Proposition 3 for a general domain $\Omega$, some modification is needed. In this case we add the lower order term $B\left(V^{\prime}, V\right)$ $={ }^{t}\left(0,0,0,0, L\left(V^{\prime}, V\right), 0\right)$ to the left side of $(10)_{1}$, where

$$
L\left(V^{\prime}, V\right)=\tilde{n}\left\{H \cdot\left(\left(u^{\prime} \cdot \nabla\right) \tilde{n}\right)-u \cdot\left(\left(H^{\prime} \cdot \nabla\right) \tilde{n}\right)\right\}
$$

and $\tilde{n}=-\nabla$ dist $(x, \Gamma)$. Then the assertion (ii) remains valid with modified (10). We owe this idea to Taira Shirota.

Proposition 4. Let $m \geq 8$ and let $V^{\prime} \in X^{m}\left(T, R_{+}^{3} ; \kappa, M_{m-1}, M_{m}\right)$. Then a solution $V \in X^{m+1}\left(T, R_{+}^{3} ; \kappa, M_{m-1}, M_{m}\right)$ of the problem (10) $)_{1,3}$ satisfies (12) $\quad\|V(t)\|_{m} \leq C\left(M_{m-1}\right)\|V(0)\|_{m} \exp \left(C\left(M_{m}\right)\right) t \quad$ for $0 \leq t \leq T$.

Here $C\left(M_{s}\right), s=m-1, m$, are positive constants depending only on $M_{s}$.
We now combine Propositions 3 -(i) and 4 with the following arguments: (i) non-characteristic regularization (see, e.g., [5]), (ii) approximation of $V^{\prime}$ by smooth functions satisfying (11) and taking the same initial value as for $V^{\prime}$. Then we have

Proposition 5. Let $m \geq 8$ and let $V^{\prime} \in X^{m}\left(T, R_{+}^{3} ; \kappa, M_{m-1}, M_{m}\right)$. Suppose that $V_{0} \in H^{m+1}\left(\boldsymbol{R}_{+}^{3}\right)$ and that $V_{0}$ satisfies conditions (5) and (6). Then the problem (10) ${ }_{1-s}$ has a unique solution $V \in X^{m}\left(T, R_{+}^{3}\right)$ with the estimate (12).

By choosing the constants $\kappa, M_{m-1}, M_{m}$, and $T$ suitably and by making use
of Propositions 5 and 3-(ii), we can show that if $V^{\prime} \in X^{m}\left(T, R_{+}^{3} ; \kappa, M_{m-1}, M_{m}\right)$, the solution $V$ of $(10)_{1-3}$ again lies in the same set. This implies that the solution of the problem (1) $)_{\mathrm{a}-\mathrm{e}}$ (2) (3) is constructed by iteration combined with smoothing of the initial data. Uniqueness of solution follows from the energy inequality (5.20) in [8].

Now we sketch the proof of Proposition 4. First we prove the following estimates by the standard energy method,

$$
\begin{gather*}
\|V(t)\|_{m, *} \leq\|V(0)\|_{m, *}+C\left(M_{m}\right) \int_{0}^{t}\left(|[v(\tau)]|_{m}+\|w(\tau)\|_{m}\right) d \tau  \tag{13}\\
\|V(t)\|_{m-1} \leq\|V(0)\|_{m-1}+C\left(M_{m-1}\right) \int_{0}^{t}\|V(\tau)\|_{m} d \tau \tag{14}
\end{gather*}
$$

for $0 \leq t \leq T$. Here

$$
\|V(t)\|_{m, *}^{2}=\sum_{|\varepsilon| \leq m}\left|\partial_{*}^{e} V(t)\right|_{0}^{2},|[v(t)]|_{m}^{2}=\sum_{k=0}^{[m / 2]} \sum_{|\ell| \leq m-2 k+1}\left|\partial_{*}^{e} \partial_{1}^{k} v(t)\right|_{0}^{2} .
$$

In deriving (13), the main terms to be estimated are the commutator parts $\left[\partial_{*}^{\ell}, A_{1}\left(V^{\prime}\right)\right] \partial_{1} V,|\ell| \leq m$, which contain the terms such as $\partial_{*}^{\nu} Q_{1}\left(V^{\prime}\right) \partial_{*}^{\ell-\nu} \partial_{1} w$, $\partial_{*}^{u} R_{1}\left(V^{\prime}\right) \partial_{*}^{\epsilon-\nu} \partial_{1} w$, with $|\nu|=1$. We deal with these terms by regarding $\partial_{*}^{\nu} Q_{1}\left(V^{\prime}\right) \partial_{*}^{\ell-\nu} \partial_{1}$ and $\partial_{*}^{\nu} R_{1}\left(V^{\prime}\right) \partial_{*}^{\ell-\nu} \partial_{1}$ as the vector fields tangential to $\partial \boldsymbol{R}_{+}^{3}$. For instance, we have $\partial_{*}^{v} Q_{1}\left(V^{\prime}\right) \partial_{*}^{\ell-\nu} \partial_{1}=g\left(V^{\prime}\right) x_{1} \partial_{*}^{\ell-\nu} \partial_{1}$ where

$$
g\left(V^{\prime}\right)=\left.\int_{0}^{1} \partial_{1} \partial_{*}^{\sim} Q_{1}\left(V^{\prime}\right)\right|_{\left(t, \theta x_{1}, x_{2}, x_{3}\right)} d \theta,
$$

because $\left.Q_{1}\left(V^{\prime}\right)\right|_{x_{1}=0}=0$. Similar argument was used in [5]. Second, we express $\partial_{1} v$ in terms of $\partial_{1} w$ and the tangential derivatives of $V$. Using this expression and Rauch's argument, we obtain
(15)

$$
|[v(t)]|_{m} \leq C\left(M_{m-1}\right)\left(\|w(t)\|_{m}+\|V(t)\|_{m-1}+\|V(t)\|_{m, *}\right)
$$

Now we observe that $w$ satisfies

$$
R_{0} \partial_{t} w+\sum_{i=1}^{3} R_{i} \partial_{i} w=-\left({ }^{t} Q_{0} \partial_{t} v+\sum_{i=1}^{3}{ }^{t} Q_{i} \partial_{i} v\right)
$$

In view of $\left.R_{1}\right|_{x_{1}=0}=0$, we conclude that

$$
\begin{equation*}
\|w(t)\|_{m} \leq\|w(0)\|_{m}+C\left(M_{m}\right) \int_{0}^{t}\left(|[v(\tau)]|_{m}+\|w(\tau)\|_{m}\right) d \tau \tag{16}
\end{equation*}
$$

for $0 \leq t \leq T$. The estimate (12) follows from (13), (14), (15), (16), and Gronwall's inequality.

## References

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