# Global Weak Solutions to a Mathematical Model Using Aptamers in Chemotherapy and Imaging 

Naïma Aïssa *

May 13, 2008


#### Abstract

We investigate the model suggested by [1] dealing with the effect of aptamers on the efficiency of anticancer drugs. The system is described by nonlinear parabolic equations. The boundary condition on the inhibitor concentration is a nonlinear Michaelis-Menton form while the others are zero flux conditions. Due to the nonlinearity of the Michaelis-Menton boundary condition, the operator associated to the problem is nonlinear so we will use theory of perturbation of maximal monotone operators by pseudomonotones ones associated with Schauder fixed point theorem to get our main result.


## 1 The model equations and main result

Many anti-cancer drugs need to penetrate the cell membrane to perform their functions. K. Boushaba, H.A. Levine and M. Nilsen Hamilton [1] propose in their recent work that the intracellular concentration and the effectiveness of a drug might be increased by the presence, inside the cell, of a means of capturing the drug and moving it through the cytoplasm. As a drug binding agent, they use an aptamer: aptamers are small single stranded nucleic acids that have been selected for tight and spedific binding to a target molecule, which in this case is the drug. The authors studied the efficiency of the aptamers by comparing the system without aptamers and with aptamers.

### 1.1 The chemical kinetics (Protein-Inhibitor)

We are given an enzyme, an inhibitor, and a product, denoted by $E, I, P$. We use the notation $E(x, t)=[E](t)$ because species are distributed in space as well as time. The mecanism

$$
E+I \rightleftharpoons \rightleftharpoons_{k-1}^{k_{1}} P
$$

[^0]is not assumed to be in equilibrium. Suppose
(1) The reaction takes place in a bounded region $\Omega$ of the three space ( the cytoplasm of a cell)
(2) The enzyme decays with rate $\mu_{e}>0$. The enzyme-inhibitor complex $P=\{E: I\}$ decays with rate $\mu_{e i} \geq 0$. The inhibitor may decay with rate $\nu \geq 0$.
(3) The inhibitor species, I, diffuses much faster than $E, P$.
(4) The cell functions as a steady source for the enzyme and the aptamer (if any), i.e. there is a nonnegative function $S_{e}(x)$ supported in a region $\Omega_{E} \subset \Omega$ for all nonnegative times that defines the cellular rate of production of the enzyme.
These assupmtions and the law of mass action lead to the following where $\Delta$ denotes the Laplacian
\[

\left\{$$
\begin{array}{l}
\partial_{t} E(t, x)=D_{e} \Delta E+k_{-1} P-k_{1} E I-\mu_{e} E+S_{e}(x), \mathbb{R}^{+} \times \Omega,  \tag{1}\\
\partial_{t} P=D_{e} \Delta P-k_{-1} P+k_{1} E I-\mu_{e i} P, \mathbb{R}^{+} \times \Omega \\
\partial_{t} I=D \Delta I+k_{-1} P-k_{1} E I-\nu I, \mathbb{R}^{+} \times \Omega \\
I(0, x)=I_{0}(x) \geq 0, P(0, x)=P_{0}(x) \geq 0, E(0, x)=E_{0}(x) \geq 0, \\
D_{e} \partial_{n} E=0, D_{e} \partial_{n} P=0, \quad-D \partial_{n} I(x, t)=\frac{K_{c}\left(I(x, t)-i_{b}\right)}{\left(K_{m}+I(x, t)\right)}+T_{e}(I-J) \quad \text { on } \partial \Omega .
\end{array}
$$\right.
\]

where $i_{b}$ is the backgroud concentration of the drug in the cell due solely to membrane transport, $T_{e}$ the membrane permeability of the inhibitor, $J$ is the concentration of the drog on the apical side of the celle membrane and is asumed to be constant, where $K_{c}, K_{m}$ may be thought of Michealis-Menten constants for transport of inhibitor through the cell membrane.
The meaning of the Michaelis-Menton boundary condition is the following: If the concentration of $I$ is larger than the threshold value $i_{b}$, the contribution to the flux out of the cell by the active transport will be positive that is, $I$ will leave the cell. If the concentration is smaller than the threshold value, this contribution to the flux out will be negative. Likewise if $I>J$, the passive diffusion will contribute positively to the flux out of the cell while if $I<J$, the passive diffusion will contribute positively to the flux into the cell. [1].

### 1.2 The chemical kinetics (Enzyme-inhibitor-Aptamer)

We consider the effect of an aptamer on the efficiency of the enzyme reaction $E+I \rightleftharpoons P=I E$. The aptamer has a single binding site for the inhibitor only and interacts with it via $A+I \rightleftharpoons_{l_{-1}}^{l_{1}}$ $Q=A I$. The aptamer source $S_{a}(x)$ is supported on $\Omega_{A} \subset \Omega$. Mass action consideration lead to

$$
\left\{\begin{array}{l}
\partial_{t} I=D \Delta I+k_{-1} P-k_{1} E I+l_{-1} Q-l_{1} A I-\nu I, \mathbb{R}^{+} \times \Omega  \tag{2}\\
\partial_{t} A=D_{a} \Delta A+l_{-1} Q-\left(l_{1} I+\nu_{a}\right) A+S_{a}(x), \mathbb{R}^{+} \times \Omega \\
\partial_{t} Q=D_{a} \Delta Q+l_{1} A I-\left(l_{-1}+\nu_{a i}\right) Q, \mathbb{R}^{+} \times \Omega \\
\partial_{t} E(t, x)=D_{e} \Delta E+k_{-1} P-\left(k_{1} I+\nu_{e}\right) E+S_{e}(x), \mathbb{R}^{+} \times \Omega \\
\partial_{t} P=D_{e} \Delta P+k_{1} E I-\left(k_{-1}+\mu_{e i}\right) P, \quad \mathbb{R}^{+} \times \Omega, \\
I(0, x)=I_{0}(x) \geq 0, \quad A(0, x)=A_{0}(x) \geq 0, \quad Q(0, x)=Q_{0}(x) \geq 0 \\
P(0, x)=P_{0}(x) \geq 0, E(0, x)=E_{0}(x) \geq 0, \\
D_{a} \partial_{n} A=D_{a} \partial_{n} Q=0, \quad D_{e} \partial_{n} E=0, D_{e} \partial_{n} P=0, \quad \text { on } \partial \Omega \\
-D \partial_{n} I(x, t)=\frac{K_{c}\left(I(x, t)-i_{b}\right)}{K_{m}+I(x, t)}+T_{e}(I-J) \quad \text { on } \partial \Omega
\end{array}\right.
$$

We assume that the molecular wheights are ordered via $M_{I} \ll M_{A} \ll M_{E}$. The smallest molecule $I$, will be the most diffusible. Therefore $D=D_{i}>D_{a}=D_{A I}=D_{a} \geq D_{e} \geq 0$.

Definition 1. Let $(I, P, Q)$ satifiying

$$
I, P, Q \in L^{2}\left(\mathbb{R}^{+} ; H^{1}(\Omega)\right), \partial_{t} I, \partial_{t} P, \partial_{t} Q \in L^{2}\left(\mathbb{R}^{+} ;\left(H^{1}(\Omega)\right)^{\prime}\right),
$$

and $(A, E)$ such that

$$
A, E \in L^{2}\left(\mathbb{R}^{+} ; H^{1}(\Omega)\right), \quad \partial_{t} A, \partial_{t} E \in L^{2}\left(\mathbb{R}^{+} ; L^{2}(\Omega)\right) .
$$

We say that $(A, E, I, P, Q)$ is a global weak solution to (2) if
1-The second and fourth equations of (2) are satisfied a.e. and the boundary conditions $\partial_{n} A=$ $\partial_{n} E=0$ are satisfied in $H^{-\frac{1}{2}}(\Omega)$.
2- $(A, E, I, P, Q)(0)=\left(A_{0}, E_{0}, I_{0}, P_{0}, Q_{0}\right) \quad$ in $\quad L^{2}(\Omega)$,
3-For all $\varphi, \psi, \eta \in V$

$$
\begin{align*}
& \frac{d}{d t} \int_{\Omega} I \varphi d x=-\int_{\Omega} \nabla I \cdot \nabla \varphi d x-K_{c} \int_{\partial \Omega} \frac{I-i_{b}}{K_{m}+I} \varphi d \sigma-T_{e} \int_{\partial \Omega}(I-J) \varphi d \sigma \\
& +\int_{\Omega}\left(k_{-1} P-k_{1} E I+l_{-1} Q-l_{1} A I-\nu I\right) \varphi d x=0, \\
& \frac{d}{d t} \int_{\Omega} Q \psi d x=-D_{a} \int_{\Omega} \nabla Q \cdot \nabla \psi d x+\int_{\Omega}\left(l_{1} A I-\left(l_{-1}+\nu_{a i}\right) Q\right) \psi d x=0,  \tag{3}\\
& \frac{d}{d t} \int_{\Omega} P \eta d x=-D_{e} \int_{\Omega} \nabla P \cdot \nabla \eta d x+\int_{\Omega}\left(k_{1} E I-\left(k_{-1}+\nu_{e i}\right)\right) \eta d x=0 .
\end{align*}
$$

Theorem 1. Assume
$S_{a}, S_{e} \geq 0, S_{a}, S_{e} \in L^{\infty}(\Omega)$,
$l_{1} \leq l_{-1}+\nu_{a i}, \quad k_{1} \leq k_{-1}+\mu_{e i}, \quad l_{-1}+k_{-1} \leq \nu, \quad l_{-1}+\left\|S_{a}\right\|_{\infty} \leq \nu_{a}, \quad k_{-1}+\left\|S_{e}\right\|_{\infty} \leq \nu_{e}$,
$0 \leq I_{0}, Q_{0}, P_{0}, E_{0}, A_{0} \leq 1, \quad\left(I_{0}, P_{0}, Q_{0}, A_{0}, E_{0}\right) \in\left(H^{1}(\Omega)\right)^{5}$,
$\left(\Delta I_{0}, \Delta P_{0}, \Delta Q_{0}, \Delta A_{0}, \Delta E_{0}\right) \in\left(L^{2}(\Omega)\right)^{2}$.
Then, there exists a global weak solution $(A, E, I, P, Q)$ to (2) such that $0 \leq A, E, I, P, Q \leq 1$.

## 2 Preliminary results

### 2.1 Step I: Parabolic problem for given aptamer and enzyme

Denote by

$$
\begin{align*}
& V=H^{1}(\Omega) ; \quad H=L^{2}(\Omega) ; \quad V^{\prime}=\left(H^{1}(\Omega)\right)^{\prime}, \\
& \mathcal{V}=L^{2}\left(0, T ; H^{1}(\Omega)\right) ; \mathcal{H}=L^{2}\left(0, T ; L^{2}(\Omega)\right),  \tag{5}\\
& \mathcal{V}^{\prime}=L^{2}\left(0, T ;\left(H^{1}(\Omega)\right)^{\prime}\right), \quad \mathcal{C}=\left\{v \in \mathcal{V}, \partial_{t} v \in \mathcal{V}^{\prime}, v(0)=0\right\} .
\end{align*}
$$

where $V^{\prime}$ is the dual of $V$, the norm in $H$ will be denoted by $\|\|$.
Let us recall the principal definitions
Definition 2. [3], [6], [5]
Let $V$ be a reflexive Banach space; consider an operator $\mathcal{A}: V \rightarrow V^{\prime}$.
$1-\mathcal{A}$ is type $M$ if $u_{n} \rightharpoonup u, \mathcal{A} u_{n} \rightharpoonup f$ and $\lim \sup \mathcal{A} u_{n}\left(u_{n}\right) \leq f(u)$ imply that $\mathcal{A} u=f$.
2-The function $\mathcal{A}$ is coercive if $\frac{\mathcal{A} u(u)}{\|u\|} \rightarrow \infty$ as $\|u\| \rightarrow \infty$.
3- $\mathcal{A}$ is hemicontinuous if for each $u, v \in V$ the real-valued function $t \rightarrow \mathcal{A}(u+t v)(v)$ is continuous.
4- $\mathcal{A}$ is bounded if $S$ bounded in $V$ implies the image $\mathcal{A}(S)$ is bounded in $V^{\prime}$.
5- $\mathcal{A}$ is monotone if $\langle\mathcal{A}(u)-\mathcal{A}(v) ; u-v\rangle \geq 0$ for all $u, v \in V$.
$6-\mathcal{A}$ is strictly monotone if $\langle\mathcal{A}(u)-\mathcal{A}(v) ; u-v\rangle>0$ for all $u, v \in V, u \neq v$.
$7-\mathcal{A}$ is strongly continuous if and only if

$$
\begin{equation*}
u_{n} \rightharpoonup u \text { as } n \rightarrow+\infty \tag{6}
\end{equation*}
$$

implies

$$
\begin{equation*}
\mathcal{A} u_{n} \rightarrow \mathcal{A} u \text { as } n \rightarrow+\infty \tag{7}
\end{equation*}
$$

8- $\mathcal{A}$ is demicontinuous if

$$
\begin{equation*}
u_{n} \rightarrow u \text { as } n \rightarrow+\infty \tag{8}
\end{equation*}
$$

implies

$$
\begin{equation*}
\mathcal{A} u_{n} \rightharpoonup \mathcal{A} u \text { as } n \rightarrow+\infty \tag{9}
\end{equation*}
$$

Let $\tilde{A}, \tilde{E} \in L_{\text {loc }}^{\infty}\left(\mathbb{R}^{+} ; L^{2}(\Omega)\right), 0 \leq \tilde{A}, \tilde{E} \leq 1$ be given and consider the nonlinear parabolic system

$$
\left\{\begin{array}{l}
\partial_{t} I=D \Delta I+k_{-1} P-k_{1} \tilde{E} I+l_{-1} Q-l_{1} \tilde{A} I-\nu I, \mathbb{R}^{+} \times \Omega  \tag{10}\\
\partial_{t} Q=D_{a} \Delta Q+l_{1} \tilde{A} I-\left(l_{-1}+\nu_{a i}\right) Q, \mathbb{R}^{+} \times \Omega \\
\partial_{t} P=D_{e} \Delta P+k_{1} \tilde{E} I-\left(k_{-1}+\mu_{e i}\right) P, \mathbb{R}^{+} \times \Omega \\
I(0, x)=I_{0}(x) \geq 0, \quad Q(0, x)=Q_{0}(x) \geq 0, P(0, x)=P_{0}(x) \geq 0 \\
D_{a} \partial_{n} Q=0, \quad D_{e} \partial_{n} P=0, \quad-D \partial_{n} I(x, t)=\frac{K_{c}\left(I(x, t)-i_{b}\right)}{K_{m}+I(x, t)}+T_{e}(I-J) \quad \text { on } \partial \Omega
\end{array}\right.
$$

Theorem 2. Assume

$$
\begin{align*}
& l_{1} \leq l_{-1}+\nu_{a i}, \quad k_{1} \leq k_{-1}+\mu_{e i} ; \quad l_{-1}+k_{-1} \leq \nu \\
& 0 \leq I_{0}, Q_{0}, P_{0} \leq 1, \quad\left(I_{0}, P_{0}, Q_{0}\right) \in V, \quad\left(\Delta I_{0}, \Delta P_{0}, \Delta Q_{0}\right) \in H \tag{11}
\end{align*}
$$

Then, there exists a unique global weak solution $(I, P, Q) \in \mathcal{V}^{3}$ such that $\partial_{t}(I, P, Q) \in\left(\mathcal{V}^{\prime}\right)^{3}$. Moreover, $0 \leq I, Q, P \leq 1$ and the following energy estimate holds

$$
\begin{align*}
& \frac{d}{2 d t} \int_{\Omega}\left(|I|^{2}+|P|^{2}+|Q|^{2}\right) d x+D \int_{\Omega}|\nabla I|^{2} d x+D_{e} \int_{\Omega}|\nabla P|^{2} d x+D_{a} \int_{\Omega}|\nabla Q|^{2} d x+\nu\|I\|^{2} \\
& +\int_{\Omega} k_{1} \tilde{E}|I|^{2}+l_{1} \tilde{A}|I|^{2} d x+\left(l_{-a}+\nu_{a i}\right)\|Q\|^{2}+\left(k_{-1}+\mu_{e i}\right)\|P\|^{2} \\
& +K_{c} \int_{\partial \Omega} \frac{I^{2}}{K_{m}+I}+T_{e} I^{2} d \sigma=\int_{\Omega} k_{-1} P I+l_{-1} Q I+l_{1} \tilde{A} I Q+k_{1} \tilde{E} I P d x \\
& +i_{b} K_{c} \int_{\partial \Omega} \frac{I}{K_{m}+I}+T_{e} I J d \sigma . \tag{12}
\end{align*}
$$

Proof. Without loss of the generality, the proof will be done in the case when $P_{0}=Q_{0}=I_{0}=0$ in order to apply approximation method of evolution operators by stationnary ones. The proof in the case when the initial conditions does not vanish is obtained in the same way by setting $\tilde{I}=I-I_{0}, \tilde{P}=P-P_{0}, \tilde{Q}=Q-Q_{0}$ and repeating the same arguments.
In order to get solutions (which are concentrations) satisfying $0 \leq I, P, Q \leq 1$, we will study, for given $0 \leq \tilde{E}, \tilde{A} \leq 1$ the following problem where $f^{+}=\max (f, 0)$

$$
\left\{\begin{array}{l}
\partial_{t} I=D \Delta I+k_{-1} P^{+}-k_{1} \tilde{E} I+l_{-1} Q^{+}-l_{1} \tilde{A} I-\nu I, \mathbb{R}^{+} \times \Omega  \tag{13}\\
\partial_{t} Q=D_{a} \Delta Q+l_{1} \tilde{A} \min \left(I^{+}, 1\right)-\left(l_{-1}+\nu_{a i}\right) Q, \mathbb{R}^{+} \times \Omega \\
\partial_{t} P=D_{e} \Delta P+k_{1} \tilde{E} \min \left(I^{+}, 1\right)-\left(k_{-1}+\mu_{e i}\right) P, \mathbb{R}^{+} \times \Omega \\
I(0, x)=0, \quad Q(0, x)=0, P(0, x)=0 \\
D_{a} \partial_{n} Q=0, \quad D_{e} \partial_{n} P=0, \quad-D \partial_{n} I(x, t)=\frac{K_{c}\left(I^{+}(x, t)-i_{b}\right)}{K_{m}+I^{+}(x, t)}+T_{e}(I-J) \quad \text { on } \partial \Omega
\end{array}\right.
$$

Let $\mathcal{A}$ be the operator defined on $H^{3}$ by

$$
\begin{align*}
& \mathcal{A} U=\left(-D \Delta u_{1}+u_{1},-D_{a} \Delta u_{2}+u_{2},-D_{e} \Delta u_{3}+u_{3}\right), \text { for } U=\left(u_{1}, u_{2}, u_{3}\right) \\
& \mathcal{D}(\mathcal{A})=\left\{U \in V^{3}, \mathcal{A} U \in H^{3}, \partial_{n} u_{2}=\partial_{n} u_{3}=0,-D \partial_{n} u_{1}=\frac{K_{c}\left(u_{1}^{+}-i_{b}\right)}{K_{m}+u_{1}^{+}}+T_{e}\left(u_{1}-J\right)\right\} \tag{14}
\end{align*}
$$

The linear operator $\mathcal{A}$ is not monotone because $u_{1}^{+}-i_{b}$ and $u_{1}-J$ may change signe on the boundary. Let $\mathcal{B}$ the operator associated with $\mathcal{A}$, that is $\mathcal{B}: \mathcal{V} \rightarrow \mathcal{V}^{\prime}$ defined by

$$
\begin{align*}
& \mathcal{B} U(W)=\int_{(0, T) \times \Omega} U W+\nabla U \cdot \nabla W d x d t+K_{c} \int_{(0, T) \times \partial \Omega} \frac{u_{1}^{+}-i_{b}}{K_{m}+u_{1}^{+}} w_{1} d \sigma d t  \tag{15}\\
& +T_{e} \int_{(0, T) \times \partial \Omega}\left(u_{1}-J\right) w_{1} d \sigma d t
\end{align*}
$$

where $U=\left(u_{1}, u_{2}, u_{3}\right) \in \mathcal{V}, W=\left(w_{1}, w_{2}, w_{3}\right) \in \mathcal{V}$ and define $K: \mathcal{C}^{3} \rightarrow \mathcal{H}^{3} \quad$ by

$$
K\left(u_{1}, u_{2}, u_{3}\right)=\left(\begin{array}{l}
-k_{-1} u_{3}^{+}+k_{1} \tilde{E} u_{1}^{+}-l_{-1} u_{2}^{+}+l_{1} \tilde{A} u_{1}+\nu u_{1}  \tag{16}\\
-l_{1} \tilde{A} \min \left(u_{1}^{+}, 1\right)+\left(l_{-1}+\nu_{a i}\right) u_{2} \\
-k_{1} \tilde{E} \min \left(u_{1}^{+}, 1\right)+\left(k_{-1}+\nu_{e i}\right) u_{3}
\end{array}\right)
$$

Lemma 1. The operator $\mathcal{B}$ is monotone.
Proof. By definition of $\mathcal{B}$, we have

$$
\begin{align*}
& \langle\mathcal{B} U-\mathcal{B} W ; U-W\rangle= \\
& \int_{\Omega_{T}} D\left|\nabla\left(u_{1}-w_{1}\right)\right|^{2}+D_{a}\left|\nabla\left(u_{2}-w_{2}\right)\right|^{2}+D_{e}\left|\nabla\left(u_{3}-w_{3}\right)\right|^{2} d x d t \\
& +\int_{\Omega_{T}}|U-W|^{2} d x d t+K_{c} \int_{\partial \Omega_{T}}\left(\frac{u_{1}^{+}-i_{b}}{K_{m}+u_{1}^{+}}-\frac{w_{1}^{+}-i_{b}}{K_{m}+w_{1}^{+}}\right)\left(u_{1}-w_{1}\right) d \sigma d t  \tag{17}\\
& +T_{e} \int_{\partial \Omega_{T}}\left(u_{1}-w_{1}\right)^{2} d \sigma d t .
\end{align*}
$$

Then

$$
\begin{align*}
& \langle\mathcal{B} U-\mathcal{B} W ; U-W\rangle= \\
& \int_{\Omega_{T}} D\left|\nabla\left(u_{1}-w_{1}\right)\right|^{2}+D_{a}\left|\nabla\left(u_{2}-w_{2}\right)\right|^{2}+D_{e}\left|\nabla\left(u_{3}-w_{3}\right)\right|^{2} d x d t \\
& +\int_{\Omega_{T}}|U-W|^{2} d x d t+K_{c}\left(i_{b}+K_{m}\right) \int_{\partial \Omega_{T}} \frac{\left(u_{1}^{+}-w_{1}^{+}\right)}{\left(K_{m}+u_{1}^{+}\right)\left(K_{m}+w_{1}^{+}\right)}\left(u_{1}-w_{1}\right) d \sigma d t  \tag{18}\\
& +T_{e} \int_{\partial \Omega_{T}}\left(u_{1}-w_{1}\right)^{2} d \sigma d t .
\end{align*}
$$

Hence, writing $u_{1}-w_{1}=\left(u_{1}^{+}-w_{1}^{+}\right)-\left(u_{1}^{-}-w_{1}^{-}\right)$we get

$$
\begin{align*}
& \langle\mathcal{B} U-\mathcal{B} W ; U-W\rangle= \\
& \int_{\Omega_{T}} D\left|\nabla\left(u_{1}-w_{1}\right)\right|^{2}+D_{a}\left|\nabla\left(u_{2}-w_{2}\right)\right|^{2}+D_{e}\left|\nabla\left(u_{3}-w_{3}\right)\right|^{2} d x d t \\
& +\int_{\Omega_{T}}|U-W|^{2} d x d t+K_{c}\left(i_{b}+K_{m}\right) \int_{\partial \Omega_{T}} \frac{\left(u_{1}^{+}-w_{1}^{+}\right)^{2}}{\left(K_{m}+u_{1}^{+}\right)\left(K_{m}+w_{1}^{+}\right)} d \sigma d t  \tag{19}\\
& +K_{c}\left(i_{b}+K_{m}\right) \int_{\partial \Omega_{T}} \frac{u_{1}^{+} w_{1}^{-}+u_{1}^{-} w_{1}^{+}}{\left(K_{m}+u_{1}^{+}\right)\left(K_{m}+w_{1}^{+}\right)} d \sigma d t \\
& +T_{e} \int_{\partial \Omega_{T}}\left(u_{1}-w_{1}\right)^{2} d \sigma d t \geq 0 .
\end{align*}
$$

Lemma 2. The operator $K: \mathcal{C}^{3} \rightarrow \mathcal{H}^{3}$ is strongly continuous.

Proof. Set $K_{1}: \mathcal{C}^{3} \rightarrow \mathcal{H}^{3}$ defined by

$$
K_{1}\left(u_{1}, u_{2}, u_{3}\right)=\left(\begin{array}{l}
-k_{-1} u_{3}+k_{1} \tilde{E} u_{1}-l_{-1} u_{2}+l_{1} \tilde{A} u_{1}+\nu u_{1}  \tag{20}\\
-l_{1} \tilde{A} u_{1}+\left(l_{-1}+\nu_{a i}\right) u_{2} \\
-k_{1} \tilde{E} u_{1}+\left(k_{-1}+\nu_{e i}\right) u_{3}
\end{array}\right)
$$

$K_{1}$ is compact and linear then $K_{1}$ is strongly continuous (Cf. [6] Proposition 26-2 p. 555). Since the mappings $\left(u_{1}, u_{2}, u_{3}\right) \rightarrow u_{i}^{+}$and $\left(u_{1}, u_{2}, u_{3}\right) \rightarrow \min \left(u_{i}, 1\right)$ are Lipschitz continous from $\mathcal{H}$ to $\mathcal{H}$ then $K$ being a composition of a strongly continuous mapping with Lipschitz mapping so $K$ is strongly continuous.

Lemma 3. The operator $\mathcal{B}+K$ is coercive.
Proof. Since

$$
\begin{equation*}
\frac{\mathcal{B} U(U)}{\|U\|_{\mathcal{V}}}=\|U\|_{\mathcal{V}}+K_{c} \frac{\int_{(0, T) \times \partial \Omega} \frac{u_{1}^{+}-i_{b}}{u_{1}^{+}+K_{m}} u_{1} d \sigma d t}{\|U\|_{\mathcal{V}}}+T_{e} \frac{\int_{(0, T) \times \partial \Omega}\left(u_{1}-J\right) u_{1} d \sigma d t}{\|U\| \mathcal{V}} \tag{21}
\end{equation*}
$$

Then

$$
\begin{equation*}
\frac{\mathcal{B} U(U)}{\|U\|_{\mathcal{V}}} \geq\|U\|_{\mathcal{V}}-K_{c} i_{b} \frac{\int_{(0, T) \times \partial \Omega} \frac{u_{1}^{+}}{u_{1}^{+}+K_{m}} d \sigma d t}{\|U\|_{\mathcal{V}}}-T_{e} J \frac{\int_{(0, T) \times \partial \Omega} u_{1} d \sigma d t}{\|U\|_{\mathcal{V}}} . \tag{22}
\end{equation*}
$$

On the other hand

$$
\begin{equation*}
\int_{(0, T) \times \Omega} K U . U d x d t \geq C\|U\|_{\mathcal{V}}^{2}-\int_{(0, T) \times \Omega} l_{-1}\left|u_{2}\right|+k_{-1}\left|u_{3}\right| d x d t . \tag{23}
\end{equation*}
$$

for some positive constant $C$ depending on the constants occuring in the system ( $C>0$ thanks to the hypothesis $\left.l_{-1}+k_{-1} \leq \nu\right)$. Hence

$$
\begin{align*}
& \frac{[(\mathcal{B}+K)] U(U)}{\|U\|_{\mathcal{V}}} \geq(1+C)\|U\|_{\mathcal{V}}-K_{c} i_{b} \frac{\int_{(0, T) \times \partial \Omega} \frac{u_{1}^{+}}{u_{1}^{+}+K_{m}} d \sigma d t}{\|U\|_{\mathcal{V}}}-T_{e} J \frac{\int_{(0, T) \times \partial \Omega} u_{1} d \sigma d t}{\|U\|_{\mathcal{V}}}  \tag{24}\\
& -\frac{\int_{(0, T) \times \Omega} l_{-1}\left|u_{2}\right|+k_{-1}\left|u_{3}\right| d x d t}{\|U\|_{\mathcal{V}}} .
\end{align*}
$$

The second term of the RHS of the last inequality goes to 0 as $\|U\|_{\mathcal{V}}$ goes to $\infty$, while the third and the fourth term of RHS are bounded. It follows that $\frac{[(\mathcal{B}+K)] U(U)}{\|U\|_{\mathcal{V}}}$ goes to $\infty$ as $\|U\|_{\mathcal{V}}$ goes to $\infty$ and $\mathcal{B}+K$ is coercive.

Lemma 4. $\mathcal{B}$ is type $M$, bounded and hemicontinuous.
Proof. Assume that

$$
\begin{equation*}
U_{n} \rightharpoonup U \text { in } \mathcal{V}, \mathcal{B} U_{n} \rightharpoonup f \text { in } \mathcal{V}^{\prime}, \quad \lim \sup \mathcal{B} U_{n}\left(U_{n}\right) \leq f(U) \tag{25}
\end{equation*}
$$

First $U_{n} \rightharpoonup U$ and $\lim \sup \mathcal{B} U_{n}\left(U_{n}\right) \leq f(U)$ imply that $U_{n}$ and $\nabla U_{n}$ are bounded in $\mathcal{H}$. Consequently $u_{n} \rightarrow u$ in $L^{2}(0, T ; \partial \Omega)$ strong. It follows from Lebesgue's convergence theorem that $\frac{u_{n}^{+}-i_{b}}{K_{m}+u_{n}^{+}} \rightarrow \frac{u^{+}-i_{b}}{K_{m}+u^{+}}$and $u_{n}-J \rightarrow u-J$ in $L^{2}((0, T) \times \partial \Omega)$ strong. Then, for all $V \in \mathcal{V}$

$$
\begin{equation*}
\int_{\partial \Omega_{T}} \frac{\left(u_{n}^{+}-i_{b}\right)}{K_{m}+u_{n}^{+}} v d \sigma d t \rightarrow \int_{\partial \Omega_{T}} \frac{\left(u^{+}-i_{b}\right)}{K_{m}+u^{+}} v d \sigma d t, \int_{\partial \Omega_{T}}\left(u_{n}-J\right) v d \sigma d t \rightarrow \int_{\partial \Omega_{T}}(u-J) v d \sigma d t \tag{26}
\end{equation*}
$$

Moreover $\mathcal{A} u_{n} \rightharpoonup f$ in $\mathcal{V}^{\prime}$ writes

$$
\begin{equation*}
\int_{\Omega_{T}} U_{n} V+\nabla U_{n} \cdot \nabla V d x d t+K_{c} \int_{\partial \Omega_{T}} \frac{u_{n}^{+}-i_{b}}{K_{m}+u_{n}^{+}} v d \sigma+T_{e} \int_{\partial \Omega_{T}}\left(u_{n}-J\right) v \rightarrow f(V) \tag{27}
\end{equation*}
$$

for all $V \in \mathcal{V}$. Then passing to the limit we get $\mathcal{B} U=f$ in $\mathcal{V}^{\prime}$ so $\mathcal{B}$ is type M.
We check easily that $\mathcal{B}$ is bounded and hemicontinuous. Consequently $\mathcal{B}$ is pseudomonotone ( Cf. [6] proposition $27-6$ p. 586). Next, we will apply the following result to prove existence of solutions to (10).

Theorem 3. [6] Corollary 32.25 p. 867

## Suppose

(H1) $C$ is a nonempty closed convex set in the real $B$-Space $X$.
(H2) The mapping $A: C \rightarrow 2^{X^{\prime}}$ is maximal monotone.
(H3) The mapping $B: C \rightarrow X^{\prime}$ is pseudomonotone, bounded, and demicontinuous.
Suppose moreover that one of the following two conditions is satisfied
(i) $C$ is bounded.
(ii) $C$ is unbounded, and $B$ is $A$-coercive, i.e. there exists $u_{0} \in C \cap \mathcal{D}(A)$ such that

$$
\begin{equation*}
\frac{\left\langle B u, u-u_{0}\right\rangle}{\|u\|} \rightarrow+\infty \text { as }\|u\| \rightarrow+\infty \text { in } C . \tag{28}
\end{equation*}
$$

Then, $R(A+B)=X^{\prime}$. That is, for each $b \in X^{\prime}$, the original problem $b \in A u+B u, u \in C$, has a solution.

Corollary 1. For each $f \in \mathcal{V}^{\prime}$, the following abstract Cauchy problem

$$
\left\{\begin{array}{l}
U^{\prime}(t)+\mathcal{B} U(t)+K U(t)=f(t), \quad 0<t<T  \tag{29}\\
u(0)=0
\end{array}\right.
$$

has a unique solution $U \in \mathcal{C}$.
Proof. Set

$$
\begin{equation*}
\left.\left.\Lambda: D(\Lambda) \rightarrow \mathcal{H} \rightarrow \mathcal{H}, \quad \Lambda u=u^{\prime}, \quad D(\Lambda)=\left\{u \in L^{2}(0, T ; V)\right), u^{\prime} \in L^{2}\left(0, T ; V^{\prime}\right)\right), u(0)=0\right\} \tag{30}
\end{equation*}
$$

We will apply the previous Theorem in the case when $\mathcal{C}=D(\Lambda), A=\Lambda, B=\mathcal{B}+K$ and $u_{0}=0$. By virtue of [6] Proposition 32.10 p. $855, A$ is maximal monotone.
As $\mathcal{B}$ is bounded, hemicontinuous and monotone then $\mathcal{B}$ is pseudomonotone (Cf. [3] p. 179 or [6] Proposition 27.6 p. 586). Since $K$ is strongly continuous then $\mathcal{B}+K$ is pseudomonotone thanks to [6] Proposition 27.6 p. 586. Moreover, by virtue of the previous lemma $\mathcal{B}+K$ is coercive. Finally, we check easily that $\mathcal{B}+K$ is demicontinuous, hence by virtue of the previous theorem where we take $u_{0}=0$, there exists at least one solution $U \in \mathcal{C}$ solution to (29).
For the uniqueness, let $U_{i} \in \mathcal{C}, i=1,2$ be two solutions to (29) and set $U=U_{1}-U_{2}$. We have

$$
\begin{equation*}
\frac{1}{2}\|U(T)\|_{H}^{2}+\left\langle\mathcal{B} U_{1}-\mathcal{B} U_{2} ; U_{1}-U_{2}\right\rangle+\left\langle K\left(U_{1}-U_{2}\right) ; U_{1}-U_{2}\right\rangle=0 \tag{31}
\end{equation*}
$$

Since $\mathcal{B}$ is monotone

$$
\begin{equation*}
\frac{1}{2}\|U(T)\|_{\mathcal{H}}^{2} \leq-\left\langle K\left(U_{1}-U_{2}\right) ; U_{1}-U_{2}\right\rangle \tag{32}
\end{equation*}
$$

As $K: H \rightarrow H$ is Lipschitz then

$$
\begin{equation*}
\|U(T)\|_{\mathcal{H}}^{2} \leq C \int_{0}^{T}\|U\|_{H}^{2}(s) d s \tag{33}
\end{equation*}
$$

Then by Gronwall's Lemma we get $\|U\|_{H}(T)=0$ for arbitrary $T>0$ then the solution is unique.

Next, we will prove that $U \geq 0$. Indeed, multiplying the first Eq of (13) by $I^{-}$

$$
\begin{align*}
& -\frac{d}{2 d t} \int_{\Omega}\left|I^{-}\right|^{2}(t) d x=D\left\|\nabla I^{-}\right\|^{2}+k_{1} \int_{\Omega} \tilde{E}\left|I^{-}\right|^{2} d x+l_{-1} \int_{\Omega} Q^{+} I^{-} d x+l_{1} \int_{\Omega} \tilde{A}\left|I^{-}\right|^{2} d x \\
& +\nu \int_{\Omega}\left|I^{-}\right|^{2} d x+i_{b} K_{c} \int_{\partial \Omega} \frac{I^{-}}{K_{m}+I^{+}} d \sigma+T_{e} \int_{\partial \Omega}\left|I^{-}\right|^{2}+J I^{-} d \sigma+k_{-1} \int_{\Omega} P^{+} I^{-} \geq 0, \tag{34}
\end{align*}
$$

hence, assuming that $I_{0} \geq 0$ we get

$$
\begin{equation*}
\int_{\Omega}\left|I^{-}\right|^{2}(t) d x \leq \int_{\Omega}\left|I_{0}^{-}\right|^{2} d x=0, \tag{35}
\end{equation*}
$$

then $I^{-}=0$ a.e. i.e. $I \geq 0$ a.e.. Moreover

$$
\begin{align*}
& -\frac{d}{2 d t} \int_{\Omega}\left|Q^{-}\right|^{2}(t) d x=D_{a}\left\|\nabla Q^{-}\right\|^{2}+l_{1} \int_{\Omega} \tilde{A} \min \left(I^{+}, 1\right) Q^{-} d x+\left(l_{-1}+\nu_{a i}\right) \int_{\Omega}\left|Q^{-}\right|^{2} d x \geq 0 \\
& -\frac{d}{2 d t} \int_{\Omega}\left|P^{-}\right|^{2}(t) d x=D_{e}\left\|\nabla P^{-}\right\|^{2}+k_{1} \int_{\Omega} \tilde{E} \min \left(I^{+}, 1\right) P^{-} d x+\left(k_{-1}+\mu_{a i}\right) \int_{\Omega}\left|Q^{-}\right|^{2} d x \geq 0 \tag{36}
\end{align*}
$$

then we proceed similarly to prove $Q \geq 0, P \geq 0$ a.e. as soon as $Q_{0} \geq 0, P_{0} \geq 0$. Next assuming hypotheses of theorem, we get $I, P, Q \leq 1$. Indeed

$$
\begin{align*}
& \frac{d}{2 d t}\left\|(1-P)^{-}\right\|^{2}(t)=-D_{e}\left\|\nabla(1-P)^{-}\right\|^{2}+k_{1} \int_{\Omega} \tilde{E} \min \left(I^{+}, 1\right)(1-P)^{-} d x \\
& -\left(k_{-1}+\mu_{e i}\right)\left\|(1-P)^{-}\right\|^{2}-\left(k_{-1}+\mu_{e i}\right) \int_{\Omega}(1-P)^{-} d x . \tag{38}
\end{align*}
$$

then

$$
\begin{align*}
& \frac{d}{2 d t}\left\|(1-P)^{-}\right\|^{2}(t) \leq k_{1} \int_{\Omega} \tilde{E} \min \left(I^{+}, 1\right)(1-P)^{-} d x-\left(k_{-1}+\mu_{e i}\right) \int_{\Omega}(1-P)^{-} d x  \tag{39}\\
& \leq\left[k_{1}-\left(k_{-1}+\mu_{e i}\right)\right] \int_{\Omega}(1-P)^{-} d x .
\end{align*}
$$

Under hypothesis $k_{1} \leq k_{-1}+\mu_{e i}$ we get $\frac{d}{2 d t}\left\|(1-P)^{-}\right\|^{2}(t) \leq 0$ hence $P \leq 1$ as soon as $P_{0} \leq 1$. Similarly

$$
\begin{align*}
& \frac{d}{2 d t}\left\|(1-Q)^{-}\right\|^{2}(t)=-D_{a}\left\|\nabla(1-Q)^{-}\right\|^{2}+l_{1} \int_{\Omega} \tilde{E} \min \left(I^{+}, 1\right)(1-Q)^{-} d x \\
& -\left(l_{-1}+\nu_{a i}\right)\left\|(1-Q)^{-}\right\|^{2}-\left(l_{-1}+\nu_{a i}\right) \int_{\Omega}(1-Q)^{-} d x \tag{40}
\end{align*}
$$

then

$$
\begin{align*}
& \frac{d}{2 d t}\left\|(1-Q)^{-}\right\|^{2}(t) \leq l_{1} \int_{\Omega} \tilde{E} \min \left(I^{+}, 1\right)(1-Q)^{-} d x-\left(l_{-1}+\nu_{a i}\right) \int_{\Omega}(1-Q)^{-} d x  \tag{41}\\
& \leq\left[l_{1}-\left(l_{-1}+\nu_{a i}\right)\right] \int_{\Omega}(1-Q)^{-} d x,
\end{align*}
$$

as $l_{1} \leq l_{-1}+\nu_{a i}$ by hypotheses, $\frac{d}{2 d t}\left\|(1-Q)^{-}\right\|^{2}(t) \leq 0$ hence $Q \leq 1$ as soon as $Q_{0} \leq 1$. Finally

$$
\begin{align*}
& \frac{d}{2 d t}\left\|(1-I)^{-}\right\|^{2}(t)=-D\left\|\nabla(1-I)^{-}\right\|^{2}+k_{-1} \int_{\Omega} P^{+}(1-I)^{-} d x-k_{1} \int_{\Omega} \tilde{E}\left|(1-I)^{-}\right|^{2} d x \\
& -k_{1} \int_{\Omega} \tilde{E}(1-I)^{-} d x l_{-1} \int_{\Omega} Q^{+}(1-I)^{-} d x-l_{1} \int_{\Omega} \tilde{A}\left|(1-I)^{-}\right|^{2} d x-l_{1} \int_{\Omega} \tilde{A}(1-I)^{-} d x  \tag{42}\\
& -\nu\left\|(1-I)^{-}\right\|^{2}-\nu \int_{\Omega}(1-I)^{-} d x,
\end{align*}
$$

hence

$$
\begin{align*}
& \frac{d}{2 d t}\left\|(1-I)^{-}\right\|^{2}(t) \leq k_{-1} \int_{\Omega} P^{+}(1-I)^{-} d x+l_{-1} \int_{\Omega} Q^{+}(1-I)^{-} d x-\nu \int_{\Omega}(1-I)^{-} d x  \tag{43}\\
& \leq\left(k_{-1}+l_{-1}-\nu\right) \int_{\Omega}(1-I)^{-} d x,
\end{align*}
$$

then assuming that $k_{-1}+l_{-1}-\nu \leq 0$ we get $I \leq 1$ as soon as $I_{0} \leq 1$.

### 2.2 Step II: Parabolic problem for given $\tilde{I}, \tilde{Q}, \tilde{P}$

For given $S_{a}, S_{e} \in L^{\infty}(\Omega), \tilde{I}, \quad \tilde{Q}, \quad \tilde{P} \in \mathcal{C}, 0 \leq \tilde{I}, \quad \tilde{P}, \quad \tilde{Q} \leq 1$, consider the linear parabolic problem

$$
\left\{\begin{array}{l}
\partial_{t} A=D_{a} \Delta A+l_{-1} \tilde{Q}-\left(l_{1} \tilde{I}+\nu_{a}\right) A+S_{a}(x), \mathbb{R}^{+} \times \Omega,  \tag{44}\\
\partial_{t} E(t, x)=D_{e} \Delta E+k_{-1} \tilde{P}-\left(k_{1} \tilde{I}+\nu_{e}\right) E+S_{e}(x), \mathbb{R}^{+} \times \Omega, \\
A(0, x)=A_{0}(x) \geq 0, E(0, x)=E_{0}(x) \geq 0, \\
D_{a} \partial_{n} A=0, D_{e} \partial_{n} E=0 \text { on } \partial \Omega .
\end{array}\right.
$$

Theorem 4. Assume that

$$
\begin{align*}
& S_{a}, S_{e} \geq 0, S_{a}, S_{e} \in L^{\infty}(\Omega) \\
& l_{-1}+\left\|S_{a}\right\|_{\infty} \leq \nu_{a}, \quad k_{-1}+\left\|S_{e}\right\|_{\infty} \leq \nu_{e} \tag{45}
\end{align*}
$$

There exists a unique mild solution $(A, E) \in \mathcal{C}^{0}\left(0, \infty ; L^{2}(\Omega) \times \mathcal{C}^{0}\left(0, \infty ; L^{2}(\Omega)\right.\right.$ to (44). Moreover, assuming that $\left(A_{0}, E_{0}\right) \in H^{1}(\Omega) \times H^{1}(\Omega),\left(\Delta A_{0}, \Delta E_{0}\right) \in L^{2}(\Omega) \times L^{2}(\Omega)$, we get $A, E \in$ $L^{2}\left(0, T ; H^{1}(\Omega)\right), \Delta A, \Delta E \in L^{2}\left(0, T ; L^{2}(\Omega)\right), \partial_{t} A, \partial_{t} E \in L^{2}\left(0, T ; L^{2}(\Omega)\right)$. In addition, if $0 \leq$ $A_{0}, E_{0} \leq 1$ then $0 \leq A, E \leq 1$.
Furthermore, there exists $C>0$ such that the following energy estimate holds

$$
\begin{align*}
& \frac{d}{2 d t}\left(\|A\|^{2}+\|E\|^{2}\right)(t)+D_{a}\|\nabla A\|^{2}(t)+D_{e}\|\nabla E\|^{2}(t)+\nu_{a}\|A\|^{2}(t)+\nu_{e}\|E\|^{2}(t) \\
& +l_{1} \int_{\Omega} \tilde{I}|A|^{2} d x+k_{1} \int_{\Omega} \tilde{I}|E|^{2} d x=l_{-1} \int_{\Omega} \tilde{Q} A d x+k_{-1} \int_{\Omega} \tilde{P} E d x+\int_{\Omega} S_{a}(x) A d x  \tag{46}\\
& +\int_{\Omega} S_{e}(x) E d x
\end{align*}
$$

Proof. Let $0 \leq \tilde{Q}, \tilde{I}, \tilde{P} \leq 1$ be given. We notice that $A$ and $E$ are independent then we can solve separately

$$
\begin{gather*}
\partial_{t} A=D_{a} \Delta A+l_{-1} \tilde{Q}-\left(l_{1} \tilde{I}+\nu_{a}\right) A+S_{a}(x), \quad A(0, x)=A_{0}(x), \quad, \partial_{n} A=0 \text { on } \partial \Omega .  \tag{47}\\
\partial_{t} E(t, x)=D_{e} \Delta E+k_{-1} \tilde{P}-\left(k_{1} \tilde{I}+\nu_{e}\right) E+S_{e}(x), E(0, x)=E_{0}(x), \partial_{n} E=0 \text { on } \partial \Omega . \tag{48}
\end{gather*}
$$

Consider the operator $\mathcal{B}$ defined on $L^{2}(\Omega)$ and $F: L^{2}(\Omega) \rightarrow L^{2}(\Omega)$ by

$$
\begin{align*}
& \mathcal{B}(u)=D_{a} \Delta u ; \mathcal{D}(\mathcal{B})=\left\{u \in H^{1}(\Omega), \mathcal{B}(u) \in L^{2}(\Omega), \partial_{n} u=0\right\} \\
& F(u)=l_{-1} \tilde{Q}-\left(l_{1} \tilde{I}+\nu_{a}\right) u+S_{a}(x) \tag{49}
\end{align*}
$$

$\mathcal{B}$ is self-adjoint and negative and $F$ is Lipschitz in $L^{2}(\Omega)$ uniformly in time then there exists a unique mild solution $u \in \mathcal{C}^{0}\left(0, \infty ; L^{2}(\Omega)\right)$ to (47). Assuming that $A_{0} \in \mathcal{D}(\mathcal{B})$, we get $A \in$ $L^{2}\left(0, T ; H^{1}(\Omega)\right), \Delta A \in L^{2}\left(0, T ; L^{2}(\Omega)\right)$ and then $\partial_{t} A \in L^{2}\left(0, T ; L^{2}(\Omega)\right)$ Cf. [2]. We proceed similarly for $E$.
Next, we will prove $A, E \geq 0$. Indeed, multiplying (47) by $A^{-}$and integrating by parts, we get

$$
\begin{equation*}
-\frac{d}{2 d t}\left\|A^{-}\right\|^{2}=D_{a}\left\|\nabla A^{-}\right\|^{2}+l_{-1} \int_{\Omega} \tilde{Q} A^{-} d x+\int_{\Omega} l_{1}\left(\tilde{I}+\nu_{a}\right)\left|A^{-}\right|^{2} d x+\int_{\Omega} S_{a}(x) A^{-} d x \tag{50}
\end{equation*}
$$

The right hand side of the last equatily is nonnegative hence $\left\|A^{-}\right\|^{2}(t) \leq\left\|A_{0}^{-}\right\|^{2}=0$ so $A \geq 0$ a.e. We procced similarly for $E \geq 0$.

Next, we will prove that $A, E \leq 1$. We have

$$
\begin{equation*}
-\partial_{t}(1-A)=-D_{a} \Delta(1-A)+l_{-1} \tilde{Q}-\left(l_{1} \tilde{I}+\nu_{a}\right) A+S_{a}(x) \tag{51}
\end{equation*}
$$

then

$$
\begin{align*}
& \partial_{t}\left\|(1-A)^{-}\right\|^{2}=-D_{a}\left\|\nabla(1-A)^{-}\right\|^{2}+l_{-1} \int_{\Omega} \tilde{Q}(1-A)^{-} d x-\int_{\Omega}\left(l_{1} \tilde{I}+\nu_{a}\right) A(1-A)^{-} d x \\
& +\int_{\Omega} S_{a}(x)(1-A)^{-} d x \tag{52}
\end{align*}
$$

in other words

$$
\begin{align*}
& \partial_{t}\left\|(1-A)^{-}\right\|^{2}=-D_{a}\left\|\nabla(1-A)^{-}\right\|^{2}+l_{-1} \int_{\Omega} \tilde{Q}(1-A)^{-} d x+\int_{\Omega}\left(l_{1} \tilde{I}+\nu_{a}\right)(1-A)(1-A)^{-} d x \\
& -\int_{\Omega}\left(l_{1} \tilde{I}+\nu_{a}\right)(1-A)^{-} d x+\int_{\Omega} S_{a}(x)(1-A)^{-} d x \tag{53}
\end{align*}
$$

then

$$
\begin{align*}
& \partial_{t}\left\|(1-A)^{-}\right\|^{2}=-D_{a}\left\|\nabla(1-A)^{-}\right\|^{2}+l_{-1} \int_{\Omega} \tilde{Q}(1-A)^{-} d x-\int_{\Omega}\left(l_{1} \tilde{I}+\nu_{a}\right)\left|(1-A)^{-}\right|^{2} d x \\
& -\int_{\Omega}\left(l_{1} \tilde{I}+\nu_{a}\right)(1-A)^{-} d x+\int_{\Omega} S_{a}(x)(1-A)^{-} d x \tag{54}
\end{align*}
$$

consequently, by the assumption on $\left\|S_{a}\right\|_{\infty}$

$$
\begin{equation*}
\partial_{t}\left\|(1-A)^{-}\right\|^{2} \leq l_{-1} \int_{\Omega}(1-A)^{-} d x-\nu_{a} \int_{\Omega}(1-A)^{-} d x+\left\|S_{a}(x)\right\|_{\infty} \int_{\Omega}(1-A)^{-} d x \leq 0 \tag{55}
\end{equation*}
$$

hence $A \leq 1$ as soon as $A_{0} \leq 1$. On the other hand

$$
\begin{equation*}
\partial_{t}\left\|(1-E)^{-}\right\|^{2}=k_{-1} \int_{\Omega} \tilde{P}(1-E)^{-} d x-\int_{\Omega}\left(k_{1} \tilde{I}+\nu_{e}\right) E(1-E)^{-} d x+\int_{\Omega} S_{e}(x)(1-E)^{-} d x \tag{56}
\end{equation*}
$$

then

$$
\begin{align*}
& \partial_{t}\left\|(1-E)^{-}\right\|^{2}=k_{-1} \int_{\Omega} \tilde{P}(1-E)^{-} d x+\int_{\Omega}\left(k_{1} \tilde{I}+\nu_{e}\right)(1-E)(1-E)^{-} d x  \tag{57}\\
& -\int_{\Omega}\left(k_{1} \tilde{I}+\nu_{e}\right)(1-E)^{-} d x+\int_{\Omega} S_{e}(x)(1-E)^{-} d x
\end{align*}
$$

hence

$$
\begin{align*}
& \partial_{t}\left\|(1-E)^{-}\right\|^{2}=k_{-1} \int_{\Omega} \tilde{P}(1-E)^{-} d x-\int_{\Omega}\left(k_{1} \tilde{I}+\nu_{e}\right)\left|(1-E)^{-}\right|^{2} d x-\int_{\Omega}\left(k_{1} \tilde{I}+\nu_{e}\right)(1-E)^{-} d x \\
& +\int_{\Omega} S_{e}(x)(1-E)^{-} d x \tag{58}
\end{align*}
$$

By virtue of hypothesis on $\left\|S_{e}(x)\right\|_{\infty}$

$$
\begin{equation*}
\partial_{t}\left\|(1-E)^{-}\right\|^{2} \leq k_{-1} \int_{\Omega}(1-E)^{-} d x-\nu_{e} \int_{\Omega}(1-E)^{-} d x+\left\|S_{e}(x)\right\|_{\infty} \int_{\Omega}(1-E)^{-} d x \leq 0 \tag{59}
\end{equation*}
$$

then $E \leq 1$ whenever $E_{0} \leq 1$.

## 3 Fixed point procedure and proof of the main theorem

We will apply Schauder's fixed point theorem to solve the nonlinear parabolic problem

$$
\left\{\begin{array}{l}
\partial_{t} I=D \Delta I+k_{-1} P-k_{1} E I+l_{-1} Q-l_{1} A I-\nu I, \mathbb{R}^{+} \times \Omega,  \tag{60}\\
\partial_{t} A=D_{a} \Delta A+l_{-1} Q-\left(l_{1} I+\nu_{a}\right) A+S_{a}(x), \mathbb{R}^{+} \times \Omega, \\
\partial_{t} Q=D_{a} \Delta Q+l_{1} A I-\left(l_{-1}+\nu_{a i}\right) Q, \mathbb{R}^{+} \times \Omega, \\
\partial_{t} E(t, x)=D_{e} \Delta E+k_{-1} P-\left(k_{1} I+\nu_{e}\right) E+S_{e}(x), \mathbb{R}^{+} \times \Omega, \\
\partial_{t} P=D_{e} \Delta P+k_{1} E I-\left(k_{-1}+\mu_{e i}\right) P, \mathbb{R}^{+} \times \Omega \\
I(0, x)=I_{0}(x) \geq 0, A(0, x)=A_{0}(x) \geq 0, Q(0, x)=Q_{0}(x) \geq 0, \\
P(0, x)=P_{0}(x) \geq 0, E(0, x)=E_{0}(X) \geq 0, \\
D_{a} \partial_{n} A=D_{a} \partial_{n} Q=0, D_{e} \partial_{n} E=0, D_{e} \partial_{n} P=0, \\
-D \partial_{n} I(x, t)=\frac{K_{c}\left(I(x, t)-i_{b}\right)}{K_{m}+I(x, t)}+T_{e}(I-J) \quad \text { on } \partial \Omega .
\end{array}\right.
$$

Set

$$
W=\left\{\begin{array}{l}
(u, v) \in L^{2}\left(0, T ; H^{1}(\Omega)\right) \times L^{2}\left(0, T ; H^{1}(\Omega)\right)  \tag{61}\\
\partial_{t}(u, v) \in L^{2}\left(0, T ; L^{2}(\Omega)\right) \times L^{2}\left(0, T ; L^{2}(\Omega)\right), \quad 0 \leq u, v \leq 1
\end{array}\right\}
$$

We check that $W$ is a convex, compact subset of $L^{2}\left(0, T ; L^{2}(\Omega)\right)$. Next, define the operator $\mathcal{L}$ on $W$ by

$$
\begin{equation*}
\mathcal{L}(A, E)=\left(A^{\star}, E^{\star}\right) \tag{62}
\end{equation*}
$$

where $\left(A^{\star}, E^{\star}\right)$ is the solution to (44) given by Theorem 4 associated to $(I, P, Q)$ which is the solution to (10) given by Theorem 2 for $\tilde{E}=E$ and $\tilde{A}=A$. It follows from Theorem 2 and 4 that $\mathcal{L}(W) \subset W$.
Our aim is to prove that $\mathcal{L}$ is Lipschitz on $L^{2}\left(0, T ; L^{2}(\Omega)\right) \times L^{2}\left(0, T ; L^{2}(\Omega)\right)$.
Let $\left(A_{i}, E_{i}\right) \in W, i=1,2$ and $\left(I_{i}, P_{i}, Q_{i}\right)$ the solution to (10) associated to $\tilde{A}=A_{i}$ et $\tilde{E}=E_{i}$. Set

$$
\begin{align*}
& \left(A_{i}^{\star}, E_{i}^{\star}\right)=\mathcal{L}\left(A_{i}, E_{i}\right), \quad A^{\star}=A_{1}^{\star}-A_{2}^{\star}, \quad E^{\star}=E_{1}^{\star}-E_{2}^{\star}  \tag{63}\\
& Q=Q_{1}-Q_{2}, \quad P=P_{1}-P_{2}, \quad I=I_{1}-I_{2}, \quad A=A_{1}-A_{2}, \quad E=E_{1}-E_{2}
\end{align*}
$$

By definition of operateur $\mathcal{L}$, we have

$$
\left\{\begin{array}{l}
\partial_{t} A_{i}^{\star}=D_{a} \Delta A_{i}^{\star}+l_{-1} Q_{i}-\left(l_{1} I_{i}+\nu_{a}\right) A_{i}^{\star}+S_{a}(x)  \tag{64}\\
\partial_{t} E_{i}^{\star}(t, x)=D_{e} \Delta E_{i}^{\star}+k_{-1} P_{i}-\left(k_{1} I_{i}+\nu_{e}\right) E_{i}^{\star}+S_{e}(x), \\
A_{i}^{\star}(0, x)=A_{0}(x) \geq 0, \quad E_{i}^{\star}(0, x)=E_{0}(x) \geq 0 \\
D_{a} \partial_{n} A_{i}^{\star}=0, \quad D_{e} \partial_{n} E_{i}^{\star}=0 \text { оп } \partial \Omega
\end{array}\right.
$$

Consequently

$$
\left\{\begin{array}{l}
\partial_{t} A^{\star}=D_{a} \Delta A^{\star}+l_{-1} Q-l_{1}\left(I_{1} A_{1}^{\star}-I_{2} A_{2}^{\star}\right)-\nu_{a} A^{\star}  \tag{65}\\
\partial_{t} E^{\star}=D_{e} \Delta E^{\star}+k_{-1} P-k_{1}\left(I_{1} E_{1}^{\star}-I_{2} E_{2}^{\star}\right)-\nu_{e} E^{\star} \\
A^{\star}(0)=0, \quad E^{\star}(0)=0, \quad \partial_{n} A^{\star}=0 ; \quad \partial_{n} E^{\star}=0 . \partial \Omega
\end{array}\right.
$$

Then

$$
\left\{\begin{array}{l}
\partial_{t} A^{\star}=D_{a} \Delta A^{\star}+l_{-1} Q-l_{1}\left(I_{1} A^{\star}+I A_{2}^{\star}\right)-\nu_{a} A^{\star}  \tag{66}\\
\partial_{t} E^{\star}=D_{e} \Delta E^{\star}+k_{-1} P-k_{1}\left(I_{1} E^{\star}-I E_{2}^{\star}\right)-\nu_{e} E^{\star} \\
A^{\star}(0)=0, \quad E^{\star}(0)=0, \quad \partial_{n} A^{\star}=\partial_{n} E^{\star}=0 ; \quad \partial \Omega
\end{array}\right.
$$

Hence

$$
\begin{align*}
\frac{d}{2 d t}\left\|A^{\star}\right\|^{2} & +D_{a}\left\|\nabla A^{\star}\right\|^{2}+\nu_{a}\left\|A^{\star}\right\|^{2}+l_{1} \int_{\Omega} I_{1}\left|A^{\star}\right|^{2} d x=l_{-1} \int_{\Omega} Q A^{\star} d x-l_{1} \int_{\Omega} I A^{\star} A_{2}^{\star} d x  \tag{67}\\
& \frac{d}{2 d t}\left\|E^{\star}\right\|^{2}+D_{e}\left\|\nabla E^{\star}\right\|^{2}+k_{1} \int_{\Omega} I_{1}\left|E^{\star}\right|^{2}+\nu_{e}\left|E^{\star}\right|^{2} d x=k_{-1} \int_{\Omega} P E^{\star} d x \\
& -k_{1} \int_{\Omega} I E_{2}^{\star} E^{\star} d x \tag{68}
\end{align*}
$$

Consequently

$$
\begin{align*}
& \frac{d}{2 d t}\left(\left\|A^{\star}\right\|^{2}+\left\|E^{\star}\right\|^{2}\right)+D_{a}\left\|\nabla A^{\star}\right\|^{2}+D_{e}\left\|\nabla E^{\star}\right\|^{2}+\nu_{a}\left\|A^{\star}\right\|^{2}+\nu_{e}\left\|E^{\star}\right\|^{2}+l_{1} \int_{\Omega} I_{1}\left|A^{\star}\right|^{2} d x \\
& +k_{1} \int_{\Omega} I_{1}\left|E^{\star}\right|^{2} d x=l_{-1} \int_{\Omega} Q A^{\star} d x+k_{-1} \int_{\Omega} P E^{\star} d x-l_{1} \int_{\Omega} I A^{\star} A_{2}^{\star} d x-k_{1} \int_{\Omega} I E_{2}^{\star} E^{\star} d x \tag{69}
\end{align*}
$$

Finally, using Gronwall's lemma, for arbitrary $T>0$, there exists $C_{T}>0$ depending on $T, l_{1}, l_{-1}, k_{1}, k_{-1}$ such that

$$
\begin{equation*}
\left(\left\|A^{\star}\right\|^{2}+\left\|E^{\star}\right\|^{2}\right)(t) \leq C_{T} \int_{0}^{T}\left(\|I\|^{2}+\|Q\|^{2}+\|P\|^{2}\right)(s) d s, \quad t \in[0, T] \tag{70}
\end{equation*}
$$

On the other hand, for given $A_{i}, E_{i}$

$$
\left\{\begin{array}{l}
\partial_{t}\left(I_{1}-I_{2}\right)=D \Delta\left(I_{1}-I_{2}\right)+k_{-1}\left(P_{1}-P_{2}\right)-k_{1}\left(E_{1} I_{1}-E_{2} I_{2}\right)+l_{-1}\left(Q_{1}-Q_{2}\right)  \tag{71}\\
+l_{1}\left(A_{1} I_{1}-A_{2} I_{2}\right)-\nu\left(I_{1}-I_{2}\right), \\
\partial_{t}\left(Q_{1}-Q_{2}\right)=D_{a} \Delta\left(Q_{1}-Q_{2}\right)+l_{1}\left(A_{1} I_{1}-A_{2} I_{2}\right)-\left(l_{-1}+\nu_{a i}\right)\left(Q_{1}-Q_{2}\right) \\
\partial_{t}\left(P_{1}-P_{2}\right)=D_{e} \Delta\left(P_{1}-P_{2}\right)+k_{1}\left(E_{1} I_{1}-E_{2} I_{2}\right)-\left(k_{-1}+\mu_{e i}\right)\left(P_{1}-P_{2}\right) \\
\left(P_{1}-P_{2}\right)(0, x)=0,\left(Q_{1}-Q_{2}\right)(0, x)=0,\left(I_{1}-I_{2}\right)(0, x)=0 \\
\partial_{n}\left(Q_{1}-Q_{2}\right)=0, \partial_{n}\left(P_{1}-P_{2}\right)=0, \\
-D \partial_{n}\left(I_{1}-I_{2}\right)=\frac{K_{c}\left(K_{m}+i_{b}\right)\left(I_{1}-I_{2}\right)}{\left(K_{m}+I_{1}\right)\left(K_{m}+I_{2}\right)}+T_{e}\left(I_{1}-I_{2}\right)
\end{array}\right.
$$

Setting $Q=Q_{1}-Q_{2}, P=P_{1}-P_{2}, I=I_{1}-I_{2}, A=A_{1}-A_{2}, E=E_{1}-E_{2}$ we get

$$
\left\{\begin{array}{l}
\partial_{t} I=D \Delta I+k_{-1} P-k_{1}\left(E I_{1}+E_{2} I\right)+l_{-1} Q+l_{1}\left(A I_{1}+A_{2} I\right)-\nu I  \tag{72}\\
\partial_{t} Q=D_{a} \Delta Q+l_{1}\left(A I_{1}+A_{2} I\right)-\left(l_{-1}+\nu_{a i}\right) Q \\
\partial_{t} P=D_{e} \Delta P+k_{1}\left(E I_{1}+E_{2} I\right)-\left(k_{-1}+\mu_{e i}\right) P \\
P(0, x)=0, Q(0, x)=0, I(0, x)=0 \\
\partial_{n} Q=0, \quad \partial_{n} P=0,-D \partial_{n} I=\frac{K_{c}\left(K_{m}+i_{b}\right) I}{\left(K_{m}+I_{1}\right)\left(K_{m}+I_{2}\right)}+T_{e} I
\end{array}\right.
$$

Hence

$$
\begin{align*}
& \frac{d}{2 d t}\left(\|I\|^{2}+\|Q\|^{2}+\|P\|^{2}\right)+D\|\nabla I\|^{2}+D_{a}\|\nabla Q\|^{2}+D_{e}\|\nabla P\|^{2}+\nu\|I\|^{2}+T_{e}\|I\|_{\partial \Omega}^{2} \\
& +\left(l_{-1}+\nu_{a i}\right)\|Q\|^{2}+\left(k_{-1}+\mu_{e i}\right)\|P\|^{2}+\int_{\partial \Omega} \frac{K_{c}\left(K_{m}+i_{b}\right)|I|^{2}}{\left(K_{m}+I_{1}\right)\left(K_{m}+I_{2}\right)} d \sigma+k_{1} \int_{\Omega} E_{2}|I|^{2} d x \\
& +l_{1} \int_{\Omega} A_{2}|I|^{2} d x=k_{-1} \int_{\Omega} P I d x+l_{-1} \int_{\Omega} Q I d x+l_{1} \int_{\Omega} A I I_{1} d x+l_{1} \int_{\Omega} A Q I_{1} d x  \tag{73}\\
& +l_{1} \int_{\Omega} A_{2} I Q d x+k_{1} \int_{\Omega} E P I_{1} d x+k_{1} \int_{\Omega} E_{2} I P d x-k_{1} \int_{\Omega} E I I_{1} d x .
\end{align*}
$$

Then, it follows from Gronwall's lemma that for $T>0$, there exsits $C_{T}^{\prime}>0$ depending on $T$ and the constants occuring in the system such that

$$
\begin{equation*}
\left(\|I\|^{2}+\|Q\|^{2}+\|P\|^{2}\right)(t) \leq C_{T}^{\prime} \int_{0}^{T}\left(\|A\|^{2}+\|E\|^{2}\right) d s, \quad t \in[0, T] \tag{74}
\end{equation*}
$$

Finally, we conclude from (70) and (74) that

$$
\begin{equation*}
\left\|\mathcal{L}\left(A_{1}, E_{1}\right)-\mathcal{L}\left(A_{2}, E_{2}\right)\right\|_{L^{2}\left(0, T ; L^{2}(\Omega)\right)} \leq C_{T}\left\|\left(A_{1}, E_{1}\right)-\left(A_{2}, E_{2}\right)\right\|_{L^{2}\left(0, T ; L^{2}(\Omega)\right)} \tag{75}
\end{equation*}
$$

for some positive constant $C_{T}>0$ depending only on $T>0$ and the constants occuring in the system. It derives from Schauder fixed point theorem that operatoe $\mathcal{L}$ has a fixed point and the proof of Theorem 1 follows.
Acknowledgements. I am very grateful to Professors H.A. Levine and M. Nilsen Hamilton for their suggestions and to Professor K. Hamdache for his hospitality.

## References

[1] K. Boushaba, H.A. Levine, M. Nilsen Hamilton, A Mathematical Feasibility Argument for the use of Aptamers in Chemotherapy and Imaging, Marrakesh International Conference and Workshop on Mathematical Biology, Marrakesh, January, 3-8, 2008 (to appear, corresponding author boushaba@iastate.edu).
[2] T. Cazenave, A. Haraux, Introduction aux problèmes d'évolution semi-linéaires Ellipses 1990.
[3] J.L. Lions Quelques méthodes de résolution des problèmes aux limites non linéaires, Dunod Gauthier-Villars (1969).
[4] A. Pazy, Semigroups of Linear Operators and Applications to Partial Differential Equations, Springer, (1983).
[5] R.E. Showalter,, Monotone operators in Banach space and Nonlinear partial differential equations,, Math. Surveys and Monographs 49 AMS, 1997.
[6] E. Zeidler Nonlinear Functional Analysis and its ApplicationsIIB Springer-Verlag (1990).
[7] Davies, R., J. Budworth, J. Riley, R. Snowden, A. Gescher, and T.W. Gant, Regulation of P-glycoprotein 1 and 2 gene expression and protein activity in two MCF7/Dox cell line subclones, Br. J. Cancer, 73, 307-15 (1996).
[8] Chen, G.K., N.J. Lacayo, G.E. Duran, Y. Wang, C.D. Bangs, S. Rea, M. Kovacs, A.M. Cherry, J.M. Brown, and B.I. Sikic, Preferential expression of a mutant allele of the amplified MDR1 (ABCB1) gene in drug-resistant variants of a human sarcoma,Genes Chromosomes Cancer, 34, 372-83, (2002).
[9] Cahill, A., T. C. Jenkins and I. N. H. White, Metabolism of 3-amino-1,2,4-benzotriazine-1,4-dioxide (SR 4233) by purified DT-diaporase unde aerobic and anaerobic conditions, Biochem. Pharmacol, 45, 321329, (1993).
[10] Ferrell, J. E.Trends, Biochem. Sci. 21 p. 460 466, (1996). Fukushima, M., A. Fujioka, J. Uchida, F. Nakagawa, and T. Takechi, Thymidylate synthase (TS) and ribonu-cleotide reductase (RNR) may be involved in acquired resistance to 5-fluorouracil (5-FU) in human cancer xenografts in vivo, Eur J Cancer, 37,1681-7, (2001).
[11] Hicks, K. O., Y. Fleming, B. G. Siim, C. J. Koch and W. R. Wilson, Extravascular diffusion of tirapazamine: effect of metabolic consumption assessed using the multicellular layer model, Int. J. Radiat. Oncol. Biol. Phy., 42, 641649, (1998).
[12] Paul, C. P., Good, D. P., winer, I., Engelke, D. R., Effective expression of small interfering RNA in human cells, Nature Biotechnology, 20, 505-508 (2002).
[13] Shibata, J., K. Aiba, H. Shibata, S. Minowa, and N. Horikoshi, Detection and quantitation of thymidylate synthase $m R N A$ in human colon adenocarcinoma cell line resistant to 5-fluorouracil by competitive PCR, Anticancer Res, 18, 1457-63, (1998).
[14] Vrana, J.A., C.K. Bieszczad, E.S. Cleaveland, Y. Ma, J.P. Park, T.K. Mohandas, and R.W. Craig, An MCL1- overexpressing Burkitt lymphoma subline exhibits enhanced survival on exposure to serum deprivation, topoisomerase inhibitors, or staurosporine but remains sensitive to 1-beta-D-arabinofuranosylcytosine, Cancer Res, 62, 892900, (2002).
[15] Tran, T. T., Mittal, A., Aldinger, T., Polli, J. W., Ayrton, A., Ellens, H., and J. Bentz, The elementary mass action rate constants for P-gp transport for a confluent monolayer of MDCKII-hMDR1 cells, Byophysical Journal, 88, p.715-738, (2005).


[^0]:    ${ }^{*}$ CMAP, Ecole Polytechnique, CNRS. 91128 Palaiseau Cedex \& USTHB, Faculté des Mathématiques, BP 32 Bab Ezzouar, 16111 Alger. Email: naima.aissa@polytechnique.edu

