# A Study of Well Posedness for Systems of Coupled Non-linear Dispersive Wave Equations 

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## THESIS

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To TATA

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## LIST OF ABBREVIATIONS

| BBM | Benjamin Bona Mahony |
| :--- | :--- |
| IVP | Initial Value Problem |
| RLW | Regularized Long-Wave |
| gBBM | Generalized Benjamin Bona Mahony |
| KdV | Korteweg-de Vries |
| gKdV | Initial Condition |
| IC | Partial Differential Equations |
| PDE | Right Hand Side |

## SUMMARY

Consideration is given to coupled systems

$$
\left\{\begin{array}{l}
u_{t}+u_{x}-u_{x x t}+P(u, v)_{x}=0 \\
v_{t}+v_{x}-v_{x x t}+Q(u, v)_{x}=0
\end{array}\right.
$$

of two evolution equations of generalized BBM-type, posed for $x \in \mathbb{R}$ and $t \geq 0$ and with $u$ and $v$ real-valued functions of $(x, t)$; the subscripts connote partial differentiation and $P$ and $Q$ are homogeneous polynomials in the variables $u$ and $v$.

A study of the initial-value problems in which $u(x, t)$ and $v(x, t)$ are specified at $t=0$, viz.

$$
u(x, 0)=u_{0}(x), \quad v(x, 0)=v_{0}(x), \quad x \in \mathbb{R},
$$

is undertake to analyze and determine conditions on the case of quadratic homogeneous polynomials $P$ and $Q$ under which these initial-value problems are globally well-posed in the $L_{2}$-based Sobolev spaces $H^{s}(\mathbb{R}) \times H^{s}(\mathbb{R})$ for $s \geq 0$. The case of cubic homogeneous polynomials $P$ and $Q$ is also investigated.

## CHAPTER 1

## INTRODUCTION

The focus of this study is systems of non-linear dispersive wave equations and the fundamental properties of their solutions. These systems arise in water wave theory, climate modeling and other situations where wave propagation is important.

The particular system of partial differential equations we are investigating is a coupled system of two evolution equations of generalized BBM-type. It has the form

$$
\left\{\begin{array}{l}
u_{t}+u_{x}-u_{x x t}+P(u, v)_{x}=0,  \tag{1.1}\\
v_{t}+v_{x}-v_{x x t}+Q(u, v)_{x}=0,
\end{array}\right.
$$

posed for $x \in \mathbb{R}$ and $t \geq 0$ and with $u$ and $v$ real-valued functions of ( $x, t$ ) where subscripts connote partial differentiation. We consider the initial-value problem (IVP henceforth) where the wave profile and the velocity are specified at a starting time taken to be 0 , viz.

$$
\begin{equation*}
u_{0}(x)=u(x, 0)=\phi(x), \quad v_{0}(x)=v(x, 0)=\psi(x), \quad x \in \mathbb{R} . \tag{1.2}
\end{equation*}
$$

For the time being, $P$ and $Q$ are arbitrary homogeneous quadratic polynomials,

$$
P(u, v)=A u^{2}+B u v+C v^{2},
$$

$$
Q(u, v)=D u^{2}+E u v+F v^{2},
$$

with coeficients $A, B, C, D, E, F \in \mathbb{R}$. Later, cubic polynomials will also be considered.

The IVP (1.1) - (1.2) is a natural generalization of the IVP for the BBM-equation or regularized long-wave (RLW) equation,

$$
\begin{equation*}
\eta_{t}+\eta_{x}+\eta \eta_{x}-\eta_{x x t}=0, \quad \eta(x, 0)=f(x) \tag{1.3}
\end{equation*}
$$

The latter was originally proposed as an alternative to the Korteweg-deVries (KdV) equation

$$
\begin{equation*}
\eta_{t}+\eta_{x}+\eta \eta_{x}+\eta_{x x x}=0 \tag{1.4}
\end{equation*}
$$

by Peregrine $(1964,1967)$ and Benjamin et al. (1972). It models unidirectional propagation of small amplitude, long wavelength water waves. This equation is written in non-dimensional, laboratory coordinates, with $x$ denoting the direction of propagation and $t$ being proportional to elapsed time. As written, it serves as a model for waves propagating in the direction of increasing values of $x$. The initial condition (IC) represents a snapshot of a disturbance already generated, without necessarily inquiring as to how the motion was initiated. The solution $\eta$ of
either (1.3) or (1.4) provides a prediction of the wave profile at future time. The model subsists on the small amplitude, long wavelength assumptions,

$$
\alpha=\frac{a}{h_{0}} \ll 1 \quad \text { and } \quad \beta=\frac{h_{0}}{\lambda} \ll 1,
$$

where $a$ is a typical amplitude, $\alpha$ a typical wavelength and $h_{0}$ is the undisturbed depth of the fluid, assumed to be constant. Consequently, the initial data $f$ can be taken in the form $f(x)=\alpha F(\beta x)$, where $F$ is independent of $\alpha$ and $\beta$, thereby enforcing the small amplitude, long wavelength presumption.

The IVP (1.3) has a distinguished history on both analytical and experimental sides. The global well-posedness in $H^{s}(\mathbb{R})$ for any $s \geq 0$ for the BBM IVP was shown in [1]. Comparisons with laboratory experiments may be found in [2].

The generalized BBM equation has the form

$$
\begin{equation*}
u_{t}+u_{x}+g(u)_{x}+L u_{t}=0 \tag{1.5}
\end{equation*}
$$

where $L$ is a Fourier-multiplier operator related to the linearized dispersion relation and $g$ is a smooth, real-valued function of a real variable. Local and global well-posedness results for generalized BBM-type equations were presented in an initial form in [3]. The results of local and global well-posedness for the gBBM-equations in weaker $L_{p}$-based spaces for appropriate values of $p$ were put forward in [4].

To model two-way propagation of waves in physical systems where nonlinear and dispersive effects are equally important, systems of partial differential equations (PDEs) have been used. The theories developed are local as well as global and use various techniques.

In the case of well-posedness for a system of equations with coupled nonlinearities, Bona, Cohen and Wang [5] improved the global existence results obtained previously by Ash, Cohen and Wang [6] for a system of two KdV equations coupled with quadratic nonlinearities. They provided conditions on the coeficients of the quadratic nonlinear terms so that the IVP is globally well-posed in the $L_{2}$-based Sobolev classes $H^{s}(\mathbb{R}) \times H^{s}(\mathbb{R})$ for any $s>-3 / 4$.

The purpose of the present work is to analyze and determine conditions under which the initial-value problem (1.1) - (1.2) is globally well-posed in the $L_{2}$-based Sobolev spaces $H^{s}(\mathbb{R}) \times H^{s}(\mathbb{R})$ for $s \geq 0$. This is accomplished by first converting the IVP into a system of integral equations. This step does not depend upon the degree of the polynomials $P$ and Q. This conversion is carried out in the first part of Chapter 2, right after Notation is fixed. In the rest of Chapter 2 and through Chapter 5, the focus is on the case when $P$ and $Q$ are homogeneous quadratic polynomials. Conservation Laws are determined in the second part of Chapter 2. These will be used later to derive the a priori bounds that lead to global well-posedness. Following the idea in [3], local well-posedness can be assured by applying the contraction-mapping principle in the spaces $H^{k}(\mathbb{R}) \times H^{k}(\mathbb{R})$ for $k \geq 1$. Chapter 3 is devoted to local well-posedness, initially deduced for the $L_{\infty} \times L_{\infty}$ case and then for $H^{k}(\mathbb{R}) \times H^{k}(\mathbb{R})$ for $k \geq 1$. Chapter 4 is concerned with the $L_{2} \times L_{2}$ case, using an approach inspired from the
work of Bona \& Tzvetkov [1]. The result obtained thus far is extended to $H^{s} \times H^{s}$ for any $s \geq 0$ by applying nonlinear interpolation theory in Chapter 5. Attention is then turned to the case when the two polynomials that comprise the nonlinearity of the model are cubic. Similar to the quadratic case, the system is locally well posed in $L_{2}$-based spaces. Conditions on the coefficients of the polynomials $P$ and $Q$ are also determined so that the invariants to be used for global well-posedness exist.

## CHAPTER 2

## PRELIMINARIES

### 2.1 Notation

The notational conventions and function-space designations used in this paper are set out here. For $1 \leq p<\infty, L_{p}=L_{p}(\mathbb{R})$ connotes the $p^{t h}$-power Lebesque-integrable functions with the usual modification for the case $p=\infty$. The norm of a function $f \in L_{r}$ with $1 \leq$ $r \leq \infty$ is written $\|f\|_{r}$ while the $L_{r} \times L_{r}$-norm of a pair $(f, g)$ of such functions is written $\|(f, g)\|_{L_{r} \times L_{r}}=\|f\|_{r}+\|g\|_{r}$. In general, if X and Y are Banach spaces, then their Cartesian product $X \times Y$ is a Banach space with a product norm defined by $\|(f, g)\|_{X \times Y}=\|f\|_{X}+\|g\|_{Y}$. The Sobolev class $H^{s}=H^{s}(\mathbb{R})$ for $s \geq 0$ is the class of measurable functions $f$ whose Fourier transform $\widehat{f(\xi)}$ is a measurable function, square integrable with respect to the measure $(1+$ $|\xi|)^{2 s} d \xi$, where

$$
\widehat{f(\xi)}=\int_{-\infty}^{\infty} f(x) e^{-i x \xi} d x
$$

We will usually use simply $H^{s}$ rather than $H^{s}(\mathbb{R})$ unless emphasis on the domain of definition of the functions is needed. The norm in $H^{s} \times H^{s}$ is denoted $\|\cdot\|_{H^{s} \times H^{s}}$. Similarly, for $s \geq 0$ an integer and $p \geq 1$, the space

$$
W_{p}^{s}(\mathbb{R})=\left\{f, f^{\prime}, \ldots, f^{(s)} \in L_{p}: \int_{-\infty}^{\infty}\left(|f(x)|^{p}+\left|f^{(s)}(x)\right|^{p}\right) d x<\infty\right\}
$$

is a Banach space endowed with the norm

$$
\|f\|_{W_{p}^{s}}=\left[\int_{-\infty}^{\infty}\left(|f(x)|^{p}+\left|f^{(s)}(x)\right|^{p}\right) d x\right]^{\frac{1}{p}} .
$$

If $X$ is any Banach space and $T>0$ given, $C(0, T ; X)$ is the class of continuous maps from $[0, T]$ into $X$ with its usual norm

$$
\|u\|_{C(0, T ; X)}=\sup _{t \in[0, T]}\|u(t)\|_{X} .
$$

The subspace $C^{1}(0, T ; X)$ of the elements of $C(0, T ; X)$ for which the limit

$$
u^{\prime}(t)=\lim _{h \rightarrow 0} \frac{u(t+h)-u(t)}{h}
$$

exists in $C(0, T ; X)$, is also a Banach space with the obvious norm. For $k \in \mathbb{N}$ the spaces $C^{k}(0, T ; X)$ are defined inductively and by analogy. For convenience and when there couldn't be any confusion created,

$$
\int_{-\infty}^{\infty} f(x) \text { is replaced by } \int f .
$$

### 2.2 Integral Equation

The analysis begins with local well-posedness in a reasonably broad set of functional classes. To accomplish this, the given system is converted to an equivalent system of integral equations. This is independent of the form and degree of the polynomials $P$ and $Q$.

Let $U$ denote the vector

$$
U=\binom{u}{v}
$$

of dependent variables, $U_{0}$ the vector

$$
U_{0}=\binom{u_{0}}{v_{0}}
$$

of initial data and $M$ the vector of non-linearities

$$
M(U)=\binom{P(u, v)}{Q(u, v)} .
$$

Then the system can be written as

$$
U_{t}+U_{x}-U_{x x t}=-M_{x},
$$

with

$$
U(x, 0)=U_{0}(x) .
$$

Rearranging this equation leads to

$$
U_{t}-U_{x x t}=-U_{x}-M_{x},
$$

or

$$
\left(I-\partial_{x}^{2}\right) U_{t}=-\partial_{x}(U+M) .
$$

Inverting the operator $\left(I-\partial_{x}^{2}\right)$ subject to boundness at $\pm \infty$ leads to

$$
\begin{equation*}
U_{t}(x, t)=-\left(I-\partial_{x}^{2}\right)^{-1} \partial_{x}[U(x, t)+M(U(x, t))] . \tag{2.1}
\end{equation*}
$$

For now, assume that $U$ is bounded or at least is not growing exponentially fast as $x \rightarrow \pm \infty$.
Under that assumption, after integrating by parts, we obtain the system

$$
\begin{equation*}
U_{t}(x, t)=\int_{-\infty}^{\infty} K(x-y)[U(y, t)+M(U(y, t))] d y \tag{2.2}
\end{equation*}
$$

where the kernel $K$ is applied to the vector $U$ componentwise and

$$
K(x)=\frac{1}{2} \operatorname{sgn}(x) e^{-|x|}
$$

Integrating the relation above with respect to $t$ leads to the integral equation

$$
\begin{equation*}
U(x, t)=U_{0}(x)+\int_{0}^{t} \int_{-\infty}^{\infty} K(x-y)[U(y, s)+M(U(y, s))] d y d s \tag{2.3}
\end{equation*}
$$

If we look at this integral formula as

$$
\begin{equation*}
U=\mathcal{A} U \tag{2.4}
\end{equation*}
$$

where the operator $\mathcal{A}$ is defined by

$$
\begin{equation*}
\mathcal{A} U(x, t):=U_{0}+\int_{0}^{t} \int_{-\infty}^{\infty} K(x-y)[U(y, s)+M(U(y, s))] d y d s \tag{2.5}
\end{equation*}
$$

the problem of finding a solution to the integral equation becomes a fixed-point problem for $\mathcal{A}$. A solution can be inferred to exist by showing that $\mathcal{A}$ is a contraction and then applying the Contraction Mapping Theorem. This is the strategy for establishing local well posedness pursued in Chapter 3.

### 2.3 Hamiltonian Structure and Conservation Laws

As we will see, local well-posedness theory doesn't depend on assumptions about the coefficients of $P$ and $Q$. To move to global in time well-posedness, we need a priori estimates for the solution. To prepare for that, observe that the system has a Hamiltonian structure which sometimes turns out to be helpful in estimating the growth of the spacial norms of the solution as functions of time. Assume that we have in hand a solution $(u, v)$ of the system which is not growing exponentially for $x \rightarrow \pm \infty$. That such solutions exist, at least locally in time, is a consequence of the theory in Chapter 3.

### 2.3.1 Energy invariant

Consider the following quadratic functional:

$$
\Omega(u, v):=\int\left(a u^{2}+b u v+c v^{2}+d u_{x}^{2}+e u_{x} v_{x}+f v_{x}^{2}\right) d x
$$

where $a, b, c, d, e, f$ are real numbers to be determined. A formal calculation where integration by parts is used with the assumption that $(u, v)$ is the solution of the system and $u$ and $v$ and their derivatives do not make any contribution at $x= \pm \infty$, leads to the formulas

$$
\begin{aligned}
\frac{d}{d t} \Omega= & \int\left(2 a u u_{t}+b u_{t} v+b u v_{t}+2 c v v_{t}+2 d u_{x t} u_{x}+e u_{x t} v_{x}+e u_{x} v_{x t}+2 f v_{x} v_{x t}\right) d x \\
= & \int\left(2 a u u_{t}+b u_{t} v+b u v_{t}+2 c v v_{t}-2 d u_{x x t} u-e u_{x x t} v-e u v_{x x t}-2 f v v_{x x t}\right) d x \\
= & 2 \int u\left(a u_{t}-d u_{x x t}\right) d x+\int v\left(b u_{t}-e u_{x x t}\right) d x+ \\
& +\int u\left(b v_{t}-e u v_{x x t}\right) d x+2 \int v\left(c v_{t}-f v_{x x t}\right) d x \\
= & 2 \int(a-d) u u_{t} d x-2 d \int u P_{x} d x+\int(b-e) v u_{t} d x-e \int v P_{x} d x+ \\
& +\int(b-e) u v_{t} d x-e \int u Q_{x} d x+2 \int(c-f) v v_{t} d x-2 f \int v Q_{x} d x
\end{aligned}
$$

To simplify things a bit assume $a=d, b=e$ and $c=f$. By expanding $P_{x}$ and $Q_{x}$, the terms in the sum above may be written in the form

$$
\begin{aligned}
\frac{d}{d t} \Omega= & -2 a \int\left(2 A u^{2} u_{x}+B u v u_{x}+B u^{2} v_{x}+2 C u v v_{x}\right) d x \\
& -b \int\left(2 A u v u_{x}+B u v v_{x}+B v^{2} u_{x}+2 C v^{2} v_{x}\right) d x \\
& -b \int\left(2 D u^{2} u_{x}+E u^{2} v_{x}+E u v u_{x}+2 F u v v_{x}\right) d x \\
& -2 c \int\left(2 D u v u_{x}+E v^{2} u_{x}+E u v v_{x}+2 F v^{2} v_{x}\right) d x \\
= & \int(-4 A a-2 D b) u^{2} u_{x} d x+\int(2 C b-4 F c) v^{2} v_{x} d x \\
& -\int\left[(2 B a+E b+2 A b+4 D c) u v u_{x}+(2 B a+E b) u^{2} v_{x}\right] d x \\
& -\int\left[(B b+2 E c+4 C a+2 F b) u v v_{x}+(B b+2 E c) v^{2} u_{x}\right] d x
\end{aligned}
$$

The first two terms of the sum above are integrals which have integrands that are x-derivatives. Thus, for smooth solutions going to 0 at $\pm \infty$, they vanish without further assumptions. The last two integrals in the sum would vanish without futher assumptions on $(u, v)$ if and only if

$$
\begin{align*}
& 2 B a+E b=2 A b+4 D c \\
& B b+2 E c=4 C a+2 F b \tag{2.6}
\end{align*}
$$

Thus $a, b$ and $c$ must solve the two equations

$$
\begin{align*}
& 2 B a+(E-2 A) b-4 D c=0,  \tag{2.7}\\
& 4 C a+(2 F-B) b-2 E c=0,
\end{align*}
$$

for it to be the case that the time derivative above is zero. This pair of equations always has a non-trivial solution.

For such values of $a, b$ and $c, \Omega$ is an invariant of the temporal evolution of smooth solutions of the given system. Furthermore one can easily see by calculation that

$$
\begin{aligned}
& \frac{\partial}{\partial u}(b P+2 c Q)=2(A b+2 D c) u+(B b+2 E c) v \\
& \frac{\partial}{\partial v}(2 a P+b Q)=(2 B a+E b) u+2(2 C a+F b) v
\end{aligned}
$$

Therefore, according to (2.7),

$$
\frac{\partial}{\partial u}(b P+2 c Q)(u, v)=\frac{\partial}{\partial v}(2 a P+b Q)(u, v) .
$$

It follows that there is a cubic polynomial $R(u, v)$ such that

$$
\begin{equation*}
\frac{\partial R}{\partial u}=2 a P+b Q \quad \text { and } \quad \frac{\partial R}{\partial v}=b P+2 c Q \tag{2.8}
\end{equation*}
$$

### 2.3.2 Second Conservation Law

For the same $a, b$ and $c$, we search for a second Conservation Law. To calculate the time derivative of the functional

$$
\int_{-\infty}^{\infty}\left(a u^{2}+b u v+c v^{2}\right) d x
$$

integration by parts and the assumption that $(u, v)$ is a solution of the system which, along with their first few derivatives, do not have any contribution at $x= \pm \infty$. Using the relations

$$
\begin{gathered}
\int u_{t}\left(u-u_{x t}+P\right)=0 \\
\int u_{t}\left(v-v_{x t}+Q\right)+v_{t}\left(u-u_{x}+P\right)=0 \\
\int v_{t}\left(v-v_{x t}+Q\right)=0
\end{gathered}
$$

it follows that

$$
\begin{aligned}
& \frac{d}{d t} \int_{-\infty}^{\infty}\left(a u^{2}+b u v+c v^{2}\right) d x=\int\left(2 a u u_{t}+b u_{t} v+b u v_{t}+2 c v v_{t}\right) d x= \\
&= 2 a \int u u_{t} d x+b \int\left(u_{t} v+u v_{t}\right) d x+2 c \int v v_{t} d x \\
&= 2 a \int\left[u_{t}\left(u-u_{x t}+P\right)\right] d x+2 a \int u_{t} u_{x t} d x-2 a \int u_{t} P d x \\
&+b \int u_{t}\left(v-v_{x t}+Q\right) d x+b \int u_{t} v_{x t} d x-b \int u_{t} Q d x \\
&+b \int v_{t}\left(u-u_{x t}+P\right) d x+b \int v_{t} u_{x t} d x-b \int v_{t} P d x \\
&+2 c \int\left[v_{t}\left(v-v_{x t}+Q\right)\right] d x+2 c \int v_{t} v_{x t} d x-2 c \int v_{t} Q d x \\
&=-\int\left[(2 a P+b Q) u_{t}+(b P+2 c Q) v_{t}\right] d x+b \int \frac{d}{d x}\left(u_{t} v_{t}\right) d x \\
&=-\int\left(\frac{\partial R}{\partial u} u_{t}+\frac{\partial R}{\partial v} v_{t}\right) d x \\
&=-\frac{d}{d t} \int R d x .
\end{aligned}
$$

Consequently, we find that

$$
\frac{d}{d t} \int_{-\infty}^{\infty}\left(a u^{2}+b u v+c v^{2}+R(u, v)\right) d x=0
$$

Therefore if the functional $\Phi(u, v)$ is defined by

$$
\begin{equation*}
\Phi(u, v):=\int_{-\infty}^{\infty}\left(a u^{2}+b u v+c v^{2}+R(u, v)\right) d x, \tag{2.9}
\end{equation*}
$$

then $\Phi(u, v)$ is also an invariant of the temporal evolution of smooth solutions of our system and it serves as a Hamiltonian for the given system as long as $4 a c-b^{2}$ is not 0 .

Both these conserved quantities $\Omega$ and $\Phi$ are going to be used further for obtaining $a$ priori bounds on solutions of the system when the quadratic form $q$ defined by

$$
\begin{equation*}
q(u, v)=a u^{2}+b u v+c v^{2}, \tag{2.10}
\end{equation*}
$$

vanishes only at origin which is the case exactly when its discriminant is negative. Of course, neither $a$ nor $c$ can be 0 if the discriminant is negative. Without loss of generality, assume both $a$ and $c$ are positive.

### 2.4 Positive Definite Condition on the Coefficients

Referring back to the system (2.6), the coefficient matrix is

$$
N=\left(\begin{array}{ccc}
2 B & E-2 A & -4 D \\
4 C & 2 F-B & -2 E
\end{array}\right) .
$$

When $\operatorname{rank}(N)=2$, then $a, b, c$ are given by

$$
\left(\begin{array}{l}
a \\
b \\
c
\end{array}\right)=\left(\begin{array}{c}
E(E-2 A)+2 D(B-2 F) \\
-2 B E+8 C D \\
B(B-2 F)+2 C(E-2 A)
\end{array}\right),
$$

up to a non-zero scalar multiple.

By calculation, it easily follows that for these values, if

$$
\begin{align*}
4 a c= & B E(E-2 A)(B-2 F)+2 C E(E-2 A)^{2}+2 B D(B-2 F)^{2} \\
& +4 C D(E-2 A)(B-2 F) \\
= & (B E+4 C D)(E-2 A)(B-2 F)+2 C E(E-2 A)^{2}+2 B D(B-2 F)^{2} \\
> & (B E-4 C D)^{2}=b^{2} ; \tag{2.11}
\end{align*}
$$

then the quadratic form $q$ defined in (2.10) is positive definite.
The case $\operatorname{rank}(N)=1$ has three subcases, the complete analysis dependents on the quantity of $(E-2 A)^{2}+8 B D$ or $(2 F-B)^{2}+8 C E$. If the bigger one is strictly positive, the system (1.1) can be decoupled, the positive definite condition holds true. If the bigger quantity is equal to 0 then both are equal to zero and the method used is not applicable. A different method would show that the global well-posedness holds true. If smallest one of the quantities is negative, the result is not true anymore. The system can be written as a single complex BBM-equation instead of two equations. Numerical methods show that the system's solution blows up in finite time. The details will be pursued later as a companion result of this thesis and omitted here. However, as a remark, a few particular examples of the coefficient choices are presented.

In the case that $B=C=D=E=0$ and $A, F \neq 0$ then both quantities $(E-2 A)^{2}+8 B D$ and $(2 F-B)^{2}+8 C E$ are strictly positive, the initial system is decoupled and by choosing $a=c=1$ and $b=0$, one obtains a time invariant $\Omega$ under the flow generated by the (1.1).

In the case that $A=\frac{1}{2}, B=C=D=F=0$, and $E=1$ then both quantities $(E-2 A)^{2}+8 B D$ and $(2 F-B)^{2}+8 C E$ are equal to zero. The system looks as follows

$$
\left\{\begin{align*}
u_{t}+u_{x}-u_{x x t}+u u_{x} & =0  \tag{2.12}\\
v_{t}+v_{x}-v_{x x t}+(u v)_{x} & =0
\end{align*}\right.
$$

The local and global well-posedness results hold true.
In the case that $A=\frac{1}{2}, B=D=F=0, C=-\frac{1}{2}$ and $E=1$ then $(E-2 A)^{2}+8 B D$ and $(2 F-B)^{2}+8 C E$ are equal to zero. The system looks as follows

$$
\left\{\begin{array}{l}
u_{t}+u_{x}-u_{x x t}+u u_{x}-v v_{x}=0  \tag{2.13}\\
v_{t}+v_{x}-v_{x x t}+(u v)_{x}=0
\end{array}\right.
$$

or equivalently, as a complex-BBM equation

$$
\begin{equation*}
(u+i v)_{t}+(u+i v)_{x}-(u+i v)_{x x t}+(u+i v)(u+i v)_{x}=0 \tag{2.14}
\end{equation*}
$$

This equation is not globally well-posed, based on numerical calculations.
With $a, b$ and $c$ determined, the next step is to find the coefficients of the cubic polynomial

$$
R(u, v):=\alpha u^{3}+\beta u^{2} v+\gamma u v^{2}+\delta v^{3}
$$

appearing in the second conservation law. A calculation yields

$$
\begin{align*}
& \alpha=\frac{1}{3}(2 a A+b D), \quad \beta=b A+2 c D, \\
& \gamma=2 a C+b F, \quad \delta=\frac{1}{3}(b C+2 c F) . \tag{2.15}
\end{align*}
$$

Because of (2.6), $\beta$ and $\gamma$ may also be written as

$$
\beta=a B+\frac{1}{2} b E, \quad \gamma=\frac{1}{2} b B+c E .
$$

## CHAPTER 3

## WELL-POSEDNESS IN INTEGER ORDER SOBOLEV SPACES

### 3.1 Local Well Posedness

Local well-posedness is proved via the Contraction Mapping Principle.

### 3.1.1 Local Well-Posedness in $L_{\infty} \times L_{\infty}$

Consider the space

$$
\begin{equation*}
X_{T}:=L_{\infty}(\mathbb{R} \times[0, T]) \times L_{\infty}(\mathbb{R} \times[0, T]) \tag{3.1}
\end{equation*}
$$

with its usual norm,

$$
\begin{equation*}
\|U-\tilde{U}\|_{X_{T}}:=\sup _{x \in \mathbb{R}, t \in[0, T]}|u(x, t)-\tilde{u}(x, t)|+\sup _{x \in \mathbb{R}, t \in[0, T]}|v(x, t)-\tilde{v}(x, t)| \tag{3.2}
\end{equation*}
$$

where

$$
U=\binom{u}{v} \quad \text { and } \quad \tilde{U}=\binom{\tilde{u}}{\tilde{v}}
$$

$\left(X_{T},\|\cdot\|_{X_{T}}\right)$ is a Banach space. To establish local existence via the Contraction Mapping Principle, the domain for the operator $\mathcal{A}$ will be restricted to the closed subset

$$
B_{R}(0):=\left\{U \in X_{T}:\|U\|_{X_{T}} \leq R\right\}
$$

where $R>0$ and $T>0$ will be determined presently.
Consider $U, \tilde{U} \in B_{R}(0)$, and look for an estimate of $\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}}$. The aim is to choose specific values of $R$ and $T$ so that $\mathcal{A}$ is a contraction on $B_{R}(0)$. Using the definition of the operator $\mathcal{A}$, one calculates thusly for $U=(u, v)$ and $\tilde{U}=(\tilde{u}, \tilde{v})$

$$
\begin{align*}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}}= & \sup _{x \in \mathbb{R}}\left|\int_{0}^{T} \int_{-\infty}^{\infty} K(x-y)[u(y, t)+P(u, v)-\tilde{u}(y, t)-P(\tilde{u} . \tilde{u})] d y d s\right| \\
& +\sup _{x \in \mathbb{R}}\left|\int_{0}^{T} \int_{-\infty}^{\infty} K(x-y)[v(y, t)+Q(u, v)-\tilde{v}(y, t)-Q(\tilde{u} . \tilde{u})] d y d s\right| \\
\leq & \sup _{x \in \mathbb{R}} \int_{0}^{T} \int_{-\infty}^{\infty}|K(x-y)| \mid(u-\tilde{u})(1+A(u+\tilde{u})+B v) \\
& +(v-\tilde{v})(C(v+\tilde{v})+B \tilde{u}) \mid d y d s \\
& +\sup _{x \in \mathbb{R}} \int_{0}^{T} \int_{-\infty}^{\infty}|K(x-y)| \mid(u-\tilde{u})(D(u+\tilde{u})+E \tilde{v}) \\
& +(v-\tilde{v})(1+F(v+\tilde{v})+E u) \mid d y d s \\
\leq & \sup _{x \in \mathbb{R}} \int_{0}^{T} \int_{-\infty}^{\infty}|K(x-y)|[|u-\tilde{u}|(1+|A|(|u|+|\tilde{u}|)+|B||v|) \\
& +|v-\tilde{v}|(|C|(|v|+|\tilde{v}|)+|B||\tilde{u}|)] d y d s \\
& +\sup _{x \in \mathbb{R}} \int_{0}^{T} \int_{-\infty}^{\infty}|K(x-y)|[|u-\tilde{u}|(|D|(|u|+|\tilde{u}|)+|E \||\tilde{v}|) \\
& +|v-\tilde{v}|(1+|F|(|v|+|\tilde{v}|)+|E||u|)] d y d s . \tag{3.3}
\end{align*}
$$

Furthermore, as $U, \tilde{U} \in B_{R}(0)$, it follows that the inequality can be continued, viz.

$$
\begin{align*}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}} \leq & c_{R}^{1} \sup _{x \in \mathbb{R}} \int_{0}^{T} \int_{-\infty}^{\infty}|K(x-y)|[|u-\tilde{u}|+|v-\tilde{v}|] d y d s \\
& +c_{R}^{2} \sup _{x \in \mathbb{R}} \int_{0}^{T} \int_{-\infty}^{\infty}|K(x-y)|[|u-\tilde{u}|+|v-\tilde{v}|] d y d s \\
\leq & \|U-\tilde{U}\|_{X_{T}} C_{R} \sup _{x, t} \int_{0}^{t} \int_{-\infty}^{\infty}|K| d y d s \\
\leq & \|U-\tilde{U}\|_{X_{T}} C_{R} T \tag{3.4}
\end{align*}
$$

where $c_{R}^{1}, c_{R}^{2}$ and $C_{R}$ are constantsonly dependant on $R$. The choices of R and T are to be made such that $\mathcal{A}$ maps $B_{R}(0)$ into itself and so that $\mathcal{A}$ is a contraction on $B_{R}(0)$. In particular, one would need to have

$$
C_{R} T<1
$$

Choose T and R so that

$$
C_{R} T=\frac{1}{2}
$$

for example. For any $U \in B_{R}(0)$, we need $\|\mathcal{A} U\|_{X_{T}} \leq R$. Since

$$
\begin{align*}
\|\mathcal{A} U\|_{X_{T}} & =\|\mathcal{A} U-0\|_{X_{T}} \leq\|\mathcal{A} U-\mathcal{A} 0\|_{X_{T}}+\|\mathcal{A} 0-0\|_{X_{T}} \\
& \leq \frac{1}{2}\|U-0\|_{X_{T}}+\left\|U_{0}\right\|_{\infty} \leq \frac{1}{2} R+\left\|U_{0}\right\|_{\infty} \leq R \tag{3.5}
\end{align*}
$$

which is satisfied by choosing $R=2\left\|U_{0}\right\|_{\infty}$. Then, the choice

$$
T=\frac{1}{2 C_{R}}
$$

allows us to apply the Contraction Mapping Theorem to the integral equation (2.3). The result is existence of a unique $U \in X_{T}$ that is solution of the given system.

### 3.1.2 Local Existence in Other Spaces

Lemma 1. Suppose that $Z$ is a space embedded in $L_{\infty}$ such that, for any $u, v \in Z$, there exist global constants $C_{1}$ and $C_{2}$ satisfying

$$
\begin{equation*}
\text { (i) } \quad\|u v\|_{Z} \leq C_{1}\|u\|_{Z}\|v\|_{Z} \quad \text { and } \quad \text { (ii) } \quad\left\|\int_{-\infty}^{\infty} K(x-y) u(y) d y\right\|_{Z} \leq C_{2}\|u\|_{Z} \tag{3.6}
\end{equation*}
$$

Then there is a unique local solution for the initial-value problem for (1.1) in the space $C(0, T ; Z \times$ Z).

Proof. The space $C(0, T ; Z \times Z)$ carries the norm

$$
\begin{equation*}
\|U-\tilde{U}\|_{C(0, T ; Z \times Z)}:=\sup _{t \in[0, T]}\|u(\cdot, t)-\tilde{u}(\cdot, t)\|_{Z}+\sup _{t \in[0, T]}\|v(\cdot, t)-\tilde{v}(\cdot, t)\|_{Z} \tag{3.7}
\end{equation*}
$$

where $U=(u, v)$ and $\tilde{U}=(\tilde{u}, \tilde{v})$.
Let $Z_{T}:=C(0, T ; Z)$. Then for the easiness of notation, the $C(0, T ; Z \times Z)$-norm and the $Z \times Z$-norm of a vector $U$ will respectively be denoted by $\|U\|_{Z_{T}}$ and $\|U\|_{Z}$ henceforth.

The collection $C(0, T ; Z \times Z)$ is a Banach space of course. Restrict the operator $\mathcal{A}$ defined in (2.5) to

$$
B_{R, Z}(0):=\left\{U \in C(0, T ; Z \times Z):\|U\|_{Z_{T}} \leq R\right\} .
$$

Consider $U, \tilde{U} \in B_{R, Z}(0)$, and look for an estimate of $\|U-\tilde{U}\|_{Z_{T}}$. The aim is to choose specific values of $R$ and $T$ so that $\mathcal{A}$ is a contraction on $B_{R, Z}(0)$. Notice that

$$
\begin{aligned}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{Z_{T}}= & \sup _{t \in[0, T]}\left\|\int_{0}^{t} \int_{-\infty}^{\infty} K(x-y)[u(y, s)+P(u, v)-\tilde{u}(y, s)-P(\tilde{u} . \tilde{u})] d y d s\right\|_{Z} \\
& +\sup _{t \in[0, T]}\left\|\int_{0}^{t} \int_{-\infty}^{\infty} K(x-y)[v(y, s)+Q(u, v)-\tilde{v}(y, s)-Q(\tilde{u} . \tilde{u})] d y d s\right\|_{Z} \\
\leq & \sup _{t \in[0, T]} \int_{0}^{t}\|K *[(u-\tilde{u})(1+A(u+\tilde{u})+B v)+(v-\tilde{v})(C(v+\tilde{v})+B \tilde{u})]\|_{Z} d s \\
& +\sup _{t \in[0, T]} \int_{0}^{t}\|K *[(u-\tilde{u})(D(u+\tilde{u})+E \tilde{v})+(v-\tilde{v})(1+F(v+\tilde{v})+E u)]\|_{Z} d s \\
\leq & T\left\|\int_{-\infty}^{\infty} K(x-y)[(u-\tilde{u})(1+A(u+\tilde{u})+B v)+(v-\tilde{v})(C(v+\tilde{v})+B \tilde{u})] d y\right\|_{Z_{T}} \\
& +T\left\|\int_{-\infty}^{\infty} K(x-y)\right\|[(u-\tilde{u})(D(u+\tilde{u})+E \tilde{v})+(v-\tilde{v})(1+F(v+\tilde{v})+E u)] d y \|_{Z_{T}} .
\end{aligned}
$$

Applying (ii) from (3.6), the triangle inequality and (3.6) part (i) yields

$$
\begin{align*}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{Z_{T}} & \leq T C_{2}\left[\|(u-\tilde{u})(1+A(u+\tilde{u})+B v)+(v-\tilde{v})(C(v+\tilde{v})+B \tilde{u})\|_{Z_{T}}\right. \\
& \left.+\|(u-\tilde{u})(D(u+\tilde{u})+E \tilde{v})+(v-\tilde{v})(1+F(v+\tilde{v})+E u)\|_{Z_{T}}\right] \\
\leq & T C_{2}\left[\|u-\tilde{u}\|_{Z_{T}}+\|(u-\tilde{u})(A(u+\tilde{u})+B v)\|_{Z_{T}}+\|(v-\tilde{v})(C(v+\tilde{v})+B \tilde{u})\|_{Z_{T}}\right. \\
& \left.+\|(u-\tilde{u})(D(u+\tilde{u})+E \tilde{v})\|_{Z_{T}}+\|v-\tilde{v}\|_{Z_{T}}+\|(v-\tilde{v})(F(v+\tilde{v})+E u)\|_{Z_{T}}\right] \\
\leq & T C_{2}\left[\|u-\tilde{u}\|_{Z_{T}}+C_{1}\|u-\tilde{u}\|_{Z_{T}}\|A(u+\tilde{u})+B v\|_{Z_{T}}+C_{1}\|v-\tilde{v}\|_{Z_{T}} \| C(v+\tilde{v})+B \tilde{u}\right) \|_{Z_{T}} \\
& \left.\left.+C_{1}\|u-\tilde{u}\|_{Z_{T}}\|D(u+\tilde{u})+E \tilde{v}\|_{Z_{T}}+\|v-\tilde{v}\|_{Z_{T}}+C_{1}\|v-\tilde{v}\|_{Z_{T}} \| F(v+\tilde{v})+E u\right) \|_{Z_{T}}\right] \\
\leq & T C_{2}\left[\|u-\tilde{u}\|_{Z_{T}}+C_{1}\|u-\tilde{u}\|_{Z_{T}}\left(|A|\left(\|u\|_{Z_{T}}+|A|\|\tilde{u}\|_{Z_{T}}+|B|\|v\|_{Z_{T}}\right)\right.\right. \\
& +C_{1}\|v-\tilde{v}\|_{Z_{T}}\left(|C|\left(\|v\|_{Z_{T}}+|C|\|\tilde{v}\|_{Z_{T}}+|B|\|\tilde{u}\|_{Z_{T}}\right)\right. \\
& +C_{1}\|u-\tilde{u}\|_{Z_{T}}\left(|D|\left(\|u\|_{Z_{T}}+|D|\|\tilde{u}\|_{Z_{T}}+|E|\|\tilde{v}\|_{Z_{T}}\right)+\|v-\tilde{v}\|_{Z_{T}}\right. \\
& \left.+C_{1}\|v-\tilde{v}\|_{Z_{T}}\left(|F|\|v\|_{Z_{T}}+|F|\|\tilde{v}\|_{Z_{T}}+|E|\|u\|_{Z_{T}}\right)\right] . \tag{3.8}
\end{align*}
$$

As $U, \tilde{U} \in B_{R, Z}(0)$, this can be further bounded, viz.

$$
\begin{aligned}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{Z_{T}} \leq & T C_{2}\left[\|u-\tilde{u}\|_{Z_{T}}\left(1+C_{1} R(2|A|+|B|+2|D|+|E|)\right)\right. \\
& \left.+\|v-\tilde{v}\|_{Z_{T}}\left(1+C_{1} R(2|C|+|B|+2|F|+|E|)\right)\right] \\
\leq & T C_{2}\left[1+C_{R}\right]\left(\|u-\tilde{u}\|_{Z_{T}}+\|v-\tilde{v}\|_{Z_{T}}\right) \\
= & T C_{2}\left[1+C_{R}\right]\|U-\tilde{U}\|_{Z_{T}}
\end{aligned}
$$

Let $R=2\left\|U_{0}\right\|_{Z}$ and demand that

$$
T C_{2}\left(1+C_{R}\right)=\frac{1}{2} .
$$

Then $\mathcal{A}$ is contractive on $B_{R, Z}(0)$ and for any $U \in B_{R}(0)$,

$$
\|\mathcal{A U}\|_{Z_{T}} \leq R .
$$

based on the same argument as in (3.5). Therefore, the Contraction Mapping Theorem can be applied to the integral equation (2.3) and the result is a unique $U$ that is solution of the given system.

This Lemma provides the local well-posedness on $H^{1}, H^{2}$ and in general in $H^{s}$, for $s>\frac{1}{2}$ as they are Banach algebras which satisfy both $(i)$ and (ii) in (3.6).

### 3.2 Smoothness and Regularity

Consider $U_{0}, V_{0} \in L_{\infty}(\mathbb{R})$ as initial data, fix $T>0, X_{T}$ defined as in (3.1) and let $U, V \in X_{T}$ be the corresponding solutions of the initial-value problem (1.1) - (1.2). To check the continuous dependence of solutions on the initial data, start locally. The overall conclusion follows from iteration of the argument below:

$$
U=\mathcal{A}_{U_{0}} U=U_{0}+\int_{0}^{t} \int_{-\infty}^{\infty} K(x-y)(U(y, s)+M) d y d s
$$

and

$$
V=\mathcal{A}_{V_{0}} V=V_{0}+\int_{0}^{t} \int_{-\infty}^{\infty} K(x-y)(V(y, s)+M) d y d s
$$

Using the fact that $\mathcal{A}$ is a contraction in $X_{T}$, at least for small $T$,

$$
\begin{aligned}
\|U-V\|_{X_{T}} & =\left\|\mathcal{A}_{U_{0}} U-\mathcal{A}_{V_{0}} V\right\|_{X_{T}} \\
& \leq\left\|\mathcal{A}_{U_{0}} U-\mathcal{A}_{U_{0}} V\right\|_{X_{T}}+\left\|\mathcal{A}_{U_{0}} V-\mathcal{A}_{V_{0}} V\right\|_{X_{T}} \\
& \leq \theta\|U-V\|_{X_{T}}+\left\|U_{0}-V_{0}\right\|_{\infty},
\end{aligned}
$$

where $0<\theta<1$. Therefore,

$$
\|U-V\|_{X_{T}} \leq \frac{1}{1-\theta}\left\|U_{0}-V_{0}\right\|_{\infty}
$$

and the mapping of the initial data $U_{0}$ to the solution $U$ is seen to be locally Lipschitz continuous. Moreover, as the initial-value problem (1.1) - (1.2) was solved locally thus far by converting it into a system of integral equations and applying a Picard iteration type of approach, there is automatically strong regularity results for the flow map. The argument for smoothness of the flow map has been presented in several contexts before (see, e.g. Bona, Sun and Zhang [7] or Bona and Tzvetkov [1]).

Notice that if $U_{0} \in C_{b}(\mathbb{R}) \times C_{b}(\mathbb{R})$, the solution determined via the last proposition lies in $C([0, T] \times \mathbb{R}) \times C([0, T] \times \mathbb{R})$.

Let $U$ be a continuous bounded solution of the initial-value problem (1.1) - (1.2). Since $U$ solves the system of integral equations it follows immediately that $U$ is differentiable with respect to $t$ and that

$$
U_{t}(x, t)=\int_{-\infty}^{\infty} K(x-y)[U(y, t)+M(U(y, t))] d y
$$

or, what is the same,

$$
U_{t}(x, t)=-\left(I-\partial_{x}^{2}\right)^{-1} \partial_{x}[U(x, t)+M(U(x, t))] .
$$

Clearly $U_{t}$ is a $C^{1}$-function in both $x$ and $t$ and therefore

$$
U_{t t}(x, t)=\int_{-\infty}^{\infty} K(x-y)\left[U_{t}(y, t)+M_{t}\right] d y .
$$

The latter is also a $C^{1}$-function in $x$ and $t$. By induction, one can conclude that $U$ is a $C^{\infty}$ function in $t$ and $\partial_{t}^{k} U$ is a $C^{1}$-function in $x$ for any $k \geq 1$.

Considering the spatial regularity, start with the fixed point $U=\mathcal{A} U$ of the system (2.3) written as

$$
\begin{equation*}
U(x, t)=U_{0}+\int_{0}^{t} K *[U(y, s)+M] d s \tag{3.9}
\end{equation*}
$$

and initial data $U_{0}$ that lies in $C_{b}^{1}(\mathbb{R}) \times C_{b}^{1}(\mathbb{R})$. This is an important assumption which will influence the conclusion drawn at the end of this paragraph. Then $U$ is also a $C^{1}$ function and direct calculation yields

$$
\begin{equation*}
\partial_{x}[K *(U+M)]=(U+M)-\frac{1}{2} \int_{-\infty}^{\infty} e^{-|x-y|}[U(y, t)+M] d y \tag{3.10}
\end{equation*}
$$

Moreover, as $U+M$ is $C^{1},(3.10)$ can be differentiated with respect to $x$ to obtain

$$
\begin{aligned}
\partial_{x}^{2}[K *(U+M)] & =(U+M)_{x}+\frac{1}{2} \int_{-\infty}^{\infty} \operatorname{sgn}(x-y) e^{-|x-y|}[U(y, t)+M] d y \\
& =(U+M)_{x}+[K *(U+M)]
\end{aligned}
$$

Therefore,

$$
\left(1-\partial_{x}^{2}\right)[K *(U+M)]=-(U+M)_{x}
$$

so $U$ will always have one derivative and the convolution on the RHS makes it even smoother. Consequently, by substituting the result above in (3.9), we find that

$$
\left(1-\partial_{x}^{2}\right) U_{t}=(U+M)_{x}
$$

Thus $U$ is a solution of the initial system (1.1) and it has reasonable regularity.
However, it is easily shown that $U$ can not have more spacial regularity than the initial data $U_{0}$. Assume that for some $t>0, U(\cdot, t)$ is a $C^{k}$-function in $x$ while $U_{0}(\cdot, t)$ is a $C^{k-1} \backslash C^{k}$ -
function in $x$. Then $(U+M)$ is a $C^{k}$-function in $x$ and the result after convolution with K , would be a $C^{k+1}$-function in $x$. Integrating this function with respect to $t$, would still preserve the $C^{k+1}$ property. Therefore from (3.9) it follows that $U_{0}(\cdot, t)$ is a $C^{k}$-function in $x$, contrary to the assumption made.

Furthermore, the smoothing associated with the temporal derivative for $U \in H^{k} \times H^{k}$ is approached just as in the $L_{\infty}$ case, to conclude that $U_{t}$ is spatially smoother and it lies in $C\left(0, \infty ; H^{k+1} \times H^{k+1}\right)$. As for the spatial regularity, the system solution has at most the same regularity as the initial data.

### 3.3 $\quad$ A priori Estimates and Global Well Posedness for Smooth Solutions

As seen in the previous sections, local well-posedness theory doesn't depend on assumptions about the coefficients of $P$ and $Q$. To move to global in time well-posedness, a priori estimates for the solution are needed. These will be derived below assuming that $u$ and $v$ are smooth solutions of the system which are not growing exponentially for $x \rightarrow \pm \infty$ and assuming the condition for the quadratic form $q$ to be positive definite is satisfied by its coefficients (see section 2.4).

In this case, the integrand of the $\Omega$ is positive definite too. Consequently, there is a $\lambda>0$ such that

$$
\begin{equation*}
\int_{-\infty}^{\infty}\left[u^{2}(x, t)+v^{2}(x, t)+u_{x}^{2}(x, t)+v_{x}^{2}(x, t)\right] d x \leq \lambda \Omega(u, v)=\lambda \Omega\left(u_{0}, v_{0}\right)=M_{1} . \tag{3.11}
\end{equation*}
$$

### 3.3.1 $\quad H^{1}$ - Bounds

Assume now that the initial condition $\left(u_{0}, v_{0}\right)$ is in $H^{1} \times H^{1}$ and let $(u, v)$ be the solution pair of our system corresponding to this data. By definition,

$$
\|(u, v)\|_{H^{1} \times H^{1}}^{2}=\int_{-\infty}^{\infty}\left(u^{2}+v^{2}+u_{x}^{2}+v_{x}^{2}\right) d x
$$

Hence, the relation (3.11) can also be written as

$$
\begin{equation*}
\|(u, v)\|_{H^{1} \times H^{1}}^{2} \leq M_{1}, \tag{3.12}
\end{equation*}
$$

so the $H^{1} \times H^{1}$ - norm of $(u, v)$ is also uniformly bounded independently of t .
Therefore the local well-posedness result obtained in section 3.1.2 can be extended to conclude the existence of globally defined solution that lies in $C\left(0, T ; H^{1} \times H^{1}\right)$ for any $T>0$. The solution is unique and depends continuously on the initial data in $H^{1} \times H^{1}$.

### 3.3.2 Differential Form of the Gronwall Lemma

This well-known result will be used going forward ([8], page 708).

Lemma 2. (i) Let $\eta(\cdot)$ be a nonnegative, absolutely continuous function on $[0, T]$ which satisfies for a.e. $t$ the differential inequality

$$
\begin{equation*}
\eta^{\prime}(t) \leq \phi(t) \eta(t)+\psi(t), \tag{3.13}
\end{equation*}
$$

where $\phi(t)$ and $\psi(t)$ are nonnegative, summable functions on $[0, T]$. Then

$$
\begin{equation*}
\eta(t) \leq e^{\int_{0}^{t} \phi(s) d s}\left(\eta(0)+\int_{0}^{t} \psi(s) d s\right) \tag{3.14}
\end{equation*}
$$

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

### 3.3.3 $\quad H^{2}$ - Bounds

Consider again the solution $(u, v)$ of the system (1.1), whose first few partial derivatives lie in $L_{2} \times L_{2}$. To estimate the growth of the $H^{2} \times H^{2}$-norm of a solution $U=(u, v)$, consider the following calculation:

$$
\begin{align*}
\frac{d}{d t} \int\left(u_{x}^{2}+v_{x}^{2}+u_{x x}^{2}+v_{x x}^{2}\right) d x= & 2 \int\left(u_{x} u_{x t}+v_{x} v_{x t}+u_{x x} u_{x x t}+v_{x x} v_{x x t}\right) d x \\
= & 2 \int\left(-u_{t} u_{x x}+u_{x x} u_{x x t}-v_{t} v_{x x}+v_{x x} v_{x x t}\right) d x \\
= & 2 \int\left(u_{x x}\left(u_{x x t}-u_{t}\right)+v_{x x}\left(v_{x x t}-v_{t}\right)\right) d x \\
= & \int\left(u_{x x}\left(u_{x}+P_{x}\right)+v_{x x}\left(v_{x}+Q_{x}\right)\right) d x \\
= & 2 \int\left(u_{x x} P_{x}+v_{x x} Q_{x}\right) d x \\
= & \int\left(2 A u u_{x} u_{x x}+B v u_{x} u_{x x}+B u v_{x} u_{x x}+2 C v v_{x} u_{x x}+\right. \\
& \left.+2 D u u_{x} v_{x x}+E u v_{x} v_{x x}+E v u_{x} v_{x x}+2 F v v_{x} v_{x x}\right) d x \tag{3.15}
\end{align*}
$$

The integrand on the RHS is a polynomial whose general monomial term is of the form $r s_{x} w_{x x}$ where $r, s, w$ are either $u$ or $v$. To obtain a bound on the $H^{2} \times H^{2}$ norm of the solution, let $M_{1}$
be the previously derived time-independent bound in $H^{1} \times H^{1}$. Then, elementary estimates imply that

$$
\begin{equation*}
\left|\int\left(r s_{x} w_{x x}\right) d x\right| \leq\|r\|_{\infty}\left\|s_{x}\right\|_{2}\left\|w_{x x}\right\|_{2} \leq\|r\|_{2}^{1 / 2}\left\|r_{x}\right\|_{2}^{1 / 2}\left\|s_{x}\right\|_{2}\left\|w_{x x}\right\|_{2} \leq M_{1}^{2}\left\|w_{x x}\right\|_{2} \tag{3.16}
\end{equation*}
$$

Let

$$
N(t)=\left[\int\left(u_{x}^{2}(x, t)+v_{x}^{2}(x, t)+u_{x x}^{2}(x, t)+v_{x x}^{2}(x, t)\right) d x\right]^{\frac{1}{2}}
$$

Applying the estimate (3.16) systematically to (3.15) yields

$$
\begin{aligned}
\frac{d}{d t} N^{2}(t) & \leq C_{1}\left(\left\|u_{x x}(\cdot, t)\right\|_{2}+\left\|v_{x x}(\cdot, t)\right\|_{2}\right) \\
& \leq C_{1} N(t)
\end{aligned}
$$

where $C_{1}$ is a constant that depends on $M_{1}$ and on the absolute values of the polynomial coefficients, but is time-independent. Hence,

$$
2 N(t) \frac{d}{d t} N(t) \leq C_{1} N(t)
$$

implies that

$$
\frac{d}{d t} N(t) \leq \frac{C_{1}}{2}
$$

It follows that

$$
\begin{equation*}
N(t) \leq N(0)+\frac{1}{2} C_{1} t \tag{3.17}
\end{equation*}
$$

for any $t>0$.

The relation (3.17) together with (3.12) lead to

$$
\|(u(\cdot, t), v(\cdot, t))\|_{H^{1} \times H^{1}}^{2}+N^{2}(t) \leq M_{1}+\left(N(0)+\frac{1}{2} C_{1} t\right)^{2}
$$

thus

$$
\|(u(\cdot, t), v(\cdot, t))\|_{H^{2} \times H^{2}}^{2} \leq M_{1}+\left(N(0)+\frac{1}{2} C_{1} t\right)^{2}
$$

Hence

$$
\|(u(\cdot, t), v(\cdot, t))\|_{H^{2} \times H^{2}} \leq C_{2}+C_{3} t
$$

where $C_{2}, C_{3}>0$ are t-independent constants that depend on $M_{1}$ and $N(0)$, so on $H^{2} \times H^{2}$ norm of $(u(\cdot, 0), v(\cdot, 0))$.

Therefore the local well-posedness result in time obtained in section 3.1 .2 can be extended to any time interval $[0, T]$. It concludes the existence of globally defined solution in $C\left(0, \infty ; H^{2} \times H^{2}\right)$. The solution is unique and depends continuously on the initial data in $H^{2} \times H^{2}$.

### 3.3.4 $H^{k}$ - Bounds

The a priori $H^{k} \times H^{k}$ bounds are determined inductively. The claim is that for any positive integer $k$, the $H^{k} \times H^{k}$ norm of the solution is bounded on bounded time intervals.

The cases $k=1$ and $k=2$ are in hand. Assuming now that for every $j$ with $2 \leq j \leq k$, there is a positive constant

$$
c_{j-1}=c_{j-1}\left(\|(u(\cdot, 0), v(\cdot, 0))\|_{H^{j-1} \times H^{j-1}} ; T\right)
$$

such that for any solution $(u, v)$ of the system (1.1), the inequality

$$
\begin{equation*}
\|\left(u(\cdot, t), v(\cdot, t) \|_{H^{j-1} \times H^{j-1}} \leq c_{j-1}\right. \tag{3.18}
\end{equation*}
$$

holds. Then there is a positive constant

$$
c_{k}=c_{k}\left(\|(u(\cdot, 0), v(\cdot, 0))\|_{H^{k} \times H^{k}} ; T\right)
$$

such that for $0 \leq t \leq T$

$$
\begin{equation*}
\|(u(\cdot, t), v(\cdot, t))\|_{H^{k} \times H^{k}} \leq c_{k} . \tag{3.19}
\end{equation*}
$$

Proof. To begin with, integration by parts and (1.1) implies

$$
\begin{aligned}
\frac{1}{2} \frac{d}{d t} \int\left[\left(\partial_{x}^{(k-1)} u\right)^{2}+\left(\partial_{x}^{(k)} u\right)^{2}\right] d x & =\int\left(\partial_{x}^{(k-1)} u \partial_{x}^{(k-1)} u_{t}+\partial_{x}^{(k)} u \partial_{x}^{(k)} u_{t}\right) d x \\
& =\int\left(\partial_{x}^{(k-1)} u \partial_{x}^{(k-1)} u_{t}-\partial_{x}^{(k-1)} u \partial_{x}^{(k+1)} u_{t}\right) d x \\
& =\int \partial_{x}^{(k-1)} u\left(\partial_{x}^{(k-1)} u_{t}-\partial_{x}^{(k-1)} u_{x x t}\right) d x \\
& =-\int \partial_{x}^{(k-1)} u \partial_{x}^{(k-1)}\left(u_{x}-P_{x}\right) d x \\
& =-\int\left(\partial_{x}^{(k-1)} u\right)\left(\partial_{x}^{(k)} P\right) d x \\
& =\int\left(\partial_{x}^{(k)} u\right)\left(\partial_{x}^{(k-1)} P\right) d x \\
\frac{1}{2} \frac{d}{d t} \int\left[\left(\partial_{x}^{(k-1)} v\right)^{2}+\left(\partial_{x}^{(k)} v\right)^{2}\right] d x & =\int\left(\partial_{x}^{(k)} v\right)\left(\partial_{x}^{(k-1)} Q\right) d x
\end{aligned}
$$

The results above will be used in the following estimate

$$
\begin{align*}
\frac{1}{2} \frac{d}{d t} \int\left[\left(\partial_{x}^{(k-1)} u\right)^{2}\right. & \left.+\left(\partial_{x}^{(k-1)} v\right)^{2}+\left(\partial_{x}^{(k)} u\right)^{2}+\left(\partial_{x}^{(k)} v\right)^{2}\right] d x= \\
& \leq \int\left(\partial_{x}^{(k)} u\right)\left(\partial_{x}^{(k-1)} P\right)+\left(\partial_{x}^{(k)} v\right)\left(\partial_{x}^{(k-1)} Q\right) d x \tag{3.20}
\end{align*}
$$

The next step is to estimate the monomials that the polynomial above is made of. These have the form

$$
\begin{align*}
\int\left|\left(\partial_{x}^{(k)} w\right)\left(\partial_{x}^{(j)} r\right)\left(\partial_{x}^{(k-1-j)} s\right)\right| d x & \leq\left\|\partial_{x}^{(k)} w\right\|_{2}\left\|\partial_{x}^{(j)} r\right\|_{\infty}\left\|\partial_{x}^{(k-1-j)} s\right\|_{2} \\
& \leq\left\|\partial_{x}^{(k)} w\right\|_{2}\left\|\partial_{x}^{(j)} r\right\|_{H^{1}}\|s\|_{H^{k-1-j}} \\
& \leq\left\|\partial_{x}^{(k)} w\right\|_{2}\|r\|_{H^{j}}\|s\|_{H^{k-1-j}} \\
& \leq\left\|\partial_{x}^{(k)} w\right\|_{2}\|r\|_{H^{j}}\|s\|_{H^{k-1-j}} \\
& \leq c_{j} c_{k-1-j}\left\|\partial_{x}^{(k)} w\right\|_{2} \tag{3.21}
\end{align*}
$$

where $r, s$ and $w$ are either $u$ or $v$. Thus,

$$
\begin{align*}
& \frac{1}{2} \frac{d}{d t} \int\left[\left(\partial_{x}^{(k-1)} u\right)^{2}+\left(\partial_{x}^{(k-1)} v\right)^{2}+\left(\partial_{x}^{(k)} u\right)^{2}+\left(\partial_{x}^{(k)} v\right)^{2}\right] d x \\
& \quad \leq c\left(\left\|\partial_{x}^{(k)} u\right\|_{2}+\left\|\partial_{x}^{(k)} v\right\|_{2}\right) \\
& \quad \leq c\left(\int\left[\left(\partial_{x}^{(k-1)} u\right)^{2}+\left(\partial_{x}^{(k-1)} v\right)^{2}+\left(\partial_{x}^{(k)} u\right)^{2}+\left(\partial_{x}^{(k)} v\right)^{2}\right] d x\right)^{\frac{1}{2}}, \tag{3.22}
\end{align*}
$$

where $c$ is a constant. Solving the ordinary differential inequality yields

$$
\left\|\partial_{x}^{(k-1)} u\right\|_{2}+\left\|\partial_{x}^{(k-1)} v\right\|_{2}+\left\|\partial_{x}^{(k)} u\right\|_{2}+\left\|\partial_{x}^{(k)} v\right\|_{2} \leq c T .
$$

This together with (3.18) indicates the following:

$$
\begin{equation*}
\|(u(., t), v(., t))\|_{H^{k} \times H^{k}} \leq\|(u(., 0), v(., 0))\|_{H^{k} \times H^{k}}+c T . \tag{3.23}
\end{equation*}
$$

Therefore, there is a positive constant

$$
c_{k}=c_{k}\left(\|(u(., 0), v(., 0))\|_{H^{k} \times H^{k}} ; T\right)
$$

such that

$$
\begin{equation*}
\|(u(., t), v(., t))\|_{H^{k} \times H^{k}} \leq c_{k} . \tag{3.24}
\end{equation*}
$$

This completes the induction.

## CHAPTER 4

## WELL-POSEDNESS IN $L_{2} \times L_{2}$

### 4.1 Local Existence

Consider the Banach space $C\left(0, T ; L_{2}(\mathbb{R}) \times L_{2}(\mathbb{R})\right)$ with its usual norm

$$
\|U-\tilde{U}\|_{C\left(0, T ; L_{2}(\mathbb{R}) \times L_{2}(\mathbb{R})\right)}:=\sup _{t \in[0, T]}\|u(\cdot, t)-\tilde{u}(\cdot, t)\|_{2}+\sup _{t \in[0, T]}\|v(\cdot, t)-\tilde{v}(\cdot, t)\|_{2}
$$

where

$$
U=\binom{u}{v} \quad \text { and } \quad \tilde{U}=\binom{\tilde{u}}{\tilde{v}}
$$

Let $X_{T}:=C\left(0, T ; L_{2}(\mathbb{R})\right)$. For simplicity, $C\left(0, T ; L_{2}(\mathbb{R}) \times L_{2}(\mathbb{R})\right)$-norm of a vector $U$ will be denoted by $\|U\|_{X_{T}}$.

To establish local existence via the Contraction Mapping Principle, the domain for the operator $\mathcal{A}$ will be restricted to the closed subset

$$
B_{R}(0):=\left\{U \in C\left(0, T ; L_{2}(\mathbb{R}) \times L_{2}(\mathbb{R})\right):\|U\|_{X_{T}} \leq R\right\}
$$

where $R>0$ and $T>0$ will be determined. Consider $U, \tilde{U} \in B_{R}(0)$, and look for an estimate of $\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}}$. The aim is to choose specific values of $R$ and $T$ so that $\mathcal{A}$ is a contraction
on $B_{R}(0)$. Using the definition of the operator $\mathcal{A}$, applying the triangle inequality and Holder's inequality one calculates that for $U=(u, v)$ and $\tilde{U}=(\tilde{u}, \tilde{v})$,

$$
\begin{align*}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}}= & \sup _{t \in[0, T]}\left\|\int_{0}^{t} K *[u(x, t)+P(u, v)-\tilde{u}(x, t)-P(\tilde{u}, \tilde{v})] d s\right\|_{2} \\
& +\sup _{t \in[0, T]}\left\|\int_{0}^{t} K *[v(x, t)+Q(u, v)-\tilde{v}(x, t)-Q(\tilde{u}, \tilde{v})] d s\right\|_{2} \\
\leq & T\|K *(u-\tilde{u})\|_{X_{T}}+T\|K *(v-\tilde{v})\|_{X_{T}} \\
& +T\|K *[(u-\tilde{u})(A(u+\tilde{u})+B v)+(v-\tilde{v})(C(v+\tilde{v})+B \tilde{u})]\|_{X_{T}} \\
& +T\|K *[(u-\tilde{u})(D(u+\tilde{u})+E \tilde{v})+(v-\tilde{v})(F(v+\tilde{v})+E u)]\|_{X_{T}} \\
\leq & 2 T\left[\|u-\tilde{u}\|_{X_{T}}\left(\frac{1}{2}+(|A|+|D|)\left(\|u\|_{X_{T}}+\|\tilde{u}\|_{X_{T}}\right)+B\|v\|_{X_{T}}+|E|\|\tilde{v}\|_{X_{T}}\right)\right. \\
& \left.+\|v-\tilde{v}\|_{X_{T}}\left(\frac{1}{2}+(|C|+|F|)\left(\|v\|_{X_{T}}+\|\tilde{v}\|_{X_{T}}\right)+|B|\|\tilde{u}\|_{X_{T}}+|E|\|u\|_{X_{T}}\right)\right] . \tag{4.1}
\end{align*}
$$

Using the fact that $U, \tilde{U} \in B_{R}(0)$, the inequality becomes

$$
\begin{align*}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}} \leq & 2 T\left[\|u-\tilde{u}\|_{X_{T}}\left(\frac{1}{2}+2 R(|A|+|D|)+|B| R+|E| R\right)\right. \\
& \left.+\|v-\tilde{v}\|_{X_{T}}\left(\frac{1}{2}+2 R(|C|+|F|)+|B| R+|E| R\right)\right] \\
\leq & 2 T\left(\frac{1}{2}+R C_{R}\right)\|U-V\|_{X_{T}} . \tag{4.2}
\end{align*}
$$

Let $R=2\left\|U_{0}\right\|_{X_{T}}$ and demand that

$$
T\left(\frac{1}{2}+R C_{R}\right)=\frac{1}{4}
$$

Then $\mathcal{A}$ is contractive on $B_{R}(0)$ and for any $U \in B_{R}(0)$,

$$
\begin{equation*}
\|\mathcal{A} U\|_{X_{T}} \leq R \tag{4.3}
\end{equation*}
$$

as in (3.5). Therefore, the Contraction Mapping Theorem can be applied to the integral equation (2.3) and the result is a unique $U$ that is solution of the given system. Moreover, using the same argument as in the previous chapter, the solution found above is seen to be continuously dependent on the initial data.

### 4.2 Long Time Solutions for Small Initial Data

Now that the $L_{2} \times L_{2}$ local existence has been established, the next step is to look for long time solutions when one starts off with $L_{2}$-small initial data. One of the forms the system (1.1) could be written in is

$$
\begin{equation*}
U_{t}(x, t)=\int_{-\infty}^{\infty} K *[U+M] d x \tag{4.4}
\end{equation*}
$$

If $R$ denotes the operator

$$
\begin{equation*}
R=-\left(I-\partial_{x}^{2}\right)^{-1} \partial_{x} \tag{4.5}
\end{equation*}
$$

then since R is a skew-adjoint operator, $\int U R(U)=0$, Young's inequality therefore implies that

$$
\begin{align*}
\int_{-\infty}^{\infty} U U_{t} d x & =\int U[R(U+M)] d x=\int U(K * M) d x= \\
& \leq\|U\|_{2}\|K * M\|_{2} \leq\|U\|_{2}\|K\|_{2}\|M\|_{1} \tag{4.6}
\end{align*}
$$

To continue the last estimate, the norm $\|M\|_{1}$ needs to be calculated and estimated in more detail. Using the triangle and the Cauchy-Schwartz inequalities, it follows that

$$
\begin{aligned}
\|M\|_{1} & =\left\|A u^{2}+B u v+C v^{2}\right\|_{1}+\left\|D u^{2}+E u v+F v^{2}\right\|_{1} \\
& \leq(|A|+|D|)\left\|u^{2}\right\|_{1}+(|B|+|E|)\|u v\|_{1}+(|C|+|F|)\left\|v^{2}\right\|_{1} \\
& \leq(|A|+|D|)\|u\|_{2}^{2}+(|B|+|E|)\|u\|_{2}\|v\|_{2}+(|C|+|F|)\|v\|_{2}^{2} \\
& \leq c\left(\|u\|_{2}^{2}+\|v\|_{2}^{2}\right) \\
& \leq c\|U\|_{2}^{2} .
\end{aligned}
$$

Therefore,

$$
\begin{align*}
\int_{-\infty}^{\infty} U U_{t} d x & \leq\|U\|_{2}\|K\|_{2} c\|U\|_{2}^{2} \\
& \leq\|U\|_{2} \tilde{c}\|U\|_{2}^{2}=\tilde{c}\|U\|_{2}^{3} \tag{4.7}
\end{align*}
$$

where c and $\tilde{c}$ are constants.

As $U \in C\left(0, T ; L_{2}\right)$ and $U_{t} \in C\left(0, T ; H^{1}\right)$ it follows that $U U_{t} \in C^{1}\left(0, T ; L_{2}\right)$ and

$$
\begin{equation*}
\int_{-\infty}^{\infty} U U_{t} d x=\frac{1}{2} \frac{d}{d t}\|U\|_{2}^{2} \tag{4.8}
\end{equation*}
$$

holds true.

The ordinary differential inequality

$$
\frac{1}{2} \frac{d}{d t}\|U\|_{2}^{2} \leq \tilde{c}\|U\|_{2}^{3}
$$

that follows from (4.7) and (4.8), implies that

$$
\frac{d}{d t}\|U\|_{2} \leq \tilde{c}\|U\|_{2}^{2}
$$

Solving the inequality, it follows that

$$
-\frac{1}{\|U\|_{2}}+\frac{1}{\left\|U_{0}\right\|_{2}} \leq \tilde{c} T \quad \text { or } \quad\|U(T)\|_{2} \leq \frac{\left\|U_{0}\right\|_{2}}{1-\tilde{c} T\left\|U_{0}\right\|_{2}}
$$

In conclusion, if the $L_{2}$-norm is small to start with, then one can infer existence of an $L_{2}$-solution over a long time.

### 4.3 Global Well-Posedness in $L_{2} \times L_{2}$

Based on the results obtained in the two preceding sections, local well-posedness will be extended to global well-posedness using a low-frequency/high-frequency decomposition method.

Fix $T>0$. If the initial data $\left(u_{0}, v_{0}\right)$ is small enough in the $L_{2} \times L_{2}$-norm, then the system (1.1) is known to have a unique solution in $C\left(0, T ; L_{2} \times L_{2}\right)$. Consider initial data of arbitrary size. Split the initial data $\left(u_{0}, v_{0}\right) \in L_{2} \times L_{2}$ into two pieces, one "small" and the other "smooth". This is possible in the following manner:

$$
\begin{aligned}
& u_{0}=u_{01}+u_{02}=\left(u_{0}-\Phi_{\epsilon} * u_{0}\right)+\Phi_{\epsilon} * u_{0}, \\
& v_{0}=v_{01}+v_{02}=\left(v_{0}-\Phi_{\epsilon} * v_{0}\right)+\Phi_{\epsilon} * v_{0},
\end{aligned}
$$

where $\Phi_{\epsilon}$ is a standard mollifier. The important properties that are to be used here are

> (i) $\Phi_{\epsilon} * w_{0}$ and $\Phi_{\epsilon} * v_{o}$ are smooth in $H^{k}$,
> (ii) $\Phi_{\epsilon} * w_{0} \rightarrow w_{0}$ in $L_{2}$ when $\epsilon \rightarrow 0$.

Hence, by choosing $\epsilon>0$ small enough, we can insure that

$$
\begin{align*}
& u_{01}=\left(u_{0}-\Phi_{\epsilon} * u_{0}\right) \text { is small in } L_{2} \text {-norm and } u_{02}=\Phi_{\epsilon} * u_{0} \text { is smooth, } \\
& v_{01}=\left(v_{0}-\Phi_{\epsilon} * v_{0}\right) \text { is small in } L_{2} \text {-norm and } v_{02}=\Phi_{\epsilon} * v_{0} \text { is smooth. } \tag{4.10}
\end{align*}
$$

Applying the previous results to the system (1.1) with initial data $\left(u_{01}, v_{01}\right)$, it is seen that there is a unique solution $(u, v)$ that lies in $C\left([0, T] ; L_{2} \times L_{2}\right)$.

Consider now the initial-value problem

$$
\left\{\begin{array}{l}
w_{t}+w_{x}-w_{x x t}+P(w, z)_{x}+(2 A u w+B u z+B v w+2 C v z)_{x}=0  \tag{4.11}\\
z_{t}+z_{x}-z_{x x t}+Q(w, z)_{x}+(2 D u w+E u z+E v w+2 F v z)_{x}=0 \\
w_{0}=u_{02}, \quad z_{0}=v_{02}
\end{array}\right.
$$

with the initial data $\left(u_{02}, v_{02}\right) \in H^{k} \times H^{k}$. If this system has a solution $\left(u^{\prime}, v^{\prime}\right)$ in $C\left([0, T] ; H^{k} \times\right.$ $\left.H^{k}\right)$ then $\left(u+u^{\prime}, v+v^{\prime}\right)$ will solve the initial system (1.1).

An a priori bound for $(w, z)$ on $H^{1} \times H^{1}$-norm is necessary to extend the local existence for (4.11) to the interval $[0, T]$. To obtain an estimation of the $H^{1} \times H^{1}$-norm, the following formal
calculation comes to our rescue. The computation is easily justified by use of the continuous dependence result and the time derivative makes sense as all the components are in $C\left(0, T ; H^{1}\right)$.

$$
\begin{align*}
\frac{1}{2} & \frac{d}{d t} \int\left(w^{2}+w_{x}^{2}+z^{2}+z_{x}^{2}\right) d x=\int\left(w w_{t}+w_{x} w_{x t}+z z_{t}+z_{x} z_{x t}\right) d x \\
= & \int\left[w\left(w_{t}-w_{x x t}\right)+z\left(z_{t}-z_{x x t}\right)\right] d x \\
= & -\int\left[w\left(w_{x}+P_{x}+(2 A u w+B u z+B v w+2 C v z)_{x}\right)\right. \\
& \left.+z\left(z_{x}+Q_{x}+(2 D u w+E u z+E v w+2 F v z)_{x}\right)\right] d x \\
= & -\int\left[\frac{B}{2}\left(w^{2}\right)_{x} z+\left(C-\frac{E}{2}\right) w\left(z^{2}\right)_{x}+A u_{x}\left(w^{2}\right)+B(u w z)_{x}\right. \\
& +(2 D-B) u z_{x} w+2 D u_{x} w z+\frac{B}{2} v_{x} w^{2}+2 C(v w z)_{x}+(E-2 C) v z w_{x} \\
& \left.+D\left(w^{2}\right)_{x} z+\frac{E}{2} u_{x} z^{2}+E v_{x} z w+F v_{x} z^{2}\right] d x \\
= & -\int\left[\left(A u_{x}+\frac{B}{2} v_{x}\right) w^{2}+\left(\frac{E}{2} u_{x}+F v_{x}\right) z^{2}+\left(\frac{B}{2}+C-\frac{E}{2}\right) w\left(z^{2}\right)_{x}\right. \\
& \left.+D\left(w^{2}\right)_{x} z+((2 D-B) u+(E-2 C) v) z w_{x}+\left(2 D u_{x}+E v_{x}\right) w z\right] d x . \tag{4.12}
\end{align*}
$$

The integrand on the RHS is a polynomial whose monomial terms are of the form $u_{x} r s, v_{x} r s$, $r\left(s^{2}\right)_{x}, u r s_{x}$ or $v r s_{x}$ where $r, s$ are either $w$ or $z$. To obtain a bound in the $H^{1} \times H^{1}$ norm of the solution $(w, z)$, the following estimates are used:

$$
\begin{equation*}
\left|\int u r_{x} s d x\right| \leq \int\left|u r_{x} s\right| d x \leq\|u\|_{2}\|s\|_{\infty}\left\|r_{x}\right\|_{2} \leq\|u\|_{2} c\|s\|_{H^{1}}\|r\|_{H^{1}} \tag{4.13}
\end{equation*}
$$

$$
\begin{align*}
\left|\int\left(u_{x} r s\right) d x\right| & =\left|\int\left(u r_{x} s+u r s_{x}\right) d x\right| \leq \int\left|u r_{x} s\right| d x+\int\left|u r s_{x}\right| d x \\
& \leq\|u\|_{2}\left(c\|s\|_{H^{1}}\|r\|_{H^{1}}+\tilde{c}\|r\|_{H^{1}}\|s\|_{H^{1}}\right) \\
& \leq c\|u\|_{2}\|s\|_{H^{1}}\|r\|_{H^{1}} \tag{4.14}
\end{align*}
$$

and

$$
\begin{align*}
\left|\int r\left(s^{2}\right)_{x} d x\right| & \leq 2\|r\|_{2}\|s\|_{\infty}\left\|s_{x}\right\|_{2} \\
& \leq 2\|r\|_{H^{1}} c\|s\|_{H^{1}}\|s\|_{H^{1}} \\
& \leq c\|r\|_{H^{1}}\left(\|s\|_{H^{1}}\right)^{2} . \tag{4.15}
\end{align*}
$$

Applying the just derived estimates to (4.12) yields

$$
\begin{equation*}
\frac{d}{d t}\|(w, z)\|_{H^{1} \times H^{1}}^{2} \leq c\|(u, v)\|_{2}\left(\|(w, z)\|_{H^{1} \times H^{1}}\right)^{2}+\tilde{c}\left(\|(w, z)\|_{H^{1} \times H^{1}}\right)^{3} \tag{4.16}
\end{equation*}
$$

where the constants $c$ and $\tilde{c}$ are independent of $t \in[0, T]$. Consequently,

$$
\begin{equation*}
\frac{d}{d t}\|(w, z)\|_{H^{1} \times H^{1}} \leq c_{1}\|(w, z)\|_{H^{1} \times H^{1}}+\tilde{c}\left(\|(w, z)\|_{H^{1} \times H^{1}}\right)^{2}, \tag{4.17}
\end{equation*}
$$

from which it follows via a Gronwall-type argument that

$$
\begin{equation*}
\|(w, z)\|_{H^{1} \times H^{1}} \leq \frac{\left\|\left(u_{02}, v_{02}\right)\right\|_{H^{1} \times H^{1}}}{1-c_{2} e^{c_{1} T}\left\|\left(u_{02}, v_{02}\right)\right\|_{H^{1} \times H^{1}}}, \tag{4.18}
\end{equation*}
$$

where the constants $c_{1}$ and $c_{2}$ are independent of $t \in[0, T]$. Hence, local well-posedness on $[0, T]$ is insured.

## CHAPTER 5

## INTERPOLATION AND WELL-POSEDNESS IN $H^{S} \times H^{S}$

To extend the results obtained for $L_{2} \times L_{2}$ and the $H^{k} \times H^{k}$ spaces, $k=1,2, \cdots$, nonlinear interpolation theory will be used. The outcome will be global well-posedness in the $L_{2}$-based Hilbert spaces $H^{s} \times H^{s}$ for all $s \geq 0$.

Nonlinear interpolation of Banach spaces and their operators goes back to the work of Peetre [9] and Peetre-Lions [10]. The following version of the theory will suffice for our purposes (see [11]). More general results of this nature are available.

Theorem 1. Let $r$ and $s$, with $r>s$ be two non-negative real numbers. Suppose that for some $T>0$, the operator $A$ is defined on both $H^{r} \times H^{r}$ and $H^{s} \times H^{s}$ and maps these spaces continuously into $C\left([0, T] ; H^{r} \times H^{r}\right)$ and $C\left([0, T] ; H^{s} \times H^{s}\right)$, respectively. Suppose in addition that $A$ respects the inequalities
(i) $\|A f-A g\|_{C\left([0, T] ; H^{s} \times H^{s}\right)} \leq c_{0}\left(\|f\|_{H^{s} \times H^{s}}+\|g\|_{H^{s} \times H^{s}}\right)\|f-g\|_{H^{s} \times H^{s}}$
and
(ii) $\|A h\|_{H^{r} \times H^{r}} \leq c_{1}\left(\|h\|_{H^{s} \times H^{s}}\right)\|h\|_{H^{r} \times H^{r}}$,
for some continuous functions $c_{0}$ and $c_{1}$.
Then, for any $b \in[s, r], A$ maps $H^{b} \times H^{b}$ continuously into $C\left([0, T] ; H^{b} \times H^{b}\right)$ and

$$
\|A f\|_{C\left([0, T] ; H^{b} \times H^{b}\right)} \leq c_{b}\left(\|f\|_{\left.H^{s} \times H^{s}\right)}\right)\|f\|_{H^{b} \times H^{b}}
$$

where, for $\gamma>0, c_{b}(\gamma)$ may be taken in the form

$$
c_{b}(\gamma)=4 c_{0}(4 \gamma)^{1-\theta} c_{1}(3 \gamma)^{\theta}
$$

with $c_{0}$ and $c_{1}$ as in (i) and (ii) and $\theta=\frac{b-s}{r-s}$.

Starting with $\left(u_{0}, v_{0}\right)$ in $H^{k} \times H^{k}$ with $k$ integer, as initial data for the system (1.1), denote by $A$ the mapping that associates to $\left(u_{0}, v_{0}\right)$ the pair $(u, v)$ in $C\left(0, T ; H^{k} \times H^{k}\right)$ that solves the system (1.1). So

$$
A: H^{k} \times H^{k} \quad \rightarrow \quad C\left([0, T] ; H^{k} \times H^{k}\right)
$$

for any $k$ non-negative integer and for any $T>0$. As seen in the previous chapter, the Contraction Mapping Theorem, which was used to prove the local existence, also assures that the mapping $A$ is locally Lipschitz in $H^{k} \times H^{k}$. The latter result can be iterated to conclude that for any $T>0$, there is continuous function $c_{k}^{T}$ such that

$$
\begin{align*}
& \left\|A\left(u_{0}, v_{0}\right)-A\left(\widetilde{u}_{0}, \widetilde{v}_{0}\right)\right\|_{C\left([0, T] ; H^{k} \times H^{k}\right)}  \tag{5.1}\\
& \quad \leq c_{k}^{T}\left(\left\|u_{0}\right\|_{H^{k}}+\left\|v_{0}\right\|_{H^{k}}+\left\|\widetilde{u}_{0}\right\|_{H^{k}}+\left\|\widetilde{v}_{0}\right\|_{H^{k}}\right)\left\|\left(u_{0}-\widetilde{u}_{0}, v_{0}-\widetilde{v}_{0}\right)\right\|_{H^{k} \times H^{k}}
\end{align*}
$$

Of course, it is probably the case that $c_{k}^{T}$ grows with $T$, but $T$ is fixed in the discussion.

Start by considering $A$ as a mapping of $L^{2} \times L^{2}$ into the space $C\left([0, T] ; L^{2} \times L^{2}\right)$ and of $H^{1} \times H^{1}$ into $C\left([0, T] ; H^{1} \times H^{1}\right)$. Because of the inequalities (3.11) and (4.2), the conditions (ii) and (i) in Theorem 1 are respectively satisfied and it is thereby concluded by applying the Theorem 1 that for any fixed $T>0$ and $s$ with $0<s<1, s \neq \frac{1}{2}, A$ maps $H^{s} \times H^{s}$ continuously into $C\left([0, T] ; H^{s} \times H^{s}\right)$ and respects the the inequality

$$
\begin{align*}
& \left\|A\left(u_{0}, v_{0}\right)\right\|_{C\left([0, T] ; H^{s} \times H^{s}\right)}  \tag{5.2}\\
& \quad \leq c_{s}^{T}\left(\left\|u_{0}\right\|_{L^{2}}+\left\|v_{0}\right\|_{L^{2}}\right)\left\|\left(u_{0}, v_{0}\right)\right\|_{H^{s} \times H^{s}},
\end{align*}
$$

where $c_{s}^{T}(y)=4 c_{0}^{T}(4 y)^{1-\theta} c_{1}^{T}(3 y)^{\theta}$ and $c_{0}^{T}$ and $c_{1}^{T}$ are as in formulas (5.1) respectively for $k=0$ and $k=1$.

Re-interpolating between $s=0$ and $s=\frac{3}{4}$, say, the same conclusion holds for $s=\frac{1}{2}$. Iteratively, one can repeat these arguments for $1<s<2,2<s<3$ and so on. It follows that the system (1.1) is globally well posed in all the $L_{2}$-based Hilbert spaces $H^{s} \times H^{s}, s \geq 0$, as long as the coefficients of the polynomials $P$ and $Q$ satisfy (2.6) and (2.11).

## CHAPTER 6

## CUBIC CASE

In this chapter, the polynomials P and Q appearing in the our system

$$
\left\{\begin{array}{l}
u_{t}+u_{x}-u_{x x t}+P(u, v)_{x}=0 \\
v_{t}+v_{x}-v_{x x t}+Q(u, v)_{x}=0
\end{array}\right.
$$

are taken to be homogeneous, cubic polynomials, viz.

$$
\begin{align*}
& \Pi(u, v)=A_{11} u^{3}+A_{12} u^{2} v+A_{13} u v^{2}+A_{14} v^{3}, \\
& \Psi(u, v)=A_{21} u^{3}+A_{22} u^{2} v+A_{23} u v^{2}+A_{24} v^{3}, \tag{6.1}
\end{align*}
$$

with all the coefficients $A_{i j} \in \mathbb{R}, i=1,2, j=1,2,3,4$. The analysis of the cubic case is a little more involved than the quadratic one.

### 6.1 Invariants

Consider the quadratic functional

$$
\Omega(u, v):=\int\left(\alpha u^{2}+\beta u v+\gamma v^{2}+\alpha u_{x}^{2}+\beta u_{x} v_{x}+\gamma v_{x}^{2}\right) d x
$$

where $\alpha, \beta, \gamma$ are real numbers to be determined. A formal calculation where integration by parts is used freely and it is presumed both $u$ and $v$ vanish at $x= \pm \infty$, leads to the formula

$$
\frac{d}{d t} \Omega=-2 \alpha \int u \Pi_{x} d x-\beta \int v \Pi_{x} d x-\beta \int u \Psi_{x} d x-2 \gamma \int v \Psi_{x} d x
$$

Writing $\Pi_{x}$ and $\Psi_{x}$ out in detail yields

$$
\begin{align*}
\frac{d}{d t} \Omega= & \int\left(-6 A_{11} \alpha-3 A_{12} \beta\right) u^{3} u_{x} d x+\int\left(-3 A_{14} \beta-6 A_{24} \gamma\right) v^{3} v_{x} d x \\
& -\int\left[\left(4 A_{12} \alpha+3 A_{11} \beta+2 A_{22} \beta+6 A_{21} \gamma\right) u^{2} v u_{x}+\left(2 A_{12} \alpha+A_{22} \beta\right) u^{3} v_{x}\right] d x \\
& -\int\left[\left(6 A_{14} \alpha+2 A_{13} \beta+3 A_{24} b+4 A_{23} \gamma\right) u v^{2} v_{x}+\left(A_{13} \beta+2 A_{23} \gamma\right) v^{3} u_{x}\right] d x \\
& -\int\left(2 A_{13} \alpha+2 A_{12} \beta+A_{23} \beta+4 A_{22} \gamma\right) u v^{2} u_{x} d x \\
& -\int\left(4 A_{13} \alpha+A_{12} \beta+2 A_{23} \beta+2 A_{22} \gamma\right) u^{2} v v_{x} d x . \tag{6.2}
\end{align*}
$$

The first two terms in the sum above are integrals which have integrands that are x -derivatives. Thus, for smooth solutions going to 0 at $\pm \infty$, they vanish without further assumptions. The last three integrals in the sum would vanish without further assumptions on $(u, v)$ if and only if

$$
\begin{align*}
2 A_{12} \alpha+A_{22} \beta & =3 A_{11} \beta+6 A_{21} \gamma, \\
A_{13} \beta+2 A_{23} \gamma & =6 A_{14} \alpha+3 A_{24} \beta, \\
2 A_{13} \alpha+A_{23} \beta & =A_{12} \beta+2 A_{22} \gamma . \tag{6.3}
\end{align*}
$$

Thus $\alpha, \beta$ and $\gamma$ must solve the three equations

$$
\left\{\begin{array}{l}
2 A_{12} \alpha+\left(A_{22}-3 A_{11}\right) \beta-6 A_{21} \gamma=0,  \tag{6.4}\\
6 A_{14} \alpha+\left(3 A_{24}-A_{13}\right) \beta-2 A_{23} \gamma=0, \\
2 A_{13} \alpha+\left(A_{23}-A_{12}\right) \beta-2 A_{22} \gamma=0,
\end{array}\right.
$$

for it to be the case that the time derivative of $\Omega(u, v)$ is zero. This system of equations has a non-trivial solution only if its determinant

$$
\Delta=\left|\begin{array}{ccc}
2 A_{12} & A_{22}-3 A_{11} & -6 A_{21}  \tag{6.5}\\
6 A_{14} & 3 A_{24}-A_{13} & -2 A_{23} \\
2 A_{13} & A_{23}-A_{12} & -2 A_{22}
\end{array}\right|
$$

is zero, which is to say,

$$
\begin{align*}
-4\left[\left(A_{12} A_{22}-A_{13} A_{21}\right)\left(3 A_{24}-A_{13}\right)\right. & +\left(A_{23}-A_{12}\right)\left(9 A_{21} A_{14}-A_{12} A_{23}\right) \\
& \left.+\left(3 A_{11}-A_{22}\right)\left(3 A_{22} A_{14}-A_{13} A_{23}\right)\right]=0 . \tag{6.6}
\end{align*}
$$

For example, one possibility to have $\Delta=0$ is that

$$
\left\{\begin{array}{c}
A_{12}=3 A_{14},  \tag{6.7}\\
A_{22}-3 A_{11}=3 A_{24}-A_{13}, \\
3 A_{21}=A_{23} .
\end{array}\right.
$$

In this case, the cubics that couple the original system's unknowns $u$ and $v$ would have the form

$$
\begin{aligned}
& \Pi(u, v)=\widetilde{A} u^{3}+3 \widetilde{B} u^{2} v+\widetilde{C} u v^{2}+\widetilde{B} v^{3} \\
& \Psi(u, v)=\widetilde{D} u^{3}+[3(\widetilde{A}+\widetilde{E})-\widetilde{C}] u^{2} v+3 \widetilde{D} u v^{2}+\widetilde{E} v^{3}
\end{aligned}
$$

where $\widetilde{A}, \widetilde{B}, \widetilde{C}, \widetilde{D}, \widetilde{E}$ are real numbers. With this choice of polynomial coefficients, the values of $\alpha, \beta$ and $\gamma$ from $\Omega$ are

$$
\left(\begin{array}{l}
\alpha \\
\beta \\
\gamma
\end{array}\right)=\left(\begin{array}{c}
(\widetilde{C}-3 \widetilde{E})(3 \widetilde{A}-\widetilde{C}+3 \widetilde{E})-9 \widetilde{D}(\widetilde{B}-\widetilde{D}) \\
-6 \widetilde{B}(\widetilde{C}-3 \widetilde{E})+18 \widetilde{A} \widetilde{B}-6 \widetilde{C} \widetilde{D} \\
\widetilde{C}(\widetilde{C}-3 \widetilde{E})-9 \widetilde{B}(\widetilde{B}-\widetilde{D})
\end{array}\right),
$$

up to a scalar multiple.

If $\alpha, \beta$ and $\gamma$ are a non-trivial solution of (6.4), it follows immediately that

$$
\begin{aligned}
\frac{\partial}{\partial u}(\beta \Pi+2 \gamma \Psi) & =\left(3 A_{11} \beta+6 A_{21} \gamma\right) u^{2}+2\left(A_{12} \beta+2 A_{22} \gamma\right) u v+\left(A_{13} \beta+2 A_{23} \gamma\right) v^{2} \\
\frac{\partial}{\partial v}(2 \alpha \Pi+\beta \Psi) & =\left(2 A_{12} \alpha+A_{22} \beta\right) u^{2}+2\left(2 A_{13} \alpha+A_{23} \beta\right) u v+\left(6 A_{14} \alpha+3 A_{24} \beta\right) v^{2}
\end{aligned}
$$

whence,

$$
\frac{\partial}{\partial u}(\beta \Pi+2 \gamma \Psi)(u, v)=\frac{\partial}{\partial v}(2 \alpha \Pi+\beta \Psi)(u, v)
$$

In consequence, there is a quartic polynomial $\Theta(u, v)$ such that

$$
\begin{equation*}
\frac{\partial \Theta}{\partial u}=2 \alpha \Pi+\beta \Psi \quad \text { and } \quad \frac{\partial \Theta}{\partial v}=\beta \Pi+2 \gamma \Psi \tag{6.8}
\end{equation*}
$$

### 6.1.1 Second Conservation Law

In case there is a non-trivial solution $(\alpha, \beta, \gamma)$ of $(6.4)$, a second Conservation Law is now contemplated. Consider the time derivative of the functional

$$
\int_{-\infty}^{\infty}\left(\alpha u^{2}+\beta u v+\gamma v^{2}\right) d x
$$

and following the same calculations as in Chapter 2 for the quadratic case to obtain

$$
\begin{equation*}
\frac{d}{d t} \int_{-\infty}^{\infty}\left(\alpha u^{2}+\beta u v+\gamma v^{2}\right) d x=-\frac{d}{d t} \int \Theta(u, v) d x \tag{6.9}
\end{equation*}
$$

If the functional $\Phi(u, v)$ is defined by

$$
\begin{equation*}
\Phi(u, v):=\int_{-\infty}^{\infty}\left(\alpha u^{2}+\beta u v+\gamma v^{2}+\Theta(u, v)\right) d x \tag{6.10}
\end{equation*}
$$

then $\Phi(u, v)$ is also an invariant of the temporal evolution of smooth solutions of our system. It serves as a Hamiltonian for the given system as long as $4 \alpha \gamma-\beta^{2}$ is not 0 .

Both these conserved quantities $\Omega$ and $\Phi$ will be used to obtain a priori bounds on solutions of the system when the quadratic form $\rho$ defined by

$$
\begin{equation*}
\rho(u, v)=\alpha u^{2}+\beta u v+\gamma v^{2}, \tag{6.11}
\end{equation*}
$$

vanishes only at origin, thus when its discriminant is negative. The details however will be pursued later as a companion result of this thesis and are omitted here.

### 6.2 Local Well Posedness

The key component to proving local well-posedness is once again the Contraction Mapping Principle.

### 6.2.1 Local Well-Posedness in $L_{\infty}$ and $C_{b}$

Consider again the Banach space

$$
X_{T}:=L_{\infty}(\mathbb{R} \times[0, T]) \times L_{\infty}(\mathbb{R} \times[0, T])
$$

endowed with its usual norm ( see 3.2). To establish local existence via the Contraction Mapping principle, the domain of the operator $\mathcal{A}$ will be restricted to the usual closed ball

$$
B_{R}(0):=\left\{U \in X_{T}:\|U\|_{X_{T}} \leq R\right\}
$$

with $R>0$ and $T>0$ to be determined.
Consider $U, \tilde{U} \in B_{R}(0)$, and look for an estimate of $\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}}$. The aim is to choose specific values of $R$ and $T$ so that $\mathcal{A}$ is a contraction on $B_{R}(0)$. As the calculations are very similar as those appearing in the quadratic case presented in section 3.1.1, the details will be omitted. Using the definition of the operator $\mathcal{A}$, for $U=(u, v)$ and $\tilde{U}=(\tilde{u}, \tilde{v})$;

$$
\begin{aligned}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}}= & \sup _{x \in \mathbb{R}}\left|\int_{0}^{T} \int_{-\infty}^{\infty} K(x-y)[u(y, s)+\Pi(u, v)-\tilde{u}(y, s)-\Pi(\tilde{u} . \tilde{u})] d y d s\right| \\
& +\sup _{x \in \mathbb{R}}\left|\int_{0}^{T} \int_{-\infty}^{\infty} K(x-y)[v(y, s)+\Psi(u, v)-\tilde{v}(y, s)-\Psi(\tilde{u} . \tilde{u})] d y d s\right| \\
\leq & \sup _{x \in \mathbb{R}} \int_{0}^{T} \int_{-\infty}^{\infty}|K(x-y)|\left[| u - \tilde { u } | \left(A_{11}\left(|u|^{2}+|u||\tilde{u}|+|\tilde{u}|^{2}\right)+A_{12}|\tilde{v}|(|u|+|\tilde{u}|)\right.\right. \\
& \left.\left.+A_{13}|\tilde{v}|^{2}\right)+|v-\tilde{v}|\left(A_{12}|u|^{2}+A_{13}|u|(|v|+|\tilde{v}|)+A_{14}\left(|v|^{2}+|v||\tilde{v}|+|\tilde{v}|^{2}\right)\right)\right] d y d s \\
& +\sup _{x \in \mathbb{R}} \int_{0}^{T} \int_{-\infty}^{\infty}|K(x-y)|\left[| v - \tilde { v } | \left(A_{24}\left(|v|^{2}+|v||\tilde{v}|+|\tilde{v}|^{2}\right)+A_{23}|\tilde{u}|(|v|+|\tilde{v}|)\right.\right. \\
& \left.\left.+A_{22}|\tilde{u}|^{2}\right)+|u-\tilde{u}|\left(A_{23}|v|^{2}+A_{22}|v|(|u|+|\tilde{u}|)+A_{21}\left(|u|^{2}+|u||\tilde{u}|+|\tilde{u}|^{2}\right)\right)\right] d y d s .
\end{aligned}
$$

Furthermore, as $U, \tilde{U} \in B_{R}(0)$, it follows that the inequality can be continued, viz.

$$
\begin{equation*}
\|\mathcal{A} U-\mathcal{A} \tilde{U}\|_{X_{T}} \leq\|U-\tilde{U}\|_{X_{T}} C_{R} T, \tag{6.12}
\end{equation*}
$$

where $C_{R}$ is a constant only dependent on $R$. The choices of R and T are to be made such that $\mathcal{A}$ maps $B_{R}(0)$ into itself through and so that $\mathcal{A}$ is a contraction. In particular, one needs to have

$$
C_{R} T<1 .
$$

As before, fix the relation

$$
C_{R} T=\frac{1}{2} .
$$

Then, as before, for any $U \in B_{R}(0),\|\mathcal{A} U\|_{X_{T}} \leq R$, provided $R=2\left\|U_{0}\right\|_{\infty}$.
This implies that

$$
T=\frac{1}{2 C_{R}} .
$$

The Contraction Mapping Theorem applies with this values of R and T to the integral equation (2.3). The result is existence of a unique $U$ that is solution of the given system when $\Pi$ and $\Psi$ are cubics.

### 6.2.2 Local Existence in Other Spaces

Lemma 1, repeated below for reader's convenience, was stated and proved in paragraph 3.1.2 for the case when $P$ and $Q$ are quadratic polynomials and contains local well-posedness results for spaces embedded in $L_{\infty}$. The proof for the cubic case is very similar thus we'll omit most of the calculations and only detail the results that are specific to the cubic polynomials.

Lemma 3. Suppose that $Z$ is a space embedded in $L_{\infty}$ such that, for any $u, v \in Z$, there exist global constants $C_{1}$ and $C_{2}$ satisfying

$$
\text { (i) }\|u v\|_{Z} \leq C_{1}\|u\|_{Z}\|v\|_{Z} \quad \text { and } \quad \text { (ii) } \quad\left\|\int_{-\infty}^{\infty} K(x-y) u(y) d y\right\|_{Z} \leq C_{2}\|u\|_{Z}
$$

Then there is a unique local solution for the initial-value problem for (1.1) in the space $X_{T}=$ $C(0, T ; Z \times Z)$.

Proof. Using the same Banach space $Z_{T}=C(0, T ; Z \times Z)$ with the norm

$$
\begin{equation*}
\|U-\tilde{U}\|_{Z_{T}}:=\sup _{t \in[0, T]}\|u(\cdot, t)-\tilde{u}(\cdot, t)\|_{Z}+\sup _{t \in[0, T]}\|v(\cdot, t)-\tilde{v}(\cdot, t)\|_{Z}, \tag{6.13}
\end{equation*}
$$

where $U=(u, v)$ and $\tilde{U}=(\tilde{u}, \tilde{v})$, restrict the operator $\mathcal{A}$ from (2.5) to

$$
B_{R, Z}(0):=\left\{U \in Z_{T}:\|U\|_{Z_{T}} \leq R\right\} .
$$

Consider $U, \tilde{U} \in B_{R, Z}(0)$. The aim is to choose specific values of $R$ and $T$ so that $\mathcal{A}$ is a contraction on $B_{R, Z}(0)$. To estimate $\|U-\tilde{U}\|_{Z_{T}}$ apply (ii) from (3.6), the triangle inequality and (3.6) part (i), and similarly with the $L_{\infty}$ case from the preceding section, yields

$$
\|\mathcal{A} U-\mathcal{A} V\|_{Z_{T}} \leq T C_{1} C_{2} C_{R}\|U-V\|_{Z_{T}}
$$

Choosing $R=2\left\|U_{0}\right\|_{Z}$ and $T$, for example, so that

$$
T=\frac{1}{2 C_{1} C_{2}^{2} C_{R}},
$$

the Contraction Mapping Theorem can be applied to the integral equation (2.3) derived from the initial system and the result is a unique $U$ that is solution of the given system. Moreover, this solution found above is seen to be continuously dependent on the initial data.

From the regularity perspective, the smoothing associated with the temporal derivative is analogous with the $L_{\infty}$ case presented in the previous section, specifically for $U \in H^{k} \times H^{k}$, $U_{t}$ is spatially smoother and it lies in $C\left(0, \infty, ; H^{k+1} \times H^{k+1}\right)$. As for the spatial regularity, the system solution has at most the same regularity as the initial data.

This Lemma provides the local well-posedness on $H^{1}, H^{2}$ and in general in $H^{k}$, for $k \geq 1$.

### 6.2.3 Smoothness and Regularity

The argument presented in Chapter 3 is independent of the polynomial degree, therefore the result applies to the cubic case as well. The strong regularity results for the flow map, the temporal smoothing and the spacial regularity being at most as good as the initial data are still valid.

## CHAPTER 7

## CONCLUSION

We have considered coupled systems

$$
\left\{\begin{array}{l}
u_{t}+u_{x}-u_{x x t}+P(u, v)_{x}=0  \tag{7.1}\\
v_{t}+v_{x}-v_{x x t}+Q(u, v)_{x}=0
\end{array}\right.
$$

of two evolution equations of generalized BBM-type, posed for $x \in \mathbb{R}$ and $t \geq 0$ and with specified initial data at $t=0, v i z$.

$$
\begin{equation*}
u(x, 0)=u_{0}(x), \quad v(x, 0)=v_{0}(x), \quad x \in \mathbb{R} . \tag{7.2}
\end{equation*}
$$

These systems were shown to be locally well-posed in the $L_{2}$-based Sobolev spaces $H^{s}(\mathbb{R}) \times$ $H^{s}(\mathbb{R})$ for $s \geq 0$ when $P$ and $Q$ are homogeneous quadratic polynomials. Moreover, the solution $(u, v)$ depends continuously on the initial data $\left(u_{0}, v_{0}\right)$. The spacial regularity of the solution $(u, v)$ is exactly as good as the initial data's. However, there is some smoothing in the temporal variable.

Furthermore, for those polynomial coefficients $A, B, \cdots, F$ such that the system

$$
\left\{\begin{aligned}
2 B a+(E-2 A) b-4 D c & =0 \\
4 C a+(2 F-B) b-2 E c & =0
\end{aligned}\right.
$$

has solutions $a, b$ and $c$ satisfying $4 a c>b^{2}$, or, what is the same, the matrix $\left(\begin{array}{cc}2 a & b \\ b & 2 c\end{array}\right)$ is positive definite, then the well-posedeness is global in time.

Similar results are sketched for the homogeneous cubic polynomials.
In future work, we plan to consider the periodic initial-value problem and the more practical case of non-homogeneous boundary conditions. We will also be interested in the question of global well-posedness in case the nonlinearities are not homogeneous.

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## EDUCATION

Ph.D. in Mathematics - University of Illinois at Chicago 2014
M.S. in Mathematics - University of Illinois at Chicago 2007
B.S. in Computer Science - West University of Timisoara, Romania 1998
B.S. in Mathematics - West University of Timisoara, Romania 1997

## WORK EXPERIENCE

- 2001 - Present Allstate Insurance Inc. (Northbrook, IL)

Information Analyst/Web Analyst/Technical Lead/Consultant (2004-Present)
Partner with functional leads on strategic planning for business intelligence and analytics for Allstate's websites and web-based applications both customer and agent-facing. Design and implement strategies working cooperatively with Business Services, Technical Services, Software Services to facilitate and coordinate the delivery of reporting and analytics across the organization. Create standard key metrics for business and technology partners, and create the corresponding presentation objects - dashboards with visual interactive analysis. Work diligently to ensure accuracy, consistency and integration between different online channels' data as well as backend data. Provide guidance on outcome analysis. Mentor and train business partners on how to best leverage data, to discover data trends and make data-driven business
decisions for site and application's optimization as well as for increase of the customer experience and satisfaction. Evaluated tools and technology solutions to determine optimum usage by end users. Managed technical aspects of contracts with vendors during negotiations. Worked with vendors of the web analytics applications on any data issues, technical implementation, tickets and other admin items. Ensured data accuracy aligns with contractual commitment and resolve any data integrity issues.

Web Developer/Professional (2001-04) - Provided defect maintenance support and made code changes for release level work for Allstate's quoting application; participated in the eXtreme Programming Pilot; provided mentoring and training during the application's transition to Allstate Northern Ireland.

- 2000-2001 Sheldon I. Dorenfest \& Associates (Chicago, IL)

Software Developer - Maintained and designed tools used for data collection, quality control and report production; constructed the data in the final format of the product to deliver to clients.

## CONFERENCES AND WORKSHOPS ATTENDED

NetConnect Annual Users' Conference in Orlando, FL - March 2004
e-Metrics Summit, Santa Barbara, CA - May 2005
e-Metrics Summit, San Francisco, CA - May 2007
Visual Sciences User's Conference, San Diego, CA - May 2007
Omniture Summit, Salt Lake City, UT - March 2009

The Workshop on "Mathematical Biology and Numerical Analysis", University of Georgia - August 2009

66th Midwest Partial Differential Equations Seminar at University of Illinois at Chicago

- November 2010

69th Midwest Partial Differential Equations Seminar at UIC - April 2012

Business Analytics Innovation Summit, Chicago - May 2012

Business Analytics Innovation Summit, Chicago - May 2013 - keynote speaker "Claims of Analytics and Analytics of Claims"

## HONORS AND MEMBERSHIPS

ComEd Scholar for the 2009-2010 Academic Year
Member, American Mathematical Society (AMS)
Member, Association for Women in Mathematics (AWM)

Member, Web Analytics Association (WAA)

