# Mathematical models based on free boundary problems 

## Citation for published version (APA):

Eijndhoven, van, S. J. L. (1990). Mathematical models based on free boundary problems. (Opleiding wiskunde voor de industrie Eindhoven : student report; Vol. 9001). Eindhoven University of Technology.

## Document status and date:

Published: 01/01/1990

## Document Version:

Publisher's PDF, also known as Version of Record (includes final page, issue and volume numbers)

## Please check the document version of this publication:

- A submitted manuscript is the version of the article upon submission and before peer-review. There can be important differences between the submitted version and the official published version of record. People interested in the research are advised to contact the author for the final version of the publication, or visit the DOI to the publisher's website.
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## REPORT 90-01

MATHEMATICAL MODELS BASED ON FREE BOUNDARY PROBLEMS

## S.J.L. van Eijndhoven

## Opleiding <br> Wiskunde voor de Industrie Eindhoven



# MATHEMATICAL MODELS BASED ON FREE BOUNDARY PROBLEMS 

S.J.L. van Eijndhoven

This report is based on a course given by Prof. A. Fasano at the Eindhoven University of Technology from December 14 until December 22, 1989.


#### Abstract

Foreword A free boundary problem for a partial (or ordinary) differential equation is characterized by the fact that the boundary of the domain in which the differential equation is to be solved is at least partly unknown. So a free boundary problem consists of determining both the solution and the unknown boundary. In case of a fixed boundary problem, when the boundary is completely described, the problem is well-posed given a set of boundary data. If part of the boundary is unknown these boundary data are insufficient for the well-posedness of the problem and additional conditions must be specified. Conditions on the free boundary are naturally called free boundary conditions. According to the above definition the term free is synonymous to unknown. Sometimes, however, a distinction is made between free and moving boundary problems referring to those cases in which the unknown boundary stays at rest or moves. In the terminology used here, there is no such difference and a moving boundary is a free boundary only if its motion is not prescribed.

In this report five examples of free boundary problems are discussed: the obstacle problem, the dam problem, the Stefan problem, the oxygen diffusion consumption problem and the flow problem of a Bingham fluid between two fixed plates. In five appendices some mathematical prerequisites are gathered.


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## CHAPTER I

## THE OBSTACLE PROBLEM

## 1. One-dimensional case

Take a rubber band and stretch it between two point $A$ and $B$.

The equilibrium configuration will be the segment $A B$. The governing differential equation is given by

$$
u^{\prime \prime}=0
$$

with boundary conditions

$$
u(a)=u(b)=0
$$

Suppose we impose the constraint

$$
u \geq \psi
$$

where $\psi$ is a $C^{1}$-function such that $\psi(a)=\psi(b)<0$ and $\psi>0$ on some interval $(c, d)$.


We have to solve the following
(i) $\quad u^{\prime \prime}=0$ on $\{u>\psi\}$
(ii) $u(0)=u(a)=0$
(iii) $\quad u=\psi$ over $\partial\{u=\psi\}$, i.e. $u\left(x_{1}\right)=\psi\left(x_{1}\right)$ and $u\left(x_{2}\right)=\psi\left(x_{2}\right)$.

The conditions (1.1.ii) and (1.1.iii) are not sufficient to determine $u, x_{1}$ and $x_{2}$. In addition we have to add

$$
\begin{equation*}
u^{\prime}\left(x_{1}\right)=\psi^{\prime}\left(x_{1}\right) \text { and } u^{\prime}\left(x_{2}\right)=\psi^{\prime}\left(x_{2}\right) \tag{iv}
\end{equation*}
$$

Conditions (I.1.1.iii and iv) constitute the free boundary conditions.

Remark. Although the differential equation is linear and the conditions at $x=a$ and $x=b$ are
linear, the problem itself is nonlinear. The nonlinearity is hidden in the free boundary conditions which involve the solution in some implicit way.

Now suppose for a moment that $\psi$ is a $C^{2}$-function. Then we see that $\psi^{\prime \prime} \leq 0$, and hence on $(a, b)$ we have the following inequalities
(I.1.2) $-u^{\prime \prime} \geq 0, u-\psi \geq 0, \quad u^{\prime \prime}(u-\psi)=0$.

The problem is put in its complementarity form. In order to generalize this form to $C^{1}$-functions $\psi$ we need a weak interpretation of (I.1.2) and seek for solutions $u$ in a wider class. To this end we use the Sobolev space $W_{0}^{2,1}([a, b])$ consisting of all absolutely continuous functions $w$ such that $\int_{a}^{b}\left|w^{\prime}(x)\right|^{2} d x<\infty$ and $w(a)=w(b)=0$. Then by $u^{\prime \prime}$ we mean the second distributional derivative of $u$ in $D^{\prime}((a, b))$. For a distribution $F \in D^{\prime}((a, b))$ we write $F \geq 0$ if $F(\phi) \geq 0$ for all positive $\phi \in D((a, b))$ (cf. Appendix C).
We remark that for $u \in W_{0}^{2,1}([a, b]), u^{\prime \prime} \geq 0$, is equivalent with

$$
\begin{equation*}
\int_{a}^{b} u^{\prime}(x) \phi^{\prime}(x) d x \geq 0, \forall \phi \in D((a, b)), \phi \geq 0 \tag{I.1.3}
\end{equation*}
$$

Further, since $D((a, b))$ is dense in $W_{0}^{2,1}([a, b]),\left(^{*}\right)$ is equivalent with

$$
\begin{equation*}
\int_{a}^{b} u^{\prime}(x) v^{\prime}(x) d x \geq 0, \forall v \in W_{0}^{2,1}([a, b]), v \geq 0 \tag{I.1.4}
\end{equation*}
$$

In the new formulation the free boundary does not appear explicitly.
In its turn the complementarity problem (I.1.2) is equivalent to the following variational inequality. Define the convex and closed set $K:=\left\{v \in W_{0}^{2,1}([a, b]) \mid v \geq \psi\right\}$ and look for $u \in K$ such that

$$
\begin{equation*}
\int_{a}^{b} u^{\prime}(x)\left(u^{\prime}(x)-v^{\prime}(x)\right) d x \leq 0 \quad, \quad \forall v \in K \tag{I.1.5}
\end{equation*}
$$

To show the equivalence we proceed as follows.

Suppose $u$ satisfies (I.1.5). Let $\phi \in D((a, b))$ with $\phi \geq 0$. Then $u+\phi \in K$ and

$$
\int_{a}^{b} u^{\prime}(x)\left(u^{\prime}(x)-\left(u^{\prime}(x)+\phi^{\prime}(x)\right)\right) d x \leq 0
$$

whence $u^{\prime \prime} \geq 0$ in weak sense. Further, let $C$ be a compact subset of $\{u>\psi\}$. Then, $u$ and $\psi$ being continuous, there exists $\delta>0$ such that $u(x)-\Psi(x) \geq \delta$ for all $x \in C$. So for all $C^{\infty}$-functions $\phi$ with support in $C$ there exists $t_{\phi}>0$ such that

$$
\forall_{t \in\left[-t_{*}, t_{*}\right]}, u+t \phi \in K
$$

whence

$$
\forall_{t \in\left[-t_{+}, t_{+}\right]}: t \int_{a}^{b} u^{\prime}(x) \phi^{\prime}(x) d x \geq 0
$$

It follows that $\int_{a}^{b} u^{\prime}(x) \phi^{\prime}(x) d x=0$ for all $\phi \in D((a, b))$ with support in $\{u>\psi\}$. Thus we conclude that $u^{\prime \prime}=0$ (weakly) on $\{u>\psi\}$.

Conversely, suppose $u$ satisfies (1.1.2) in its weak interpretation. Then for all $v \in K, u-v \leq u-\psi=0$, on $\{u=\psi\}$ and hence by (I.1.4)

$$
\int_{\{u=\psi\}} u^{\prime}(x)\left(u^{\prime}(x)-v^{\prime}(x)\right) d x \leq 0
$$

Moreover, on $\{u>\psi\}$ we have $u^{\prime \prime}=0$ weakly, i.e.

$$
\forall_{0 \in D((a, b))}: \int_{\{u>\psi\}} u^{\prime}(x) \phi^{\prime}(x) d x=0
$$

which yields, again because $D((a, b))$ is dense in $W_{0}^{2,1}([a, b])$,

$$
\forall_{v \in K}: \int_{\{u>\psi\}} u^{\prime}(x)\left(u^{\prime}(x)-v^{\prime}(x)\right) d x=0 .
$$

Uniqueness of $u$ can be established straightforwardly from (1.1.5). Indeed suppose $u_{1}$ and $u_{2}$ satisfy (1.1.5). Then

$$
\int_{a}^{b} u_{1}^{\prime}(x)\left(u_{1}^{\prime}(x)-u_{2}^{\prime}(x)\right) d x \leq 0
$$

ánd

$$
\int_{a}^{b} u_{2}^{\prime}(x)\left(u_{2}^{\prime}(x)-u_{1}^{\prime}(x)\right) d x \leq 0
$$

so that

$$
\int_{a}^{b}\left(u_{1}^{\prime}(x)-u_{2}^{\prime}(x)\right)^{2} d x \leq 0
$$

Hence $u_{1}^{\prime}=u_{2}^{\prime}$. Since $u_{j}(a)=u_{j}(b)=0, j=1,2$, the result follows.

Finally, the variational problem (I.1.5) is equivalent with the problem of determining $u \in K$ for which the quadratic form

$$
\begin{equation*}
J(w)=\int_{a}^{b}\left|w^{\prime}(x)\right|^{2} d x, w \in K \tag{I.1.6}
\end{equation*}
$$

is minimal. To see this, observe first that

$$
\begin{align*}
J(\bar{w})-J(w)= & \int_{a}^{b} \bar{w}^{\prime}(x)\left(\bar{w}^{\prime}(x)-w^{\prime}(x)\right) d x  \tag{*}\\
& +\int_{a}^{b} w^{\prime}(x)\left(\bar{w}^{\prime}(x)-w^{\prime}(x)\right) d x .
\end{align*}
$$

Let $u \in K$ satisfy (1.1.5).
Then for $\bar{w}, w \in K$ with $\tilde{w} \neq u$

$$
\int_{a}^{b} \tilde{w}^{\prime}(x)\left(\bar{w}^{\prime}(x)-w^{\prime}(x)\right) d x \geq 0
$$

Hence both summands on the right hand side of $\left(^{*}\right)$ are positive if we take $\tilde{w}=v \in K, v \neq u$, and $\boldsymbol{w}=u$.
It follows that $J(v) \geq J(u)$ for all $v \in K$.
Conversely, if $u$ minimizes $J$ over the convex set $K$, then

$$
\left.\frac{d}{d t} J(u+t(v-u))\right|_{t=0} \geq 0
$$

from which (I.1.5) results.
Now $J(w)^{1 / 2}$ is the norm of $w$ in $W_{0}^{2,1}([a, b])$ corresponding to the inner product

$$
\left(w_{1}, w_{2}\right)=\int_{a}^{b} w_{1}^{\prime}(x) w_{2}^{\prime}(x) d x
$$

Since the set $K$ is closed and convex in $W_{0}^{2,1}([a, b])$ a classical result from Hilbert space theory says that there exist a unique $u \in K$ such that

$$
J(u)=\min \{J(w) \mid w \in K\}=\operatorname{dist}^{2}(0, K) .
$$

Thus both existence and uniqueness of $u$ is established.

## 2. Two-dimensional case

As we have seen the obstacle problem in one dimension has a simple solution. In two dimensions the problem becomes less trivial. In this case we consider a membrane stretched over a profile, e.g. the boundary of a given domain $\Omega$ at which $u=0$, and with the same constraint that $u \geq \psi$ where now $\psi<0$ on $\partial \Omega$. The classical formulation of this problem is rather complicated.


Look for a function $u$ on $\Omega_{0}=\Omega \backslash I$ such that

$$
u \in C^{2}\left(\Omega_{0}\right) \cap C^{1}\left(\bar{\Omega}_{0}\right)
$$

and such that $u$ satisfies the following relations

$$
\Delta u=0 \text { in } \Omega_{0}
$$

(ii) $u=0$ on $\partial \Omega$
(iii) $\quad u=\psi$ on $\gamma$
(iv) $\frac{\partial u}{\partial n}=\frac{\partial \psi}{\partial n}$ on $\gamma$.

Here $I$ denotes the coincidence set $\{u=\psi\}$ and $\gamma$ its $C^{1}$-boundary. Further, $\frac{\partial}{\partial n}$ denotes the normal derivative at $\gamma$, and the relation

$$
\frac{\partial u}{\partial n}=\frac{\partial \psi}{\partial n}
$$

expresses that the membrane is tangent to the obstacle. Under the assumption that $\psi \in C^{2}(\Omega)$ we can prove again that the problem (I.2.1) is equivalent to the complementarity problem valid in the whole domain $\Omega$
(1.2.2) $-\Delta u \geq 0, u-\psi \geq 0,-\Delta u(u-\psi)=0$.

In order to interpret problem (1.2.2) for a more general class of functions $\psi$ we take $\psi \in W^{2.1}(\Omega)$ and seek for solutions $u \in W_{0}^{2.1}(\Omega)$, i.e. the closure of $D(\Omega)$ in $W^{2.1}(\Omega)$. Thus (I.2.2) is given the following weak interpretation.

$$
\begin{equation*}
-\Delta u \geq 0 \text { means } \forall_{\varphi} \in D(\Omega), \phi \geq 0: \int_{\Omega}(\nabla u \cdot \nabla \phi) d x \geq 0 \tag{I.2.3}
\end{equation*}
$$

$\Delta u(u-\psi)=0$ means $\int_{\Omega} \nabla u \cdot \nabla \phi=0$ for $\phi \in D(\Omega)$ with support contained in $\{u>\psi\}$.
(Observe that we have applied Green's first identity.)
From the weakly interpreted complementarity form we arrive at the following variational inequality: Introduce the convex and closed set $K \subset W_{0}^{2.1}(\Omega)$ by

$$
K=\left\{v \in W_{0}^{2.1}(\Omega) \mid v \geq \psi\right\}
$$

Then (1.2.2) is equivalent with the problem of searching $u \in K$ for which

$$
\begin{equation*}
\int_{\Omega} \nabla u \cdot(\nabla u-\nabla v) d x \leq 0, \quad \forall v \in K \tag{1.2.4}
\end{equation*}
$$

It can be proved that the norm in $W_{0}^{2.1}(\Omega)$ is equivalent with

$$
\|u\|_{2,1}^{0}=\left(\int_{\Omega}|\nabla u|^{2} d x\right)^{1 / 2} .
$$

Now (I.2.4) is, in its turn, equivalent with searching $u \in K$ that minimizes the functional

$$
J(w)=\int_{\Omega}|\nabla w|^{2} d x=\left(\|w\|_{2,1}^{0}\right)^{2}, \quad w \in K .
$$

Since $K$ is closed and convex such a unique $u \in K$ exists.

## 3. A comparable problem

In comparison with the obstacle problem we present the problem of refining metal surfaces using electrolytic processes. The workpiece and the tool are respectively the anode and the cathode in an electrolytic circuit. The current flows
 in the electrolytic solution under the action of a potential difference between the electrodes and causes dissolution of the anode surface which allows for micrometric machining of the workpiece.
Clearly the free boundary in this problem is the anode surface. Keeping the potential difference between the electrodes constant we can use a non dimensional potential $\phi$ with value $\phi=0$ on the cathode surface $\Gamma$ and $\phi=1$ on the moving anode surface $\gamma_{i}$. The function $\phi$ is harmonic in the region $\Omega_{t}$ occupied by the electrolyte. So we have the following formulation of the problem.
Given the cathode surface $\Gamma$ and the initial configuration $\gamma_{0}$ of the anode surface, find the pair ( $\gamma_{l}, \phi$ ) such that

$$
\begin{equation*}
\Delta \phi=0 \text { in } \Omega_{t}, t>0 \tag{1.3.1}
\end{equation*}
$$

(ii) $\left.\quad \phi\right|_{\Gamma}=0, t>0$
(iii) $\left.\quad \phi\right|_{\gamma_{t}}=1, t>0$
(iv) $\quad \gamma_{i=0}=\gamma_{0}$.

Of course, we need an additional free boundary condition relating the local dissolution rate, i.e. the normal component $v_{n}$ of the velocity of the anode surface, to the value of the electric field. We take
(v) $\quad v_{n}=f\left[\frac{\partial \phi}{\partial n}\right]$ on $\gamma_{t}$.

A realistic form of $f$ exhibits the presence of a threshold current below which no or very little machining occurs.

(I.3.2)
(i)

$$
\Delta \phi=0 \text { in } \Omega
$$

(ii) $\left.\quad \phi\right|_{\Gamma}=0$
(iii) $\left.\quad \phi\right|_{\gamma=1}$
(iv) $\left.\quad \frac{\partial \phi}{\partial n}\right|_{\gamma}=\lambda$.

The free boundary will be steady if

$$
\frac{\partial \phi}{\partial n} \leq \lambda .
$$

Therefore the threshold current model has a limiting steady state ( $\phi, \gamma$ ) satisfying the equations

Problem (I.3.2) can be interpreted as a membrane equilibrium problem: Given the profile $\Gamma$ on $\phi=0$ look for a profile $\gamma$ on $\phi=1$ such that the slope of the membrane complies with the condition (iv).


This problem is much harder to solve than the previous one since it cannot be made variational.

## CHAPTER II

## THE DAM PROBLEM

In this chapter we consider the steady filtration of a fluid through a porous dam. The problem is sketched in the following figure.


We assume complete saturation, there is no capillarity.

Recall that a fluid flow in a (saturated) porous medium is governed by Darcy's law (cf. Appendix B)

$$
\vec{v}=-k \nabla(p+\rho g y)
$$

where $\vec{v}$ is the volumetric velocity,
$k$ the hydraulic conductivity,
$\rho$ the fluid density,
$g$ gravity, and
y upwardly directed vertical coordinate.
Incompressibility of the fluid implies that $\operatorname{div} \vec{v}=0$ and so that $\Delta p=0$ on the region $\{p>0$ \}.

We have to find $p$ in the saturated region and the free boundary $\Gamma$ which we describe by $y=\phi(x), 0 \leq x \leq c$. First we derive the conditions on the fixed boundaries $\Gamma_{0}, \Gamma_{1}, \Gamma_{2}$ and $\Gamma_{3}$. We normalize such that $k=\rho g=1$.

$$
\begin{array}{ll}
\text { on } \Gamma_{1}: & p(0, y)=y_{1}-y, \\
\text { on } \Gamma_{2}: & 0 \leq y \leq y_{1}, \\
\text { on } \Gamma_{0}: & p(c, y)=y_{2}-y, 0 \leq y \leq y_{2}, \\
, ~ & y_{2} \leq y \leq \phi(c) .
\end{array}
$$

For the boundary condition on $\Gamma_{3}$ we apply Darcy's law

$$
v_{x}=-\frac{\partial p}{\partial x} \quad, \quad v_{y}=-\left[\frac{\partial p}{\partial y}+1\right]=0
$$

and so

$$
\text { on } \Gamma_{3}: \frac{\partial p}{\partial y}(x, 0)=-1, \quad 0 \leq x \leq c .
$$

On the free boundary $\Gamma$ we have

$$
p(x, \phi(x))=0, \quad 0 \leq x \leq c
$$

and the additional free boundary condition

$$
\left.\frac{\partial(p+y)}{\partial n}\right|_{\Gamma}=0
$$

Explanation: At $\Gamma$ the fluid flows tangent to the curve $\Gamma$. It follows that $\vec{v} \cdot \vec{n}=0$ and hence by Darcy's law

$$
\nabla(p+y) \cdot \vec{n}=0 \quad \text { on } \Gamma \text {. }
$$

We have the following observations.
(II.1) Because of the maximum principle $p$ cannot be negative in $\Omega$ since otherwise $p$ would be strictly negative in a point of $\Gamma_{3}$ which yields a contradiction. In fact even the stronger assertion that $p>0$ in $\Omega$ is valid.
(II.2) We have $\frac{\partial p}{\partial y}>-1$ in $\Omega$.


Indeed, the curve $\Gamma$ is a level line of $p$ since

$$
p(x, \phi(x))=0, \quad 0 \leq x \leq c
$$

So $\nabla p \mid \Gamma$ is normal to $\Gamma$ and hence

$$
\vec{v} \cdot \nabla p=0 \text { on } \Gamma .
$$

$$
\nabla(p+y) \cdot \nabla p=0 \text { on } \Gamma .
$$

This means that

$$
\left[\frac{\partial p}{\partial x}\right]^{2}+\left[\frac{\partial p}{\partial y}+\frac{1}{2}\right]^{2}=\frac{1}{4} \text { on } \Gamma
$$

Since $\Delta\left[\frac{\partial p}{\partial y}\right]=0$ the result follows from the maximum principle.
(II.3) $\phi(x)$ is decreasing.

* Assume


Applying Gauss' divergence theorem to $G$ yields

$$
0=\int_{\partial G}(\nabla p \cdot \vec{n}) d s=\int_{K_{1}} \nabla p \cdot \vec{n} d s-\int_{K_{2}} \nabla p \cdot \vec{n} d s .
$$

Now

$$
\nabla p \cdot \vec{n}=-|\nabla p| \text { on } K_{1}:(t, \phi(t)), x_{1} \leq t \leq x_{2}
$$

and

$$
\nabla p \cdot \vec{n}=-\frac{\partial p}{\partial y} \text { on } K_{2}:\left(t, \phi\left(x_{1}\right)\right), x_{1} \leq t \leq x_{2} .
$$

It follows that

$$
\int_{K_{1}}|\nabla p| d s=\int_{K_{2}} \frac{\partial p}{\partial y} d s<0
$$

which yields a contradiction.

* Assume


Then with Gauss' divergence theorem

$$
\int_{K_{1}}|\nabla p| d s=\int_{K_{2}} \frac{\partial p}{\partial y} d s+\int_{K_{3}} \frac{\partial p}{\partial x} d s<0
$$

(observe that $v_{x}>0$ on $K_{3}$ )
and again we arrive at a contradiction.

* Assume $\phi$ is constant on $\left(x_{1}, x_{2}\right)$ with $\phi\left(x_{1}\right)=y^{\prime}$.

Since $\nabla p$ is normal to the curve $(t, \phi(t))$ we have

$$
\begin{aligned}
& \frac{\partial p}{\partial x}=0 \text { on } K_{1} \\
& \frac{\partial p}{\partial y}=-1 \text { on } K_{1}
\end{aligned}
$$

and of course also

$$
p=0 \text { on } \mathrm{K}_{1} .
$$

The solution of the problem

$$
\begin{aligned}
& \Delta \bar{p}=0 \text { in }\left(x_{1}, x_{2}\right) \times\left(0, y^{\prime}\right) \\
& \frac{\partial \tilde{p}}{\partial x}=0 \quad, \quad \frac{\partial \tilde{p}}{\partial y}=-1 \text { and } \tilde{p}=0 \text { on } K_{1} \\
& \frac{\partial \tilde{p}}{\partial y}=-1 \text { on }\left(x_{1}, x_{2}\right)
\end{aligned}
$$

is given by

$$
\tilde{p}(x, y)=y^{\prime}-y .
$$

This solution can be extended in a unique way to a harmonic function in the rectangle $(0, c) \times\left(0, y^{\prime}\right)$. Since the looked for solution $p$ agrees with $\tilde{p}$ on $\left(x_{1}, x_{2}\right) \times\left(0, y^{\prime}\right)$ and is harmonic in $\Omega$ we must have $p=\tilde{p}$ on $\Omega \cap(0, c) \times\left(0, y^{\prime}\right)$. Thus we get a contradiction with the boundary conditions on $\Gamma_{1}$ or on $\Gamma_{2}$.

$$
\begin{equation*}
\frac{\partial p}{\partial x}<0 \text { in } \Omega . \tag{II.4}
\end{equation*}
$$

Indeed on $\Gamma_{1}$ and $\Gamma_{2}$ we have $\frac{\partial p}{\partial y}=-1$ and so $\frac{\partial}{\partial x}\left[\frac{\partial p}{\partial x}\right]=0$ on $\Gamma_{1}$ and $\Gamma_{2}$. The fluid flows tangent to $\Gamma_{3}$, so $\frac{\partial}{\partial y}\left[\frac{\partial p}{\partial x}\right]=0$ on $\Gamma_{3}$. Moreover $\frac{\partial p}{\partial x}<0$ on $\Gamma_{0}$, because $v_{x}>0$ on $\Gamma_{0}$.
It is clear that $\frac{\partial p}{\partial x} \leq 0$ on $\Gamma$.
Now apply the maximum principle:

$$
\begin{aligned}
& \frac{\partial}{\partial x}\left(\frac{\partial p}{\partial x}\right)<0 \quad \Delta\left(\frac{\partial p}{\partial x}\right)=0 \quad \begin{array}{l}
\Gamma: \frac{\partial p}{\partial x} \leqslant 0 \\
\Gamma_{c}: \frac{\partial p}{\partial x}<c \\
\Gamma_{2}: \frac{\partial}{\partial x}\left(\frac{\partial p}{\partial x}\right)<0
\end{array} \\
& \Gamma_{3}: \frac{\partial}{\partial y}\left(\frac{\partial p}{\partial x}\right)=c
\end{aligned}
$$

(II.5)


Consider the region $G$ as indicated in the figure.

Since $\operatorname{div} \vec{v}=0$, Gauss' divergence theorem yields

$$
\int_{\partial G}(\vec{v} \cdot \vec{n}) d s=0 .
$$

Now $(\vec{v} \cdot \vec{n})=0$ on $\Gamma$ and $\Gamma_{3}$, what yields

$$
\int_{0}^{\left(x_{1}\right)} v_{x}\left(x_{1}, y\right) d y=\int_{0}^{\phi\left(x_{2}\right)} v_{x}\left(x_{2}, y\right) d y=q \quad \text { (same discharges) }
$$

So $q c=\int_{0}^{c}\left(\int_{0}^{\Theta(x)}\left(-\frac{\partial p}{\partial x}(\xi, \eta)\right) d \eta\right) d \xi=\frac{1}{2}\left(y_{1}^{2}-y_{2}^{2}\right)$.
(II.6) The curve $\Gamma_{0}$ really exists.

We have $\vec{v}=(u, v), \vec{v}=-\nabla(p+y), \operatorname{div}(\vec{v})=0$.
So $u=-\frac{\partial p}{\partial x}$ and $v=-\frac{\partial p}{\partial y}-1$.

From $\operatorname{div} \vec{v}=0$ we obtain

$$
\left.\begin{array}{l}
\frac{\partial u}{\partial y}=\frac{\partial v}{\partial x} \\
\frac{\partial u}{\partial x}=-\frac{\partial v}{\partial y}
\end{array}\right\} \text { Cauchy-Riemann equations. }
$$

Put $z=x+i y$ and $f(z)=u(x, y)+i v(x, y)$. Then $f$ is holomorphic in $\Omega \subset \mathbb{C}$ and $f$ maps the region $\Omega$ in the $(x, y)$-plane onto the region $\Omega^{*}$ in the ( $u, v$ )-plane with boundary segments $\Gamma_{0}^{*}, \Gamma_{1}^{*}$, $\Gamma_{2}^{*}, \Gamma_{3}^{*}$ and $\Gamma_{4}^{*}$.
In the latter plane we have



## Explanation:

As we have seen on $\Gamma$ we have $\left[\frac{\partial p}{\partial x}\right]^{2}+\left[\frac{\partial p}{\partial y}+\frac{1}{2}\right]^{2}=\frac{1}{4}$ implying that

$$
u^{2}+\left(v+\frac{1}{2}\right)^{2}=\frac{1}{4}
$$

So $\Gamma^{*}=\left\{(u, v) \left\lvert\, u^{2}+\left(v+\frac{1}{2}\right)^{2}=\frac{1}{4}\right., u \geq 0\right\}$.
Further on $\Gamma_{1}, \Gamma_{2}$ and $\Gamma_{3}$ we have $\frac{\partial p}{\partial y}=-1$, i.e. $v=0$.
Now if $D$ were equal to $B$, then $D^{*}$ should be equal to $B^{*}$, which cannot be the case.
On $\Gamma_{0}$ we have $\frac{\partial p}{\partial y}=0$, i.e. $v=-1$. It follows that the intersection point $B$ of $\Gamma_{0}$ and $\Gamma_{2}$ is send to infinity by the mapping $f$. Hence $u$ has a singular point in $B$.
Intuitively, water is coming down from $D$ to $B$ and at $\Gamma_{2}$ the $y$-component of the velocity is zero. So there would be water accumulation if the $x$-component of the velocity at $B$ were not zero.

Having done some qualitative analysis for the problem, the next question is how to prove existence of a solution. For convenience we state the problem again
(i) $\quad \Delta p=0$ in $\Omega:\{(x, y) \mid 0<y<\phi(x), 0<x<c\}$
(ii) $\quad p(0, y)=y_{1}-y, 0 \leq y \leq y_{1}$,
(iii) $\quad p(c, y)=y_{2}-y, 0 \leq y \leq y_{2}$,
(iv) $\quad p(c, y)=0, y_{2} \leq y \leq \phi(c)$,
(v)

$$
\frac{\partial p}{\partial y}(x, 0)=-1, \quad 0 \leq x \leq c
$$

(vi)

$$
\left.p\right|_{\Gamma=0}
$$

$$
\begin{equation*}
\left.\frac{\partial}{\partial n}(p+y)\right|_{\Gamma=0,} \Gamma:\{(x, \phi(x)) \mid 0 \leq x \leq c\} . \tag{vii}
\end{equation*}
$$

Enclose the region $\Omega$ in a sufficiently large rectangle $D=[0, c] \times[0, \tilde{y}] \bar{y}>y_{1}$. The pressure $p$ is zero in $D \backslash \Omega$. For $\psi \in C_{c}^{\infty}(D)$ consider the expression

$$
\psi \Delta(p+y)
$$

Then by Green's first identity

$$
\int_{\Omega} \psi \Delta(p+y) d x=-\int_{\Omega}(\nabla \psi \cdot \nabla(p+y)) d x+\int_{\partial \Omega} \psi \frac{\partial(p+y)}{\partial n} d s
$$

Due to the boundary conditions the integral $\int_{\partial \Omega} \cdots d s$ is zero. So we end up with the following weak formulation of the problem

$$
\begin{aligned}
0 & =\int_{\Omega}(\nabla \psi \cdot \nabla(p+y)) d x= \\
& =\int_{\Omega} \frac{\partial \psi}{\partial y} d x+\int_{\Omega}(\nabla \psi \cdot \nabla p) d x \\
& =\int_{D} \chi_{\Omega} \frac{\partial \psi}{\partial y} d x+\int_{\Omega}(\nabla \psi \cdot \nabla p) d x
\end{aligned}
$$

Using distributional derivatives we get the distributional equation
(II.8) $-\Delta p=\frac{\partial \chi_{\Omega}}{\partial y}$ in $D$.

Next we apply the so called Baiocchi transform

$$
\begin{equation*}
w(x, y)=\int_{y}^{\tilde{y}} p(x, \eta) d \eta, \frac{\partial w}{\partial y}=-p \tag{II.9}
\end{equation*}
$$

It follows that

$$
-\Delta p=\frac{\partial}{\partial y}(\Delta w)
$$

Since $w(x, \bar{y}) \equiv 0$ and $\chi_{\Omega}(\bar{y})=0$ it follows that
(II.10) $\Delta w=\chi_{\Omega}$ in $D$.

We look for a solution $w \in W^{2,1}(D)$ with $w \geq 0$ and values on the boundary $\partial D$ given by

$$
w(x, \bar{y})=0,0 \leq x \leq c
$$

$$
\begin{aligned}
& w(0, y)=\int_{y}^{\tilde{y}}\left(y_{1}-\eta\right)^{+} d \eta, \quad 0<y<\tilde{y} \\
& w(c, y)=\int_{y}^{\tilde{y}}\left(y_{2}-\eta\right)^{+} d \eta, \quad 0<y<\tilde{y} .
\end{aligned}
$$

For the boundary condition at $y=0$ we have

$$
q=\int_{0}^{\bar{y}} v_{x}(x, \eta) d \eta=-\frac{d}{d x} \int_{0}^{\tilde{y}} p(x, \eta) d \eta
$$

so that

$$
\int_{0}^{\bar{y}} p(x, \eta) d \eta=c-q x
$$

with

$$
c=\int_{0}^{\bar{y}} p(0, \eta) d \eta=\int_{0}^{y_{1}}\left(y_{1}-\eta\right) d \eta=\frac{1}{2} y_{1}^{2} .
$$

It follows that

$$
w(x, 0)=\frac{1}{2} y_{1}^{2}-\frac{1}{2 c}\left(y_{1}^{2}-y_{2}^{2}\right) x, \quad 0 \leq x \leq c
$$

where we used (II.5).
So for all $u \in W^{2,1}(D)$ such that $\left.u\right|_{\partial D}=w \mid \partial D$ and $u \geq 0$

$$
\int_{D} \nabla w \cdot \nabla(u-w) d x=-\int_{D} \chi_{\Omega}(u-w) d x
$$

yielding the integral inequality

$$
\int_{D} \nabla w \nabla(u-w) d x \geq-\int_{D}(u-w) d x .
$$

We can conclude from this that there exists exactly one solution. In fact we have to minimize the functional

$$
I(v)=\int_{D}|\nabla v|^{2} d x+2 \int_{D} v d x
$$

over the closed convex set

$$
K=\left\{v \in W^{2,1}(D)|v \geq 0, v|_{\partial D}=\left.w\right|_{\partial D}\right\}
$$

Next we discuss a generalization of the dam problem.

The saturation $\sigma$ of a porous medium is defined as follows

$$
\begin{equation*}
\sigma=\frac{\text { volume occupied by the fluid }}{\text { total available volume }}, 0 \leq \sigma \leq 1 . \tag{II.11}
\end{equation*}
$$

Until now we considered the situation $\sigma=0$ (dry) or $\sigma=1$ (complete saturation). It is clear that a complete saturation requires a certain pressure $p_{s}$ :

$$
p \geq p_{s} \Rightarrow \sigma=1
$$

In this more general situation time has to be taken into account and instead of the incompressibility condition $\operatorname{div} \vec{v}=0$ we have

$$
\begin{align*}
& \frac{\partial \sigma}{\partial t}+\operatorname{div} \vec{v}=0  \tag{II.12}\\
& \vec{v}=-k(\sigma) \nabla(p+\rho g y) .
\end{align*}
$$

A combination of the above relations gives

$$
\frac{\partial \sigma}{\partial t}-\nabla \cdot(k(\sigma) \nabla(p+\rho g y))=0
$$

If gravity can be neglected, then $p-\sigma, 0<\sigma<1$, and we end up with a heat equation of type

$$
\frac{\partial p}{\partial t}-\nabla \cdot(a(p) \nabla p)=0
$$

We consider the problem with $a(p)$ constant and for only one space variable. So we look at the differential equation

$$
\frac{\partial p}{\partial t}-a \frac{\partial^{2} p}{\partial x^{2}}=0
$$

valid in the unsaturated region $0 \leq x \leq s(t), s(t)$ denoting the free boundary point dependent on $t$. Also, we rescale such that $p_{s}=0$.

Consider the following schematic plot


Statement of the problem
(II.13)
(i)

$$
\frac{\partial^{2} p}{\partial x^{2}}=0, s(t)<x<L
$$

(ii)

$$
\frac{\partial p}{\partial t}=a \frac{\partial^{2} p}{\partial x^{2}}, \quad 0<x<s(t)
$$

(iii)

$$
p(0, t)=p_{0}(t), t \geq 0
$$

(iv)

$$
p(L, t)=p_{1}(t), t \geq 0
$$

(v) $\quad p(s(t), t)=0$ (first condition on the free boundary).

So in the saturated region we have

$$
p(x, t)=p_{1}(t) \frac{(x-s(t))}{(L-s(t))} \quad, \quad s(t) \leq x \leq L .
$$

Consequently
(vi) $\quad \frac{\partial p}{\partial x}(s(t), t)=\frac{p_{1}(t)}{L-s(t)}$ (second condition on the free boundary).

## CHAPTER III

## THE STEFAN PROBLEM

We consider a heat conducting medium occupying a given domain $\Omega$ in $\mathbb{R}^{n}$ in which phase change from solid to liquid or from liquid to solid is taken place. By $\Omega_{s}$ we denote the solid part of the domain and by $\Omega_{l}$ the liquid part. The crucial point of the mathematical scheme is the description of what happens at the interface $\Gamma$. The temperature in each phase obeys the heat equation. So

$$
\begin{equation*}
\rho_{s} c_{s} \frac{\partial \theta_{s}}{\partial t}-k_{s} \Delta \theta_{s}=0 \text { in } \Omega_{s} \tag{III.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{l} c_{l} \frac{\partial \theta_{l}}{\partial t}-k_{l} \Delta \theta_{l}=0 \text { in } \Omega_{l} \tag{III.2}
\end{equation*}
$$

Here $\rho_{s}, c_{s}$ and $k_{s}$ are respectively the density, the specific heat and the thermal conductivity of the solid phase. Similarly, $\rho_{l}, c_{l}$ and $k_{l}$ are defined.
At the interface $\Gamma$ the temperature must be equal to the phase change temperature $\theta_{m}$ which is normally a constant. So we derive the condition

$$
\begin{equation*}
\theta_{s}(x, t)=\theta_{l}(x, t)=\theta_{m}, \quad x \in \Gamma(t), \quad t \geq 0 . \tag{III.3}
\end{equation*}
$$

This is not the only condition for $\theta_{s}$ and $\theta_{l}$ on $\Gamma$. An additional condition is derived from the heat balance at the interface.
Let $\vec{n}$ be the unit normal vector at $\Gamma$ pointing towards the solid phase and let $\vec{v}$ be the velocity of a point of $\Gamma$. The the normal component of $\vec{v}$, i.e. $\vec{v} \cdot \vec{n}$, represents the local rate of melting or of
 solidification if positive or negative, respectively. If $L$ denotes the heat absorbed (or released) for melting (or solidifying) a unit volume of the material, then $L \vec{v} \cdot \vec{n}$ is the local rate of heat absorption (or heat release) in the process. Further, the heat coming to the interface from the liquid phase equals $k_{l} \frac{\partial \theta_{l}}{\partial n}$ and the heat flux leaving it through the solid phase equals $k_{s} \frac{\partial \theta_{s}}{\partial n}$. There is balance of heat whenever

$$
\begin{equation*}
L \vec{v} \cdot \vec{n}=-k_{l} \frac{\partial \theta_{l}}{\partial n}+k_{s} \frac{\partial \theta_{s}}{\partial n} \text { on } \Gamma . \tag{III.4}
\end{equation*}
$$

Let the interface $\Gamma$ be described by the equation

$$
\begin{equation*}
S(x, t)=0 \tag{III.5}
\end{equation*}
$$

where $S$ is a continuously differentiable function. Then we have

$$
\begin{equation*}
\vec{n}(x, t)= \pm \nabla_{x} S(x, t) /\left|\nabla_{x} S(x, t)\right| \tag{III.6}
\end{equation*}
$$

From (III.5) taking the total time derivative we obtain

$$
\begin{equation*}
0=\frac{d}{d t} S(x(t), t)=\vec{v} \cdot \nabla_{x} S(x, t)+\frac{\partial S}{\partial t}(x, t) \tag{III.7}
\end{equation*}
$$

Thus we arrive at the Stefan condition
(III.8) $\quad-L \frac{\partial S}{\partial t}=\left[-k_{l} \nabla_{x} \theta_{l}+k_{s} \nabla_{x} \theta_{s}\right] \cdot \nabla_{x} S$.

The problem of determining $\left(S, \theta_{l}, \theta_{s}\right)$ is called a Stefan problem. In its classical formulation $S$ is required to be $C^{1}$ for $t$ in some interval $(0, T), \theta_{l}$ and $\theta_{s}$ must be $C^{2,1}$ in $\Omega_{l} \times(0, T)$ and $\Omega_{s} \times(0, T)$, respectively. Moreover, the temperature $\theta$ composed of $\theta_{s}$ and $\theta_{s}$, must be continuous in $\bar{\Omega} \times[0, T]$.

## 1. One-dimensional case

In the one-dimensional problem $x$ is a scalar variable and the free boundary can be expressed by the equation
(III.1.1) $\quad S(x, t)=x-s(t)=0$.

The Stefan condition (III.8) takes the form
(III. 1.2)

$$
L \dot{s}(t)=-k_{l} \frac{\partial \theta_{l}}{\partial x}+k_{s} \frac{\partial \theta_{s}}{\partial x}
$$

which can be written as

$$
L \dot{s}(t)=\llbracket-k \frac{\partial \theta}{\partial x} \rrbracket_{s}^{l}
$$

where $\llbracket f \rrbracket_{s}^{l}$ denotes the jump of $f$ from the liquid side to the solid side at the interface. Whenever either $\theta_{s}$ or $\theta_{l}$ is identically equal to $\theta_{m}$ the Stefan condition simplifies to
(i) $L \dot{s}(t)=-k_{l} \frac{\partial \theta_{l}}{\partial x} \quad$ (melting, liquid phase problem).
(ii) $\quad L \dot{s}(t)=k_{s} \frac{\partial \theta_{s}}{\partial x} \quad$ (solidifying, solid phase problem).

We have a so called one phase problem.

Let us consider the liquid phase problem, which becomes after rescaling.
(III.1.4)
(i) $\quad \frac{\partial^{2} u}{\partial x^{2}}-\frac{\partial u}{\partial t}=0$
(ii) $u(x, t)=0, x \geq s(t), t>0$
(iii) $\quad \dot{s}(t)=-\frac{\partial u}{\partial x}(s(t), t), t>0$.

Observe that $u$ denotes the rescaled temperature such that $u=0$ is the melting temperature. The situation is sketched in the following plot


First we construct explicit solutions using self-similar solutions of the heat equation.
Therefore, take $u(x, t)=f(a(t) x)$. Then we have

$$
\begin{aligned}
& \frac{\partial^{2} u}{\partial x^{2}}(x, t)=(a(t))^{2} f^{\prime \prime}(a(t) x) \\
& \frac{\partial u}{\partial t}(x, t)=\dot{a}(t) f^{\prime}(a(t) x) .
\end{aligned}
$$

Taking $\eta=a(t) x$ it follows that

$$
\eta \dot{a}(t) f^{\prime}(\eta)=(a(t))^{3} f^{\prime \prime}(\eta) .
$$

Separation of variables yields

$$
\frac{\dot{a}}{a^{3}}=\lambda \quad \text { and } \quad f^{\prime \prime}=\lambda \eta f^{\prime}
$$

what for $\lambda=-2$ leads to the solution

$$
a(t)=\frac{1}{2 \sqrt{t}} \quad \text { and } \quad f^{\prime}(\eta)=A e^{-\eta^{2}}
$$

Thus the following self-similar solutions have been constructed
(III.1.5)

$$
u(x, t)=A\left[\operatorname{erf}\left(\frac{1}{2} \alpha\right)-\operatorname{erf}\left[\frac{x}{2 \sqrt{t}}\right]\right]
$$

where

$$
\operatorname{erf}(y)=\frac{2}{\sqrt{\pi}} \int_{-\infty}^{y} e^{-\eta^{2}} d \eta .
$$

Introducing III.1.5 into the free boundary condition

$$
u(s(t), t)=0
$$

yields
(III.1.6) $s(t)=\alpha \sqrt{t}, t>0$
whence

$$
\dot{s}(t)=\frac{\alpha}{2 \sqrt{t}}, t>0 .
$$

Moreover, $\frac{\partial u}{\partial x}=\frac{-2 A}{\sqrt{\pi}} e^{-x^{2} / 4 t} \frac{1}{2 \sqrt{t}}$.

So the second free boundary condition

$$
\dot{s}(t)=-\frac{\partial u}{\partial x}(s(t), t)
$$

is fulfilled whenever
(III.1.7) $\quad A=\frac{1}{2} \sqrt{\pi} \alpha e^{\alpha^{2} / 4}$.

It follows that the free boundary can be any parabola. The result is the so called Neumann solution

$$
u(x, t)=\alpha e^{\alpha^{2} / 4} \int_{x / 2 \sqrt{t}}^{\alpha / 2} \exp \left(-\eta^{2}\right) d \eta \quad, x<\alpha \sqrt{t} .
$$

We have found a one-parameter family ( $\alpha$ ) of solutions.


At $t=0$ the half space $x<0$ is occupied by the liquid and the half space $x>0$ is occupied by the solid at zero temperature. We can find the initial value of $u, u_{0}(x)$, for $x<0$ by letting $t$ tend to zero. This yields

$$
\begin{equation*}
u_{0}(x)=\alpha e^{\alpha^{2} / 4} \int_{-\infty}^{\alpha / 2} e^{-\eta^{2}} d \eta \quad, x<0 \tag{III.1.8}
\end{equation*}
$$

We also see that the temperature at $x=0$ remains constant

$$
u(0, t)=\alpha e^{\alpha^{2} / 4} \int_{0}^{\omega / 2} e^{-\eta^{2}} d \eta, t>0
$$

This way we can solve problems with constant initial data $u_{0}$ or constant boundary data $u_{1}$. Therefore we have to find $\alpha$ such that

$$
u_{0}=\alpha e^{\alpha^{2 / 4}} \int_{-\infty}^{\alpha / 2} e^{-\eta^{2}} d \eta
$$

or

$$
u_{1}=\alpha e^{\alpha^{2} / 4} \int_{0}^{\alpha / 2} e^{-\eta^{2}} d \eta
$$

It can be shown that for each $u_{0}>-1$ there exists a unique $\alpha$. Moreover the corresponding solution $u$ has the same sign as $u_{0}$. So for $-1<u_{0}<0$ we are dealing with a supercooled fluid

$$
-1<u(x, t)<0, x<s(t), t>0 .
$$

There appears a discontinuity at $(x, t)=(0,0)$ since the temperature jumps from $u_{0}$ to 0 .
Remark. In general (so not only for similarity solutions) the following can be said

- The temperature jump at the origin cannot exceed -1.
- If there is some neighbourhood of $x=0$ in which $u_{0}(x)>-1$ then the supercooled problem has one unique solution.
- If $u_{0}(x) \leq-1$ in some neighbourhood of $x=0$ then no solution exists.

We return to the one-phase problem in which we attach boundary data at $x=0$ and initial data for $t=0$.


We assume the following: $h$ and $f$ are continuous, $h(b)=0$, $0 \leq h(x) \leq H(b-x)$
and $f$ is bounded
with $0 \leq f(t) \leq H b$.
(III.1.9) Formulation of the problem.

$$
\begin{equation*}
\frac{\partial^{2} u}{\partial x^{2}}-\frac{\partial u}{\partial t}=0,0<x<s(t), t>0 \tag{i}
\end{equation*}
$$

(ii) $u(x, 0)=h(x), 0<x<b$
(iii) $u(0, t)=f(t), t>0$
(iv) $\quad s(0)=b>0$
(v) $\quad u(s(t), t)=0, t>0$
(vi) $\quad \frac{\partial u}{\partial x}(s(t), t)=-\dot{s}(t), t>0$.

## (III. 1.10) Theorem.

There exists precisely one solution $u$ (global in time) of (III.1.9) with $0 \leq \dot{s} \leq H$.

## Proof.

We use Schauder's fixed point theorem to prove existence.

## Consider the following family of curves

$$
\begin{array}{r}
S(T, A)=\left\{s \in C((0, T)) \mid s(0)=b \wedge \forall_{0<t_{1}, t_{2}<T}:\right. \\
\left.0 \leq \frac{s\left(t_{1}\right)-s\left(t_{2}\right)}{t_{1}-t_{2}} \leq A\right\} .
\end{array}
$$

Take a fixed $s \in S(T, A)$ and consider the problem
(III.1.11)

$$
\left[\begin{array}{l}
\frac{\partial u}{\partial t}-\frac{\partial^{2} u}{\partial x^{2}}=0,0<x<s(t), 0<t<T \\
u(s(t), t)=0,0<t<T \\
u(x, 0)=h(x), 0<x<b \\
u(0, t)=f(t), 0<t<T
\end{array}\right.
$$

Since the curve $s$ is Lipschitz continuous, $\frac{\partial u}{\partial x}$ is continuous up to the curve $s$. (This is an old theorem of Gevrey, 1913.) Next solve the equation

$$
\left[\begin{array}{l}
\dot{\sigma}(t)=-\frac{\partial u}{\partial x}(s(t), t) \\
\sigma(0)=b
\end{array}\right.
$$

where $u$ denotes the solution of (III.1.11). This generates a new curve $\sigma$. The described procedure leads to an operator $\mathfrak{t}$ mapping a curve $s$ on a curve $\sigma$. If there is a constant $A$ such that $\mathfrak{t}$ maps $S(T, A)$ into ifself, then a fixed point of $t$ yields a solution of the free boundary Stefan problem.
To find such a constant $A$ we first observe that we have assumed that

* $h(x)<H(b-x)$
* $f(t)<H b$
* $s$ is increasing.

Take a fixed $t_{0}, 0<t_{0}<T$.


Let $v$ be defined by $v(x)=H\left(s\left(t_{0}\right)-x\right)$. Then $v(x) \geq h(x)$, $\nu(0) \geq f(t)$, $v(s(t))>0,0<t<t_{0}$, and
$v$ is a solution of the heat equation.

Now let $w=v-u$. Then we have

$$
\begin{aligned}
& \frac{\partial w}{\partial t}=\frac{\partial^{2} w}{\partial x^{2}}, 0<x<s(t), 0<t<t_{0} \\
& w(0, t)>0,0<t<t_{0} \\
& w(x, 0)>0,0<x<b
\end{aligned}
$$

$$
\begin{aligned}
& w(s(t), t)>0,0<t<t_{0} \\
& w\left(s\left(t_{0}\right), t_{0}\right)=0
\end{aligned}
$$

So the maximum principle says that $w \geq 0$ for $0 \leq x \leq s(t)$ and $0 \leq t \leq t_{0}$, which means

$$
\begin{aligned}
& v(x) \geq u(x, t), \quad 0 \leq x \leq s(t), 0 \leq t \leq t_{0} \\
& v\left(s\left(t_{0}\right)\right)=u\left(s\left(t_{0}\right), t_{0}\right)=0 .
\end{aligned}
$$

It follows that

$$
\frac{\partial u}{\partial x}\left(s\left(t_{0}\right), t_{0}\right) \geq-H
$$

(see picture).


As a simple consequence of the maximum principle, yielding $u>0$, we also must have

$$
\frac{\partial u}{\partial x}\left(s\left(t_{0}\right), t_{0}\right) \leq 0
$$

So $t_{0}$ being arbitrary, we obtain

$$
\forall_{t>0}:-H \leq \frac{\partial u}{\partial x}(s(t), t) \leq 0 .
$$

Thus we find $0 \leq \dot{\sigma} \leq H$ and we can therefore take $A=H$, i.e. $\mathbf{t}$ maps $S(T, H)$ into $S(T, H)$.

Now we are in the following position

- $\quad \mathbf{t}$ is a mapping from $S(T, H)$ into $S(T, H)$.
- $S(T, H)$ is a closed convex bounded subset of the Banach space $C([0, T])$ with $\|f\|=\max _{0 \leq i \leq T}|f(t)|, f \in C([0, T])$.
In order to be able to apply Schauder's fixed point theorem we prove that
I : $S(T, H)$ is compact in $C([0, T])$
II : $\mathbf{t}: S(T, H) \rightarrow S(T, H)$ is continuous .
I. By definition the set $S(T, H)$ is equicontinuous. So by Ascoli's theorem the set $S(T, H)$ is compact in $C([0, T])$.
II. Proving continuity of $\mathbf{t}$ requires a more lengthy proof.

$0 \leq h(x) \leq H(b-x)$, $0 \leq f(t) \leq H b$, $s$ is increasing with $s(0=b$
Lipschitz constant $\leq H$
$\Rightarrow s(t) \leq H t$.

Consider the vector field

$$
\vec{v}(x, t)=\left(x u, x \frac{\partial u}{\partial x}-u\right) .
$$

Then by Green's integral identity

$$
\int_{G^{*}}\left[\frac{\partial v_{2}}{\partial x}-\frac{\partial v_{1}}{\partial t}\right] d \sigma=\int_{\partial G^{*}}(\vec{v}, \vec{t}) d s
$$

It follows that

$$
\begin{align*}
0=\int_{0}^{b} x h(x) d x & -\int_{0}^{s\left(t^{*}\right)} x u\left(x, t^{*}\right) d x+\int_{0}^{t^{*}} f(\tau) d \tau  \tag{*}\\
& +\int_{0}^{t^{*}} s(\tau) \frac{\partial u}{\partial x}(s(\tau), \tau) d \tau
\end{align*}
$$

(observe that $u(s(\tau), \tau)=0$ ).
We have defined $\dot{\sigma}(t)=-\frac{\partial u}{\partial x}(s(t), t), \sigma(0)=b$.
Take $s_{1}, s_{2} \in S(T, H)$ and use $\left({ }^{*}\right)$ to get

$$
\begin{aligned}
0= & -\int_{0}^{s_{1}} x u_{1}\left(x, t^{*}\right) d x+\int_{0}^{s_{2}} x u_{2}\left(x, \tau^{*}\right) d x \\
& -\int_{0}^{t^{*}} s_{1}(\tau) \dot{\sigma}_{1}(\tau) d \tau+\int_{0}^{t^{*}} s_{2}(\tau) \dot{\sigma}_{2}(\tau) d \tau .
\end{aligned}
$$

This yields
(**)

$$
\begin{aligned}
& \int_{0}^{t^{*}} s_{1}(\tau)\left(\dot{\sigma}_{1}(\tau)-\dot{\sigma}_{2}(\tau)\right) d \tau=-\int_{0}^{i^{*}} \dot{\sigma}_{2}(\tau)\left(s_{1}(\tau)-s_{2}(\tau)\right) d \tau \\
& -\int_{0}^{\min \left(s_{1}\left(\tau^{*}\right), s_{2}\left(t^{*}\right)\right)} x\left(u_{1}\left(x, t^{*}\right)-u_{2}\left(x, t^{*}\right)\right)+(-1)^{j} \int_{\min \left(s_{1}\left(i^{*}\right), s_{2}\left(t^{*}\right)\right)}^{\max \left(s_{1}\left(t^{*}\right), s_{2}\left(t^{*}\right)\right)} x u_{j}\left(x, t^{*}\right) d x .
\end{aligned}
$$

Integrating the left hand side by parts in a slightly generalized sense we get
$\left({ }^{* * *}\right) \quad s_{1}\left(t^{*}\right)\left(\sigma_{1}\left(t^{*}\right)-\sigma_{2}\left(t^{*}\right)\right)=\int_{0}^{t^{*}} \dot{s}_{1}(\tau)\left[\sigma_{1}(\tau)-\sigma_{2}(\tau)\right] d \tau+\cdots$
Put $\|z\|_{t^{*}}=\max \left\{|z(x, \tau)| \mid 0 \leq x \leq s(\tau), 0<\tau<t^{*}\right\}$.
Then $\left({ }^{* *}\right)$ and $\left({ }^{* * *}\right)$ gives
$(* * * *) \quad b\left|\sigma_{1}\left(t^{*}\right)-\sigma_{2}\left(t^{*}\right)\right| \leq s_{1}\left(t^{*}\right)\left|\sigma_{1}\left(t^{*}\right)-\sigma_{2}\left(t^{*}\right)\right| \leq$
$\leq H t^{*}\left\|\sigma_{1}-\sigma_{2}\right\|_{t^{*}}+H t^{*}\left\|s_{1}-s_{2}\right\|_{i^{*}}+$
$+\frac{1}{2}\left(b+H t^{*}\right)^{2}\left\|u_{1}-u_{2}\right\|_{t^{*}}+\frac{1}{6} H\left\|s_{1}-s_{2}\right\|_{i^{*}}^{3}$
where we used that $u_{j}\left(x, t^{*}\right) \leq H\left(s_{j}\left(t^{*}\right)-x\right), 0 \leq x \leq s_{j}\left(t^{*}\right)$.
By the maximum principle

$$
\left\|u_{1}-u_{2}\right\|_{i^{*}} \leq \sup _{0<\tau<i^{*}}\left|u_{1}-u_{2}\right|_{x=\min \left(s_{1}(\tau), s_{2}(\tau)\right)} \leq H\left\|s_{1}-s_{2}\right\|_{i^{*}}
$$

Now take $T^{*}=\frac{1}{2} b / H$ and conclude that
$(* * * * *) \quad \frac{1}{2} b\left\|\sigma_{1}-\sigma_{2}\right\|_{T^{*}} \leq C\left\|s_{1}-s_{2}\right\|_{T^{*}}$
with $C$ a constant depending on $b$ and $H$.
Hence $\mathbf{t}$ is a continuous mapping from $S\left(T^{*}, H\right)$ into $S\left(T^{*}, H\right)$.

Remark. We have obtained a solution up to $t=T^{*}$. However we know that $0 \leq u\left(x, T^{*}\right) \leq H\left(s\left(T^{*}\right)-x\right)$. So taking

$$
h^{*}(x)=u\left(x, T^{*}\right)
$$

we can solve the problem

$$
\begin{aligned}
& \frac{\partial u^{*}}{\partial t}=\frac{\partial^{2} u^{*}}{\partial x^{2}}, 0<x<s(t), T^{*}<t<T^{* *} \\
& u^{*}(0, t)=f(t), T^{*}<t<T^{* *} \\
& u^{*}(x, 0)=h^{*}(x), 0<x<s\left(T^{*}\right) \\
& u^{*}(s(t), t)=0, T^{*}<t<T^{* *} \\
& \dot{s}(t)=\frac{-\partial u}{\partial x}(s(t), t), T^{*}<t<T^{* *}
\end{aligned}
$$

Thus the solution can be extended up to $t=T^{* *}=\frac{1}{2} s\left(T^{*}\right) / H$.
In other words the solution exists for all $T>0$.

Our next aim is to prove uniqueness.

## (III.1.12) Lemma.

For any solution $(s, u)$ we have the following identity

$$
\frac{1}{2}(s(t))^{2}=\frac{1}{2} b^{2}+\int_{0}^{b} x h(x) d x-\int_{0}^{s(t)} x u(x, t) d x+\int_{0}^{t} f(\tau) d \tau
$$

## Proof.

As we have seen in the above proof

$$
\begin{aligned}
0=\int_{0}^{b} x h(x) d x & -\int_{0}^{s(t)} x u(x, t) d x+\int_{0}^{b} f(\tau) d \tau+ \\
& +\int_{0}^{t} s(\tau) \frac{\partial u}{\partial x}(s(\tau), \tau) d \tau
\end{aligned}
$$

Since $-\dot{s}(\tau)=\frac{\partial u}{\partial x}(s(\tau), \tau)$ the result follows.

The following result is on the monotone dependence of the free boundary on the data.
(III.1.13) Theorem.

Let ( $s_{1}, u_{1}$ ) and ( $s_{2}, u_{2}$ ) be solutions corresponding to the respective data ( $b_{1}, h_{1}, f_{1}$ ) and $\left(b_{2}, h_{2}, f_{2}\right)$ satisfying the requirements

$$
\begin{aligned}
& 0 \leq h_{j}(x) \leq H_{j}\left(b_{j}-x\right) \\
& 0 \leq f_{j}(t) \leq H_{j} b_{j}, \quad j=1,2
\end{aligned}
$$

If $b_{1} \geq b_{2}, h_{1} \geq h_{2}$ and $f_{1} \geq f_{2}$ then $s_{1} \geq s_{2}$.

## Proof.

Suppose first that $b_{1}>b_{2}$.
We shall show that $s_{1}(t)>s_{2}(t)$ for all $t \in[0, T]$.
Assume this were not true. Then there exists $t_{0} \in(0, T] l$ such that $s_{1}\left(t_{0}\right)=s_{2}\left(t_{0}\right)$ while $s_{1}(t)>s_{2}(t)$ for $0 \leq t<t_{0}$.
The function $u_{1}-u_{2}$ satisfies

$$
\begin{gathered}
\left(u_{1}-u_{2}\right)(0, t)=f_{1}(t)-f_{2}(t) \geq 0 \\
\left(u_{1}-u_{2}\right)(x, 0)=h_{1}(x)-h_{2}(x) \geq 0, \\
0 \leq x \leq b_{2} \\
\left(u_{1}-u_{2}\left(s_{2}(t), t\right)=u_{1}\left(s_{2}(t), t\right) \geq 0,\right.
\end{gathered}
$$



## $0 \leq t \leq t_{0}$.

So by the maximum principle $u_{1} \geq u_{2}$ on $0 \leq x \leq s_{2}(t), 0 \leq t \leq t_{0}$.
Since $u_{1}\left(s_{1}\left(t_{0}\right), t_{0}\right)=u_{2}\left(s\left(t_{0}\right), t_{0}\right)=0$ it also follows that

$$
\frac{\partial u_{1}}{\partial x}\left(s_{1}\left(t_{0}\right), t_{0}\right)<\frac{\partial u_{2}}{\partial x}\left(s_{2}\left(t_{0}\right), t_{0}\right)
$$

because of the so called boundary point principle.
Hence $\dot{s}_{1}\left(t_{0}\right)>\dot{s}_{2}\left(t_{0}\right)$ which yields a contradiction.

Next consider the case $b_{1}=b_{2}=b$.
In this case take a family of solutions $\left(s^{\delta}, u^{5}\right)$ to the problem with data

$$
\begin{aligned}
& s^{\delta}(0)=b+\delta \\
& u^{\delta}(x, 0)=h_{1}(x) \quad, \quad 0 \leq x<b \\
& u^{\delta}(x, 0)=0, b \leq x \leq b+\delta \\
& u^{\delta}(0, t)=f_{1}(t) .
\end{aligned}
$$

From the above arguments we see that $s_{1}<s^{\delta}$ and $s_{2}<s^{\delta}$ for all $t$. For the difference $s^{\delta}-s_{1}$ we have by the previous lemma

$$
\begin{aligned}
& \frac{1}{2}\left(s^{\delta}(t)-s_{1}(t)\right)\left(s^{\delta}(t)+s_{1}(t)\right)= \\
& =\frac{1}{2} \delta(2 b+\delta)-\int_{0}^{s_{1}(t)} x\left(u^{\delta}(x, t)-u_{1}(x, t)\right) d x \\
& -\int_{s_{1}(t)}^{s^{\circ}(t)} x u^{\delta}(x, t) d x .
\end{aligned}
$$

Since $u^{\delta} \geq 0$ and $u^{\delta} \geq u_{1}$ we have

$$
\frac{1}{2}\left(s^{\delta}(t)-s_{1}(t)\right)\left(s^{\delta}(t)+s_{1}(t)\right) \leq \frac{1}{2} \delta(b+\delta) .
$$

Moreover, since $s^{\delta}(t)+s_{1}(t) \geq 2 b+\delta$, we have

$$
s^{\delta}(t)-s_{1}(t) \leq \delta
$$

and using $s_{2}(t)<s^{\delta}$,

$$
s_{2}(t)<s_{1}(t)+\delta .
$$

The result follows by letting $\delta$ tend to zero.

## (III.1.14) Corollary.

Problem (III.1.9) has exactly one solution ( $s, u$ ) with $s$ in the class $S(T, H)$.

## Proof.

If $\left(s_{1}, u_{1}\right)$ and $\left(s_{2}, u_{2}\right)$ are solutions corresponding to the same set of data, then by the previous theorem both $s_{1} \geq s_{2}$ and $s_{2} \leq s_{1}$. Hence $s_{1}=s_{2}$ and a posteriori $u_{1}=u_{2}$.

## 2. More dimensional Stefan problem

Throughout we assume that no supercooling or superheating is allowed. So each phase is characterized by the sign of the temperature, taking zero as the melting point. The thermal energy (enthalpy) stored in a unit volume of the solid can be taken equal to $\rho_{s} c_{s} \theta$ while in the liquid the thermal energy can be taken equal to $\rho_{l} c_{l} \theta+L, L$ denoting the latent heat. So we can define

$$
\begin{align*}
& E(\theta)=\rho_{s} c_{s} \theta \text { for } \theta<0  \tag{II.2.1}\\
& E(\theta)=\rho_{l} c_{l} \theta+L \text { for } \theta>0,
\end{align*}
$$

leaving aside for the moment the question of defining $E$ for $\theta=0$.
The heat equation can be written in the following form

$$
\begin{align*}
& \frac{\partial E}{\partial t}=k_{s} \Delta \theta, \quad \theta<0,  \tag{II.2.2}\\
& \frac{\partial E}{\partial t}=k_{l} \Delta \theta, \quad \theta>0,
\end{align*}
$$

where $k_{s}$ and $k_{l}$ are supposed to be constant.
Consider the following space of test functions

$$
\begin{equation*}
V=\left\{\phi \in C^{\infty}(\bar{\Omega} \times[0, T]) \mid \phi=0 \text { on } \partial \Omega \times(0, T) \wedge \phi(x, T)=0\right\} \tag{II.2.3}
\end{equation*}
$$

With this test function space a weak formulation of the Stefan problem will be derived.

$\Omega_{0}=\Omega \times(0)$
$\Gamma_{l}$ : outer boundary with external normal $\vec{n}_{l}$
$k_{l} \theta \mid \Gamma_{l}=f(x, t)>0$
$\Gamma_{s}$ : inner boundary with external
normal $\vec{n}_{s}$
$\left.k_{s} \theta\right|_{\Gamma_{s}}=g(x, t)<0$
initial condition
$E(x, 0)=E_{0}(x)$ in $\Omega$.
$\vec{v}$ normal to the free boundary $\Gamma_{0}$ in $\mathbb{R}^{n+1}$ (directed towards the solid)
$\vec{\mu}$ projection of $\vec{v}$ on $\mathbb{R}^{n}$
$\vec{n}$ normal to $\Gamma_{0} \cap \Omega \times\{t\}$ in $\mathbb{R}^{n}$
$Q=\Omega \times(0, T), Q_{l}\left(Q_{s}\right)$ subset of $Q$ occupied by the liquid (solid).

Observe that $\vec{v}(x, t)=\left(v_{1}(x, t), \ldots, v_{n}(x, t), v_{0}(x, t)\right)$.
Since $\frac{\partial}{\partial t}(\phi E)=\phi \frac{\partial E}{\partial t}+E \frac{\partial \phi}{\partial t}, \phi \in V$, we find by Gauss' divergence theorem applied to the region $Q_{l}$
(i) $\quad \int_{Q_{1}} \phi \frac{\partial E}{\partial t} d x d t=-\int_{Q_{1}} E \frac{\partial \phi}{\partial t} d x d t+\int_{\Omega_{a}} \phi E_{0} d \sigma+\int_{\Gamma_{0}} L \cdot \phi \cdot v_{0} d \sigma$.

Similarly for the region $Q_{s}$

$$
\begin{equation*}
\int_{Q_{2}} \phi \frac{\partial E}{\partial t} d x d t=-\int_{Q_{0}} E \frac{\partial \phi}{\partial t} d x d t+\int_{\Omega_{0}} \phi E_{0} d \sigma . \tag{ii}
\end{equation*}
$$

Further, applying Gauss' integral identity again,

$$
\begin{aligned}
& \int_{Q_{l}} \nabla_{x} \cdot\left(\phi \nabla_{x} \theta-\theta \nabla_{x} \phi\right) d x d t= \\
& \quad=\int_{\Gamma_{1}}\left(\phi \frac{\partial \theta}{\partial n_{l}}-\theta \frac{\partial \phi}{\partial n_{l}}\right) d \sigma+\int_{\Gamma_{0}}\left(\left.\phi \nabla_{x} \theta\right|_{l}-\left.\theta \nabla_{x} \phi\right|_{l}\right) \cdot \vec{\mu} d \sigma .
\end{aligned}
$$

(Observe that $\nabla_{x} \phi, \nabla_{x} \theta \in \mathbb{R}^{n} \times\{0\}$.)

It yields
(iii)

$$
\begin{aligned}
\int_{Q_{l}} k_{l} \phi \Delta \theta d x d t= & -\int_{\Gamma_{l}} f \frac{\partial \phi}{\partial n_{l}} d \sigma+\left.\int_{\Gamma_{0}} \phi k_{l} \nabla_{x} \theta\right|_{l} \cdot \vec{\mu} d \sigma \\
& +\int_{Q_{l}} k_{l} \theta \Delta \phi d x d t
\end{aligned}
$$

Similarly,
(iv)

$$
\begin{aligned}
\int_{Q_{土}} k_{s} \phi \Delta \theta d x d t= & -\int_{\Gamma_{a}} g \frac{\partial \phi}{\partial n_{s}} d \sigma-\left.\int_{\Gamma_{0}} \phi k_{s} \nabla_{x} \theta\right|_{s} \cdot \vec{\mu} d \sigma \\
& +\int_{Q_{a}} k_{s} \theta \Delta \phi d x d t
\end{aligned}
$$

Combining the relations (i)-(iv) we get

$$
\int_{Q}\left(E \frac{\partial \phi}{\partial t}+k(\theta) \theta \Delta \phi\right) d x d t=
$$

$$
\begin{aligned}
& =\int_{\Omega_{0}} \phi E_{0} d x+\int_{\Gamma_{i}} f \frac{\partial \phi}{\partial n_{l}} d \sigma+\int_{\Gamma_{s}} g \frac{\partial \phi}{\partial n_{s}} d \sigma+ \\
& +\int_{\Gamma_{0}} \phi\left\{L v_{0}-\left.k_{l} \nabla_{x} \theta\right|_{l} \cdot \vec{\mu}-\left.k_{s} \nabla_{x} \theta\right|_{s} \cdot \vec{\mu}\right\} d \sigma
\end{aligned}
$$

where

$$
k(\theta)= \begin{cases}k_{s} & , \quad \theta<0 \\ k_{l} & , \quad \theta>0\end{cases}
$$

Describing the free boundary $\Gamma_{0}$ by $S(x, t)=0$ we have

$$
v_{0}(x, t)=(\partial S / \partial t) /\left(\left|\nabla_{x} S\right|^{2}+(\partial S / \partial t)^{2}\right)^{1 / 2}
$$

and

$$
\vec{\mu}(x, t)=\nabla_{x} S /\left(\left|\nabla_{x} S\right|^{2}+(\partial S / \partial t)^{2}\right)^{1 / 2}
$$

So because of the Stefan condition

$$
L v_{0}=\left.k_{l} \nabla_{x} \theta\right|_{l} \cdot \vec{\mu}+\left.k_{s} \nabla_{x} \theta\right|_{s} \cdot \vec{\mu}
$$

Thus we arrive at the following weak formulation of the Stefan problem.
(II.2.4) $\quad \int_{Q}\left(E \frac{\partial \phi}{\partial t}+k(\theta) \theta \Delta \phi\right) d x d t=$

$$
=\int_{\Omega_{0}} \phi E_{0} d x+\int_{\Gamma_{i}} f \frac{\partial \phi}{\partial n_{l}} d \sigma+\int_{\Gamma_{i}} g \frac{\partial \phi}{\partial n_{s}} d \sigma
$$

We note that $\theta=\beta(E)$ is a single valued function of $E$ :


We define a weak solution of the Stefan problem as a measurable bounded function $E$, satisfying the below integral relation for all test function $\phi$ in $V$.

$$
\begin{align*}
& \int_{Q}\left(E \frac{\partial \phi}{\partial t}+k(\beta(E)) \beta(E) \Delta \phi\right) d x d t=  \tag{II.2.5}\\
& \quad=\int_{\Omega_{0}} \phi E_{0} d x+\int_{\Gamma_{l}} f \frac{\partial \phi}{\partial n_{l}} d \sigma+\int_{\Gamma_{i}} g \frac{\partial \phi}{\partial n_{s}} d \sigma .
\end{align*}
$$

Remark. We have $\{(x, t) \in Q \mid \theta(x, t)=0\}=\{(x, t) \in Q \mid 0<E(x, t)<L\}$.

It may happen that this set has nonzero measure. In this case we have a mushy region, i.e. an intermediate phase between pure states.

We want to prove the existence of a weak solution. First we define the function $E$ and $K$ as

$$
E(\lambda)= \begin{cases}\rho_{s} c_{s} \lambda, & \lambda<0, \\ \rho_{l} c_{l} \lambda+L, & \lambda>0\end{cases}
$$



$$
K(\lambda)= \begin{cases}\lambda k_{s}, & \lambda<0, \\ \lambda k_{l}, & \lambda>0 .\end{cases}
$$



We introduce two sequences of smooth functions: a sequence $E_{m}$ with $E_{m}{ }^{\prime}>0$ and $E_{m}(\lambda)=E(\lambda),|\lambda| \geq \frac{1}{m}$, a sequence $K_{m}$ with $K_{m}{ }^{\prime}>0$ and $K_{m}(\lambda)=K(\lambda),|\lambda| \geq \frac{1}{m}$, and $K_{m} \leq K$. Define for each $m \in \mathbb{N} \quad \tilde{E}_{m}$ by

$$
\tilde{E}_{m}\left(K_{m}(\lambda)\right)=E_{m}(\lambda),
$$

i.e. $\bar{E}_{m}=E_{m} \circ K_{m}^{\leftarrow}$. Then $\tilde{E}_{m}$ is monotoneously increasing and smooth.

Let $\bar{\theta}_{m}$ denote the solution of the non-linear initial value problem

$$
\begin{aligned}
& \frac{\partial \tilde{E}_{m}\left(\bar{\theta}_{m}\right)}{\partial t}-\Delta \tilde{\theta}_{m}=0 \text { in } Q \\
& \left.\bar{\theta}_{m}\right|_{\Gamma_{l}}=f(x, t),\left.\quad \tilde{Q}_{m}\right|_{\Gamma_{s}}=g(x, t) \\
& \bar{\theta}_{m}(x, 0)=\bar{\theta}_{m 0}(x)=K_{m}\left(\theta_{0}(x)\right) .
\end{aligned}
$$

## Remark. In fact we solve the problem

$$
\begin{aligned}
& \frac{\partial E_{m}\left(\theta_{m}\right)}{\partial t}-\Delta K_{m}\left(\theta_{m}\right)=0 \\
& \left.k_{l} \theta_{m}\right|_{\Gamma_{l}}=f,\left.\quad k_{s} \theta_{m}\right|_{\Gamma_{s}}=g
\end{aligned}
$$

$$
\theta_{m}(x, 0)=\theta_{0}(x) \text { with } E_{0}(x)=E\left(\theta_{0}\right) .
$$

It can be proved on the basis of a (generalized) maximum principle that

## (III.2.6)

- $\quad \bar{\theta}_{m}$ is uniformly bounded in $Q$,i.e.

$$
\sup _{m} \sup _{(x, t) \in Q} \tilde{\theta}_{m}(x, t)<\infty
$$

- $\frac{\partial \bar{\theta}_{m}}{\partial n_{s}}, \frac{\partial \bar{\theta}_{m}}{\partial n_{l}}$ are uniformly bounded on $\Gamma_{s}$ and $\Gamma_{l}$, respectively.

Since there exists $c>0$ such that $\tilde{E}_{m}^{\prime} \geq c>0$ we have the estimation

$$
\begin{aligned}
& c \int_{Q_{r}}\left|\frac{\partial \bar{\theta}_{m}}{\partial t}\right|^{2} d x d t \leq \int_{Q_{r}} \bar{E}_{m}^{\prime}\left(\tilde{\theta}_{m}\right)\left|\frac{\partial \bar{\theta}_{m}}{\partial t}\right|^{2} d x d t= \\
& \quad=\int_{Q_{r}} \frac{\partial \bar{\theta}_{m}}{\partial t} \Delta \bar{\theta}_{m} d x d t .
\end{aligned}
$$

Now

$$
\frac{\partial \tilde{\theta}_{m}}{\partial t} \Delta \bar{\theta}_{m}=\nabla_{x} \cdot\left[\frac{\partial \tilde{\theta}_{m}}{\partial t} \nabla_{x} \tilde{\theta}_{m}\right]-\frac{1}{2} \frac{\partial}{\partial t}\left|\nabla_{x} \theta_{m}\right|^{2}
$$

so that the right hand side is equal to

$$
\int_{Q_{r}} \nabla_{x} \cdot\left[\frac{\partial \bar{\theta}_{m}}{\partial t} \nabla_{x} \bar{\theta}_{m}\right] d x d t-\frac{1}{2} \int_{0}^{t^{*}}\left(\frac{d}{d t} \int_{\Omega_{1}}\left|\nabla_{x} \theta_{m}\right|^{2} d x\right) d t
$$

(Observe that $Q_{l^{*}}=\Omega \times\left[0, t^{*}\right]$ and $\Omega_{t}=\Omega \times\{t\}$.)

## By Gauss

$$
\int_{\boldsymbol{e}_{r}} \nabla_{x} \cdot\left[\frac{\partial \bar{\theta}_{m}}{\partial t} \nabla_{x} \bar{\theta}_{m}\right] d x d t=\int_{\Gamma_{l, r}} \frac{\partial \bar{\theta}_{m}}{\partial n_{l}} \cdot \frac{\partial f}{\partial t} d \sigma+\int_{\mathrm{r}_{t, r}} \frac{\partial \bar{\theta}_{m}}{\partial n_{s}} \frac{\partial g}{\partial t} d \sigma .
$$

So consequently

$$
\begin{aligned}
& c \int_{Q_{r}}\left|\frac{\partial \bar{\theta}_{m}}{\partial t}\right|^{2} d x d t+\frac{1}{2} \int_{\Omega\left(t^{*}\right)}\left|\nabla_{x} \bar{\theta}_{m}\right|^{2} d x \leq \\
& \leq \int_{\Gamma_{1}} \frac{\partial \bar{\theta}_{m}}{\partial n_{l}} \frac{\partial f}{\partial t} d \sigma+\int_{\Gamma_{1}} \frac{\partial \bar{\theta}_{m}}{\partial t} d \sigma+\int_{\Omega_{0}}\left|\nabla_{x} \bar{\theta}_{m}\right|^{2} d x .
\end{aligned}
$$

Since the above inequality is valid for all $t^{*} \in(0, T]$ it follows that

$$
\begin{aligned}
& \int_{Q}\left(\left|\frac{\partial \tilde{\theta}_{m}}{\partial t}\right|^{2}+\left|\nabla \tilde{\theta}_{m}\right|^{2}\right) d x d t \\
& \quad \leq \frac{T}{\min \left(c T, \frac{1}{2}\right)}\left\{\int_{\Gamma_{l}} \frac{\partial \bar{\theta}_{m}}{\partial n_{l}} \frac{\partial f}{\partial t} d \sigma+\int_{\Gamma_{,}} \frac{\partial \bar{\theta}_{m}}{\partial n_{s}} \frac{\partial g}{\partial t} d \sigma+\int_{\Omega_{0}}\left|\nabla \tilde{\theta}_{m}\right|^{2} d x\right\} .
\end{aligned}
$$

Now

$$
\begin{aligned}
\int_{\Omega_{0}}\left|\nabla \bar{\theta}_{m}\right|^{2} d x & =\int_{\Omega}\left|K_{m}^{\prime}\left(\theta_{0}\right)\right|^{2}\left|\nabla \theta_{0}\right|^{2} d x \leq \\
& \leq \max \left(k_{s}^{2}, k_{l}^{2}\right) \int_{\Omega}\left|\nabla \theta_{0}\right|^{2} d x
\end{aligned}
$$

So applying (II.2.6) there exists a constant $C_{1}$ independent of $m$ such that
(II.2.7.i) $\quad \forall_{m}: \int_{Q}\left(\left|\frac{\partial \bar{\theta}_{m}}{\partial t}\right|^{2}+\left|\nabla \bar{\theta}_{m}\right|^{2}\right) d x d t \leq C_{1}$.

Further since the sequence ( $\bar{\theta}_{m}$ ) is uniformly bounded in $Q$ there exist constants $C_{2}>0$ and $C_{3}>0$ such that
(II.2.7.ii)

$$
\begin{array}{ll}
\text { (II.2.7.ii) } & \forall_{m}: \int_{Q}\left|\bar{\theta}_{m}\right|^{2} d x d t<C_{2} \\
\text { (II.2.7.iii) } & \forall_{m}: \int_{Q}\left|\tilde{E}_{m}\left(\tilde{\theta}_{m}\right)\right|^{2} d x d t<C_{3} .
\end{array}
$$

Now use the following well-known theorem.
Let $\left(\psi_{m}\right)$ be a bounded sequence in a Hilbert space $H$ with inner product $(\cdot, \cdot)_{H}$. Then there exists a subsequence ( $\psi_{m_{*}}$ ) and $\psi \in H$ such that for all $\phi \in H$

$$
\lim _{k \rightarrow \infty}\left(\psi_{m_{4}}, \phi\right)_{H}=(\psi, \phi)_{H} \quad \text { (weak convergence !). }
$$

We conclude from (II.2.7.i) and (II.2.7.ii) that there exists a subsequence $\left(\bar{\theta}_{m_{2}}\right)$ which tends weakly to some $\bar{\theta}$ in the Hilbert space $W^{2,1}(Q)$ and for which also by (II.2.7.iii) the sequence ( $\bar{E}_{m_{4}}\left(\tilde{\theta}_{m_{*}}\right)$ ) tends to some $\hat{E}$ weakly in the Hilbert space $L_{2}(Q)$. Thus we derive the following

$$
\begin{aligned}
\forall_{\phi \in V}: & \lim _{k \rightarrow \infty} \int_{a}\left(\bar{E}_{m_{2}}\left(\bar{\theta}_{m_{2}}\right) \frac{\partial \phi}{\partial t}+\bar{\theta}_{m_{2}} \Delta \phi\right) d x d t= \\
& \left.=\int_{Q} \hat{E} \frac{\partial \phi}{\partial t}+\tilde{\theta} \Delta \phi\right) d x d t .
\end{aligned}
$$

On the other hand it can be proved (cf. page 33) that

$$
\begin{aligned}
& \int_{Q}\left(\bar{E}_{m}\left(\bar{\theta}_{m_{2}}\right) \frac{\partial \phi}{\partial t}+\bar{\theta}_{m_{t}} \Delta \phi\right) d x d t= \\
& \quad=\int_{\Omega_{0}} \phi \tilde{E}_{m_{t}}\left(\bar{\theta}_{m_{t}}\right) d x+\int_{\Gamma_{l}} f \frac{\partial \phi}{\partial n_{l}} d \sigma+\int_{\Gamma_{s}} g \frac{\partial \phi}{\partial n_{s}} d \sigma .
\end{aligned}
$$

Now $\tilde{E}_{m_{t}}\left(\tilde{\theta}_{m_{1}}\right)=\tilde{E}_{m_{t}}\left(K_{m_{t}}\left(\theta_{0}\right)\right)=E_{m_{t}}\left(\theta_{0}\right) \rightarrow E_{0}(k \rightarrow \infty)$.
It follows that $\hat{E}$ is a weak solution if we can prove that $\hat{E}=E(\theta)$ for $\theta=K^{\leftarrow}(\bar{\theta})$.
To show this we observe that the canonical injection from $W^{2,1}(Q)$ into $L_{2}(Q)$ is compact. So the sequence $\left(\bar{\theta}_{m_{*}}\right)$ tends to $\bar{\theta}$ in the norm of $L_{2}(Q)$. Consequently, taking a subsequence if necessary, $\tilde{\theta}_{m_{t}} \rightarrow \tilde{\theta}$ almost uniform in $Q$. Hence

$$
\tilde{E}_{m_{4}}\left(\bar{\theta}_{m_{*}}\right) \rightarrow \tilde{E}(\bar{\theta})=E(\theta) .
$$

Also uniqueness of the weak solution can be proved. In short, we sum some arguments for this.

Let $E_{1}$ and $E_{2}$ be weak solutions and suppose that the set $\Psi=\left\{(x, t) \in Q \mid E_{1}(x, t) \neq E_{2}(x, t)\right\}$ has non-zero measure. It follows that

$$
\forall_{\varrho \in V}: \int_{\Psi}\left[\left(E_{1}-E_{2}\right) \frac{\partial \phi}{\partial t}+\left(K\left(\beta\left(E_{1}\right)\right)-K\left(\beta\left(E_{2}\right)\right) \Delta \phi\right] d x d t=0 .\right.
$$

Put

$$
\sigma=\frac{K\left(\beta\left(E_{1}\right)\right)-K\left(\beta\left(E_{2}\right)\right)}{E_{1}-E_{2}},(x, t) \in \Psi .
$$

Then $0 \leq \sigma \leq \sigma_{0}$ for some constant $\sigma_{0}$.
We extend $\sigma$ to the whole of $Q$ as a bounded measurable function, whence

$$
\forall_{\phi \in V}: \int_{Q}\left(E_{1}-E_{2}\right)\left(\frac{\partial \phi}{\partial t}+\sigma \Delta \phi\right) d x d t=0 .
$$

Next consider a sequence of smooth functions in $Q$ such that $\sigma_{n} \geq \frac{1}{n}$ and $\sigma_{n} \rightarrow \sigma$. Then it can be proved that for each $\psi \in C^{\infty}(Q)$ there exists $\phi_{n} \in V$ such that

$$
\frac{\partial \phi_{n}}{\partial t}+\sigma_{n} \Delta \phi_{n}=\psi
$$

It follows that

$$
\int_{\boldsymbol{Q}}\left[\left(E_{1}-E_{2}\right) \psi+\left(E_{1}-E_{2}\right)\left(\sigma-\sigma_{n}\right) \Delta \phi_{n}\right] d x d t=0 .
$$

Taking the limit $n \rightarrow \infty$,

$$
\int_{Q}\left(E_{1}-E_{2}\right) \psi d x d t=0 .
$$

Being valid for all $\psi \in C_{c}^{\infty}(Q)$ we must have $E_{1}=E_{2}$.

## CHAPTER IV

## REACTION-DIFFUSION PROBLEMS

A substance diffuses through a medium and at the same time undergoes a chemical reaction which may involve heat absorption or heat release. The concentration of the diffusing substance is denoted by $c$ and the absolute temperature by $T$. Then $c$ and $T$ are usually described by the following system of (nonlinear) partial differential equations

$$
\begin{equation*}
\frac{\partial c}{\partial t}-d \Delta c=-A c^{m} \exp (-E / R T) \tag{IV.1}
\end{equation*}
$$

in the set $\{c>0\}$ and

$$
\begin{equation*}
C \frac{\partial T}{\partial t}-k \Delta T=Q A c^{m} \exp (-E / R T) \tag{IV.2}
\end{equation*}
$$

throughout the medium.
The right hand side of equation (IV.1) describes the rate at which the chemical is used in the reaction. Here $A$ is a constant, the so called pre-exponential factor, $m \geq 0$ is the order of the reaction, $E \geq 0$ the activation energy and $R$ the universal gas constant. The right hand side of the second equation is the rate at which heat is released ( $Q>0$ ) or absorbed ( $Q<0$ ). If $Q=0$ the reaction is isothermal, i.e. the temperature remains constant.
Let $\Omega$ be the domain occupied by the medium. A fundamental problem is whether $c$ can vanish identically over a subset $D$ of $\Omega, D$ is called a dead core. In this case, besides $c$ and $T$ one should determine the evolution of the free boundary $\partial D$.
We investigate the case of a one-dimensional stationary isothermal reaction-diffusion in a slab $0 \leq x \leq a$ with prescribed boundary conditions

$$
c(0, t)=c(a, t)=c_{0}>0 .
$$

Using the non-dimensional variable $u=c / c_{0}$, the differential equation reduces to

$$
\begin{equation*}
\frac{\partial^{2} u}{\partial x^{2}}=\lambda u_{+}^{m} \quad, \quad(\lambda>0) \tag{IV.3}
\end{equation*}
$$

where $u(0)=u(a)=1$.

As a preliminary we consider the problem in the half space $x>0$ imposing the conditions

$$
u(0)=1 \text { and } \lim _{x \rightarrow \infty} u(x)=0 .
$$

From (IV.3) we see that $\frac{d^{2} u}{d x^{2}}$ is nonnegative whence $\frac{d u}{d x}$ is nondecreasing.
Since $\lim _{x \rightarrow \infty} u(x)=0$, it follows that also $\lim _{x \rightarrow \infty} u^{\prime}(x)=0$.
Multiplying by $\frac{d u}{d x}$ and integrating over $(0, x)$ yields as long as $u(x)>0$

$$
\begin{equation*}
\frac{1}{2}\left(\left(u^{\prime}(x)\right)^{2}-\left(u^{\prime}(0)\right)^{2}\right)=\lambda\left((u(x))^{m+1}-1\right) /(m+1) . \tag{IV.4}
\end{equation*}
$$

We see that there are two possibilities


From (IV.4) we deduce by taking $x \rightarrow \infty$ (extending the validity of (IV.4) to the dead core in case (b))

$$
u^{\prime}(0)=-[2 \lambda / m+1]^{1 / 2}
$$

so that (IV.4) takes the form

$$
u^{\prime}(x)=-(2 \lambda / m+1)^{1 / 2}\left(u_{+}(x)\right)^{m+1 / 2}
$$

As long as $u>0$ separation, of variables is permitted and yields

$$
\begin{aligned}
\lambda^{1 / 2} x & =\left[\frac{m+1}{2}\right]^{1 / 2} \int_{u}^{1} z^{-(m+1) / 2} d z \\
& =\left[\frac{m+1}{2}\right]^{1 / 2}\left[\begin{array}{ll}
\log u \text { as } m=1 \\
\frac{2}{m-1}\left(u^{-(m-1) / 2}-1\right)
\end{array} \text { as } m \neq 1 .\right.
\end{aligned}
$$

Taking the limit $u \downarrow 0$ we see that $x \rightarrow \infty$ as $m \geq 1$ and $x \rightarrow \lambda^{-1 / 2}(2(m+1))^{1 / 2} /(1-m)$ as $m<1$. It follows that for $m \geq 1$ the solution has to be positive everywhere, while for $m<1$ a dead core appears at
(IV.5) $\quad x_{d}=\lambda^{-1 / 2}(2(m+1))^{1 / 2} /(1-m)$.

Coming back to our original problem for a slab $0 \leq x \leq a$ we can now say that a dead core can be expected if and only if $a>2 x_{d}$.

In case $m=0$ we deal with the so called oxygen-diffusion consumption problem. The term $c_{+}^{m}$ has to be replaced by the Heaviside function $H(c)$ taking the value 1 if $c>0$ and the value 0 otherwise.
The non-stationary problem is classically described as follows
(IV.6)
(i)

$$
\frac{\partial c}{\partial t}-\frac{\partial^{2} c}{\partial x^{2}}=-1,0<x<s(t), 0<t<T
$$

(ii) $c(x, 0)=u_{0}(x) \geq 0,0<x<1$,
(iii) $\frac{\partial c}{\partial x}(0, t)=0,0<t<T$,
(iv) $\quad s(0)=1$,
(v) $\quad c(s(t), t)=0,0<t<T$,
(vi) $\quad \frac{\partial c}{\partial x}(s(t), t)=0,0<t<T$.

For the corresponding stationary problem we know that a dead core appears at $x_{d}=(2 / \lambda)^{1 / 2}=\sqrt{2} \quad(\lambda=1)$.


The weak formulation consists in writing the equation $\frac{\partial c}{\partial t}-\frac{\partial^{2} c}{\partial x^{2}}=-H(c)$ in the whole domain, and it is equivalent to a variational inequality.

Put $u=\frac{\partial c}{\partial t}$. Then $\frac{\partial^{2} c}{\partial x^{2}}=u+1$ for $0<x<s(t)$ and $0<t<T$. It can be checked that $u$ satisfies the following equations
(IV.7)
(i)

$$
\frac{\partial u}{\partial t}-\frac{\partial^{2} u}{\partial x^{2}}=0,0<x<s(t), 0<t<T,
$$

(ii) $u(x, 0)=u_{0}(x)=c_{0}^{\prime \prime}(x)-1$,
(iii) $\frac{\partial u}{\partial x}(0, t)=0$,
(iv) $\quad s(0)=1$
(v) $\quad u(s(t), t)=0$
(vi) $\quad \frac{\partial u}{\partial x}(s(t), t)=-\dot{s}(t)$.

## Explanation.

* $\left.\quad 0=\frac{d}{d t}(c(s(t), t))=\dot{s}(t) \frac{\partial c}{\partial x}(x, t) \right\rvert\, x=s(t)+\frac{\partial c}{\partial t}(s(t), t)$

$$
=u(s(t), t)
$$

* $\left.\quad 0=\frac{d}{d t}\left(\frac{\partial c}{\partial x}(s(t), t)\right)=\dot{s}(t) \frac{\partial^{2} c}{\partial x^{2}}(x, t) \right\rvert\, x=s(t)+\frac{\partial}{\partial x}\left(\frac{\partial c}{\partial t}(s(t), t)\right)$

$$
\begin{aligned}
& =\dot{s}(t)(u(x, t)+1) \left\lvert\, x=s(t)+\frac{\partial u}{\partial x}(s(t), t)=\right. \\
& =\dot{s}(t)+\frac{\partial u}{\partial x}(s(t), t)
\end{aligned}
$$

If $c_{0}{ }^{\prime \prime}-1<0$, then $u_{0}<0$ and it follows that $u<0$. In this case the equations for $u$ corresponds to the Stefan model for the solidification of a supercooled fluid.

We have

$$
\begin{aligned}
c(x, t) & =\int_{x}^{s(t)} d \xi \int_{\xi}^{s(t)} \frac{\partial^{2} c}{\partial x^{2}}(\eta, t) d \eta= \\
& =\int_{x}^{s(t)} d \xi \int_{\xi}^{s(t)}[u(\eta, t)+1] d \eta,
\end{aligned}
$$

whence

$$
c_{0}(x)=\int_{x}^{1} d \xi \int_{\xi}^{1}\left[u_{0}(\eta)+1\right] d \eta
$$

and

$$
c_{0}^{\prime}(0)=-\int_{0}^{1}\left[u_{0}(\eta)+1\right] d \eta=-Q
$$

We distinguish three cases
A.

There exists a global solution for $0<x<s(t), t>0$.
B.


There is extinction in a finite time: $\exists_{T}: s(T)=0$.

C.

There is a blow up
$\exists_{T}: \dot{s}(t) \rightarrow-\infty, t \rightarrow T$
It can be proved that


$$
\begin{aligned}
& A \Rightarrow Q>0 \\
& B \Rightarrow Q=0 \\
& C \Leftarrow Q<0
\end{aligned}
$$

Suppose a negativity set is formed. Of course then we are dealing with a wrong model because $c$ has to be nonzero. But a mathematical meaning is there.


The following statements are valid.
(1) The negativity set expands.
(2) The negativity set is bound to meet the free boundary in finite time.
(3) The meeting point is a point of essential blow up.
(4) There are no other cases of essential blow up.

Physically, essential blow ups do not occur.
When $c$ becomes zero in a point, one has to introduce a new dead core and solve a new free boundary problem.

## CHAPTER V

## BINGHAM FLUIDS

Consider a Newtonian fluid flow in some region of $\mathbb{R}^{2}$ with velocity field $\vec{v}=\vec{v}(x, t)$. Let $p=p(x, t)$ denote the pressure of the fluid. Then $p$ and $\vec{v}$ satisfy the Navier-Stokes equations

$$
\begin{equation*}
\rho \frac{\partial \vec{v}}{\partial t}+\frac{1}{2}(\nabla \vec{v})^{2}+(\operatorname{curl} \vec{v}) \times \vec{v}=-\nabla p+\eta \Delta \vec{v} \tag{V.1}
\end{equation*}
$$

where $\rho$ denotes the homogeneous mass density and $\eta$ the viscosity of the fluid. If the fluid is incompressible we have $\operatorname{div} \vec{v}=0$.

Consider a flow given by $\vec{v}=(v(x, y, t), 0)$ and assume that the fluid is incompressible. Then the Navier-Stokes equations reduce to

$$
\begin{align*}
\rho \frac{\partial v}{\partial t} & =-\frac{\partial p}{\partial x}+\eta \frac{\partial^{2} v}{\partial y^{2}}  \tag{V.2}\\
\frac{\partial p}{\partial y} & =0 \\
\frac{\partial v}{\partial x} & =0
\end{align*}
$$

It follows that $v(x, y, t)=v(y, t)$ and $p(x, y, t)=p_{1}(t) x+p_{0}(t)$.

For a Newtonian fluid the stress tensor is related to the gradient of the velocity field in the following way

$$
\begin{equation*}
\mathrm{T}=\eta\left(\nabla \vec{v}+(\nabla \vec{v})^{T}\right) . \tag{V.3}
\end{equation*}
$$

So in our particular situation

$$
\begin{aligned}
& \tau_{x x}=2 \eta \frac{\partial v}{\partial x}=0, \quad \tau_{y y}=0 \\
& \tau_{x y}=\eta \frac{\partial v}{\partial y} \equiv \eta \sigma .
\end{aligned}
$$

Here $\tau \equiv \tau_{x y}$ is the shear stress and $\sigma \equiv \frac{\partial v}{\partial y}$ is the strain rate.

## Example.

Consider the stationary flow of an incompressible Newtonian fluid between two (infinitely long) plates under the assumption that there are no volume forces.


For the velocity field we take $\vec{v}(x, y)=(v(x, y), 0)$, i.e. a Poiseuille flow. Adding no-slip conditions at the boundary, i.e. $v(x, L)=v(x,-L)=0$ we obtain the solution

$$
\begin{aligned}
& p(x, y)=p_{0} x+p_{1} \\
& v(x, y)=\frac{p_{0}}{2 \eta}\left(L^{2}-y^{2}\right) .
\end{aligned}
$$

Observe that in this case $\tau=\tau(y)=-p_{0} y$.
A Bingham fluid is a non-Newtonian fluid characterized by the presence of a threshold value $\tau_{0}$ for the shear stress, such that if the shear stress $\tau$ is less than $\tau_{0}$ the fluid behaves like a rigid body, while for $\tau>\tau_{0}$ it behaves as a fluid where the relationship between shear stress and strain rate is linear, i.e.
(V.4) $\quad \tau=\tau_{0}+\eta \sigma$.

The dynamics of a Bingham fluid, described by a velocity field $\vec{v}$ obeys the Navier-Stokes equations in the region $\left\{\tau>\tau_{0}\right\}$, while on the boundary with the rigid core (the free boundary) we have $\tau=\tau_{0}$, i.e. zero strain rate. Another free boundary condition results from the balance of momentum.

Here we consider an incompressible Bingham fluid flowing between two parallel plates. Again $x$ denotes the coordinate along the direction of motion and $y$ the coordinate in the direction perpendicular to the plates. So in the representative $x y$-plane the velocity has the form $\vec{v}=(v(y, t), 0)$ and the equation of motion in the viscous region equals

$$
\begin{equation*}
\rho \frac{\partial v}{\partial t}=-\frac{\partial p}{\partial x}+\eta \frac{\partial^{2} v}{\partial y^{2}}, \frac{\partial p}{\partial y}=0 . \tag{V.5}
\end{equation*}
$$

Consequently,

$$
\begin{equation*}
\frac{\partial p}{\partial x}=-f(t), f(t)>0 \tag{V.6}
\end{equation*}
$$

which we assume to be given.


Let $y= \pm s(t)$ be the equation for the free boundary of the rigid core. Then the velocity field satisfies the parabolic equation

$$
\rho \frac{\partial v}{\partial t}-\eta \frac{\partial^{2} v}{\partial y^{2}}=f(t)
$$

in the regions $-L<y<-s(t)$ and $s(t)<y<L$. By symmetry we only need to consider the upper half layer. We impose the no-slip condition

$$
v(L, t)=0, t>0
$$

and some initial condition

$$
v(y, 0)=v_{0}(y)
$$

such that $v_{0}(y)$ is constant for $0<y<s(0), v_{0}(L)=0$.

Since $\sigma=0$ at the free boundary

$$
\begin{equation*}
\left.\frac{\partial v}{\partial y}\right|_{x=s(t)}=0 \tag{V.7}
\end{equation*}
$$

For the second free boundary condition we apply Newton's law. Consider a portion of the rigid core situated between two unit squares parallel to the plates


The driving force equals $2 s(t) f(t)-2 \tau_{0}$. The mass of the portion equals $2 s(t) \rho$. Hence
(V.8) $\left.\quad 2 s(t) \rho \frac{\partial v}{\partial t}\right|_{x=s(t)}=2 s(t) f(t)-2 \tau_{0}$
or equivalently

$$
\left.\frac{\partial v}{\partial t}\right|_{x=s(t)}=\frac{1}{\rho} f(t)-\frac{\tau_{0}}{\rho s(t)}
$$

The free boundary conditions are neither of Cauchy nor of Stefan type.

As a special case we consider stationary solutions for $f(t)=f_{0}$ with $f_{0}>\tau_{0} / L$. Then $\frac{\partial v}{\partial t}=0$ and so by (V.8), $s(t)=\tau_{0} / f_{0}$. The equations reduces to

$$
\begin{align*}
& \frac{d^{2} v}{d y^{2}}=\frac{1}{\eta} f_{0}, \frac{\tau_{0}}{f_{0}}<y<L  \tag{V.9}\\
& v(L)=0 \\
& \frac{d v}{d y}\left[\frac{\tau_{0}}{f_{0}}\right]=0 .
\end{align*}
$$

Hence

$$
v(y)=-\frac{f_{0}}{\eta}(L-y)\left(L+y-\frac{2 \tau_{0}}{f_{0}}\right), \frac{\tau_{0}}{y_{0}}<y<L .
$$

For the non-stationary problem we take $w=\frac{\partial v}{\partial y}$. Then for $w$ we get
(i) $\quad \rho \frac{\partial w}{\partial t}-\eta \frac{\partial^{2} w}{\partial y^{2}}=0$,
(ii) $\quad \frac{\partial w}{\partial y}(L, t)=-\frac{1}{\eta} f(t)$,
(iii) $\quad w(y, 0)=v_{0}^{\prime}(y)$,
(iv) $\quad w(s(t), t)=0$,
(v) $\quad \eta \frac{\partial w}{\partial y}(s(t), t)=-\tau_{0} / s(t)$.
(This is a free boundary problem with Cauchy data on the free boundary.)

## Explanation.

We have $\rho \frac{\partial v}{\partial t}-\eta \frac{\partial^{2} v}{\partial y^{2}}=f(t)$. So at $y=s(t)$

$$
f(t)-\frac{\tau_{0}}{s(t)}-\left.\eta \frac{\partial w}{\partial y}\right|_{y=s(t)}=f(t)
$$

so that

$$
\left.\eta \frac{\partial w}{\partial y}\right|_{y=s(t)}=-\frac{\tau_{0}}{s(t)}
$$

## Another transformation is the following

$$
z=\frac{\partial v}{\partial t}=\frac{\eta}{\rho} \frac{\partial w}{\partial y}+\frac{1}{\rho} f
$$

yielding the Stefan type problem
(V.11)
(i)

$$
\rho \frac{\partial z}{\partial t}-\eta \frac{\partial^{2} z}{\partial y^{2}}=\dot{f}, s(t)<y<L, t>0
$$

(ii) $z(L, t)=0, t>0$,
(iii) $\quad z(y, 0)=\frac{\eta}{\rho} v_{0}{ }^{\prime \prime}+\frac{1}{\rho} f(0), s(0)<y<L$
(iv) $\quad z(s(t), t)=\left(f(t)-\frac{\tau_{0}}{s(t)}\right) / \rho, t>0$
(v) $\quad \frac{\partial z}{\partial y}(s(t), t)=\frac{1}{\eta} \frac{\tau_{0}}{s(t)} \dot{s}(t), t>0$.

## Appendix A

## Some concepts and theorems of vector analysis in $\mathbb{R}^{2}$ and $\mathbb{R}^{3}$

## A.1. Curve.

$$
\underline{x}=\underline{x}(t)=\left(x_{1}(t), x_{2}(t), x_{3}(t)\right) .
$$

The curve is called smooth if $x_{1}(t), x_{2}(t)$ and $x_{3}(t)$ are differentiable. A smooth curve has a tangent

$$
\underline{\dot{x}}(t)=\left(\dot{x}_{1}(t), \dot{x}_{2}(t), \dot{x}_{3}(t)\right) .
$$

## A.2. Surface.

$$
\underline{x}=\underline{x}=\left(x_{1}(u, v), x_{2}(u, v), x_{3}(u, v)\right), \underline{t}=\frac{\dot{x}(t)}{|\dot{x}(t)|} .
$$

The surface is said to be smooth if $x_{1}(u, v), x_{2}(u, v)$ and $x_{3}(u, v)$ are differentiable. The tangent surface at $\underline{x}=\underline{a}$ is given by

$$
\underline{y}=\underline{a}+\lambda \frac{\partial x}{\partial u}(\underline{a})+\mu \frac{\partial x}{\partial v}(\underline{a})
$$

or correspondingly

$$
\left.\underline{( }, \frac{\partial x}{\partial u}(a) \times \frac{\partial x}{\partial v}(a)\right)=\left(\underline{a}, \frac{\partial x}{\partial u}(\underline{a}) \times \frac{\partial x}{\partial v}(\underline{a})\right) .
$$

So the normal of the surface at $\underline{x}=\underline{a}$ equals

$$
\underline{n}(\underline{a})= \pm \frac{\frac{\partial x}{\partial u}(\underline{a}) \times \frac{\partial x}{\partial v}(\underline{a})}{\left|\frac{\partial x}{\partial u}(\underline{a}) \times \frac{\partial x}{\partial v}(\underline{a})\right|} .
$$

Special case $z=z(x, y), \underline{x}=(x, y, z(x, y))$.
Take $u=x$ and $v=y$. Then

$$
\frac{\partial x}{\partial u}=\left(1,0, \frac{\partial z}{\partial x}\right), \frac{\partial x}{\partial v}=\left(0,1, \frac{\partial z}{\partial y}\right) .
$$

So

$$
\frac{\partial x}{\partial u} \times \frac{\partial x}{\partial v}=\left[-\frac{\partial z}{\partial x},-\frac{\partial z}{\partial y}, 1\right] .
$$

## A.3. Scalar field.

Let $\Omega \subset \mathbb{R}^{3}$. A function $\phi$ from $\Omega$ in $\mathbb{R}$ is said to be a scalar field. Notation: $\phi(\underline{x})$ or $(\underline{x}, \phi(\underline{x})$ ). Fysical examples: temperature, mass density, pressure. Surfaces of the form $\phi(\underline{x})=c$ are called equiscalar (or equipotential) surfaces.
A scalar field is said to be differentiable at $\underline{x} \in \Omega$ if there exists a linear functional $L(\underline{x}): \mathbb{R}^{3} \rightarrow \mathbb{R}$ such that

$$
\phi(\underline{x}+\underline{h})-\phi(\underline{x})=L(\underline{x}) \underline{h}+o(|\underline{h}|) .
$$

If $(\cdot, \cdot)$ denotes the Euclidean inner product we can write

$$
L(\underline{x}) \underline{h}=\nabla \phi(\underline{x}) \cdot \underline{h}
$$

for some vector $\nabla \phi(x) \in \mathbb{R}^{3}$. The mapping

$$
\nabla \phi: \Omega \rightarrow \mathbb{R}^{3}
$$

is called the gradient of $\phi$. We have in cartesian coordinates

$$
(\nabla \phi)(\underline{a})=\left[\frac{\partial \phi}{\partial x_{1}}(\underline{a}), \frac{\partial \phi}{\partial x_{2}}(\underline{a}), \frac{\partial \phi}{\partial x_{3}}(\underline{a})\right], \underline{a} \in \Omega .
$$

Let $\underline{v} \in \mathbb{R}^{3}$ with $|\underline{v}|=1$. Then the directional derivative

$$
\lim _{t \rightarrow 0} \frac{\phi(x+t v)-\phi(x)}{t}
$$

equals

$$
\nabla \phi(\underline{x}) \cdot \underline{v} .
$$

It follows that the directional derivative is maximal or minimal if $\underline{v}=\lambda \nabla \phi(\underline{x})$. So $\nabla \phi \underline{x})$ points in the direction of maximal increase or decrease.
Consider the equiscalar surface $\phi(\underline{x})=\phi(\underline{a})$. Taking $z=z(x, y)$ we obtain

$$
\frac{\partial \phi}{\partial x}+\frac{\partial \phi}{\partial z} \frac{\partial z}{\partial x}=0, \frac{\partial \phi}{\partial y}+\frac{\partial \phi}{\partial z} \frac{\partial z}{\partial y}=0 .
$$

So the tangent plane at $\underline{x}=\underline{a}$ is spanned by

$$
\left[-\frac{\partial \phi}{\partial z}(\underline{a}), 0, \frac{\partial \phi}{\partial x}(\underline{a})\right] \text { and }\left[0,-\frac{\partial \phi}{\partial z}(\underline{a}), \frac{\partial \phi}{\partial y}(\underline{a})\right]
$$

with normal $\nabla \phi(\underline{a}) /|\nabla \phi(\underline{a})|$.
So the equation for the tangent plane is given by

$$
\nabla \phi(\underline{a}) \cdot(\underline{y}-\underline{a})=0 .
$$

## A.4. Vector field.

Let $\Omega \subset \mathbb{R}^{3}$. A vector field is a function from $\Omega$ into $\mathbb{R}^{3}$. Notation $\underline{v}(\underline{x})$ or $(\underline{x}, \underline{v}(\underline{x}))$. The second notation suggests that at each point of $\Omega$ a vector is attached. Fysical examples: electro-magnetic field, velocity field of a fluid. In particular for each scalar field $\phi$, the gradient $\nabla \phi$ is a vector field. A curve $\underline{x}(t)$ with the property that

$$
\underline{\dot{x}}(t)=\lambda \underline{v}(\underline{x}(t))
$$

is called a stream line or field line.
A vector field $\underline{v}$ is said to be differentiable at $\underline{x} \in \Omega$ if there exists a linear mapping $A(\underline{x}): \mathbb{R}^{3} \rightarrow \mathbb{R}^{3}$ such that

$$
\underline{v}(\underline{x}+\underline{h})-\underline{v}(\underline{x})=A(\underline{x}) \underline{h}+o(|\underline{h}|) .
$$

## A.5. Operations on vector fields.

Given a vector field $\underline{v}=\left(v_{1}, v_{2}, v_{3}\right)$ which is differentiable in a region $\Omega \subset \mathbb{R}^{3}$.
curl $\underline{v} \mapsto \operatorname{curl} \underline{v}$

$$
\operatorname{curl} \underline{v}=\nabla \times \underline{v}=\left[\frac{\partial v_{3}}{\partial x_{2}}-\frac{\partial v_{2}}{\partial x_{3}}, \frac{\partial v_{1}}{\partial x_{3}}-\frac{\partial v_{3}}{\partial x_{1}}, \frac{\partial v_{2}}{\partial x_{1}}-\frac{\partial v_{1}}{\partial x_{2}}\right]
$$

$-\quad \operatorname{curl}(\underline{v}+\underline{w})=\operatorname{curl} \underline{v}+\operatorname{curl} \underline{w}$
$-\quad \operatorname{curl}(\phi \underline{v})=\phi \operatorname{curl} \underline{v}+\nabla \phi \times \underline{v}$
$-\quad \operatorname{curl}(\nabla \phi)=0$
Divergence. $\underline{v} \rightarrow \operatorname{div} \underline{v}$

$$
\operatorname{div} \underline{v}=\nabla \cdot \underline{v}=\frac{\partial v_{1}}{\partial x_{1}}+\frac{\partial v_{2}}{\partial x_{2}}+\frac{\partial v_{3}}{\partial x_{3}}
$$

$-\quad \operatorname{div}(\underline{v}+\underline{w})=\operatorname{div} \underline{v}+\operatorname{div} \underline{w}$
$-\quad \operatorname{div}(\phi \underline{v})=\phi \operatorname{div} \underline{v}+\nabla \phi \cdot \underline{v}$
$-\quad \operatorname{div}(\underline{v} \times \underline{w})=\underline{w} \cdot \operatorname{curl} \underline{v}-\underline{v} \cdot \operatorname{curl} \underline{w}$

- $\quad \operatorname{div} \operatorname{curl} \underline{v}=0$.

Further for a twice differentiable scalar field (vector field) $\phi(\underline{v})$ we have

$$
\begin{aligned}
& \Delta \phi=\operatorname{div}(\nabla \phi)=\frac{\partial^{2} \phi}{\partial x_{1}^{2}}+\frac{\partial^{2} \phi}{\partial x_{2}^{2}}+\frac{\partial^{2} \phi}{\partial x_{3}^{2}} \\
& \Delta \underline{v}=\left(\Delta v_{1}, \Delta v_{2}, \Delta v_{3}\right)
\end{aligned}
$$

$-\quad \nabla(\operatorname{div} \underline{v})=\operatorname{curl}(\operatorname{curl} \underline{v})+\Delta \underline{v}$
$-\quad \nabla(\Delta \phi)=\Delta(\nabla \phi)$
$-\operatorname{curl}(\Delta \underline{v})=\Delta(\operatorname{curl} \underline{v})$
$-\quad \operatorname{div}(\Delta \underline{v})=\Delta(\operatorname{div} \underline{v})$.

## A.6. Line-integrals, surface integrals.

Given a smooth curve $K, \underline{x}=\underline{x}(t), a \leq t \leq b$


$$
\text { Then } \begin{aligned}
\Delta s & \approx|x(t+\Delta t)-x(t)| \\
& \approx|\dot{x}(t)| \Delta t
\end{aligned}
$$

In the limit $\Delta t \downarrow 0$ we get

$$
d s=|\dot{x}(t)| d t
$$

length of $K: \int_{K} d s=\int_{a}^{b}|\dot{x}(t)| d t$

For a scalar field $\phi$ we define

$$
\int_{K} \phi d s=\int_{a}^{b} \phi(\underline{x}(t)) \underline{\dot{x}}(t) d t
$$

Let $\underline{t}(\underline{x})$ denote the normalized tangent vector at $\underline{x} \in K, \underline{t}=\frac{\dot{\dot{x}}}{|\underline{\dot{x}}|}$. Then for a vector field $\underline{v}$ we have by taking $\phi(x)=(\underline{v}(\underline{x}), \underline{t}(\underline{x}))$

$$
\int_{K}(\underline{v} \cdot \underline{t}) d s=\int_{a}^{b}(\underline{v}(\underline{x}(t)) \cdot \underline{\dot{x}}(t)) d t
$$

## Important special case.

Take $\underline{v}=\nabla \phi$. Then

$$
((\nabla \phi) \underline{x}(t) \cdot \underline{\dot{x}}(t))=\frac{d}{d t}(\phi(\underline{x}(t))
$$

So that

$$
\int_{K}(\nabla \phi \cdot \underline{t}) d s=\phi\left(x_{\ell}\right)-\phi\left(x_{s}\right)
$$

where $\underline{x}_{e}=\underline{x}(b)$ is the endpoint of $K$ and $\underline{x}_{s}=\underline{x}(a)$ its startingpoint.
In particular, if $K$ is a closed curve then

$$
\oint_{K}(\nabla \phi \cdot t) d s=0 .
$$

Gradient fields are said to be conservative. For instance if a force field $\mathbf{F}=\nabla \boldsymbol{\phi}$, the energy needed to go from $\underline{p}$ to $\underline{q}$ does not depend on the road being followed.

Let there be given a smooth surface $S$,

$$
\underline{x}=\underline{x}(u, v),(u, v) \in P \subset \mathbb{R}^{2}
$$



$$
d \sigma=\left|\frac{\partial x}{\partial u} \times \frac{\partial x}{\partial v}\right| d u d v
$$

For a scalar field $\phi$ we define

$$
\iint_{S} \phi d \sigma=\iint_{P} \phi(\underline{x}(u, v))\left|\frac{\partial x}{\partial u} \times \frac{\partial \underline{x}}{\partial v}\right| d u d v .
$$

If the surface $S$ is given by $z=z(x, y)$, we get

$$
\iint_{S} \phi d \sigma=\iint_{P} \phi(x, y, z(x, y)) \sqrt{1+\left[\frac{\partial z}{\partial x}\right]^{2}+\left[\frac{\partial z}{\partial y}\right]^{2}} d x d y .
$$

Let $\underline{w}$ be a vector field. Suppose the surface $S$ is orientable, i.e. the normal $\underline{n}: S \rightarrow \mathbb{R}^{3}$

$$
\underline{n}= \pm \frac{\partial x}{\partial u} \times \frac{\partial x}{\partial v} /\left|\frac{\partial x}{\partial u} \times \frac{\partial x}{\partial v}\right|
$$

can be taken continuous on $S$. So the plus or minus sign is fixed by giving the direction of the normal in- one point of $S$. Let $S$ thus be given an orientation. Then the scalar field ( $\underline{w} \cdot \underline{n}$ ) is well defined on $S$. We have

$$
\iint_{S}(\underline{w} \cdot \underline{n}) d \sigma= \pm \iint_{P}\left(\underline{w} \cdot \frac{\partial x}{\partial u} \times \frac{\partial \underline{x}}{\partial v}\right) d u d v
$$

$\uparrow$ depending on the orientation of $S$.

## A.7. Gauss' integral relation.

Let $\Omega$ be a bounded domain in $\mathbb{R}^{3}$ with piecewise smooth, orientable boundary $\partial \Omega$. Let $\underline{n}$ denote the outwardly directed normal on $S$. Then for a differentiable vector field $\underline{\boldsymbol{w}}: \Omega \rightarrow \mathbb{R}^{\mathbf{3}}$

$$
\iint_{\Omega} \int_{\operatorname{div}} \underline{w} d \underline{x}=\iint_{\partial \Omega}(\underline{w} \cdot \underline{n}) d \sigma .
$$

## Consequences.

$$
\begin{array}{ll} 
& \iiint_{\Omega} \frac{\partial w_{j}}{\partial x_{j}} d x=\iint_{\Omega} w_{j} n_{j} d \sigma \\
& \iiint_{\Omega} \nabla \phi d x=\iint_{\partial \Omega} \phi \underline{n} d \sigma \\
& \iiint_{\Omega} \operatorname{rot} \underline{w} d x=\iint_{\partial \Omega}(\underline{n} \times \underline{w}) d \sigma \\
& \iiint_{\Omega} \Delta \phi d x=\iiint_{\partial \Omega}(\nabla \phi \cdot \underline{n}) d \sigma .
\end{array}
$$

Gauss formula also holds for bounded regions $\Omega \subset \mathbb{R}^{2}$ and differentiable vector fields $\underline{w}: \Omega \rightarrow \mathbb{R}^{2}$.
To this end define

$$
\begin{aligned}
& \bar{\Omega}=\Omega \times[0,1] \subset \mathbb{R}^{3} \\
& \underline{\bar{w}}=\left(w_{1}, w_{2}, 0\right) .
\end{aligned}
$$

Then $\underline{n}=\left(\underline{n}, \underline{n}_{2}, 0\right)$ and

$$
\begin{aligned}
& \iint_{\Omega}\left[\frac{\partial w_{1}}{\partial x_{1}}+\frac{\partial w_{2}}{\partial x_{2}}\right] d \underline{x}=\iint_{\Omega} \operatorname{div} \bar{w} d \underline{x}= \\
& \quad=\int_{\partial \Omega}(\tilde{w} \cdot \bar{n}) d \sigma=\int_{\partial \Omega}(\underline{w} \cdot \underline{n}) d s .
\end{aligned}
$$

## Example.

Consider a fluid with mass density $\rho(\underline{x}, t)$ which flows with velocity $\underline{v}(\underline{x}, t)$ through a surface $S$ which is the boundary of a bounded region $\Omega$.
The amount of fluid streaming out of $\Omega$ during a unit of time $\Delta t$ can be approximated by

$$
\left(\int_{S}(\rho \underline{v} \cdot \underline{n}) d \sigma\right) \Delta t .
$$

Further, the change of mass of the volume $\Omega$

$$
\iiint_{\Omega}\left(\rho(\underline{x}, t+\Delta t)-\rho(\underline{x}, t) d \underline{x} \approx\left[\iint_{\Omega} \frac{\partial \rho}{\partial t} d \underline{x}\right] \Delta t .\right.
$$

Thus we find

$$
\iiint_{\Omega} \frac{\partial \rho}{\partial t} d \underline{x}=-\iint_{S}(\rho \underline{v} \cdot \underline{n}) d \sigma
$$

and applying Gauss divergence theorem yields

$$
\iiint_{\Omega}\left(\frac{\partial \rho}{\partial t}+\operatorname{div} \rho \underline{v}\right) d \underline{x}=0
$$

Since $\Omega$ is an arbitrarily taken region in $\mathbb{R}^{3}$ we end up with the continuity equation

$$
\frac{\partial \rho}{\partial t}+\operatorname{div} \rho \underline{v}=0
$$

In the special case that $\rho$ is constant we deal with an incompressible fluid yielding $\operatorname{div} \underline{\boldsymbol{v}}=\mathbf{0}$.

## A.8. Stokes' integral relation.

Let $S$ be a piecewise smooth, orientable surface in $\mathbb{R}^{3}$ with a piecewise smooth, simply connected, closed curve $K$ as its boundary. Let $\underline{v}$ be a differentiable vector field on $S$

$$
\iint_{S}(\operatorname{rot} \underline{v} \cdot \underline{n}) d \sigma=\oint_{K}(\underline{a} \cdot \underline{t}) d s
$$

where $\underline{t}$ and $\underline{n}$ fit together in a counterclockwise manner. If we take $S \subset \mathbb{R}^{2}$, we get Green's integral formula

$$
\left.\oint_{K}\left(v_{1} t_{1}+v_{2} t_{2}\right) d s=\iint_{S} \frac{\partial a_{2}}{\partial x_{1}}-\frac{\partial a_{1}}{\partial x_{2}}\right] d x_{1} d x_{2} .
$$

## Example.

An electromagnetic field has two components
the electric field $\underline{E}(\underline{x}, t)$
the magnetic induction $\underline{B}(\underline{x}, t)$.
In a region $\Omega$ without currents, $\underline{E}$ and $\underline{B}$ satisfy the following integral relation

$$
\begin{aligned}
& \oint_{K}(\underline{E} \cdot \underline{t}) d s=-\frac{d}{d t} \iint_{S}(\underline{B} \cdot \underline{n}) d \sigma \quad \text { (Faraday's law) } \\
& \oint_{K}(\underline{B} \cdot \underline{t}) d s=\varepsilon_{0} \mu_{0} \frac{d}{d t} \iint_{S}(\underline{E} \cdot \underline{n}) d \sigma \quad \text { (Maxwell's law) }
\end{aligned}
$$

where $\varepsilon_{0}$ denotes the permittivity and $\mu_{0}$ the permeability of vacuum. Further any piecewise smooth, orientable surface $S$ in $\Omega$ with piecewise smooth closed boundary $K$ may be taken. Applying Stokes we obtain Maxwell's equations

$$
\operatorname{rot} \underline{E}=-\frac{\partial B}{\partial t}, \operatorname{rot} \underline{B}=\varepsilon_{0} \mu_{0} \frac{\partial E}{\partial t} \text { in } \Omega .
$$

## A.9. Green's identities.

Let $\Omega$ be a bounded region in $\boldsymbol{R}^{3}$ such that $\partial \Omega$ is piecewise smooth and orientable. By $\underline{n}$ the outwardly directed normal on $\partial \Omega$ is denoted.
Let $u \in C^{2}(\Omega) \cap C^{1}(\bar{\Omega})$, i.e. $u$ is a scalar field from $\bar{\Omega}$ in $\mathbb{R}$ such that $u: \Omega \rightarrow \mathbb{R}$ is twice differentiable and $u: \partial \Omega \rightarrow \mathbb{R}$ is differentiable. Applying Gauss divergence theorem to $\Delta u$
yields

$$
\iint_{\Omega} \Delta u d \underline{x}=\iint_{\partial \Omega}(\nabla u \cdot \underline{n}) d \sigma .
$$

In stead of $\nabla u \cdot \underline{n}$ one writes $\frac{\partial u}{\partial n}$, called the normal derivative of $u$ on $\partial \Omega$.

## Green's first identity.

Let $u \in C^{2}(\Omega) \cap C^{1}(\bar{\Omega})$ and $v \in C^{1}(\Omega)$.

$$
\iint_{\Omega}(v \Delta u+(\nabla u \cdot \nabla v)) d \underline{x}=\iint_{\partial \Omega} v \frac{\partial u}{\partial n} d \sigma .
$$

(Observe that $\nabla \cdot(v \nabla u)=v \Delta u+\nabla u \cdot \nabla v$.

## Green's second identity.

Let $u, v \in C^{2}(\Omega) \cap C^{1}(\bar{\Omega})$.

$$
\iiint_{\Omega}(v \Delta u-u \Delta v) d x=\iint_{\partial \Omega}\left[v \frac{\partial u}{\partial n}-u \frac{\partial v}{\partial n}\right] d \sigma .
$$

For a function $u \in C^{2}(\Omega) \cap C^{1}(\bar{\Omega})$ which is harmonic, i.e. $\Delta u=0$ in $\Omega$, it follows that

$$
\iint_{\partial \Omega} \frac{\partial u}{\partial n} d \sigma=0 \text { and } \iint_{\Omega}|\nabla u|^{2} d \tau=\iint_{\partial \Omega} u \frac{\partial u}{\partial n} d \sigma .
$$

We apply these results to the boundary problems:
Dirichlet: $\Delta u=-f$ in $\Omega, u=g$ on $\partial \Omega$
Neumann: $\Delta u=-f$ in $\Omega, \frac{\partial u}{\partial n}=h$ on $\partial \Omega$
with respect to the uniqueness of their solutions.

- $\Delta u=0$ in $\Omega$ and $u=0$ on $\partial \Omega$ :

$$
\iint_{\Omega}|\nabla u|^{2} d x=0
$$

whence $\nabla u=0$, i.e. $u$ is constant in $\Omega$.
Since $u=0$ on $\partial \Omega$, we obtain $u \equiv 0$.

- $\Delta u=0$ in $\Omega, \frac{\partial u}{\partial n}=0$ on $\partial \Omega$.

We derive similarly that $u \equiv c, c$ a constant. So the solution of the Neumann problem is unique up to a constant.

## Appendix B

## Flow of a Newtonian fluid

Each continuum theory is based on two systems of laws

- universal laws of balance, e.g. balance of mass, balance of momentum.
- constitutive laws or material laws, which are specific for the considered material, e.g. Hooke's law for a linearly elastic medium, Ohm's law for a conductor.


## B.1. Balance of mass.

Consider a material body $B$ which occupies at each time $t$ a bounded region $\Omega(t)$ of $\mathbb{R}^{3}$. Each point $P \in B$ is at time $t$ described by a vector $\underline{x}=\underline{x}(P, t)$. Now conservation of mass says

$$
\int_{\Omega(t)} \rho(\underline{x}, t) d \underline{x}=\int_{\Omega(\tau)} \rho(\underline{x}, \tau) d \underline{x}
$$

where $\rho(\underline{x}, t)$ denotes the mass density of $B$ at time $t$. So we see that

$$
\frac{d}{d t}\left(\int_{\Omega(t)} \rho(\underline{x}, t) d \underline{x}\right)=0 .
$$

From this equation one can deduce the so called continuity equation

$$
\frac{d \rho}{\partial t}+\operatorname{div} \rho \underline{v}=0
$$

where $\underline{v}=\underline{\dot{x}}$ denotes the velocity $(\underline{v}(P, t)=\underline{\dot{x}}(P, t))$. In the special case of an incompressible homogeneous fluid, where $\rho(\underline{x}, t)=\rho_{0}$, a constant, we obtain the incompressibility condition

$$
\operatorname{div} \underline{v}=0 .
$$

For an incompressible medium both mass and volume (not the shape of the volume) are conserved.

Heuristic two-dimensional explanation.


$$
\underline{v}=\left(v_{x}, v_{y}\right)
$$

At time $t$ we have $P_{1}=(x, y), P_{2}=(x+d x, y), P_{3}=(x, y+d y) P_{4}=(x+d x, y+d y)$.

At time $t+d t$ we have $Q_{1} \simeq\left(x+v_{x}(x, y) d t, y+v_{y}(x, y) d t\right)$

$$
\begin{aligned}
& Q_{2} \approx\left(x+d x+v_{x}(x+d x, y) d t, y+v_{y}(x+d x, y) d t\right), \\
& Q_{3} \approx\left(x+v_{x}(x, y+d y) d t, y+d y+v_{y}(x y+d y) d t\right), \\
& Q_{4} \approx\left(x+d x+v_{x}(x+d x, y+d y) d t, y+d y+v_{y}(x+d x, y+d y) d t\right) .
\end{aligned}
$$

Volume at time $t: d x d y$
Volume at time $t+d t:=\left(Q_{2} \overrightarrow{Q_{1}}\right)_{x}\left(Q_{3} \overrightarrow{Q_{1}}\right)_{y}=$

$$
\begin{aligned}
& =\left(d x+\left(v_{x}(x+d x, y)-v_{x}(x, y)\right) d t\right) \cdot\left(d y+\left(v_{y}(x, y+d y)-v_{y}(x, y)\right) d t\right. \\
& \approx d x d y+\left[\frac{\partial v_{x}}{\partial x}+\frac{\partial v_{y}}{\partial y}\right] d t+O\left(d t^{2}\right)
\end{aligned}
$$

Conservation of volume yields

$$
\frac{\partial v_{x}}{\partial x}+\frac{\partial v_{y}}{\partial y}=0
$$

B.2.

Balance of momentum.


At time $t$, let $\underline{k}$ denote the total force on a material unit volume and let $\underline{p}$ denote the total momentum of the volume. Then we have

$$
\underline{k}=\dot{p} \quad \text { (Newton's law })
$$

which can be worked out to the following form

$$
\operatorname{div} T+\underline{b}=\rho \underline{\dot{v}}
$$

the equation of motion.

## One dimensional explanation.

$$
-60-
$$

Consider a bar on which on both end points a force $F$ is imposed

$\sigma_{1}$ : force per unit surface (= stress) which is exerted on part II by part I
$\sigma_{2}$ : force per unit surface which is exerted on part I by part II.
So balance of forces yields

$$
\sigma_{1} S=P, \quad \sigma_{2} S=P
$$

and so

$$
\sigma_{1}=\sigma_{2}=\frac{P}{S}
$$

$\sigma$ is called the normal stress:

(action = reaction!)

Consider a one-dimensional "unit-volume"

$\longrightarrow e_{x}$
$\sigma:$ stress, $b:$ volume force (e.g. mass), $a=\dot{v}$.
Total force in the $x$-direction on a unit element $d x$

$$
K=-\sigma(x) S+\sigma(x+d x) S+b S d x .
$$

According to Newton's law ( $K=m a$ ) we have

$$
(\sigma(x+d x)-\sigma(x)) S+b S d x=m a=\rho S d x \cdot a
$$

from which we derive the equation of motion

$$
\frac{d \sigma}{d x}+b=\rho a .
$$

Suppose

- volume force is zero (or negligible) : $b=0$
- motion is stationary (or quasi-stationary) : $a=0$

Then we have

$$
\sigma^{\prime}(x)=\frac{d \sigma}{d x}=0 .
$$

Two dimensional analogue of the previous equation ( $\underline{b}=\underline{a}=\underline{0}$ )


Stresses are forces per unit-surface, e.g. total force on the $B C$-surface

$$
\left(t_{x x} \underline{e}_{x}+t_{y x} \underline{e}_{y}\right) d y d z
$$

From the balance of moments we obtain the symmetry relation

$$
t_{x y}=t_{y x} .
$$

Balancing the forces in the $x$-direction yields

$$
\begin{aligned}
& t_{x x}(x+d x, y) d y d z-t_{x x}(x, y) d y d z+ \\
& \quad+t_{x y}(x, y+d y) d x d z-t_{x y}(x, y) d x d z= \\
& \quad=\left[\frac{\partial t_{x x}}{\partial x}+\frac{\partial t_{x y}}{\partial y}\right] d x d y d z=0
\end{aligned}
$$

and similarly in the $y$-direction

$$
\left[\frac{\partial t_{y x}}{\partial x}+\frac{\partial t_{y y}}{\partial y}\right] d x d y d z=0
$$

Thus we obtain the balance equations

$$
\left[\begin{array}{l}
\frac{\partial t_{x x}}{\partial x}+\frac{\partial t_{x y}}{\partial y}=0 \\
\frac{\partial t_{y x}}{\partial x}+\frac{\partial t_{y y}}{\partial y}=0
\end{array}\right.
$$

$t_{x x}, t_{y y}$ are called normal stresses $t_{x y}=t_{y z}$ are called shear stresses.

## Recapitulation.

For a two-dimensional flow of an incompressible fluid with a homogeneous mass density without volume forces or acceleration there are the unknowns

$$
v_{x}, v_{y}, t_{x x}, t_{x y}, t_{y y} \text { and } t_{z z}
$$

satisfying the equation

$$
\begin{aligned}
& \frac{\partial v_{x}}{\partial x}+\frac{\partial v_{y}}{\partial y}=0 \\
& \frac{\partial t_{x x}}{\partial x}+\frac{\partial t_{x y}}{\partial y}=0 \\
& \frac{\partial t_{x y}}{\partial x}+\frac{\partial t_{y y}}{\partial y}=0 .
\end{aligned}
$$

In order to determine the unknowns we need three more equations.

## B.3. Constitutive (= material) equations.

These equations characterise the specific medium. They are often determined on experimental/empirical grounds.
The starting point for a Newtonian fuid is that the stresses are proportional to the velocity gradients:

$$
t_{i j}-\frac{\partial v_{i}}{\partial x_{j}}+\frac{\partial v_{j}}{\partial x_{i}} .
$$

(We have to take the symmetric part of the velocity gradient since $t_{i j}=t_{j i}$.)
A fluid can be taken to be incompressible. A compressible medium under hydrostatic pressure ( $t_{x x}=t_{y y}=t_{z z}=-p$ ) gets a smaller volume ( $p$ increases). For a compressible medium the hydrostatic pressure is related to the normal stresses as follows

$$
\begin{equation*}
p=-\frac{1}{3}\left(t_{x x}+t_{y y}+t_{z z}\right) \tag{}
\end{equation*}
$$

However for an incompressible we do not have this constitutive relation: hydrostatic pressure is an additional essential unknown.
Therefore the stresses are replaced by the so called deviatoric stresses $\tau_{i j}$ :

$$
\begin{aligned}
& t_{x x}=-p+\tau_{x x} \\
& t_{y y}=-p+\tau_{y y} \\
& t_{z z}=-p+\tau_{z z} \text { and } t_{x y}=\tau_{x y}
\end{aligned}
$$

Since $\tau_{z z}=-\left(\tau_{x x}+\tau_{y y}\right)$, we deal with the unknowns

$$
v_{x}, v_{y}, \tau_{x x}, \tau_{y y}, \tau_{x y} \text { and } p
$$

We need three constitutive equations. For a Newtonian fluid they have the following simple form

$$
\begin{aligned}
& \tau_{x x}=2 \eta \frac{\partial v_{x}}{\partial x}, \quad \tau_{y y}=2 \eta \frac{\partial v_{y}}{\partial y} \\
& \tau_{x y}=\eta\left[\frac{\partial v_{x}}{\partial y}+\frac{\partial v_{y}}{\partial x}\right]
\end{aligned}
$$

where $\eta$ is the viscosity (= material constant).

## Recapitulation.

For the unknowns $v_{x}, v_{y}, \tau_{x x}, \tau_{x y}, \tau_{y y}$ and $p$ we have found the following equations

$$
\begin{aligned}
& \frac{\partial v_{x}}{\partial x}+\frac{\partial v_{y}}{\partial y}=0 \\
& -\frac{\partial p}{\partial x}+\frac{\partial \tau_{x x}}{\partial x}+\frac{\partial \tau_{x y}}{\partial y}=0 \\
& \frac{\partial \tau_{x y}}{\partial x}-\frac{\partial p}{\partial y}+\frac{\partial \tau_{y y}}{\partial y}=0 \\
& \tau_{x x}=2 \eta \frac{\partial v_{x}}{\partial x} \\
& \tau_{y y}=2 \eta \frac{\partial v_{y}}{\partial y} \\
& \tau_{x y}=\eta\left[\frac{\partial v_{x}}{\partial y}+\frac{\partial v_{y}}{\partial x}\right]
\end{aligned}
$$

B. 4 .

Flow between two plates.


Given two plates at $y=0$ and $y=d$, the plate at $y=0$ remains at rest while the other plate at $y=d$ move in the $x$-direction with a constant velocity $V$. Between these plates there is a flow of a Newtonian fluid with viscosity $\eta$ and with a homogeneous mass density $\rho_{0}$.

## Assumptions.

- the flow is stationary
- the flow is laminar
- the volume forces are negligible
- no $z$-dependence
- the pressure gradient has no $x$-component

Then the first two assumptions yield $\underline{v}(x, y)=v(x, y) \underline{e}_{x}$, i.e. $v_{x}=v(x, y)$ and $v_{y}=0$. So from the incompressibility condition we derive

$$
\frac{\partial v_{x}}{\partial x}=0 \Rightarrow v_{x}=v(y)
$$

The constitutive equations are given by

$$
\begin{aligned}
& \tau_{x x}=\tau_{y y}=0 \quad\left[\frac{\partial v_{x}}{\partial x}=\frac{\partial v_{y}}{\partial y}=0\right] \\
& \tau_{x y}=\eta v^{\prime}(y) .
\end{aligned}
$$

The equations of motion are given by

$$
\begin{aligned}
& \eta v^{\prime \prime}(y)=0 \quad\left[\frac{\partial p}{\partial x}=0\right] \\
& \frac{d p}{d y}=0 \Rightarrow p=p_{0} \text { (constant). }
\end{aligned}
$$

We obtain

$$
v(y)=C_{1} y+C_{2} .
$$

Introducing the no slip conditions

$$
v_{x}(0)=0 \quad \text { and } \quad v_{x}(d)=V
$$

it follows that $C_{1}=\frac{V}{d}$ and $C_{2}=0$. Hence

$$
v(y)=\frac{V}{d} y
$$

For the corresponding stresses we find

$$
\tau_{x y}=t_{x y}=\eta \frac{V}{d}, \quad \tau_{x x}=\tau_{y y}=0 \Rightarrow t_{x x}=t_{y y}=-p_{0} .
$$

Remark. This kind of simple results do not hold any longer for more complex (non-linear) constitutive equations.

## B.5. Poiseuille flow, Darcy's law.

Darcy's law describes the flow of a fluid through a porous medium. Although the law possesses a more general validity, we restrict here to a rigid porous medium. Clearly, if the fluid flows through the pores of the medium it encounters resistance which is due to the viscosity of the fluid and, more importantly, due to the surface tensions in the fluid. For our very heuristic derivation of Darcy's law we consider first the Poiseuille flow through a circular pipe.


$$
\begin{aligned}
& \text { As a solution we obtain } \\
& \quad \underline{v}=v(r) \underline{e}_{2} \\
& \text { met } \\
& \quad v(r)=A\left(R^{2}-r^{2}\right) .
\end{aligned}
$$

Let $\bar{v}$ denote the mean value of the velocity $v(r)$, then we have

$$
A=\frac{2 \bar{v}}{R^{2}} \Rightarrow v(r)=\frac{2 \bar{v}}{R^{2}}\left(R^{2}-r^{2}\right)
$$

So for the stresses we obtain

$$
\frac{\partial v_{2}}{\partial r}=v^{\prime}(r)=-\frac{4 \bar{v}}{R^{2}} r \Rightarrow \tau_{r 2}=\tau_{2 r}=-\frac{4 \eta}{R^{2}} \bar{v} r .
$$

The equations of motion are in cylindrical coordinates

$$
-\frac{\partial p}{\partial z}+\frac{1}{r} \frac{d}{d r} r t_{r z}+b_{z}=\rho a_{z}=0
$$

$$
\frac{\partial p}{\partial r}=0 \Rightarrow p=p(z)
$$

Hence

$$
\frac{-d p}{d z}+b_{2}=-\frac{1}{r} \frac{d}{d r}\left(r \tau_{2 r}\right)=\frac{8 \eta}{R^{2}} \bar{v} .
$$

Consider a verticle tube. dV


Then

$$
b_{z}=-p g
$$

and herewith

$$
-\frac{d p}{d z}-\rho g=\frac{8 \eta}{R^{2}} \bar{v}
$$

Define $\frac{8 \eta}{R^{2}}=\frac{1}{S}=\frac{\eta}{k}$
with $S$ the permeability ( $S$ increases if $\eta$ decreases or $R$ increases) and with $k$ the porosity of the medium.

Then we obtain

$$
\bar{v}=-\frac{k}{\eta} \frac{d}{d z}(p(z)+\rho g z) .
$$

This equation can be generalized to Darcy's law

$$
\underline{v}=-\frac{k}{\eta} \nabla(p+\rho g z) .
$$

## Appendix C

## Sobolev spaces

Consider the closed interval $[0,1]$.
By $C([0,1])$ we mean the vector space of all continuous functions on the interval $[0,1]$. Defining for every $f \in C([0,1])$

$$
\|f\|_{\infty, 0}=\max _{t \in[0,1]}|f(t)|=\sup _{t \in\langle 0,1)}|f(t)|
$$

$C([0,1])$ becomes a Banach space, i.e. a complete normed space. By $C^{0}([0,1])$ we denote the vector space of all restrictions to the open interval $(0,1)$ of functions in $C([0,1])$. Put differently, $C^{0}([0,1])$ consists of all uniformly continuous functions on $(0,1)$.
For every $k \in \mathbb{N}$ define $C^{k}([0,1])$ as the vector space of all $k$ times continuously differentiable functions on ( 0,1 ) for which $f^{(k)} \in C^{0}([0,1])$. For $f \in C^{k}([0,1])$ we have $f^{(j)} \in C^{0}([0,1])$ $j=0,1, \ldots, k$ because

$$
f(x)=f(0)+x f^{\prime}(0)+\cdots+\frac{x^{k-1}}{(k-1)!} f^{(k-1)}(0)+\frac{1}{(k-1)!} \int_{0}^{x}(x-t)^{k-1} f^{(k)}(t) d t
$$

So a suitable norm on $C^{k}([0,1])$ is given by

$$
\|f\|_{\infty, k}=\sum_{j=0}^{k} \sup _{t \in(0,1)}\left|f^{(j)}(t)\right|
$$

By introducing other norms on $C^{0}([0,1])$ we arrive at completions of $C^{0}([0,1])$, e.g. Banach spaces of (equivalence classes of) measurable functions. Classical are the $L_{p}$-spaces where the norm is given by

$$
\|f\|_{p, 0}=\left(\int_{0}^{1}|f(t)|^{p} d t\right)^{1 / p}, 1 \leq p<\infty .
$$

The corresponding completions are denoted by $L_{p}([0,1])$. They do not consists of functions, although we often treat them as if they do.

Correspondingly we introduce on $C^{k}([0,1])$ the norms

$$
\|f\|_{p, k}=\sum_{j=0}^{k}\left(\int_{0}^{1}\left|f^{(j)}(t)\right|^{p} d t\right)^{1 / p} \quad, \quad 1 \leq p<\infty .
$$

The corresponding completions are the Sobolev spaces $W^{p, i}([0,1])$. We have the Sobolevimbedding

$$
W^{p, k}([0,1]) \hookrightarrow C^{k-1}([0,1])
$$

i.e.

$$
W^{p, k}\left([0,1] \subset C^{k-1}([0,1])\right.
$$

ànd

$$
\exists_{C>0} \forall_{f \in W^{r^{2}}([0,1])}:\|f\|_{\infty, k-1} \leq C\|f\|_{p, k} .
$$

The Banach space $W^{p, k}([0,1])$ admits the following characterisation:
$W^{p, k}([0,1])$ is the Banach of all functions
$f \in C^{k-1}([0,1])$ for which the $(k-1)$-th derivative $f^{(k-1)}$
is absolutely continuous and has a generalized derivative
As an explartath belone dedqe (Staldment we consider the case $k=1$, the general case can then the obtained by induction.

## Definition.

A function $f \in C^{0}([0,1])$ is said to be absolutely continuous if there is a (Lebesgue) integrable function $g$ on $(0,1)$ such that

$$
f(x)=f(0)+\int_{0}^{x} g(t) d t
$$

$g$ is called the generalized derivative of $f$ and is determined up to a function which is zero almost everywhere.

## Lemma.

Let $f \in C^{1}([0,1])$. Then for every $x \in(0,1)$

$$
|f(x)| \leq 2\|f\|_{p, 1} .
$$

## Proof.

Since $f(x)=f(0)+\int_{0}^{x} f(t) d t$, we get the estimate

$$
\begin{aligned}
|f(0)| & \leq \int_{0}^{1}|f(x)| d x+\int_{0}^{1}\left(\int_{0}^{x}\left|f^{\prime}(t)\right| d t\right) d x \\
& { }^{\text {H}} \underline{C l d e r}^{1}\left(\int_{0}^{1}|f(x)|^{p} d x\right)^{1 / p}+\left(\int_{0}^{1}\left|f^{\prime}(t)\right|^{p} d t\right)^{1}=\|f\|_{p, 1} .
\end{aligned}
$$

So

$$
|f(x)| \leq|f(0)|+\int_{0}^{x}\left|f^{\prime}(t)\right| d t \leq
$$

$$
\leq\|f\|_{p, 1}+\|f\|_{p, 1}=2\|f\|_{p, 1} .
$$

## Corollary.

For any $f \in C^{1}([0,1])$ we have $\|f\|_{\infty, 0} \leq 2\|f\|_{p, 1}$.

Now let $\left(f_{n}\right)$ be a Cauchy sequence in $C^{1}([0,1])$ with respect to the norm $\|\cdot\|_{p, 1}$. Then $\left(f_{n}\right)$ is a Cauchy sequence in $C^{0}([0,1])$ because

$$
\left\|f_{n}-f_{m}\right\|_{\infty, 0} \leq 2\left\|f_{n}-f_{m}\right\|_{p, 1} .
$$

So there is $f \in C^{0}([0,1])$ such that

$$
\left\|f_{n}-f\right\|_{\infty, 0} \rightarrow 0 \quad(n=>\infty) .
$$

Since

$$
\left\|f_{n}^{\prime}-f_{m}^{\prime}\right\|_{p, 0} \leq\left\|f_{n}-f_{m}\right\|_{p, 1}
$$

the sequence $\left(f_{n}{ }^{\prime}\right)$ is Cauchy in $L_{p}([0,1])$ and hence converges to some $g \in L_{p}([0,1])$.
We have

$$
f_{n}(x)=f_{n}(0)+\int_{0}^{x} f_{n}^{\prime}(t) d t
$$

and so in the limit $n \rightarrow \infty$

$$
f(x)=f(0)+\int_{0}^{x} g(t) d t .
$$

It follows that $f$ is absolutely continuous with generalized derivative $g \in L_{p}([0,1])$. Moreover

$$
\left\|f_{n}-f\right\|_{p, 1}=\left\|f_{n}-f\right\|_{p, 0}+\left\|f_{n}^{\prime}-g\right\|_{p, 0} .
$$

The Sobolev space $W^{p, k}([0,1])$ can be introduced in a different way.
To this end we introduce the space $C_{c}^{\infty}((0,1))$ consisting of all infinitely differentiable functions on $(0,1)$ with a compact support within $(0,1) . C_{c}^{\infty}((0,1))$ is a vector space with a topology defined such that

$$
\phi_{n} \rightarrow \phi \text { in } C_{c}^{\infty}((0,1))
$$

means that there exists a compact set $K \subset(0,1)$ such that

$$
\forall_{n}: \operatorname{supp}\left(\phi_{n}\right) \subset K
$$

ànd

$$
\forall_{j \in \mathbb{N}}: \sup _{t \in K}\left|\phi_{n}^{(j)}(t)-\phi^{(j)}(t)\right| \rightarrow 0 \quad(n \rightarrow \infty)
$$

In the sequel we write $D((0,1))$ instead of $C_{c}^{\infty}((0,1))$.
By $D^{*}((0,1))$ we denote the dual space of $D((0,1))$, i.e. the vector space of all continuous linear functionals on $D((0,1))$. The elements of $D^{*}((0,1))$ are called distributions or generalized functions.
We introduce the concept of distributional derivative.
For each $j \in \mathbb{N}$ the linear mapping

$$
D^{j}: \phi \mapsto \phi^{(j)}
$$

is continuous on $D((0,1))$. Consequently for every continuous linear functional $L$ on $D((0,1))$, $L \circ D^{j}$ is a continuous linear functional. Now we define

$$
\hat{D}^{j} L=(-1)^{j} L \circ D^{j}
$$

Let $f$ be an integrable function on $[0,1]$. Then $f$ determines a continuous linear functional $\hat{f}$ on $D((0,1))$ through

$$
\hat{f}(\phi)=\int_{0}^{1} f(t) \phi(t) d t
$$

and so

$$
\left(\hat{D}^{j} \hat{f}\right)(\phi)=(-1)^{j} \int_{0}^{1} f(t) \phi^{(j)}(t) d t
$$

In particular, for $f \in C^{j}([0,1])$

$$
\begin{aligned}
\left(\hat{D}^{j} \hat{f}\right)(\phi) & =(-1)^{j} \int_{0}^{1} f(t) \phi^{(j)}(t) d t \\
& =\int_{0}^{1} f^{(j)}(t) \phi(t) d t
\end{aligned}
$$

and therefore

$$
\hat{D}^{j} \hat{f}=\left(D^{j} f\right)^{n}
$$

We conclude that the mapping $\hat{D}^{j}$ yields a generalization of the classical differential operator $D^{j}$. In this new terminology we have the following characterization result

$$
\begin{aligned}
f \in W^{p, k}([0,1]) \Rightarrow & f \in L_{p}([0,1]) \\
& \text { and } \\
& \exists_{g_{1}, \ldots, g_{k} \in L_{r}([0,1])}: \hat{D}^{j} \hat{f}=\hat{g}_{j}
\end{aligned}
$$

or put differently

$$
(-1)^{j} \int_{0}^{1} f(t) \phi^{(j)}(t) d t=\int_{0}^{t} g_{j}(t) \phi(t) d t
$$

The Sobolev spaces $W^{2, k}((0,1))$ are Hilbert spaces with inner products

$$
(f, g)_{k}=\sum_{j=0}^{k} \int_{0}^{1} f^{(j)}(t) \overline{g^{(j)}(t)} d t
$$

where $f^{(k)}$ denote the $k$-th generalized derivative of $f$.

Consider the space $C_{\text {per }}^{1}([0,1])$ consisting of all $f \in C^{1}([0,1])$ with $f(0)=f(1)$ and $f^{\prime}(0)=f^{\prime}(1)$. Each $f \in C_{p e r}^{1}([0,1])$ has an absolutely convergent Fourier series

$$
f(x)=\sum_{n \in \mathbb{Z}} a_{n} e^{2 \pi i n x} \quad, \quad x \in \mathbb{R}
$$

with

$$
a_{n}=\int_{0}^{1} f(t) e^{-2 \pi i n t} d t \quad, \quad n \in \mathbb{Z}
$$

Put

$$
e_{n}(x)=\left(1+4 \pi^{2} n^{2}\right)^{-\frac{1}{2}} e^{2 \pi i n x} \quad, \quad x \in \mathbb{R}, n \in \mathbb{Z}
$$

Then $\left\{e_{n}\right\}_{n \in \mathbb{Z}}$ is an orthonormal set in $W^{2,1}([0,1])$. Consequently we have

## Corollary.

Let $W_{p e r}^{2,1}([0,1])$ denote the closed subspace of $W^{2,1}([0,1])$ in which $\left\{e_{n}\right\}_{n \in \mathbb{Z}}$ is an orthonormal basis. Then $f \in W_{p e r}^{2,1}([0,1])$ if and only if $f \in W^{2,1}([0,1])$ and $f(0)=f(1)$.

Let $l$ denote the continuous linear functional

$$
l(f)=f(1)-f(0), f \in W^{2,1}([0,1])
$$

Then we have

$$
W_{p e r}^{2,1}([0,1])=\operatorname{ker}(l)
$$

Since $W^{2,1}([0,1])$ is a Hilbert space, there exists $f_{0} \in W^{2,1}([0,1])$ such that

$$
l(f)=\left(f, f_{0}\right)_{2,1}
$$

It follows that

$$
W^{2,1}([0,1])=W_{p e r}^{2,1}([0,1]) \oplus\left\langle f_{0}\right\rangle
$$

## Lemma.

$$
f_{0}(x)=\frac{\sinh \left(x-\frac{1}{2}\right)}{\sinh (1 / 2)} .
$$

Proof.
By $F_{0}$ we denote the primitive of $f_{0}$ satisfying $F_{0}(0)=F_{0}(1)=0$. Then

$$
f(1)-f(0)=\int_{0}^{1} f^{\prime}(x) d x
$$

and

$$
\begin{aligned}
f(1)-f(0) & =\int_{0}^{1} f_{0}(x) f(x)+f_{0}^{\prime \prime}(x) f^{\prime}(x) d x \\
& =\int_{0}^{1}\left[-F_{0}(x)+F_{0}^{\prime \prime}(x)\right] f^{\prime}(x) d x
\end{aligned}
$$

so that

$$
F_{0}{ }^{\prime \prime}-F_{0}=-1 \quad, \quad F_{0}(0)=F_{0}(1)=0 .
$$

Hence

$$
F_{0}(x)=\frac{\cosh 1-1}{\sinh 1} \sinh x-\cosh x+1
$$

and

$$
F_{0}{ }^{\prime}(x)=f_{0}(x)=\frac{\sinh \left(x-\frac{1}{2}\right)}{\sinh \frac{1}{2}} .
$$

Now we turn to the more dimensional case.

Let $\Omega$ be a bounded region in $\mathbb{R}^{n}$.

## Definition.

$C^{0}(\bar{\Omega})$ is the vector space consisting of all restrictions to $\Omega$ of continuous functions on $\bar{\Omega}$. Put differently, $C^{0}(\bar{\Omega})$ consists of all (bounded) uniformly continuous functions on $\Omega$. The norm in $C^{0}(\bar{\Omega})$ is defined by

$$
\|f\|_{\infty, 0}=\sup _{x \in \Omega}|f(x)|,
$$

$C^{0}(\bar{\Omega})$ is a Banach space.
Let $D_{j}=\frac{\partial}{\partial x_{j}}$. Then for every multi-index $\alpha \in \mathbb{N}_{0}^{n}$,

$$
D^{\alpha}=D_{1}^{\alpha_{1}} \cdots D_{n}^{\alpha_{\alpha}}
$$

Further we write $|\alpha|=\alpha_{1}+\cdots+\alpha_{n}$.

## Definition.

$C^{k}(\bar{\Omega})$ is the vector space of all functions $f$ for which for all $\alpha$ with $|\alpha| \leq k$ the derivative $D^{\alpha} f$ exists with

$$
D^{\alpha} f \in C^{0}(\bar{\Omega})
$$

The norm

$$
\|f\|_{\infty, k}=\sum_{|\alpha| \leq k}\left\|D^{\alpha} f\right\|_{\infty, 0}
$$

turns $C^{k}(\bar{\Omega})$ into a Banach space.
Besides we introduce

$$
C^{\infty}(\bar{\Omega})=\bigcap_{k=1}^{\infty} C^{k}(\bar{\Omega})
$$

and

$$
D(\Omega)=C_{c}^{\infty}(\Omega)=\left\{\phi \in C^{\infty}(\bar{\Omega}) \mid \operatorname{supp} \phi \subset \Omega \text { compact }\right\}
$$

With the aid of the space $D(\Omega)$ we can introduce weak (generalized) derivatives for nondifferentiable functions. First we introduce the topology of $D(\Omega)$.

- A sequence $\left(\phi_{m}\right)$ in $D(\Omega)$ is said to be convergent to $\phi \in D(\Omega)$ if

$$
\begin{aligned}
& \underset{K \subset \Omega}{\exists K \text { compact }} \forall_{n \in \mathbb{N}}: \operatorname{supp} \phi_{m} \subset K \text { and } \\
& \qquad \forall_{\alpha \in \mathbb{N}_{0}^{*}}: \sup _{x \in K}\left|\left(D^{\alpha} \phi_{m}\right)(x)-\left(D^{\alpha} \phi\right)(x)\right| \rightarrow 0 .
\end{aligned}
$$

Every continuous linear functional on $D(\Omega)$ is called a distribution. For instance, the functional

$$
l_{\alpha, x_{0}}: \phi \mapsto\left(D^{\alpha} \phi\right)\left(x_{0}\right)
$$

with $\alpha \in \mathbb{N}_{0}^{n}$ and $x_{0} \in \Omega$ is continuous on $D(\Omega)$. This functional is heuristically written as

$$
l_{\alpha, x_{0}}(\phi)=(-1)^{|\alpha|} \int_{\Omega} \phi(x) \delta^{(\alpha)}\left(x-x_{0}\right) d x
$$

and $\delta\left(x-x_{0}\right)$ is called a delta function.
Also, let $p \geq 1$ and let $f \in L_{p}(\Omega)$. Then the linear functional $\hat{f}$ is defined by

$$
\hat{f}(\phi)=\int_{\Omega} \phi(x) f(x) d x, \quad \phi \in D(\Omega) .
$$

The functional $\hat{f}$ is continuous. Thus for any $p \geq 1, L_{p}(\Omega)$ is imbedded into the dual space $D^{*}(\Omega)$. For $L \in D^{*}(\Omega)$ we write $L \in L_{p}(\Omega)$ meaning that $L=\hat{f}$ for some $f \in L_{p}(\Omega)$.

For every distribution $L$ in $D^{*}(\Omega)$ its distributional derivative $\hat{D}^{\alpha} L$ is defined by

$$
\left(\hat{D}^{\alpha} L\right)(\phi)=(-1)^{|\alpha|} L\left(D^{\alpha} \phi\right)
$$

In particular for $f \in L_{p}(\Omega)$

$$
\left(\hat{D}^{\alpha} \hat{f}\right)(\phi)=(-1)^{|\alpha|} \int_{\Omega}\left(D^{\alpha} \phi\right)(x) f(x) d x .
$$

The distribution $\hat{D}^{\alpha} \hat{f}$ is called the weak derivative of order $\alpha$ of $f$. It is a natural extension of the classical notion of derivative:

Take $\psi \in C^{\infty}(\bar{\Omega})$. Then

$$
\left(\hat{D}^{\alpha} \hat{\psi}\right)(\phi)=\int_{\Omega} \phi(x)\left(D^{\alpha} \psi\right)(x) d x
$$

Often one writes $D^{\alpha} f$ in stead of $\hat{D}^{\alpha} \hat{f}$ although $f$ is not differentiable.
We come to the definition of the Sobolev spaces $W^{p, k}(\Omega)$.

## Definition.

$$
\begin{aligned}
f \in W^{p, k}(\Omega): \Leftrightarrow & f \in L_{p}(\Omega) \text { and } \\
& \forall_{\alpha \in \mathbb{N}_{0}^{k},|\alpha| \leq k}: D^{\alpha} f \in L_{p}(\Omega)
\end{aligned}
$$

$\uparrow$ weak interpretation.
The norm in $W^{p, k}(\Omega)$ is given by

$$
\|f\|_{p, k}=\left(\sum_{|\alpha| \leq k}\left\|D^{\alpha} f\right\|_{p, 0}^{p}\right)^{1 / p} .
$$

$W^{p, k}(\Omega)$ is a Banach space, in particular $W^{2, k}(\Omega)$ is a Hilbert space with inner product

$$
(f, g)_{2, k}=\sum_{|\alpha| \leq k}\left(D^{\alpha} f, D^{\alpha} g\right)_{2}
$$

For $k=1$ we have

$$
(f, g)_{2,1}=\int_{\Omega}[f(x) g(x)+(\nabla f \cdot \nabla g)(x)] d x
$$

Remark. It can be proved (Serrin and Morrey, 1964) that $W^{p, k}(\Omega)$ is the completion of $C^{k}(\bar{\Omega})$ with respect to the norm $\|\cdot\|_{p, k}$.

Finally some remarks on the so called Sobolev embedding theorems.
For the proof of these theorems conditions are needed on the shape of the region $\Omega$. Most classical configurations such as balls, cones, cylinders and ellipsoides satisfy these conditions. We have

$$
W^{k, p}(\Omega) \hookrightarrow C^{l}(\bar{\Omega}) \text { with } l<k-\frac{n}{p}
$$

i.e.

$$
W^{k, p}(\Omega) \subset C^{l}(\bar{\Omega})
$$

and

$$
\exists C_{k_{k}>0}\|f\|_{\infty, l} \leq C_{k, l}\|f\|_{p, k} .
$$

## Appendix D

## The notion: weak solution

Let $K$ be a linear second order differential operator given by

$$
K u=\sum_{i, j=1}^{n} a_{i j}(x) D_{i} D_{j} u+\sum_{i=1}^{n} b_{i}(x) D_{i} u+c(x) u
$$

where $D_{i}=\frac{\partial}{\partial x_{i}}$ and where $u$ belongs to $C^{2}\left(\mathbb{R}^{n}\right)$.
Assuming that the coefficients $a_{i j}(x), b_{i}(x)$ and $c(x)$ belong to $C^{\infty}\left(\mathbb{R}^{n}\right), K$ is a continuous linear mapping from $D\left(\mathbb{R}^{n}\right)=C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ into itself. We can extend $K$ to the distribution space $D^{*}\left(\mathbb{R}^{n}\right)$ in the following natural way

For every $L \in D^{*}\left(\mathbb{R}^{n}\right)$ define $\hat{K} L$ by

$$
\hat{K} L=\sum_{i, j=1}^{n} a_{i j}(x) \hat{D}_{i} \hat{D}_{j} L+\sum_{i=1}^{n} b_{i}(x) \hat{D}_{i} L+c(x) L
$$

i.e.

$$
(\hat{K} L)(\phi)=L\left(K^{*} \phi\right), \phi \in D\left(\mathbb{R}^{n}\right)
$$

where $K^{*}$ denotes the differential operator.
The operator $K^{*}$ is called the formal adjoint of $K$. If $K=K^{*}$ then $K$ is called (essentially) selfadjoint. E.g. the Laplacian $\Delta$ in $\mathbb{R}^{n}$ is self-adjoint.
Let $u, v \in C^{2}\left(\mathbb{R}^{n}\right)$. Then we have

$$
v(K u)-u\left(K^{*} v\right)=\operatorname{div} P=\sum_{i=1}^{n} \frac{\partial P_{i}}{\partial x_{i}}
$$

where the vector field $P$ is defined by

$$
P_{i}(x)=\sum_{j=1}^{n}\left[a_{i j}(x) v D_{j} u-u D_{j}\left(a_{i j}(x) v\right)\right]+b_{i}(x) u v
$$

So for a bounded region $\Omega \subset \mathbb{R}^{n}$ with piecewise smooth orientable boundary $\partial \Omega$ we have by Gauss' divergence theorem

$$
\int_{\Omega}\left(v K u-u K^{*} u\right) d x=\int_{\partial \Omega}(P \cdot n) d \sigma
$$

This relation is called Green's formula. If $K=\Delta$. Green's formula become the second identity of Green

$$
\int_{\Omega}(v \Delta u-u \Delta v) d x=\int_{\partial \Omega}\left(v \frac{\partial u}{\partial n}-u \frac{\partial v}{\partial n}\right) d \sigma .
$$

Consider the linear second order partial differential equation

$$
\begin{equation*}
K u \equiv \sum_{i, j=1}^{n} a_{i j}(x) D_{i} D_{j} u+\sum_{i=1}^{n} b_{i}(x) D_{i} u+c(x) u=f(x) \tag{*}
\end{equation*}
$$

Then $u$ is called a classical or strong solution of (*) in a region $\Omega \subset \mathbb{R}^{n}$ if
i) $u \in C^{2}(\Omega)$,
ii) $K u=f$ in $\Omega$.

A distribution $L \in D^{*}(\Omega)$ is called a generalized or weak solution of $\left({ }^{*}\right)$ if

$$
\hat{K} L=\hat{f}
$$

i.e.

$$
\begin{aligned}
& \text { for all } \phi \in D(\Omega): \\
& L\left(K^{*} \phi\right)=\hat{f}(\phi)=\int_{\Omega} f(x) \phi(x) d x .
\end{aligned}
$$

## Theorem.

Any strong solution of $\left(^{*}\right)$ is a weak solution.

Proof.
Let $u \in C^{2}(\Omega)$ with $K u=f$. Then for all $\phi \in D(\Omega)$

$$
\begin{aligned}
(\hat{K} \hat{u})(\phi) & =\int_{\Omega} u\left(K^{*} \phi\right) d x \\
& =\int_{\Omega}(K u) \phi d x-\int_{\partial \Omega}(P, n) d \sigma .
\end{aligned}
$$

Since $\phi \in D(\Omega)$ we have $\left.P\right|_{\partial \Omega}=0$.
Conversely,

## Theorem.

Let $u \in C^{2}(\Omega)$ be a weak solution of (*) with $f \in C^{0}(\Omega)$. Then $u$ is a strong solution.

Let for every $y \in \Omega$ the distribution $\delta_{y} \in D^{*}(\Omega)$ be defined by $\delta_{y}(\phi)=\phi(y), \phi \in D(\Omega)$. A distribution $S_{y} \in D^{*}(\Omega)$ is called a fundamental solution of the differential operator $K$ in $\Omega$ if

$$
\hat{K} S_{y}=\delta_{y}
$$

i.e.

$$
S_{y}\left(K^{*} \phi\right)=\phi(y) \quad, \phi \in D(\Omega) .
$$

In many concrete cases $S_{y}$ is represented by a smooth function with only a singularity at $x=y$. It is clear that $S_{y}$ is in general not uniquely determined. Therefore boundary conditions have to be added.
Suppose $S_{y}=\hat{\delta}_{y}$, we mean

$$
S_{y}(\phi)=\int_{\Omega} g_{y}(x) \phi(x) d x .
$$

Then under certain conditions the solution of $K u=f$ is given by

$$
u(x)=\int_{\Omega} g_{y}(x) \phi(y) d y .
$$

Example. Take $K=-\Delta$. Then a fundamental solution in $\mathbb{R}^{n}$ is given by

$$
g_{y}(x)=g(x-y)
$$

with

$$
\begin{aligned}
& g(x)=\frac{1}{2 \pi} \log \frac{1}{|x|}, n=2 \\
& g(x)=\frac{1}{(n-2) \sigma_{n}} \frac{1}{|x|^{n-2}}, n \geq 3
\end{aligned}
$$

where

$$
\sigma_{n}=\frac{n \pi^{1 / n n}}{\Gamma\left(\frac{1}{2} n+1\right)}
$$

So a solution of the Poisson equation (in case $n=3$ )

$$
\Delta u=-f(x), x \in \Omega
$$

is given by

$$
u(x)=\frac{1}{4 \pi} \int_{\Omega} \frac{f(y)}{|x-y|} d y \quad, \quad x \in \Omega .
$$

In order to satisfy the possible boundary conditions a harmonic function can be added to $u$.

## Appendix E

## Some fixed point theorems

## Banach's contraction theorem.

Let $(X, d)$ be a complete metric space and $T: X \rightarrow X$ a contractive mapping, i.e.

$$
\begin{aligned}
& \exists_{\alpha, 0<\alpha<1} \forall_{x, y \in X, x \neq y}: \\
& \quad d(T x, T y) \leq \alpha d(x, y)
\end{aligned}
$$

Then there exists $x_{0} \in X$ with $T x_{0}=x_{0}$.

## Proof.

Let $x \in X$ and define the sequence $x_{n}=T^{n} x$. If we can prove that the sequence $\left(x_{n}\right)$ is convergent in $X$ with limit $x_{0}$, then $x_{0}$ is a fixed point of $T$. Indeed,

$$
x_{0}=\lim _{n \rightarrow \infty} T^{n} x=\lim _{n \rightarrow \infty} T\left(T^{n-1} x\right)-T x_{0}
$$

where in the last step the continuity of $T$ is used.
To prove that $\left(x_{n}\right)$ is a convergence we show that it is a Caucy sequence.
Let $n>m$. Then

$$
\begin{align*}
& d\left(x_{n}, x_{m}\right)=d\left(T^{n} x, T^{m} x\right) \leq \\
& \quad \leq d\left(T^{n} x, T^{n-1} x\right)+\cdots+d\left(T^{m+1} x, T^{m} x\right) \leq \\
& \quad \leq\left(\alpha^{n-1}+\alpha^{n-2}+\cdots+\alpha^{m}\right) d(T x, x) \\
& \quad \leq \frac{\alpha^{m}}{1-\alpha} d(T x, x)
\end{align*}
$$

Hence $d\left(x_{n}, x_{m}\right) \rightarrow 0$ if $n, m \rightarrow \infty$.

## Theorem.

Let $(X, d)$ be a compact metric space and $T: X \rightarrow X$ a semi-contractive mapping, i.e.

$$
\forall_{x, y \in X, x \neq y}: d(T x, T y)<d(x, y) .
$$

Then $T$ has a fixed point.

## Proof.

The mapping $T$ is continuous, and hence also $\phi: X \rightarrow \mathbb{R}^{+}$defined by

$$
\phi(x)=d(x, T x) \quad, \quad x \in X,
$$

is continuous:

$$
|d(x, T x)-d(y, T y)| \leq d(x, y)+d(T x, T y)
$$

So the set $\phi(X)=\{d(x, T x) \mid x \in X\}$ is compact in $\mathbb{R}^{+}$and so has a minimum. Suppose

$$
d\left(x_{0}, T x_{0}\right)=\min \{d(x, T x) \mid x \in X\}
$$

If $x_{0} \neq T x_{0}$, then

$$
d\left(T x_{0}, T^{2} x_{0}\right)<d\left(x_{0}, T x_{0}\right)
$$

which gives a contradiction. Conclusion $x_{0}=T x_{0}$.

## Brouwer's fixed point theorem.

Let $B$ denote the closed unit ball in $\mathbb{R}^{n}$ and $f: B \rightarrow B$ a continuous mapping. Then $f$ has a fixed point in $B$.

This theorem has the following consequence.

Corollary.
Let $K$ be a non-empty compact convex subset of a finite dimensional normed space $X$ and let $f: K \rightarrow K$ be a continuous mapping. Then $f$ possesses a fixed point in $K$.

## Proof.

We can as well assume that $X=\mathbb{R}^{d}$. If $K=\left\{x \in \mathbb{R}^{d} \mid\|x\| \leq r\right\}$ then define $f_{r}(x)=\frac{1}{r} f\left(\frac{x}{r}\right)$ and apply Brouwer's theorem.
Otherwise, take $r>0$ so large that $K \subseteq B_{r}=\left\{x \in \mathbb{R}^{d} \mid\|x\| \leq r\right\}$.
Define $\phi: B \rightarrow K$ by

$$
\phi(x)=y \text { with }\|x-y\|=\operatorname{dist}(x, K) .
$$

(Remark: $y$ is uniquely determined.)

Then $\phi$ is continuous and $\phi(x)=x$ for all $x \in K$. So $f \circ \phi$ from $B_{r}$ into $K \subset B_{r}$ is continuous. According to Brouwer's theorem $f \circ \phi$ has a fixed point $x \in B_{r}$. Since $(f \circ \phi)(x)=x, x \in K$ so that $\phi(x)=x$ and $f(x)=x$.

## Schauder's fixed point theorem.

Let $E$ denote a closed bounded convex subset of a normed space $X$, and let $f: E \rightarrow E$ be a mapping with the property that

- $\quad f(E)$ is relatively compact.

Then $f$ has a fixed point in $E$.

Remark: If $f: X \rightarrow X$ is a compact mapping, then condition is satisfied for all $E$.

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