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# A Quadrature Eulerian-Lagrangian WENO Scheme for Reservoir Simulation 

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# A Quadrature Eulerian-Lagrangian WENO Scheme for Reservoir Simulation 

by

Jamie Micayla Pool, B.S.

## DISSERTATION

Presented to the Faculty of the Graduate School of The University of Texas at Austin
in Partial Fulfillment
of the Requirements
for the Degree of

DOCTOR OF PHILOSOPHY

In loving memory of Bill and Mary Ann Pool, who inspire me to live with a sense of purpose and wonder.

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# A Quadrature Eulerian-Lagrangian WENO Scheme for Reservoir Simulation 

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This dissertations focuses on solving the advection problem with the motivation of simulating transport in porous media. A quadrature based Eulerian-Lagrangian scheme is developed to solve the nonlinear advection problem in multiple spatial dimensions. The schemes combines the ideas of Lagrangian traceline methods with high order WENO reconstructions to compute the mass that flows into a given cell over a time step. These schemes are important since they have a relaxed CFL constraint, and can be run in parallel.

In this thesis we provide two improvements to Eulerian-Lagrangian schemes. To do this an integration based WENO (IWENO) interpolation technique is derived by reconstructing the primitive function and differentiating. This technique gives a high order reconstruction of the mass at an arbitrary point.

This WENO scheme is used to solve the linear advection problem. A scheme is derived by backwards tracing of quadrature points located on mesh elements. The mass at these tracepoints is used to compute the mass in the trace region, without resolving its boundary. This process defines a high order quadrature Eulerian-Lagrangian WENO (QEL-WENO) scheme that solves the multi-dimensional problem without the need for a spatial splitting technique.

The second improvement is for solving the nonlinear advection problem using an approximate velocity field. The velocity field is used to transport mass in the manner of a standard Eulerian-Lagrangian scheme. Then a flux correction is applied to compute the flow across the tracelines. The contribution is to use a variation of the IWENO technique to reduce the stencil size of this computation.

Numerical results are presented demonstrating the capabilities of the scheme. An application to two-phase flow in porous media is provided.

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## Chapter 1

## Introduction

We consider the basic physical phenomena of fluid transport in a velocity field. This type of problem naturally arises in a diverse variety of fields such as continuum mechanics and thermodynamics. While the range of application areas of this dissertation are broad, we will restrict our focus to problems arising in porous media simulation. The primary goal is to look at transport of a miscible tracer and of two phases, oil and water, in a petroleum reservoir system.

For the reservoir problem, the porous media are natural materials such as rocks, soils and sands. These media are characterized by porosity, which measures the void space in the material, and permeability, which measures the ability of a fluid to flow through the void network. Reservoir systems have large spatial scales, with dimensions in terms of kilometers and potentially large time scales. For instance, the geological process of $\mathrm{CO}_{2}$ sequestration lasts for hundreds of years. For this reason, computational simulation of these phenomena are run on modern parallel supercomputers, and hence we seek to develop numerical algorithms that are suitable in this setting.

### 1.1 Problem

One of the simplest cases for transport in a porous media is the situation with one substance, the tracer, being transported in another. Conservation of the tracer mass gives

$$
\begin{equation*}
\frac{\partial(\phi c)}{\partial t}+\nabla \cdot(c u-D \nabla c)=q_{c}(c) \tag{1.1}
\end{equation*}
$$

where the unknown tracer concentration $c(x, t)$ is transported in a Darcy velocity field $u(x, t)$. The parameter $\phi(x)$ is the porosity of the medium, $D$ is a tensor describing the molecular diffusion and macroscopic dispersion of the particles based on the pore scale geometry of the medium and $q_{c}(c)$ is an external source or sink function, possibly describing a well in a reservoir.

This problem contains a parabolic part describing mixing through diffusion; this process smooths data. Equations of parabolic type have an infinite speed of information propagation, so implicit numerical methods must be used. There is also an advective or hyperbolic part, describing particle movement in the velocity field. In the presence of a smooth velocity field and initial condition, this equations permits the development of discontinuous solutions, or shocks. Advection equations propagate information at finite speed and explicit solvers may be used. The differences between the advection and diffusion processes, makes it natural to use a splitting technique to solve for them individually. The first step to solve (1.1) using a splitting technique is to solve the advection equation

$$
\begin{equation*}
\frac{\partial(\phi c)}{\partial t}+\nabla \cdot(c u)=q_{c}(c) \tag{1.2}
\end{equation*}
$$

Its solution is then used as an initial condition to solve the diffusion equation

$$
\begin{equation*}
\frac{\partial(\phi c)}{\partial t}-\nabla \cdot(D \nabla c)=0 \tag{1.3}
\end{equation*}
$$

For the purpose of this dissertation, we will focus on the advection equation. However, there are many numerical solvers for the diffusion equation, including ones based on finite element [12], mixed [13, 33] and discontinuous Galerkin (DG) [2] methods.

The various types of methods to solve the advection problem can be categorized by the type of mesh that is used. Eulerian methods, or fixed grid methods, look to accurately reconstruct the flux across the boundaries of the elements. These methods have had a good deal of success as they are accurate while minimizing numerical dispersion errors. Examples of these schemes include weighted essentially non-oscillatory (WENO) [41, 42, 49, 54, 60] Godnov [15, 27, 50] and DG [7, 65]. These methods all must satisfy the Courant Friedrichs Levy (CFL) condition which states that the length of a time step size is limited by the mesh spacing divided by the maximal magnitude of the velocity field. This condition can require a large number of time steps to be taken, leading to increased computation time and the accumulation of numerical dispersion errors over time.

A Lagrangian method, or a moving grid method, attempts to alleviate the CFL constraint by allowing the grid points to move with the velocity field over time. These types of methods are challenging to extend to higher dimensions where multiple shock interaction provides a challenge to mesh adaptabil-
ity. An example of this type of method is given in [73] by Stockie, Mackenzie and Russell. In order to concentrate grid points in optimal locations, they employ a technique that couples a moving mesh equation with the physical equation.

The final type of method is Eulerian-Lagrangian, or Semi-Lagrangian, where at every time step there is a projection between and Lagrangian and Eulerian frameworks. This type of method was first made viable by Douglas and Russell in 1982 with their development of the Modified Method of Characteristics (MMOC) [32]. Further discussion on the history of these types of methods is provided in Chapter 3.

### 1.2 Dissertation Outline

The dissertation is organized as follows. In Chapter 2, we give the background of the development of WENO schemes. Followed by a detailed development of a new integration based WENO reconstruction technique, with numerical results illustrating formal convergence rates.

Chapter 3 provides a review of characteristic methods for advection problems. We then define a new Eulerian-Lagrangian scheme for the linear advection problem which has the capability of solving the multi-dimensional problem without the aid of a spatial splitting technique.

In Chapters 4 and 5, we define a new algorithm for solving the nonlinear advection problem. This algorithm has a compact local stencil for com-
puting the flux integral. Numerical results demonstrating the capabilities of the scheme are presented.

An application to two phase flow is presented in Chapter 6.
This dissertation is finalized in Chapter 7 by giving some conclusions and suggesting directions for future research.

## Chapter 2

## Integration-based WENO (IWENO)

### 2.1 History of WENO Schemes

Essentially Non-Oscillatory (ENO) schemes were introduced in 1987 in a paper by Harten, Enquist, Osher and Chakravarthy [41] as a way to develop a shock capturing method to solve hyperbolic conservation laws. The design of these schemes involves construction of multiple polynomials over select stencils. These polynomials have average values on grid intervals that agree with the cell averaged data coming from the scheme at the given time step. The ENO scheme adaptively chooses the polynomial on a stencil that avoids shocks in the data for use in the forward time stepping scheme.

Weighted Essentially Non-Oscillatory (WENO) schemes were later developed in $[49,54]$ using a convex combination of all the candidate stencils as opposed to the single stencil used in ENO schemes. These schemes achieve a higher order of accuracy in smooth regions, while in regions of discontinuity weight towards the same stencil as the ENO scheme would adaptively choose. In addition WENO schemes are computationally preferable to ENO schemes since they rely on a linear combination of data, instead of the logic statements in ENO schemes. This decreases the total computation time to solve
the problem.
There is an extensive literature detailing the many uses of ENO/WENO schemes for both conservation laws and Hamilton-Jacobi equations. Topics discussed include the use of ENO/WENO schemes in higher dimensions, on triangular grids and theory on stencil choice, stability, and choice of basis functions which can be wavelet, trigonometric or polynomial. In 1997 Shu wrote a set of self contained lecture notes [71] which serve as a more in-depth overview of ENO and WENO schemes.

### 2.1.1 Review of WENO Interpolation

Our interest in WENO schemes is in the interpolation process. To this end we give a brief overview of the typical WENO reconstruction scheme. We are interested in computing $u\left(x^{*}, t^{n}\right)$, the value of $u$ at the point $x^{*}$ located in the grid cell $I_{0}$, given that we have the average value of $u$ on the cells. We are given $n$, the degree of the lower order polynomial used in the approximation. We take a computational grid with points $x_{-n}<\cdots<x_{n+1}$, with corresponding intervals $I_{i}=\left[x_{i}, x_{i+1}\right]$ for $i=-n, \ldots, n$. We further assume that we have cell averaged data for a function $u(x)$, given by

$$
\begin{equation*}
\bar{u}_{i}=\frac{1}{x_{i+1}-x_{i}} \int_{I_{i}} u(x) d x, \quad i=-n, \ldots, n . \tag{2.1}
\end{equation*}
$$

In order to approximate the value of $u\left(x^{*}\right)$ we define $n+1$ polynomials, $p_{0}, \ldots, p_{n}$ of degree $n$ so that they agree with the cell averaged data, i.e.,

$$
\begin{equation*}
\int_{I_{i}} p_{k} d x=\left(x_{i+1}-x_{i}\right) \bar{u}_{i}, \quad i=-n+k, \ldots, k \tag{2.2}
\end{equation*}
$$



Figure 2.1: WENO Reconstruction Polynomials. Illustration of WENO reconstruction polynomials $p_{k}(x)$ which are used to approximate the value of $u$ at a point $x^{*} \in I_{0}$.
as shown in Figure 2.1. Taylor series analysis shows that for smooth $u(x)$ the approximation $p_{k}\left(x^{*}\right)$ is an approximation of $u\left(x^{*}\right)$ of order $n+1$, i.e.,

$$
\begin{equation*}
p_{k}\left(x^{*}\right)-u\left(x^{*}\right)=O\left(h^{n+1}\right) \tag{2.3}
\end{equation*}
$$

where $h=\max \left|x_{i+1}-x_{i}\right|$. There is also a polynomial of degree $2 n$ defined with cell averages agreeing with all the candidate cells that is a higher order approximation to $u\left(x^{*}\right)$. WENO schemes attempt to find weights $c_{k}$ such that the linear combination of the polynomials and weights achieves this higher order

$$
\begin{equation*}
\sum_{k=0}^{n} c_{k} p_{k}\left(x^{*}\right)-u\left(x^{*}\right)=O\left(h^{2 n+1}\right) \tag{2.4}
\end{equation*}
$$

However, there are situations where such weights cannot be found. A well known example of this is found in attempting to derive a 3rd order central WENO scheme. Here the reconstruction point $x^{*}$ is the center of the interval $I_{0}$. As seen in Figure 2.2, the linear reconstructions agree at this point.


Figure 2.2: The Third Order WENO Midpoint Reconstruction. The linear interpolation polynomials $p_{k}$ have the same value at the midpoint.

Since any linear combination of these values will yield the same value at the midpoint, there does not exist coefficients for the linear combination to agree with the quadratic interpolation. Several different techniques have been used to get around this difficulty. Levy, Puppo and Russo [53] use the quadratic polynomial in addition to the linear interpolations in order to achieve this higher order approximation in smooth regions. Huang, Arbogast and Hung [47] define a re-averaging technique which converts cell averages on one grid to cell averages on another. The re-averaging puts the midpoint on a grid boundary where linear weights are guaranteed to exist. However, this scheme contains a large implicit stencil. This technique will be further addressed in Chapter 4.

For the transport schemes we will develop later, we need a high order approximation at any point $x^{*}$ in our mesh. We also want to avoid using a larger stencil, since this can lead to numerical diffusion in the transport
scheme. The rest of this chapter will derive an Integration-based WENO (IWENO) scheme which addresses this problem.

### 2.2 IWENO Scheme

We are motived by the work of Carlini, Ferretti and Russo in [19], where they show that for problems involving nodal interpolation, as opposed to cell average values, high order reconstructions can be found for any point in the domain. Furthermore, the weighting functions are simply polynomials of the position $x^{*}$.

We define polynomials $P_{k}$ of degree $n+1$ that interpolate nodal values of the primitive function,

$$
\begin{equation*}
U(x)=\int_{x_{0}}^{x} u(\xi) d \xi \tag{2.5}
\end{equation*}
$$

This interpolation is shown in Figure 2.3. Results from [19] state that there are coefficient polynomials $c_{k}(x)$ of degree $n-1$ such that for $x^{*} \in I_{0}$ there is a high order reconstruction of the primitive function, i.e.,

$$
\begin{array}{ll}
U(x)-P_{k}(x)=O\left(h^{r+1}\right), & \forall x \in I_{0} \\
U(x)-\sum_{k=0}^{n} c_{k}(x) P_{k}(x)=O\left(h^{2 r}\right), & \forall x \in I_{0} \\
\sum_{k=0}^{n} c_{k}(x)=1 \tag{2.8}
\end{array}
$$

Here $r+1=(n+1)+1$ is the order of polynomial interpolation of the primitive function.


Figure 2.3: IWENO Reconstruction Polynomials. Illustration of IWENO reconstruction polynomials $P_{k}(x)$ which are used to approximate the value of the primitive function $U(x)$ at a point $x^{*} \in I_{0}$.

The Lagrange interpolation satisfies

$$
\begin{equation*}
U(x)=P_{k}(x)+\frac{U^{(r+1)} \gamma}{(r+1)!} \prod_{j=k-r+1}^{k+1}\left(x-x_{j}\right) \tag{2.9}
\end{equation*}
$$

for some $\gamma(x) \in\left(x_{-k-r+1}, x_{k+1}\right)$. If $u(x) \in C\left(\left[x_{-k+r+1}, x_{k+1}\right]\right)$ then for any $x_{-k-r+1} \leq \alpha<\beta \leq x_{k+1}$

$$
\begin{equation*}
\frac{1}{\beta-\alpha} \int_{\alpha}^{\beta} u(x) d x=\frac{P_{k}(\beta)-P_{k}(\alpha)}{\beta-\alpha}+O\left(h^{r}\right) . \tag{2.10}
\end{equation*}
$$

Which comes from (2.9) and the mean value theorem. The limit as $\alpha \rightarrow \beta$ gives

$$
\begin{equation*}
u(\alpha)=P_{k}^{\prime}(\alpha)+O\left(h^{r}\right) \tag{2.11}
\end{equation*}
$$



Figure 2.4: IWENO Nodal Values. Cell average data is used to compute the value of the primitive function at grid nodes.

Using this relation the differentiation of equations (2.6) - (2.8) gives

$$
\begin{array}{ll}
u(x)-p_{k}(x)=O\left(h^{r}\right), & \forall x \in I \\
u(x)-\sum_{k=0}^{n}\left(c_{k}^{\prime}(x) P_{k}(x)+c_{k}(x) p_{k}(x)\right)=O\left(h^{2 r-1}\right), & \forall x \in I \\
\sum_{k=0}^{n} c_{k}^{\prime}(x)=0, & \tag{2.14}
\end{array}
$$

where $p_{k}(x)$ is the derivative of $P_{k}(x)$. Showing the reconstruction of $u(x)$,

$$
\begin{equation*}
\sum_{k=0}^{n}\left(c_{k}^{\prime}(x) P_{k}(x)+c_{k}(x) p_{k}(x)\right) \tag{2.15}
\end{equation*}
$$

is of order $h^{2 r-1}$. The sum (2.15) is the IWENO reconstruction using linear weighting coefficients.

### 2.2.1 Implementation Details

To implement this scheme we first need to convert our cell averages to nodal values by using the definition of the primitive function. An illustration
of this process is seen in Figure 2.4. The nodal values of the primative function are sums of the cell averages scaled by their corresponding element size, so

$$
\begin{align*}
& U\left(x_{i}\right)=-\int_{x_{i}}^{x_{0}} u(\xi) d \xi=-\prod_{j=i}^{-1}\left(x_{j+1}-x_{j}\right) \bar{u}_{j}, \quad i<0  \tag{2.16}\\
& U\left(x_{0}\right)=\int_{x_{0}}^{x_{0}} u(\xi) d \xi=0  \tag{2.17}\\
& U\left(x_{i}\right)=\int_{x_{0}}^{x_{i}} u(\xi) d \xi=\prod_{j=0}^{i-1}\left(x_{j+1}-x_{j}\right) \bar{u}_{j}, \quad i>0 . \tag{2.18}
\end{align*}
$$

In order to compute the polynomials as shown in Figure 2.3, a polynomial interpolation can be evaluated using Lagrange interpolation polynomials. The polynomial $P_{k}$ is defined on the stencil $x_{-n+k}, \ldots, x_{1+k}$ and is reconstructed as

$$
\begin{align*}
& P_{k}(x)=\sum_{i=-n+k}^{1+k} U\left(x_{i}\right) \mathcal{L}_{(k, i)}(x)  \tag{2.19}\\
& \text { where } \quad \mathcal{L}_{(k, i)}(x)=\prod_{\substack{j=-n+k \\
j \neq i}}^{j=1+k} \frac{\left(x-x_{j}\right)}{\left(x_{i}-x_{j}\right)} \tag{2.20}
\end{align*}
$$

The Lagrange interpolation polynomials, $\mathcal{L}_{(k, i)}$, have values at nodes given by the Kroneker delta, i.e.,

$$
\mathcal{L}_{(k, i)}\left(x_{j}\right)=\delta_{i j}=\left\{\begin{array}{ll}
1 & i=j  \tag{2.21}\\
0 & i \neq j
\end{array} .\right.
$$

Differentiation of polynomials $P_{k}$ gives the polynomials $p_{k}$

$$
\begin{align*}
& p_{k}(x)=\frac{d P_{k}}{d x}=\sum_{i=-n+k}^{1+k} U\left(x_{i}\right) \frac{d \mathcal{L}_{k, i}}{d x},  \tag{2.22}\\
& \text { where } \frac{d \mathcal{L}_{k, i}}{d x}=\sum_{\substack{m=-n+k \\
m \neq i}}^{1+k}\left(\frac{1}{\left(x-x_{m}\right)} \prod_{\substack{j=-n+k \\
j \neq i}}^{1+k} \frac{\left(x-x_{j}\right)}{\left(x_{i}-x_{j}\right)}\right) . \tag{2.23}
\end{align*}
$$

```
\(\overline{\text { Algorithm } 1 \text { for the computation of the primitive function and its derivative }}\) on an arbitrary grid.
\[
P=0
\]
\[
d P=0
\]
\[
\text { for } i=-n+k, i<2+k \text { do }
\]
\[
\operatorname{Lag}=U_{i}
\]
\[
\text { Denom }=1.0
\]
\[
\text { for } j=0, j<4, j \neq i \text { do }
\]
\[
\operatorname{Lag} *=\left(x-x_{j}\right)
\]
\[
\text { Denom } *=\left(x_{i}-x_{j}\right)
\]
end for
\[
P+=L a g / D e n o m
\]
\[
d \operatorname{Lag}=0
\]
\[
\text { for } m=-n+k, m<k+2, m \neq i \text { do }
\]
\[
d L=U_{i}
\]
\[
\text { for } j=-n+k, j<k+2, j \neq m, j \neq i \text { do }
\]
\[
d L *=\left(x-x_{j}\right)
\]
end for
\[
d \operatorname{Lag}+=d L
\]
end for
\[
d P+=d L a g / D e n o m
\]
```

end for
Computationally, an algorithm for this process is given by Algorithm

1. However, this is sensitive to rounding errors, when $x$ is near a grid point.

For irregular or adaptive grids using a Newton divided difference algorithm to compute these polynomials minimizes this problem. For uniform meshes computing the coefficients of the polynomials ahead of time reduces this error. The $3^{\text {rd }}$ order reconstruction of $P_{k}(x)$ and $p_{k}(x)$ on a grid with uniform spacing $h$ are

$$
\begin{align*}
& P_{k}\left(x_{l o c, k}\right)=\left(2 U_{-1+k}\right.+x_{l o c, k}\left(-3 U_{-1+k}+4 U_{k}-U_{1+k}\right) \\
&\left.+x_{l o c, k}^{2}\left(U_{-1+k}-2 U_{k}+U_{1+k}\right)\right) / 2  \tag{2.24}\\
& p_{k}\left(x_{l o c, k}\right)=\left(1 / 2\left(3 \bar{u}_{-1+k}-\bar{u}_{k}\right)-\left(\bar{u}_{-1+k}-\bar{u}_{k}\right) x_{l o c, k}\right) / h,  \tag{2.25}\\
& \text { where } \quad x_{l o c, k}=(1-k)+\alpha \quad \text { and } \quad k=0,1
\end{align*}
$$

and the $5^{\text {th }}$ order scheme has polynomials

$$
\begin{align*}
& P_{k}\left(x_{l o c, k}\right)=\left(U_{-2+k}\right.+x_{l o c, k}\left(-11 U_{-2+k}+18 U_{-1+k}-9 U_{k}+2.0 U_{k+1}\right) \\
&+3 x_{l o c, k}^{2}\left(2 U_{-2+k}-5 U_{-1+k}+4 U_{k}-U_{k+1}\right) \\
&\left.+x_{l o c, k}^{3}\left(-U_{-2+k}+3 U_{-1+k}-3 U_{k}+U_{k+1}\right)\right) / 6,  \tag{2.26}\\
& p_{k}\left(x_{l o c, k}\right)=\left(\begin{array}{rl}
(11 \quad & \left.\bar{u}_{-2+k}-7 \bar{u}_{-1+k}+2 \bar{u}_{k}\right) / 6 \\
& x_{l o c, k}\left(2 \bar{u}_{-2+k}-3 \bar{u}_{-1+k}+\bar{u}_{k}\right) \\
& \left.-x_{l o c, k}^{2}\left(-\bar{u}_{-2+k}+2 \bar{u}_{-1+k}-\bar{u}_{k}\right) / 2\right) / h
\end{array}\right. \\
& \text { where } \quad x_{l o c, k}=(2-k)+\alpha \quad \text { and } \quad k=0,1,2 .
\end{align*}
$$

These polynomials are evaluated at the local reference position $x_{l o c, k}$ defined through a scaling factor $\alpha \in[0,1]$ describing the relative position of $x^{*} \in I_{0}$ namely $\alpha=\left(x^{*}-x_{0}\right) / h$.

Carlini, Ferretti and Russo [19] give the algorithm to compute the linear weighting coefficients. To do this first set

$$
\begin{equation*}
W_{k}^{m}=\prod_{j=-n}^{-n+k-1}\left(x_{m+1}-x_{j}\right) \prod_{j=2+k}^{n+1}\left(x_{j}-x_{m+1}\right) \tag{2.28}
\end{equation*}
$$

for $k, m=0, \ldots, n-1$, then solve the lower triangular matrix problem

$$
\left[\begin{array}{cccc}
W_{0}^{0} & 0 & \cdots & 0  \tag{2.29}\\
W_{1}^{0} & W_{1}^{1} & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
W_{n-1}^{0} & W_{n-1}^{1} & \cdots & W_{n-1}^{n-1}
\end{array}\right]\left[\begin{array}{c}
\hat{c}_{0} \\
\hat{c}_{1} \\
\vdots \\
\hat{c}_{n-1}
\end{array}\right]=\left[\begin{array}{c}
1 \\
1 \\
\vdots \\
1
\end{array}\right] .
$$

On a uniform mesh the solution is

$$
\begin{equation*}
\hat{c}_{k}=\binom{n}{k} \frac{n!}{(2 n+1)!} h^{n} . \tag{2.30}
\end{equation*}
$$

Finally, the linear weights are computed from $W_{k}^{m}$ as

$$
\begin{align*}
& c_{k}(x)=\hat{c}_{k} \prod_{j=-n}^{-n+k-1}\left(x-x_{j}\right) \prod_{j=2+l}^{2 n-1}\left(x_{j}-x\right) \quad l=0,1, \ldots, n-1  \tag{2.31}\\
& c_{n}(x)=1-\sum_{l=0}^{n-1} c_{k}(x) . \tag{2.32}
\end{align*}
$$

For a scheme with uniform mesh size we compute these terms directly. For the $3^{\text {rd }}$ order scheme the coefficients are

$$
\begin{array}{ll}
c_{0}(x(\alpha))=(2-\alpha) / 3, & c_{0}^{\prime}(x(\alpha))=-1 /(3 h), \\
c_{1}(x(\alpha))=(1+\alpha) / 3, & c_{1}^{\prime}(x(\alpha))=1 /(3 h) . \tag{2.34}
\end{array}
$$

The $5^{\text {th }}$ order scheme has coefficients

$$
\begin{array}{ll}
c_{0}(x(\alpha))=\left(6-5 \alpha+\alpha^{2}\right) / 20, & c_{0}^{\prime}(x(\alpha))=(-5+2 \alpha) /(20 h), \\
c_{1}(x(\alpha))=\left(6+\alpha-\alpha^{2}\right) / 10, & c_{1}^{\prime}(x(\alpha))=(1-2 \alpha) /(10 h), \\
c_{2}(x(\alpha))=\left(2+3 \alpha+\alpha^{2}\right) / 20, & c_{2}^{\prime}(x(\alpha))=(3+2 \alpha) /(20 h) . \tag{2.37}
\end{array}
$$

The IWENO scheme in 1D using linear coefficients on a uniform grid is implemented by solving equation (2.15) using polynomials given in equation (2.24) - (2.27) and coefficients listed in equations (2.33) - (2.37).

### 2.2.2 Rounding Errors

In this section we demonstrate the performance of the IWENO algorithm. We will first show that the algorithm is sensitive to rounding error accumulation. This can be overcome by adjusting the precision of the code, we will outline these changes, and compare numerical results.

In order to test IWENO we use the $5^{\text {th }}$ order reconstruction polynomials. The point $x^{*}$ and relative location $\alpha \in[0,1]$ are given. These parameters define the center cell

$$
\begin{equation*}
I_{0}=\left[x^{*}-\alpha \Delta x, x^{*}+(1-\alpha) \Delta x\right] \tag{2.38}
\end{equation*}
$$

where $\Delta x$ is the size of the interval, and we assume uniform spacing. We perform successive grid refinements $\Delta x=\Delta x / 2$ and analyze the errors.

The function reconstructed is $u(x)=\sin (\pi x)$. The reconstruction is computed at the point $x^{*}=0.716$ with relative location $\alpha=0.334$. This test is run with both double and long double precision. The results are in Table 2.1. Double precision should exhibit precision at the $16^{\text {th }}$ digit. However the IWENO scheme uses a process of differentiation, where we scale out by a factor of $\Delta x$, this reduces the total number of digits we see for the computation.

For problems where high precision is necessary, using the long double

| $\Delta x^{-1}$ | Error | Order |
| :---: | :---: | :---: |
| 10 | $1.49 \mathrm{E}-5$ |  |
| 20 | $4.52 \mathrm{E}-7$ | 5.04 |
| 40 | $1.39 \mathrm{E}-8$ | 5.03 |
| 80 | $4.29 \mathrm{E}-10$ | 5.01 |
| 160 | $1.33 \mathrm{E}-11$ | 5.01 |
| 320 | $4.20 \mathrm{E}-13$ | 4.99 |
| 640 | $4.50 \mathrm{E}-14$ | 3.22 |
| 1280 | $8.21 \mathrm{E}-15$ | 2.45 |
| 2560 | $1.13 \mathrm{E}-14$ | -0.46 |
| 5120 | $5.53 \mathrm{E}-13$ | -5.61 |$\quad$| $\Delta x^{-1}$ | Error | Order |  |
| :---: | :---: | :---: | :---: |
|  |  | $1.49 \mathrm{E}-5$ |  |
| 40 | $4.52 \mathrm{E}-7$ | 5.04 |  |
| 80 | $4.38 \mathrm{E}-8$ | 5.03 |  |
| 160 | $1.33 \mathrm{E}-10$ | 5.01 |  |
| 320 | $4.16 \mathrm{E}-13$ | 5.01 |  |
| 5.00 |  |  |  |
| 640 | $1.28 \mathrm{E}-14$ | 5.02 |  |
| 1280 | $2.65 \mathrm{E}-16$ | 5.59 |  |
| 2560 | $2.00 \mathrm{E}-16$ | 0.41 |  |
| 5120 | $1.28 \mathrm{E}-16$ | 0.63 |  |

Table 2.1: Rounding Errors. The error in computing $\sin (0.716 \pi)$ using the IWENO scheme. The results on the left use double precision, on the right is long double precision.
precision adds the memory necessary. This is implemented in the scheme by changing the precision of variables $p_{k}$ and $P_{k}$ to be of type long double, however to minimize total storage memory the grid variable $x, \alpha, c(x)$, and $c^{\prime}(x)$ are of type double. The numerical results provided were computed in a code that is implemented in $\mathrm{C}++$. Here we note the use of the math library <cmath> since it overloads mathematical operations for variable type long double.

### 2.2.3 Non-Linear Weighting

The IWENO reconstruction detailed previously gives a high order accurate reconstruction technique for smooth data. However when the initial data is not smooth it is preferable to use a polynomial interpolation whose stencil does not contain the shock in the data. WENO schemes adjust the weights of the linear coefficients by an appropriate choice of smoothness indicator, IS.

The non-linear weighting polynomials, $w_{k}(x)$, are defined as a function of these indicators and the linear polynomials $c_{k}(x)$ as

$$
\begin{equation*}
w_{k}=\frac{\sigma_{k}}{\sum_{i=0}^{n} \sigma_{i}} \quad \text { with } \quad \sigma_{i}=\frac{c_{i}}{\left(\epsilon+\mathrm{IS}_{i}\right)^{p}} \quad \text { for } \quad k=0, \ldots, n \tag{2.39}
\end{equation*}
$$

Here $p \geq 1$ is an integer and $0<\epsilon \ll 1$ is a small number to avoid dividing by 0 . The smoothness indicators $I S_{i}$ are actually in some sense 'roughness indicators' since larger values should indicate that there is a shock in the $i^{\text {th }}$ stencil.

The reconstruction approximation defining the IWENO scheme is

$$
\begin{equation*}
u(x) \approx \sum_{k=0}^{n}\left(w_{k}^{\prime}(x) P_{k}(x)+w_{k}(x) p_{k}(x)\right) \tag{2.40}
\end{equation*}
$$

The linear coefficients, $c_{k}$ are choosen such that scheme achieves a high order of accuracy. We will next show that as long as $c_{k}-w_{k}=O\left(h^{r-1}\right)$, we maintain the accuracy presented in equations (2.6) - (2.14).

Theorem 2.2.1. Let $c_{k}(x), p_{k}(x)$, and $P_{k}(x)$ satisfy equations (2.6) - (2.14) and $w_{k}(x)$ be such that $c_{k}-w_{k}=O\left(h^{r-1}\right)$, then $U(x)-\sum w_{k} P_{k}=O\left(h^{2 r}\right)$ and $u(x)-\sum\left(w_{k} p_{k}+w_{k}^{\prime} P_{k}\right)=O\left(h^{2 r-1}\right)$.

Proof. First we consider the primitive function

$$
\begin{aligned}
U(x)- & \sum w_{k} P_{k} \\
& =\left(U(x)-\sum c_{k} P_{k}\right)+\sum\left(c_{k}-w_{k}\right) P_{k} \\
& =O\left(h^{2 r}\right)+\sum\left(c_{k}-w_{k}\right)\left(P_{k}-U\right) \quad \text { since } \quad U \sum\left(c_{k}-w_{k}\right)=U(1-1)=0 \\
& =O\left(h^{2 r}\right)+O\left(h^{r-1}\right) O\left(h^{r+1}\right)=O\left(h^{2 r}\right)
\end{aligned}
$$

Then considering $u(x)$ we have

$$
\begin{aligned}
u(x)- & \sum\left(w_{k} p_{k}+w_{k}^{\prime} P_{k}\right) \\
& =\left(u(x)-\sum\left(c_{k} p_{k}+c_{k}^{\prime} P_{k}\right)\right)+\sum\left(\left(c_{k}-w_{k}\right)^{\prime} P_{k}+\left(c_{k}-w_{k}\right) p_{k}\right) \\
& =O\left(h^{2 r-1}\right)+\sum\left(c_{k}-w_{k}\right)^{\prime} P_{k}+\sum\left(c_{k}-w_{k}\right) p_{k} \\
& =O\left(h^{2 r-1}\right)+\sum\left(c_{k}-w_{k}\right)^{\prime}\left(P_{k}-U\right)+\sum\left(c_{k}-w_{k}\right)\left(p_{k}-u\right)
\end{aligned}
$$

since $U \sum\left(c_{k}^{\prime}-w_{k}^{\prime}\right)=U(0-0)=0 \quad$ and $\quad u \sum\left(c_{k}-w_{k}\right)=u(1-1)=0$

$$
=O\left(h^{2 r-1}\right)+O\left(h^{r-2}\right) O\left(h^{r+1}\right)_{k}+O\left(h^{r-1}\right) O\left(h^{r}\right)=O\left(h^{2 r-1}\right) .
$$

In the classic paper [49], Jiang and Shu show that choosing smoothness indicators of the form

$$
\begin{equation*}
\sum_{l=1}^{r-1} \int_{I} h^{2 l-1}\left(p_{k}^{(l)}\right)^{2} d x \tag{2.41}
\end{equation*}
$$

satisfy this requirement, at least for $r=2$ or 3 , that is schemes of $3^{\text {rd }}$ and $5^{\text {th }}$ order. For a uniform grid the smoothness indicators for the $3^{\text {rd }}$ order scheme are

$$
\begin{align*}
& I S_{0}=\left(\bar{u}_{-1}-\bar{u}_{0}\right)^{2},  \tag{2.42}\\
& I S_{1}=\left(\bar{u}_{0}-\bar{u}_{1}\right)^{2} \tag{2.43}
\end{align*}
$$

and for the $5^{\text {th }}$ order scheme they are

$$
\begin{array}{r}
I S_{0}=4 \bar{u}_{-2}^{2}+25 \bar{u}_{-1}^{2}-31 \bar{u}_{-1} \bar{u}_{0}+10 \bar{u}_{0}^{2}+\bar{u}_{-2}\left(-19 \bar{u}_{-1}+11 \bar{u}_{0}\right), \\
I S_{1}=4 \bar{u}_{-1}^{2}+13 \bar{u}_{0}^{2}-13 \bar{u}_{0} \bar{u}_{1}+4 \bar{u}_{1}^{2}+\bar{u}_{-1}\left(-13 \bar{u}_{0}+5 \bar{u}_{1}\right), \\
I S_{2}=10 \bar{u}_{0}^{2}+25 \bar{u}_{2}^{2}-19 \bar{u}_{1} \bar{u}_{2}+4 \bar{u}_{2}^{2}+\bar{u}_{0}\left(-31 \bar{u}_{1}+11 \bar{u}_{2}\right) . \tag{2.46}
\end{array}
$$

The IWENO scheme solves equation (2.40) for a uniform grid we use reconstruction polynomials (2.24) - (2.27) with weights defined in (4.33), using coefficients in equations (2.33) - (2.37) and smoothness indicators in equations (2.42) - (2.46).

### 2.2.4 Numerical Results in 1D

In this section we will briefly demonstrate the performance the of the IWENO algorithm. We proceed with the same set up as in section 2.2.2 using long double precision.

We reconstruct the function $u(x)=\sin (\pi x)$ at the point $x=0.716$ with the relative location $\alpha=0.334$ in $I_{0}$. The results for the $3^{\text {rd }}$ and $5^{\text {th }}$ order schemes are given in Table 2.2. These serve to verify that non-linear weighting preserves the order of the scheme.

This example is repeated to illustrate the performance at the midpoint where $\alpha=0.5$. Numerical results are given in Table 2.3. On the coarse mesh we have smaller errors than when $\alpha=0.334$. Also, as the mesh size is reduced there is a super convergence effect where our third order scheme converges at fourth order, and our fifth order scheme converges at sixth order. This variability will influence the convergence rate of the transport schemes derived later.

| $\Delta x^{-1}$ | Error | Order | $\Delta x^{-1}$ | Error | Order |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 10 | 7.16 E-4 |  | 10 | $1.49 \mathrm{E}-5$ |  |
| 20 | 8.61 E-5 | 3.06 | 20 | $4.52 \mathrm{E}-7$ | 5.04 |
| 40 | $1.05 \mathrm{E}-5$ | 3.03 | 40 | 1.39 E-8 | 5.03 |
| 80 | $1.30 \mathrm{E}-6$ | 3.02 | 80 | 4.29 E-10 | 5.01 |
| 160 | 1.62 E-7 | 3.01 | 160 | 1.33 E-11 | 5.01 |
| 320 | $2.02 \mathrm{E}-8$ | 3.00 | 320 | 4.16 E-13 | 5.00 |
| 640 | 2.51 E-9 | 3.00 | 640 | $1.28 \mathrm{E}-14$ | 5.02 |
| 1280 | $3.14 \mathrm{E}-10$ | 3.00 | 1280 | 2.65 E-16 | 5.59 |
| 2560 | $3.92 \mathrm{E}-11$ | 3.00 | 2560 | $2.00 \mathrm{E}-16$ | 0.41 |
| 5120 | 4.91 E-12 | 3.00 | 5120 | $1.28 \mathrm{E}-16$ | 0.63 |

Table 2.2: IWENO Reconstruction. The error in computing $\sin (\pi x)$ at $x=0.716$ at a relative cell location of 0.334 using the IWENO scheme. The results on the left are for the $3^{\text {rd }}$ order scheme and on the right is the $5^{\text {th }}$ order scheme.

| $\Delta x^{-1}$ | Error | Order |
| :---: | :---: | :---: |
| 10 | $3.53 \mathrm{E}-5$ |  |
| 20 | $2.22 \mathrm{E}-6$ | 3.99 |
| 40 | $1.39 \mathrm{E}-7$ | 4.00 |
| 80 | $8.68 \mathrm{E}-8$ | 4.00 |
| 160 | $5.42 \mathrm{E}-10$ | 4.00 |
| 320 | $3.39 \mathrm{E}-11$ | 4.00 |
| $x^{-1}$ | Error | Order |
| 10 | $5.16 \mathrm{E}-7$ |  |
| 20 | $8.13 \mathrm{E}-9$ | 5.99 |
| 40 | $1.27 \mathrm{E}-10$ | 6.00 |
| 80 | $1.99 \mathrm{E}-12$ | 6.00 |
| 160 | $3.10 \mathrm{E}-14$ | 6.00 |
| 320 | $6.51 \mathrm{E}-16$ | 5.57 |

Table 2.3: IWENO Midpoint Reconstruction. The error in computing $\sin (\pi x)$ at $x=0.716$ located at the midpoint of the cell, using the IWENO scheme. The results on the left are for the $3^{\text {rd }}$ order scheme and on the right is the $5^{\text {th }}$ order scheme.

### 2.2.5 Reconstruction in Higher Dimensions

The purpose of this section is to outline how this algorithm extends to multiple spatial dimensions. To do this we first define the primitive function as

$$
\begin{equation*}
U(x)=\int_{x_{1}}^{x_{1}} \cdots \int_{x_{d_{0}}}^{x_{d}} u(\xi) d \xi_{d} \ldots d \xi_{1}, \quad x=\left(x_{1}, \cdots, x_{d}\right) . \tag{2.47}
\end{equation*}
$$

The point $x_{0}=\left(x_{1_{0}}, \ldots, x_{d_{0}}\right)$ is a vertex of the center element. Here elements are of the form
$E=\left[x_{1_{i_{1}}}, x_{1_{i_{1}+1}}\right] \times \cdots \times\left[x_{2_{i_{d}}}, x_{2_{i_{d}+1}}\right], \quad$ for $\quad i_{j}=-n, \ldots, n \quad$ and $\quad j=0, \ldots, d$. The function $u(x)$ is the multicomponent derivative of the primitive function, i.e.,

$$
\begin{equation*}
u(x)=\frac{d}{d x_{1}} \cdots \frac{d}{d x_{d}} U(x) . \tag{2.48}
\end{equation*}
$$

For indexing simplicity, the remainder of this section will be presented for three spatial dimensions. The extension to higher dimensions is straightforward. The function $U(x)$ is reconstructed as a tensor product of the 1D reconstruction polynomials

$$
\begin{array}{r}
\sum_{k, l, m=0}^{n} c_{k}\left(x_{1}\right) P_{k}\left(x_{1}\right) \\
c_{l}\left(x_{2}\right) P_{l}\left(x_{2}\right) c_{m}\left(x_{3}\right) P_{m}\left(x_{3}\right)  \tag{2.49}\\
=\sum_{k, l, m=0}^{n} c_{k l m}(x) P_{k l m}(x)
\end{array}
$$

where the polynomials $c_{k l m}(x)$ and $P_{k l m}(x)$ are

$$
\begin{align*}
c_{k l m}(x) & =c_{k}\left(x_{1}\right) c_{l}\left(x_{2}\right) c_{m}\left(x_{3}\right)  \tag{2.50}\\
P_{k l m}(x) & =P_{k}\left(x_{1}\right) P_{l}\left(x_{2}\right) P_{m}\left(x_{3}\right) \tag{2.51}
\end{align*}
$$

The function $u(x)$ is computed by

$$
\begin{equation*}
\frac{d}{d x_{1}} \frac{d}{d x_{2}} \frac{d}{d x_{3}} \sum_{k, l, m=0}^{n} c_{k l m}(x) P_{k l m}(x) \tag{2.52}
\end{equation*}
$$

using the product rule for differentiation. These are computed using the one dimensional equations given in (2.24) - (2.27) and (2.33) - (2.37).

The nonlinear weighting is defined in the usual way as

$$
\begin{equation*}
w_{k l m}(x)=\frac{\sigma_{k l m}\left(x_{i}\right)}{\sum_{a b c} \sigma_{a b c}(x)} \quad \text { with } \quad \sigma_{k l m}(x)=\frac{c_{k l m}(x)}{\left(\epsilon+\mathrm{IS}_{k l m}\right)^{p}} \tag{2.53}
\end{equation*}
$$

The smoothness indicator, $I S_{k l m}$ must contain information from the entire polynomial $P_{k, l, m}(x)$ in order to maintain the high order reconstruction. Theorem 2.2.1 generalizes to the following result.

Theorem 2.2.2. Let $c_{k l m}(x)$ and $P_{k l m}(x)$ satisfy the following relations

$$
\begin{aligned}
U-P_{k l m} & =O\left(\Delta x_{1}^{r+1}+\Delta x_{2}^{r+1}+\Delta x_{3}^{r+1}\right), \\
U-\sum P_{k l m} c_{k l m} & =O\left(\Delta x_{1}^{2 r}+\Delta x_{2}^{2 r}+\Delta x_{3}^{2 r}\right) \\
\sum_{\gamma=0}^{n} c_{\gamma} & =1
\end{aligned}
$$

and $w_{k l m}(x)$ be such that the $\sum\left(c_{k l m}-w_{k l m}\right)=O\left(\Delta x_{1}^{r-1}+\Delta x_{2}^{r-1}+\Delta x_{3}^{r-1}\right)$. Then

$$
\begin{aligned}
U-\sum w_{k l m} P_{k l m} & =O\left(\Delta x_{1}^{2 r}+\Delta x_{2}^{2 r}+\Delta x_{3}^{2 r}\right) \\
u-\frac{d}{d x_{1}} \frac{d}{d x_{2}} \frac{d}{d x_{3}} \sum w_{k l m} P_{k l m} & =O\left(\Delta x_{1}^{2 r-1}+\Delta x_{2}^{2 r-1}+\Delta x_{3}^{2 r-1}\right)
\end{aligned}
$$

Proof. Follow exactly as in the proof for Theorem 2.2.1.

The code implements the standard set of smoothness indicators for a WENO scheme. In multiple spatial dimensions these are

$$
\begin{equation*}
\mathrm{IS}_{k l m}=\sum_{|\alpha|=1}^{r-1} \int_{\mathrm{E}} h^{2|\alpha|-1}\left(p_{k l m}^{(\alpha)}\right)^{2} d x \tag{2.54}
\end{equation*}
$$

### 2.2.6 Numerical Results in 2D

In this section we give an example of the performance of the $5^{\text {th }}$ order IWENO scheme in two spatial dimensions. To do this we define the point $\left(x^{*}, y^{*}\right)$ that we wish to reconstruct along with a relative locations $\alpha, \beta \in[0,1]$ inside cell $I_{0}$ that $\left(x^{*}, y^{*}\right)$ is located. $I_{0}$ is the rectangular element

$$
\begin{equation*}
\left[x^{*}-\alpha \Delta x, x^{*}+(1-\alpha) \Delta x\right] \times\left[y^{*}-\beta \Delta y, y^{*}+(1-\beta) \Delta y\right] \tag{2.55}
\end{equation*}
$$

The function we reconstruct is $u(x, y)=\sin (\pi x+\pi y)$. The reconstruction is computed for the point $\left(x^{*}, y^{*}\right)=(0.489,0.372)$ at the relative locations $(\alpha, \beta)=(0.123,0.746)$. We use a uniform grid with $\Delta x=\Delta y$, and analyze the errors as we take successive grid refinements $\Delta x=\Delta x / 2$. These are presented in Table 2.4.

Note that the rounding error accumulation for these is more than that for the 1D problem, the reason is taking a multicomponent derivative of the primitive function corresponds with dividing out by $\Delta x$ twice. If increased precision is desired quadratic precision can be implemented.

Furthermore the results using non-linear weighting can develop oscillations in convergence rates. In the example provided, we see what appears to

| $\Delta x^{-1}$ | Error | Order |
| :---: | :---: | :---: |
| 10 | $1.40 \mathrm{E}-5$ |  |
| 20 | $3.98 \mathrm{E}-7$ | 5.13 |
| 40 | $1.17 \mathrm{E}-8$ | 5.08 |
| 80 | $3.56 \mathrm{E}-10$ | 5.04 |
| 160 | $1.17 \mathrm{E}-11$ | 4.92 |
| 320 | $1.64 \mathrm{E}-12$ | 2.83 |
| 640 | $1.54 \mathrm{E}-11$ | -3.23 |


| $\Delta x^{-1}$ | Error | Order |
| :---: | :---: | :---: |
| 10 | $1.40 \mathrm{E}-5$ |  |
| 20 | $3.98 \mathrm{E}-7$ | 5.13 |
| 40 | $1.17 \mathrm{E}-8$ | 5.08 |
| 80 | $3.56 \mathrm{E}-10$ | 5.04 |
| 160 | $1.09 \mathrm{E}-11$ | 5.02 |
| 320 | $3.41 \mathrm{E}-13$ | 5.00 |
| 640 | $1.37 \mathrm{E}-14$ | 4.64 |


| $\Delta x^{-1}$ | Error | Order |
| :---: | :---: | :---: |
| 10 | $7.74 \mathrm{E}-4$ |  |
| 20 | $3.43 \mathrm{E}-5$ | 4.50 |
| 40 | $4.52 \mathrm{E}-7$ | 6.24 |
| 80 | $9.04 \mathrm{E}-10$ | 8.97 |
| 160 | $1.22 \mathrm{E}-11$ | 6.21 |
| 320 | $1.64 \mathrm{E}-12$ | 2.89 |
| 640 | $1.54 \mathrm{E}-11$ | -3.23 |


| $\Delta x^{-1}$ | Error | Order |
| :---: | :---: | :---: |
| 10 | $7.74 \mathrm{E}-4$ |  |
| 20 | $3.43 \mathrm{E}-5$ | 4.50 |
| 40 | $4.52 \mathrm{E}-7$ | 6.24 |
| 80 | $9.04 \mathrm{E}-10$ | 8.97 |
| 160 | $1.15 \mathrm{E}-11$ | 6.30 |
| 320 | $3.42 \mathrm{E}-13$ | 5.06 |
| 640 | $1.36 \mathrm{E}-14$ | 4.65 |

Table 2.4: IWENO 2D Reconstruction. The error in computing $\sin (\pi x+$ $\pi y$ ) using the $5^{\text {th }}$ order IWENO scheme. On the top are results using linear weighting coefficients, the results on the bottom use non-linear weighting. The left side is double precision, and the right is long double precision.
be a super convergence effect, however this is due to the larger errors on the coarse meshes. For the results using the non-linear weighting and long double precision, the convergence rate for points corresponding to $\Delta x^{-1}=80, \ldots, 640$ is 5.31 , as found through a linear regression analysis.

## Chapter 3

## A Quadrature Based Eulerian-Lagrangian Scheme for Linear Advection

### 3.1 History of Characteristic Methods

We overview the history of the development of characteristic methods to solve the advection problem,

$$
\begin{equation*}
u_{t}+\nabla \cdot f(u)=q . \tag{3.1}
\end{equation*}
$$

This equation models the transport of a conserved quantity $u$, such as mass, energy or momentum. The fact that $u$ is conserved means that assuming no mass is created or destroyed over the boundaries of the domain $\Omega$ that the quantity $\int_{\Omega} u d x$ is constant in time. The flux function $f(u)$ describes the local movement of $u$ in the domain and $q$ is an external source term. One method to theoretically analyze conservation equations is to look at the characteristics curves $\xi(t)$. These are solutions of the ordinary differential equation $(\mathrm{ODE}) \xi^{\prime}(t)=f^{\prime}(u)$. The solution of the advection problem restricted to these characteristic curves is constant, i.e.,

$$
\begin{align*}
\frac{d}{d t} u(\xi(t), t) & =\frac{\partial}{\partial t} u(\xi(t), t)+\nabla u(\xi(t), t) \cdot \xi^{\prime}(t) \\
& =u_{t}+\nabla \cdot f(u) \\
& =0 \tag{3.2}
\end{align*}
$$

The linear advection problem is defined by setting $f(u)=v(x, t) u$. In this case the characteristics curves are the same as the streamlines or particle velocity paths. Characteristic numerical methods are developed by taking advantage of the solution being constant along these curves as a way to relieve the restrictive CFL constraint.

Characteristic schemes were first made viable with the modified method of characteristics (MMOC) defined by Douglas and Russell in [32]. In order to solve the linear problem on a time interval $\left[t^{n}, t^{n+1}\right]$, they develop a method to backward trace along characteristic curves. Their approach traces node points $x$ at $t^{n+1}$ to traceback points $\check{x}$ at $t^{n}$. These values are used to approximate the derivative in the characteristic direction as

$$
\begin{equation*}
u_{\tau}(x, t) \equiv \frac{1}{\left(1+v^{2}\right)^{1 / 2}}\left(u_{t}+v \cdot \nabla u(x, t)\right) \approx \frac{u\left(x, t^{n+1}\right)-u\left(\check{x}, t^{n}\right)}{\Delta t} \tag{3.3}
\end{equation*}
$$

resulting in the finite difference approximation

$$
\begin{equation*}
u\left(x, t^{n+1}\right)=u\left(\check{x}, t^{n}\right)+\Delta t q . \tag{3.4}
\end{equation*}
$$

Since this method is based on individual points, it fails to be either mass or volume conservative. This work was furthered by the introduction of the modified method of characteristics with adjusted advection (MMOCAA) [30]. This method provides a correction for the global mass balance; however, there is no local mass conservation.

Eulerian-Lagrangian schemes were then developed in terms of a local mass balance. These types of methods include various Eulerian-Lagrangian
localized adjoint methods (ELLAM) $[20,26,78,77,75]$ which include characteristic mixed methods (CMM) [9]. Instead of tracing node points, an entire grid element $E$ is traced along the characteristic to a traceback element $\check{E}$ creating a tessellation of the domain. This gives local mass conservation, since all the mass in $\check{E}$ is numerically transported forward into $E$. These methods are not locally volume conservative, since in multiple dimensions it is not possible to trace the entire boundary of $E$. Methods were developed tracing a finite subset of the boundary points and connecting them with a piecewise polynomial interpolation of the boundary. Here the volume of the approximated traceback element is different than the volume of the theoretical traceback element. Arbogast and Huang developed the volume corrected characteristic mixed method (VCCMM) [3, 8] by adjusting the locations of the trace points in a way that ensured the method was both locally mass and volume conservative.

In 2011 Qiu, Shu and Christlieb [61, 62] derived semi-Lagrangian WENO finite difference schemes as a way to develop a high order transport scheme. Their scheme only works for constant velocity. In 2012 Arbogast and Huang created the Eulerian-Lagrangian WENO (EL-WENO) finite volume scheme [48] for the linear advection problem. This method does not naturally extend to higher spatial dimensions without the aid of a splitting technique.

The nonlinear problem is challenging since tracing either characteristic or particle paths is a function of the solution over the traceline. In 2012 Arbogast, Huang, and Russell [6] solve the nonlinear two-phase flow system
by tracing particle paths approximately and then adjusting them by solving an expensive minimization problem. The work later developed in [46] by Arbogast and Huang extends EL-WENO schemes to the nonlinear problem by solving the related problem

$$
\begin{equation*}
\left(u_{t}+v u_{x}\right)+\left(f(u)_{x}-v u_{x}\right)=q \tag{3.5}
\end{equation*}
$$

where $v$ is a chosen as an approximation to either the particle velocity field $f(u) / u$ or the wave velocity field $f^{\prime}(u)$. The linear advection problem is solved using an EL-WENO technique, leaving a flux correction term which subjects the method to a relaxed CFL constraint. This method alleviates the need to solve the minimization problem, however the technique used to solve the flux correction at a high order accuracy involves a large spatial stencil.

The remainder of this dissertation will be focused on providing improvements to Eulerian-Lagrangian schemes, with the aid of the IWENO technique derived in Chapter 2. This chapter will be focused on developing a method to solve the multi-dimensional linear advection equation.

### 3.2 Quadrature Based Scheme

Let $u(x, t)$ be the mass which is transported in a velocity field $v(x, t)$ as described by

$$
\begin{array}{ll}
u_{t}+\nabla \cdot(v u)=0, & x \in \mathbb{R}^{d}, \quad t \in \mathbb{R}^{+} \\
u(x, 0)=u_{0}(x), & x \in \mathbb{R}^{d} . \tag{3.6}
\end{array}
$$



Figure 3.1: Eulerian-Lagrangian Scheme. Illustration of a theoretical Eulerian-Lagrangian scheme. All the points in $E$ are traced through the velocity field $v$ in order to find a traceback region $\dot{E}$. This process also defines a space-time region $\mathcal{E}$ for the time step $\left[t^{n}, t^{n+1}\right]$.

We will numerically solve (3.6) using an Eulerian-Lagrangian schemes where the mass in an element $E$ is given by the mass located in a region $\check{E}$ defined through backward characteristic tracing as shown in Figure 3.1. The tracing is defined at every point in $E$ as the curve $\check{x}(t)$ satisfying

$$
\begin{align*}
& \frac{d \check{x}}{d t}=v(\check{x}, t), \quad t^{n} \leq t<t^{n+1}  \tag{3.7}\\
& \check{x}\left(t^{n+1}\right)=x
\end{align*}
$$

This tracing defines a map $F(x): E \rightarrow \check{E}$. The volume swept out by this tracing is called the space-time traceback region and is denoted $\mathcal{E}$. Furthermore, the divergence theorem gives the mass in the traceback region $\check{E}$ is exactly the same as the mass in the cell $E$ at the future time.

$$
\begin{equation*}
\int_{E} u\left(x, t^{n+1}\right) d x=\int_{\check{E}} u\left(x, t^{n}\right) d x . \tag{3.8}
\end{equation*}
$$

There is no flow across the boundary since the velocity field is tangent to the
boundary of $\mathcal{E}$.
Computationally it is not possible in spatial dimensions greater than one to trace all the points on the boundary of $E$. Meaning that the trace back region $\check{E}$ has unknown geometry. In order to compute the integral over $\check{E}$ a coordinate transformation is applied to the mass integral. This transformation is defined through the map $F$ and the integral after the coordinate transformation is

$$
\begin{equation*}
\int_{\check{E}} u\left(x, t^{n}\right) d x=\int_{E} u\left(F(x), t^{n}\right)\left|\frac{\partial F}{\partial x}\right| d x . \tag{3.9}
\end{equation*}
$$

Where $\left|\frac{\partial F}{\partial x}\right|$ is the determinant of the Jacobian matrix associated with the transformation. This integral is computed with a quadrature rule as

$$
\begin{equation*}
\sum_{k} w_{k} u\left(F\left(x_{k}\right), t^{n}\right)\left|\frac{\partial F}{\partial x}\right|_{x=x_{k}} \tag{3.10}
\end{equation*}
$$

where $x_{k}$ and $w_{k}$ are quadrature points and weights. This evaluation requires evaluating two terms. The first is the reconstruction of $u$ at time $t^{n}$ and location $F\left(x_{k}\right)$ which is computed using the IWENO scheme detailed in Chapter 2. The second term is the determinant of the Jacobian matrix evaluated at the quadrature point. In section 3.2.1 we develop a system of ODEs that solve for the Jacobian determinant.

### 3.2.1 Jacobian Matrix

This section looks at the determinant of the Jacobian Matrix,

$$
\left|\frac{\partial F}{\partial x}\right|=\left|\begin{array}{ccc}
\frac{\partial F_{1}}{\partial x_{1}} & \cdots & \frac{\partial F_{1}}{\partial x_{d}}  \tag{3.11}\\
\vdots & \ddots & \vdots \\
\frac{\partial F_{d}}{\partial x_{1}} & \cdots & \frac{\partial F_{d}}{\partial x_{d}}
\end{array}\right| .
$$

The map $F: E \rightarrow E$ is defined pointwise by particle tracing through the velocity field; that is, the components of the Jacobian are of the form

$$
\begin{equation*}
\frac{\partial F_{i}}{\partial x_{j}}=\frac{\partial \check{x}_{i}}{\partial x_{j}} \tag{3.12}
\end{equation*}
$$

The value of these terms at $t=t^{n+1}$ is understood since $\check{x}\left(t^{n+1}\right)=x$ or $\left(\check{x}_{1}\left(t^{n+1}\right), \ldots, \check{x}_{d}\left(t^{n+1}\right)\right)=\left(x_{1}, \ldots, x_{d}\right)$ and

$$
\begin{equation*}
\left.\frac{\partial \check{x}_{i}}{\partial x_{j}}\right|_{t=t^{n+1}}=\delta_{i j} \tag{3.13}
\end{equation*}
$$

Here $\delta_{i j}$ is the Kroneker delta, associated with the identity matrix. In order to get the values of the components at $t^{n}$, we take the time derivative, which can be related to the gradient of the velocity field as

$$
\begin{equation*}
\frac{d}{d t} \frac{\partial \check{x}_{i}}{\partial x_{j}}=\frac{\partial}{\partial x_{j}} \frac{d \check{x}_{i}}{d t}=\frac{\partial v_{i}\left(\check{x}_{i}, t\right)}{\partial x_{j}}=\nabla v_{i} \cdot\left(\frac{\partial \check{x}_{1}}{\partial x_{j}}, \cdots, \frac{\partial \check{x}_{d}}{\partial x_{j}}\right) . \tag{3.14}
\end{equation*}
$$

This sets up a system of ODEs

$$
\begin{align*}
& \frac{d}{d t} \sigma_{i j}=\nabla v_{i}(\check{x}, t) \cdot\left(\sigma_{1 j}, \ldots, \sigma_{d j}\right),  \tag{3.15}\\
& \frac{d}{d t} \check{x}=v  \tag{3.16}\\
& \sigma_{i, j}\left(t^{n+1}\right)=\delta_{i j}  \tag{3.17}\\
& \check{x}\left(t^{n+1}\right)=x \tag{3.18}
\end{align*}
$$

where $\sigma_{i j}=\frac{\partial \check{x}_{i}}{\partial x_{j}}$. The solution of the system at $t=t^{n}$ gives the components of the Jacobian Matrix at any point. This is evaluated for the finite subset of quadrature points and used to compute the value of the determinant at these locations.

An alternative method is to solve

$$
\begin{equation*}
\frac{d}{d t}\left|\frac{\partial \check{x}}{\partial x}\right|=\left|\frac{\partial \check{x}}{\partial x}\right|(\nabla \cdot v) \tag{3.19}
\end{equation*}
$$

for the determinant without resolving the individual components as derived in [39].

### 3.2.2 Local Conservation

Previous work with Eulerian-Lagrangian schemes guaranteed a local mass conservation since the elements tessellated the domain, and mass was directly transported forward. Since the quadrature scheme only traces select points and not the entire domain, it is neither locally mass nor volume conservative. In this section, we will derive a volume conservation equation and demonstrate how to determine the number of quadrature points to use in the scheme.

The volume difference of an element $E$ and that of its traceback element $\check{E}$ is the integral over $\mathcal{E}$ of the divergence of the velocity field, i.e.

$$
\begin{equation*}
\int_{\mathcal{E}} \nabla_{x} \cdot v d V=|E|-|\check{E}| \tag{3.20}
\end{equation*}
$$

This is verified in two spatial dimensions by noting

$$
\begin{aligned}
\int_{\mathcal{E}} \nabla_{x} \cdot v d V= & \int_{\mathcal{E}} \nabla_{x, t} \cdot\binom{v}{1} d V \\
= & \int_{\partial \mathcal{E}}\binom{v}{1} \cdot \eta_{x, t} d A \\
= & \int_{E}\binom{v}{1} \cdot\binom{0}{1} d x+\int_{\check{E}}\binom{v}{1} \cdot\binom{0}{-1} d x \\
& +\sum_{i=1}^{4} \int_{S_{i}}\binom{v}{1} \cdot \eta_{x, t} d x \\
= & \int_{E} 1 d x+\int_{\check{E}}-1 d x+\sum_{i=1}^{4} \int_{S_{i}} 0 d x \\
= & |E|-|\check{E}|
\end{aligned}
$$

here $S_{i}$ denotes a face of the traceback region. A normal vector to a given surface $S_{i}$ must be orthogonal to the tangent vector on the surface which is $(\check{x}(t), t)^{\prime}=(v, 1)$. This explains why the the dot product over these surfaces in the direction of the normal is zero. A scheme is called volume conservative if equation (3.20) is satisfied. When $\nabla \cdot v=0$, we say the flow is divergence free, and in this case volume conservation is simply $|E|=|\check{E}|$.

For the quadrature based scheme, the volume conservation equation is used to evaluate a local volume error. The right hand side of (3.20) is evaluated using the coordinate transformation with respect the the map $F$
and with quadrature points, $x_{k}$, and weights, $w_{k}$, as

$$
\begin{align*}
|E|-|\check{E}| & =|E|-\int_{\check{E}} 1 d A \\
& =|E|-\int_{E}\left|\frac{\partial F}{\partial x}\right|_{t^{n}} d A  \tag{3.21}\\
& =|E|-\sum_{k} w_{k}\left|\frac{\partial F}{\partial x}\right|_{\left(x_{k}, t^{n}\right)}
\end{align*}
$$

The left hand side is similarly computed as

$$
\begin{align*}
\int_{\mathcal{E}} \nabla_{x} \cdot v d V & =\int_{t^{n}}^{t^{n+1}} \int_{E(t)} \nabla_{x} \cdot v d A d t \\
& =\sum_{j} \tilde{w}_{j} \int_{E\left(t_{j}\right)} \nabla_{x} \cdot v d A \\
& =\sum_{j} \tilde{w}_{j} \int_{E}\left(\nabla_{x} \cdot v\right)\left|\frac{\partial F}{\partial x}\right|_{t_{j}} d A  \tag{3.22}\\
& =\sum_{j} \tilde{w}_{j} \sum_{k} w_{k}\left(\nabla_{x} \cdot v\right)\left|\frac{\partial F}{\partial x}\right|_{\left(x_{k}, t_{j}\right)}
\end{align*}
$$

where $\tilde{w}_{j}$ and $t_{j}$ are the quadrature weights and points in time. The local volume error estimate is the difference between equations (3.21) and (3.22),

$$
\begin{equation*}
|E|-\left(\sum_{k} w_{k}\left|\frac{\partial F}{\partial x}\right|_{\left(x_{k}, t_{n}\right)}\right)-\left(\sum_{j} \tilde{w}_{j} \sum_{k} w_{k}\left(\nabla_{x} \cdot v\right)\left|\frac{\partial F}{\partial x}\right|_{\left(x_{k}, t_{j}\right)}\right) \tag{3.23}
\end{equation*}
$$

The code adaptively chooses the number of quadrature points such that the volume error on a given element is less than a preset tolerance.

Local mass conservation means that the mass on element $\check{E}$ is forward transported to $E$ that is $\int_{E} u d A=\int_{\check{E}} u d A$. Since quadrature is used to solve for the mass on $\check{E}$ this requirement is not satisfied. In order to minimize this error a fifth order IWENO reconstruction scheme is used to compute
$u\left(F\left(x_{k}\right), t^{n}\right)$ and the selection of the quantity of quadrature points helps to reduce the error associated with computing $\left|\frac{\partial F}{\partial x}\right|_{x_{k}}$. In addition, a global mass error is monitored by checking the difference in total mass at the initial state with the total mass at a given time step.

### 3.2.3 Algorithm

The problem is solved on a uniform grid with spacing $\Delta x$ in the $x$ direction and $\Delta y$ in the $y$-direction, as is required by the IWENO interpolation scheme to construct the primitive function. The following procedure gives the algorithm to compute the mass at $t^{n+1}$ for a given element.

1. Start with an initial number of quadrature points $n_{q}$. For the two dimensional problem we use a tensor product of quadrature points; that is, in both spatial directions we use $n_{q}$ points for a total of $n_{q}^{2}$ quadrature points. Then solve the system of ODE's (3.15)-(3.18) to obtain the values of $\check{x}, \check{y}$ and $\left|\frac{\partial F}{\partial x}\right|$ for each quadrature point. The numerical solver implemented in our code is the fifth order components of the RungeKutta 45 scheme [35] proposed by Fehlberg in 1969, it is defined by the Butcher tableau in Table 3.1. For complex velocity fields and large time steps it is necessary to take multiple micro time steps in the Runge-Kutta scheme.
2. Use the values of the determinant of the Jacobian $\left|\frac{\partial F}{\partial x}\right|$ to compute the local volume error in equation (3.23). If the error is greater than the set

| 0 |  |  |  |  |  |  |  |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $1 / 4$ | $1 / 4$ | $9 / 32$ |  |  |  |  |  |
| $3 / 8$ | $3 / 32$ | $9296 / 2197$ |  |  |  |  |  |
| $12 / 13$ | $1932 / 2197$ | $-7200 / 2197$ | $72968 / 513$ | $-845 / 4104$ |  |  |  |
| 1 | $439 / 216$ | -8 | 3680 |  |  |  |  |
| $1 / 2$ | $-8 / 27$ | 2 | $-3544 / 2565$ | $1859 / 4104$ | $-11 / 40$ |  |  |
|  | $16 / 135$ | 0 | $6656 / 12825$ | $28561 / 56430$ | $-9 / 50$ | $2 / 55$ |  |

Table 3.1: Runge-Kutta Butcher Tableau. The fifth order Runge-Kutta scheme used to solve the ODE system (3.15)-(3.18).
tolerance, the number of quadrature points is incremented by one and step 1 is repeated to solve for $\check{x}, \check{y}$ and $\left|\frac{\partial F}{\partial x}\right|$ associated with the new set of quadrature points.
3. Evaluate $u\left(\check{x}, \check{y}, t^{n}\right)$ using the IWENO scheme given in Chapter 2.2.5. The values $u\left(\check{x}, \check{y}, t^{n}\right)$ and $\left|\frac{\partial F}{\partial x}\right|$ are used to solve equation 3.10 which approximates the total mass in the element at $t^{n+1}$.

### 3.2.4 Numerical Results in 1D

In this section numerical results for a variety of linear advection problems are provided. The problem is set up by giving the initial condition $u_{0}(x)$ and the velocity field $v$ and its Jacobian $J=\frac{d v}{d x}$ that the mass will transport in. A periodic boundary condition is implemented for all of the examples and the tolerance for the volume computation is set at $1 \mathrm{E}-9$. The global mass loss is verified to be less in magnitude than the computed error. These errors
are presented using the discrete $L^{1}$ norm

$$
\begin{equation*}
\left\|u^{n}-\bar{u}^{n}\right\|_{1}:=\sum_{i}\left|\frac{1}{h} \int_{E_{i}} u\left(x, t^{n}\right) d x-\bar{u}_{i}^{n}\right| h, \tag{3.24}
\end{equation*}
$$

and the discrete $L^{\infty}$ norm

$$
\begin{equation*}
\left\|u^{n}-\bar{u}^{n}\right\|_{\infty}:=\sup _{i}\left|\frac{1}{h} \int_{E_{i}} u\left(x, t^{n}\right) d x-\bar{u}_{i}^{n}\right| . \tag{3.25}
\end{equation*}
$$

The quadrature rule use for the numerical results is Gaussian Quadrature [45] which on the reference element $[-1,1]$ has points and weights listed in Table 3.2. The order of the approximation is $\Delta x^{2 n_{q}-1}$. In order to maintain fifth order accuracy the scheme initializes the number of quadrature points at $n_{q}=3$. There are other natural choices for the quadrature rule such as Gauss-Lobatto which includes points on the boundary; however, this particular quadrature rule did not produce significantly different results from the Gaussian quadrature and requires a greater number of points to achieve a given order of accuracy.

### 3.2.4.1 Example: Constant Speed Transport

We first test our scheme with the simple case of constant speed transport with velocity $v(x, t)=1$ where $v_{x}=0$. The initial condition is $u_{0}(x)=$ $\sin (\pi x)$ for $x \in[0,2]$. The exact solution is $u(x, t)=u_{0}(x-t)$. The trace points and determinant of the Jacobian are computed using the Runge-Kutta scheme using 1 time step. This is more than sufficient since the trace points are of the form $\check{x}=x-\Delta t$ and $\left|\frac{\partial \check{x}}{\partial x}\right|=1$ and are solved exactly using the single Runge-Kutta time step.

| Gaussian Quadrature |  |  | Gauss-Lobatto Quadrature |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  | $n_{q}$ | $x_{i}$ | $w_{i}$ |
| $n_{q}$ | $x_{i}$ | $w_{i}$ | 2 | $\pm 1$ | 1 |
| 1 | 0 | 2 |  | 0 | $\frac{4}{3}$ |
| 2 | $\pm \sqrt{\frac{1}{3}}$ | 1 | 3 | $\pm 1$ | $\frac{1}{3}$ |
| 3 | 0 | $\frac{8}{9}$ | 4 | $\frac{1}{5} \sqrt{5}$ | $\frac{5}{6}$ |
|  | $\pm \sqrt{\frac{3}{5}}$ | $\frac{5}{9}$ |  | $\pm 1$ | $\frac{1}{6}$ |
| 4 | $\pm \sqrt{\frac{3}{7}-\frac{2}{7} \sqrt{\frac{6}{5}}}$ | $\frac{18+\sqrt{30}}{36}$ | 5 | 0 | $\frac{32}{45}$ |
|  | $\pm \sqrt{\frac{3}{7}+\frac{2}{7} \sqrt{\frac{6}{5}}}$ | $\frac{18-\sqrt{30}}{36}$ |  | $\pm \frac{1}{7} \sqrt{21}$ | $\frac{49}{90}$ |
| 5 | 0 | $\frac{128}{225}$ |  | $\pm 1$ | $\frac{1}{10}$ |
|  | $\pm \frac{1}{3} \sqrt{5-2 \sqrt{\frac{10}{7}}}$ | $\frac{322+13 \sqrt{70}}{900}$ | 6 | $\pm \sqrt{\frac{1}{21}(7-2 \sqrt{7})}$ | $\frac{1}{30}(14+\sqrt{7})$ |
|  | $\pm \frac{1}{3} \sqrt{5+2 \sqrt{\frac{10}{7}}}$ | $\frac{322-13 \sqrt{70}}{900}$ |  | $\pm \sqrt{\frac{1}{21}(7+2 \sqrt{7})}$ | $\frac{1}{30}(14-\sqrt{7})$ |
|  |  |  |  | $\pm 1$ | $\frac{1}{15}$ |

Table 3.2: Quadrature Rules. The quadrature points $x_{i}$ and weights $w_{i}$ on a reference element $[-1,1]$. On the left is Gaussian quadrature, on the right is Gauss-Lobatto quadrature.

Numerical results in Table 3.3 are at $t=1.5$ at three different time steps. The first is for $\Delta t=0.1$ here elements will trace exactly to other elements in the domain. The results show the error to be within the rounding error of the scheme. The second choice of time step is $\Delta t=0.12$ the results show $5^{\text {th }}$ order convergence. For the final test $\Delta t=4.5 \Delta x$ this is 4.5 times the CFL time limit of a fixed grid WENO scheme where

$$
\Delta t_{\mathrm{CFL}}=\Delta x / \max |v|=\Delta x
$$

. The results show a larger error accumulation in comparison to taking a small fixed number of time steps.

### 3.2.4.2 Example: Time Dependent Velocity

This example uses the velocity $v=\sin (t)$ with $v_{x}=0$. The initial condition is again $u_{0}(x)=\sin (\pi x)$ for $x \in[0,2]$. The analytical solution to this problem is $u(x, t)=u_{0}(x+1+\cos (t))$. The trace points and determinant of the Jacobian are computed using the Runge-Kutta scheme with a single step. The error results at $t=1.5$ using 10 time steps are listed in Table 3.4 show the $5^{\text {th }}$ order convergence rate.

### 3.2.4.3 Example: Spatially Dependent Velocity

The most challenging test we run in 1D is with the velocity $v=\sin (x)$ with $v_{x}=\cos (x)$. The initial condition is $u_{0}(x)=1$ for $x \in[0,2 \pi]$. The

| $N$ | Error $L^{1}$ | Order | Error $L^{\infty}$ | Order |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\Delta t=0.1$ |  |  |  |  |  |  |  |
| 20 | $4.62 \mathrm{E}-15$ | - | $6.23 \mathrm{E}-15$ | - |  |  |  |
| 40 | $1.45 \mathrm{E}-14$ | - | $1.84 \mathrm{E}-14$ | - |  |  |  |
| 80 | $4.82 \mathrm{E}-14$ | - | $1.09 \mathrm{E}-13$ | - |  |  |  |
| 160 | $9.00 \mathrm{E}-14$ | - | $3.74 \mathrm{E}-13$ | - |  |  |  |
| 320 | $2.09 \mathrm{E}-13$ | - | $1.82 \mathrm{E}-12$ | - |  |  |  |
| 640 | $3.74 \mathrm{E}-13$ | - | $6.28 \mathrm{E}-12$ | - |  |  |  |
| $\Delta t=0.12$ |  |  |  |  |  |  |  |
| 10 | $4.93 \mathrm{E}-3$ | - | $3.81 \mathrm{E}-3$ | - |  |  |  |
| 20 | $1.42 \mathrm{E}-4$ | 5.11 | $1.11 \mathrm{E}-4$ | 5.10 |  |  |  |
| 40 | $2.54 \mathrm{E}-6$ | 5.80 | $2.00 \mathrm{E}-6$ | 5.79 |  |  |  |
| 80 | $1.35 \mathrm{E}-7$ | 4.23 | $1.06 \mathrm{E}-7$ | 4.23 |  |  |  |
| 160 | $2.25 \mathrm{E}-9$ | 5.91 | $1.77 \mathrm{E}-9$ | 5.91 |  |  |  |
| 320 | $1.32 \mathrm{E}-10$ | 4.09 | $1.04 \mathrm{E}-10$ | 4.09 |  |  |  |
| 640 | $2.25 \mathrm{E}-12$ | 5.87 | $6.37 \mathrm{E}-12$ | 4.02 |  |  |  |
| LR |  |  |  |  |  | 5.12 | 4.92 |
| $\Delta t=4.5 \Delta x$ |  |  |  |  |  |  |  |
| 10 | $7.89 \mathrm{E}-4$ | - | $5.69 \mathrm{E}-4$ | - |  |  |  |
| 20 | $3.06 \mathrm{E}-5$ | 4.69 | $3.90 \mathrm{E}-5$ | 3.87 |  |  |  |
| 40 | $1.66 \mathrm{E}-6$ | 4.21 | $1.94 \mathrm{E}-6$ | 4.33 |  |  |  |
| 80 | $6.31 \mathrm{E}-8$ | 4.71 | $1.09 \mathrm{E}-7$ | 4.15 |  |  |  |
| 160 | $8.19 \mathrm{E}-9$ | 2.95 | $6.64 \mathrm{E}-9$ | 4.04 |  |  |  |
| 320 | $2.25 \mathrm{E}-10$ | 4.08 | $\mathrm{E}-10$ | 4.03 |  |  |  |
| 640 | $5.93 \mathrm{E}-11$ | 4.63 | $\mathrm{E}-11$ | 3.14 |  |  |  |
| LR |  | 4.03 |  | 4.00 |  |  |  |

Table 3.3: Constant Transport Test 1D. Error and convergence order for transport with velocity $v=1$ at time $t=1.5$ with $\Delta t=0.1,0.12$, and $4.5 \Delta x$. LR stands for linear regression, and the order is the slope of the linear regression line through all the data points.

| N | Error $L_{1}$ | Order $L_{1}$ | Error $L_{\infty}$ | Order $L_{\infty}$ |
| :---: | :---: | :---: | :---: | :---: |
| 10 | $2.92 \mathrm{E}-3$ | - | $2.26 \mathrm{E}-3$ | - |
| 20 | $9.71 \mathrm{E}-5$ | 4.91 | $7.65 \mathrm{E}-5$ | 4.89 |
| 40 | $8.20 \mathrm{E}-7$ | 6.89 | $6.42 \mathrm{E}-7$ | 6.90 |
| 80 | $3.90 \mathrm{E}-8$ | 4.39 | $3.06 \mathrm{E}-8$ | 4.39 |
| 160 | $3.83 \mathrm{E}-9$ | 3.35 | $3.00 \mathrm{E}-9$ | 3.35 |
| LR |  | 5.04 |  | 5.03 |

Table 3.4: Time Dependent Transport Test 1D. Error and convergence order for transport with velocity $v=\sin (t)$ at time $t=1.5$ with $\Delta t=0.15$.
analytical solution to this problem is

$$
\begin{equation*}
u(x, t)=\frac{\sin \left(2 \arctan \left(e^{-t} \tan (x / 2)\right)\right)}{\sin (x)} . \tag{3.26}
\end{equation*}
$$

The trace points and determinant of the Jacobian are computed using the Runge-Kutta scheme using $N_{\mathrm{RK}}=10$ micro time steps. This value $N_{\mathrm{RK}}$ was sufficient to keep the volume conservation equation within the required tolerance, and experimentation using a larger value of $N_{\mathrm{RK}}$ did not improve the error results. These results are presented in Table 3.5 at $t=1$ with $\Delta t=0.1$.

### 3.2.4.4 Example: Shu's Linear Test

The final example in 1D is the standard test case Shu's linear test [14]. In this case $v(x, t)=1$ and $v_{x}=0$, and the analytical solution is $u(x, t)=u_{0}(x-t)$. The initial condition is

| N | Error $L_{1}$ | Order $L_{1}$ | Error $L_{\infty}$ | Order $L_{\infty}$ |
| :---: | :---: | :---: | :---: | :---: |
| 10 | $6.39 \mathrm{E}-2$ | - | $3.27 \mathrm{E}-3$ | - |
| 20 | $3.63 \mathrm{E}-3$ | 4.13 | $7.65 \mathrm{E}-5$ | 3.61 |
| 40 | $2.12 \mathrm{E}-4$ | 4.10 | $6.42 \mathrm{E}-7$ | 3.75 |
| 80 | $7.13 \mathrm{E}-6$ | 4.89 | $3.06 \mathrm{E}-8$ | 3.98 |
| 160 | $1.85 \mathrm{E}-7$ | 5.27 | $3.00 \mathrm{E}-9$ | 5.32 |
| LR |  | 4.58 |  | 4.10 |

Table 3.5: Spatially Dependent Transport Test 1D. Error and convergence order for transport with velocity $v=\sin (x)$ at time $t=1$ with $\Delta t=0.1$.

$$
u_{0}(x)= \begin{cases}\frac{1}{6}(G(x, \beta, z-\delta)+G(x, \beta, z+\delta)+4 G(x, \beta, z)), & 0.2 \leq x \leq 0.4 \\ 1, & 0.6 \leq x \leq 0.8 \\ 1-|10(x-1.1)|, & 1 \leq x \leq 1.2 \\ \frac{1}{6}(F(x, \alpha, a-\delta)+F(x, \alpha, a+\delta)+4 F(x, \alpha, a)), & 1.4 \leq x \leq 1.6 \\ 0, & \text { otherwise }\end{cases}
$$

where

$$
G(x, \beta, z)=e^{-\beta(x-z)^{2}} \quad \text { and } \quad F(x, \alpha, a)=\sqrt{\max \left(1-\alpha^{2}(x-a)^{2}, 0\right)}
$$

The constants are set to $a=0.5, z=0.3, \delta=0.005, \alpha=10$, and $\beta=$ $\log (2) /\left(36 \delta^{2}\right)$. This test is designed to see the resolution properties of a given scheme. One time step in used in the Runge-Kutta scheme solving the ODE system. Two tests are run, the first is for the number of elements $N=200$ the second is for $N=201$. The problem is solved after 4 rotations at time $t=8$ using 20 time steps. The solutions at the final time are in Figure 3.2. The first case exhibits similar behavior to the first test in Section 3.2.4.1. There is no visible diffusion and the global mass error is within rounding error. The


Figure 3.2: Shu's Linear Test 1D. Numerical solutions for Shu's Linear Test at time $t=8$ with 20 time steps. On the left there are 200 grid elements and on the right there are 201 elements.
second test we run with $\mathrm{N}=201$ this time trace element are no longer the same as grid elements. In this case some diffusion accumulates, and the global mass error is of the magnitude $1 \mathrm{E}-4$.

### 3.2.5 Numerical Results in 2D

In 2D the the linear advection equation is

$$
\begin{equation*}
u_{t}+\left(a_{1} u\right)_{x}+\left(a_{2} u\right)_{y}=0, \tag{3.27}
\end{equation*}
$$

where the velocity is $v=\left(a_{1}, a_{2}\right)$ with Jacobian $J(v)=\left(\begin{array}{ll}\left(a_{1}\right)_{x} & \left(a_{1}\right)_{y} \\ \left(a_{2}\right)_{x} & \left(a_{2}\right)_{y}\end{array}\right)$. The discrete $L^{1}$ norm is

$$
\begin{equation*}
\left.\left\|u^{n}-\bar{u}^{n}\right\|_{1}:=\sum_{i, j} \left\lvert\, \frac{1}{h^{2}} \int_{E_{i, j}} u\left(x, t^{n}\right) d A-\bar{u}_{i, j}^{n}\right.\right) \mid h^{2}, \tag{3.28}
\end{equation*}
$$

and the discrete $L^{\infty}$ norm is

$$
\begin{equation*}
\left.\left\|u^{n}-\bar{u}^{n}\right\|_{\infty}:=\sup _{i, j} \left\lvert\, \frac{1}{h^{2}} \int_{E_{i, j}} u\left(x, t^{n}\right) d A-\bar{u}_{i, j}^{n}\right.\right) \mid . \tag{3.29}
\end{equation*}
$$

All the numerical examples in this section use a tolerance of $1 \mathrm{E}-9$ for the volume computation, and use periodic boundary conditions on the uniform grid with $\Delta x=\Delta y$.

| $N$ | Error $L^{1}$ | Order | Error $L^{\infty}$ | Order |
| :---: | :---: | :---: | :---: | :---: |
| 10 | $2.34 \mathrm{E}-2$ | - | $9.02 \mathrm{E}-3$ | - |
| 20 | $4.30 \mathrm{E}-4$ | 5.76 | $1.67 \mathrm{E}-4$ | 5.75 |
| 40 | $2.10 \mathrm{E}-5$ | 4.36 | $8.24 \mathrm{E}-6$ | 4.32 |
| 80 | $4.06 \mathrm{E}-7$ | 5.69 | $1.59 \mathrm{E}-7$ | 5.69 |
| 160 | $2.04 \mathrm{E}-8$ | 4.32 | $8.00 \mathrm{E}-9$ | 4.32 |
| LR |  | 5.03 |  | 5.02 |

Table 3.6: Constant Transport Test 2D. Error and convergence order for transport with velocity $v=(1,1)$ at time $t=1.2$ with $\Delta t=0.12$.

### 3.2.5.1 Example: Constant Speed Transport

The first example we give is for the constant velocity field $v=(1,1)$ with Jacobian $J(v)=\left(\begin{array}{ll}0 & 0 \\ 0 & 0\end{array}\right)$. The initial condition is $u_{0}(x, y)=\sin (\pi(x+y))$ for $(x, y) \in[0,2]^{2}$. The analytic solution is $u(x, y, t)=\sin (\pi(x+y-2 t))$. The ODE system is solved using one Runge-Kutta time step, which is sufficient to solve for the trace points and Jacobian exactly. The error at time $t=1.2$ using 10 time steps is given in Table 3.6 which shows the theoretical $5^{\text {th }}$ order convergence rate.

### 3.2.5.2 Example: Rigid Body Rotation

Two dimensional rigid body rotation is described by

$$
\begin{equation*}
u_{t}-((y-1) u)_{x}+((x-1) u)_{y} \tag{3.30}
\end{equation*}
$$

for $(x, y) \in[0,2]^{2}$. The Jacbobian of the velocity field is

$$
\left(\begin{array}{cc}
0 & -1  \tag{3.31}\\
1 & 0
\end{array}\right)
$$

| $N$ | Error $L^{1}$ | Order | Error $L^{\infty}$ | Order |
| :---: | :---: | :---: | :---: | :---: |
| 10 | $1.23 \mathrm{E}-3$ | - | $1.49 \mathrm{E}-3$ | - |
| 20 | $1.33 \mathrm{E}-4$ | 3.2 | $1.46 \mathrm{E}-4$ | 3.35 |
| 40 | $5.44 \mathrm{E}-6$ | 4.6 | $8.15 \mathrm{E}-6$ | 4.16 |
| 80 | $1.23 \mathrm{E}-7$ | 5.5 | $2.07 \mathrm{E}-7$ | 5.3 |
| 160 | $2.90 \mathrm{E}-8$ | 2.1 | $1.33 \mathrm{E}-8$ | 4.0 |
| LR |  | 4.08 |  | 4.3 |

Table 3.7: Rigid Body Rotation Test 2D. Error and convergence order for transport with velocity $v=(1-y, x-1)$ at time $t=1$ with $\Delta t=0.1$.

The initial condition is a smooth radial bump function

$$
\begin{align*}
u(x, y, 0) & =\frac{2}{5}[\psi(1+r(x, y)) \psi(1-r(x, y))+1]  \tag{3.32}\\
r(x, y) & =\sqrt{(x-1)^{2}+(y-1)^{2}} \tag{3.33}
\end{align*}
$$

with $\psi(s)=e^{-1 / s^{2}}$ for $s>2$ and $\psi(s)=0$ otherwise. The Runge-Kutta scheme uses one time step. The numerical results at time $t=1$ using 10 time steps are presented in Table 3.7.

### 3.2.5.3 Example: Swirling and Deforming Velocity

The final and most severe test is obtained using a swirling deformation flow where the velocity field is given by

$$
\begin{align*}
& a_{1}(x, y)=\sin ^{2}\left(\frac{\pi x}{2}\right) \sin (\pi y) g(t)  \tag{3.34}\\
& a_{2}(x, y)=-\sin ^{2}\left(\frac{\pi y}{2}\right) \sin (\pi x) g(t) \tag{3.35}
\end{align*}
$$

The function $g(t)$ introduces a time dependance to the flow field

$$
\begin{equation*}
g(t)=2 \cos (\pi t / T) \tag{3.36}
\end{equation*}
$$



Figure 3.3: Swirling Deforming Flow 2D. On the left is the initial condition, on the right is the numerical solution at $t=T / 2$.
this equation solved on a time interval $t \in[0, T]$. This velocity field has flow that slows down and reverses directions so that the initial condition is recovered at the final time $T$. The initial condition is on the domain $[0,2]^{2}$ and is a radial bump function centered at the point (0.7, 0.7). Figure 3.3 illustrates the initial condition and the solution at $t=T / 2$. This test has a Jacobian given by

$$
\begin{align*}
& \left(a_{1}\right)_{x}=2 \pi \cos (\pi x / 2) \sin (\pi x / 2) \sin (\pi y) \cos (\pi t / T)  \tag{3.37}\\
& \left(a_{1}\right)_{y}=2 \pi \sin (\pi x / 2)^{2} \cos (\pi y) \cos (\pi t / T)  \tag{3.38}\\
& \left(a_{2}\right)_{x}=-2 \pi \sin (\pi y / 2)^{2} \cos (\pi x) \cos (\pi t / T)  \tag{3.39}\\
& \left(a_{2}\right)_{y}=-2 \pi \cos (\pi y / 2) \sin (\pi y / 2) \sin (\pi x) \cos (\pi t / T) \tag{3.40}
\end{align*}
$$

Numerical results in Table 3.8 are for the solution at $T=1$ using $\Delta t=0.1$ and 10 micro time steps for the Runge-Kutta scheme.

| $N$ | Error $L^{1}$ | Order | Error $L^{\infty}$ | Order |
| :---: | :---: | :---: | :---: | :---: |
| 10 | $2.95 \mathrm{E}-3$ | - | $2.03 \mathrm{E}-2$ | - |
| 20 | $8.99 \mathrm{E}-4$ | 1.7 | $4.87 \mathrm{E}-3$ | 2.06 |
| 40 | $9.22 \mathrm{E}-5$ | 3.29 | $8.34 \mathrm{E}-4$ | 2.55 |
| 80 | $4.58 \mathrm{E}-6$ | 4.33 | $9.12 \mathrm{E}-5$ | 3.19 |
| 160 | $1.17 \mathrm{E}-7$ | 5.28 | $3.48 \mathrm{E}-6$ | 4.71 |
| LR |  | 3.82 |  | 3.08 |

Table 3.8: Swirling and Deforming Flow 2D. Error and convergence order for transport with a time and spatially dependent velocity field at time $t=1$ with $\Delta t=0.1$.

## Chapter 4

## EL-WENO for Nonlinear Advection in 1D

We follow the method of Arbogast and Huang in [46] solving the nonlinear problem $u_{t}+f(u)_{x}=0$ by solving the related problem

$$
\begin{equation*}
\left(u_{t}+(v u)_{x}\right)+(f(u)-v u)_{x}=0 \tag{4.1}
\end{equation*}
$$

for an appropriately chosen velocity $v$. After integrating over the space-time region $\mathcal{E}(v)$ their method uses an EL-WENO scheme to solve the linear advection problem and uses a quadrature rule to solve for the flux integral

$$
\begin{equation*}
\int_{\mathcal{E}(v)}(f(u)-v u)_{x} d A=\int_{\partial \mathcal{E}(v)}\binom{f(u)-v u}{0} \cdot \eta_{x, t} d x \tag{4.2}
\end{equation*}
$$

This requires the evaluation of $u$ along the boundary of the space-time region. The $5^{\text {th }}$ order scheme of Arbogast and Huang computes this using a RungeKutta method that requires a 17 point stencil.

Arbogast and Huang were motivated by the CWENO schemes of Levy Puppo and Russo [51, 53, 52]. This chapter first reviews CWENO schemes as a way to understand the stencil for computing the integral of the flux term. This is a reasonable framework since it is easier to understand in a fixed mesh setting and the EL-WENO method generalizes to a fixed mesh WENO method when the velocity field $v$ is zero.

The purpose of this chapter is to provide a new algorithm for computing the flux term (4.2) using a compact local stencil. Concluding with numerical results directly comparing these methods and demonstrating formal convergence rates for the new scheme.

### 4.1 CWENO Flux Stencil

Central WENO schemes solve for the staggered cell averages

$$
\begin{align*}
\bar{u}_{i+1 / 2}^{n+1}= & 1 / h \int_{x_{i}}^{x_{i+1}} R\left(x, t^{n}\right) d x \\
& +\frac{\Delta t}{\Delta x} \sum_{l=0}^{m} \gamma_{l}\left[f\left(u\left(x_{i}, t^{n}+\beta_{l} \Delta t\right)\right)-f\left(u\left(x_{i+1}, t^{n}+\beta_{l} \Delta t\right)\right)\right] \tag{4.3}
\end{align*}
$$

where $R$ is a piecewise polynomial reconstruction of the cell averaged data. The values $\gamma_{l}, \beta_{l}$ are quadrature weights and points describing the intermediate times $t^{n}+\beta_{l} \Delta t$ where $u$ is computed. An illustration of this scheme is shown in Figure 4.1.

A Natural Continuous Extension (NCE) [85] of a Runge Kutta scheme is used to evolve the mass $u$ at the fixed position $x_{i}$ in time. The concentration changes in time according to the nonlinear advection equation as $u_{t}=\left.f(u)_{x}\right|_{x_{i}}$. The NCE Runge-Kutta scheme provides uniform accuracy for the solution in the time interval $\left[t^{n}, t^{n+1}\right]$. Each $v$-stage Runge Kutta method of order $p$ has a NCE of order $d \leq p$. That is there exist polynomials $b_{i}(\theta), i=1, . ., v$ of degree


Figure 4.1: Central WENO Scheme. An illustration of a time step in a central WENO scheme. The staggered grid element average $\bar{u}_{i+1 / 2}^{n+1}$ at the advanced time $t^{n+1}$ is computed from the reconstructed solution $R(x)$ over the interval $\left[x_{i}, x_{i+1}\right]$ at the time $t^{n}$ and a computation of the flux across the midpoints over the time interval $\left[t^{n}, t^{n+1}\right]$.
at most $d$ where

$$
\begin{align*}
u\left(t^{n}+\theta \Delta t\right) & =u\left(t^{n}\right)+\Delta t \sum_{i=1}^{v} b_{i}(\theta) g^{(i)}  \tag{4.4}\\
g^{(i)} & =F\left(t^{n}+\Delta t c_{i}, u^{n}+\Delta t \sum_{j=1}^{v} a_{i j} g^{(j)}\right)  \tag{4.5}\\
c_{i} & =\sum_{j} a_{i j}  \tag{4.6}\\
F(\tau, u(\tau)) & =-f_{x}\left(u\left(t^{n}+\tau\right)\right) \tag{4.7}
\end{align*}
$$

defines the NCE. The underlying Runge-Kutta scheme is determined by the selection of matrix $a$ and vector $b$. When matrix $a$ is lower triangular the scheme is explicit. These schemes are useful since they evaluate $u$ at multiple times along the line $x=x_{i}$ without solving for the parameters $g$ at each time separately.

The $5^{\text {th }}$ order CWENO method uses a 4-stage Runge Kutta scheme. At


Figure 4.2: Implicit Stencil for CWENO5 Flux Computation A depiction of the 4 -stage Runge Kutta stencil using a 5 point WENO reconstruction for computing the derivative $f\left(u\left(x_{i}, t\right)\right)_{x}$.
time $t^{n}$ the values of $u\left(x_{i+m}\right)$ are evaluated at the 17 points with $m=-8, \ldots, 8$. Then a 5 point WENO reconstruction scheme is used to get $f\left(u\left(x_{i+m}\right)\right)_{x}$ for $m=-6, \ldots, 6$. The value of $u$ at these 13 points are evolved to the second Runge-Kutta stage. They are used in the WENO scheme to evaluate $f\left(u\left(x_{i+m}\right)\right)_{x}$ for $m=-4, \ldots, 4$. These value of $u$ at these 9 points is computed at stage 3 , where the value of $f\left(u\left(x_{i+m}\right)\right)_{x}$ is reconstructed for 5 points with $m=-2, \ldots, 2$. At the fourth stage these values are used to compute $f\left(u\left(x_{i}\right)\right)_{x}$. This process is illustrated in Figure 4.2. In a standard CWENO scheme this stencil is implicitly defined since at each stage all the grid points are evolved at the same time.

This method was used by Arbogast and Huang in their EL-WENO scheme solving the nonlinear problem except for each traceline $\check{x}(t)$ they temporarily freeze the velocity field and explicitly compute these terms for each grid point separately. There is a detailed explanation of this procedure in their
paper [46]. The contribution in this chapter is to use a variation of the IWENO scheme to develop an alternative method for computing the flux integral with a compact local stencil.

### 4.2 A Flux Corrected EL-WENO Scheme

This section will develop a complete algorithm for solving the flux corrected EL-WENO scheme for the nonlinear advection problem. This is done by solving the related problem (4.1) in an Eulerian-Lagrangian setting. In order to define the velocity field $v$ the characteristic equation

$$
\begin{align*}
& \frac{d \xi}{d t}=\tilde{v}(\check{x}, t), \quad t^{n} \leq t<t^{n+1}  \tag{4.8}\\
& \xi\left(t^{n+1}\right)=x \tag{4.9}
\end{align*}
$$

is solved at node points. Here $\tilde{v}$ is an approximation to either the particle velocity field $f(u) / u$ or the wave velocity field $f^{\prime}(u)$. The velocity field $v$ at a grid point is the constant value $v_{i}=\frac{x_{i}-\xi_{i}}{\Delta t}$ and the velocity on element $E_{i}=\left[x_{i}, x_{i+1}\right]$ is the linear interpolation

$$
\begin{equation*}
v\left(x_{i+\alpha}\right)=(1-\alpha) v_{i}+\alpha v_{i+1} \tag{4.10}
\end{equation*}
$$

with $x_{i+\alpha}=x_{i}+\alpha \Delta x$, where $\alpha \in[0,1]$. The tracelines at grid points are

$$
\begin{equation*}
\check{x}_{i}(t)=x_{i}\left(t^{n+1}\right)-v_{i}\left(t^{n+1}-t\right) . \tag{4.11}
\end{equation*}
$$

This velocity also defines the space-time region $\mathcal{E}(v)$ and the trace region $\check{E}(v)$ for every element.

The scheme is defined by integrating over the space-time region $\mathcal{E}(v)$

$$
\begin{equation*}
\int_{\mathcal{E}(v)} \nabla_{x, t} \cdot\binom{v u}{u} d A+\int_{\mathcal{E}(v)} \nabla_{x, t} \cdot\binom{f(u)-v u}{0} d A=0 \tag{4.12}
\end{equation*}
$$

and applying the divergence theorem

$$
\begin{equation*}
\oint_{\partial \mathcal{E}(v)}\binom{v u}{u} \cdot\binom{\eta_{x}}{\eta_{t}} d S+\oint_{\partial \mathcal{E}(v)}\binom{f(u)-v u}{0} \cdot\binom{\eta_{x}}{\eta_{t}} d S=0 . \tag{4.13}
\end{equation*}
$$

The boundary of the space-time region $\partial \mathcal{E}(v)$ has 4 sides. They are $E_{i}$ with outward unit normal vector $\eta=\binom{0}{1}, \check{E}_{i}$ with $\eta=\binom{0}{-1}, S_{i}$ with $\eta=\binom{-1}{v_{i}} / \sqrt{1+v_{i}^{2}}$ and $S_{i+1}$ with $\eta=\binom{1}{-v_{i+1}} / \sqrt{1+v_{i+1}^{2}}$. This is depicted in Figure 4.3. Substituting these in the equation yields

$$
\begin{gather*}
\int_{E_{i}} u^{n+1} d x=\int_{\check{E}_{i}} u^{n} d x+\int_{S_{i}}(f(u)-v u) \frac{d \sigma}{\sqrt{1+v_{i}^{2}}} \\
-\int_{S_{i+1}}(f(u)-v u) \frac{d \sigma}{\sqrt{1+v_{i+1}^{2}}} \tag{4.14}
\end{gather*}
$$

In order to solve the right hand side we apply a two step approach The first step is a reconstruction $R(x)$ of the integral of the concentration $u$ on the trace element $\check{E}$ using WENO techniques. The second step is to solve for the flux integral approximately using a quadrature rule

$$
\begin{equation*}
\int_{S_{i}} f(u)-v u \frac{d \sigma}{\sqrt{1+v_{i}^{2}}}=\sum_{k=1}^{n_{q}} w_{k} \Delta t\left(f\left(u\left(x_{k}, t_{k}\right)\right)-v_{i} u\left(x_{k}, t_{k}\right)\right) \tag{4.15}
\end{equation*}
$$

where $n_{q}$ is the number of quadrature points, and $\left(x_{k}, t_{k}\right)$ and $w_{k}$ are the quadrature locations and weights. Since the point $\left(x_{k}, t_{k}\right)$ is on $S_{i}$ it satisfies the relation

$$
\begin{equation*}
x_{k}=\check{x}_{i}\left(t^{n}\right)+v_{i}\left(t_{k}-t^{n}\right) . \tag{4.16}
\end{equation*}
$$



Figure 4.3: Flux Corrected EL-WENO Scheme. An illustration of a time step in a flux corrected EL-WENO scheme. The grid element average $\bar{u}_{i}^{n+1}$ at the advanced time $t^{n+1}$ is computed from the reconstructed solution, $R(x)$, over the interval $\left[\check{x}_{i}, \check{x}_{i+1}\right]$ at the time $t^{n}$, and a correction for flux values across the trace lines.

The evolution of $u$ along $S_{i}$ is computed using a NCE Runge-Kutta scheme that uses the differentiation technique presented for the IWENO scheme.

This method can be extend to two spatial dimensions by using a Strang splitting technique where the conservation equation

$$
\begin{array}{ll}
u_{t}+f_{1}(u)_{x}+f_{2}(u)_{y}=0, & \text { for } \quad(x, y) \in \mathbb{R}^{2}, t>0, \\
u(x, y, 0)=u_{0}(x, y), & \text { for } \quad(x, y) \in \mathbb{R}^{2}, \tag{4.18}
\end{array}
$$

is solved by decoupling to a series of one dimensional problems. We do not provide details for this sort of algorithm, instead in Chapter 5 this flux correction technique is combined with the QEL-WENO scheme of Chapter 3 to provide a method for the nonlinear problem in multiple spatial dimensions.

### 4.2.1 Mass Reconstruction

There are two choices of methods to compute the mass integral

$$
\begin{equation*}
\int_{\check{E}_{i}} u\left(x, t^{n}\right) d x . \tag{4.19}
\end{equation*}
$$

The first is the QEL-WENO scheme given in Chapter 3. The second choice, which is the one implemented in the code for the numerical results presented in this chapter, is the standard EL-WENO method. The reason for this choice is it allows direct comparison of the computation of the flux term for the two different techniques. In the paper [48] developing the EL-WENO technique, Arbogast and Huang provide a detailed explanation of the mass integration. In this section I will briefly outline their method and give the details to implement it for schemes of $3^{\text {rd }}$ and $5^{\text {th }}$ order.

Trace regions on a given element are explicitly defined as $\check{E}_{j}^{n}=\left[\check{x}_{j}, \check{x}_{j+1}\right]$. They are viewed as a union of subintervals of grid elements $E_{i}$ of the form

1. $\left[x_{i}, x_{i+1}\right]$, when the whole $E_{i}$ is contained in $\check{E}_{j}^{n}$;
2. $\left[x_{i}, x_{i+\alpha}\right], 0<\alpha<1$, when $\check{x}_{j+1}^{n}$ intersects $E_{i}$ but not $\check{x}_{j}^{n}$;
3. $\left[x_{i+\beta}, x_{i+1}\right], 0<\beta<1$, when $\check{x}_{j}^{n}$ intersects $E_{i}$ but not $\check{x}_{j+1}^{n}$;
4. $\left[x_{i+\beta}, x_{i+\alpha}\right], 0<\beta<\alpha<1$, when $\check{E}_{j}^{n}$ is contained in $E_{i}$.

High order reconstruction of $u$ over these types of subintervals can be computed as integrals over type 1 and type 2 subintervals by noticing that type 3 and 4
subintervals are combinations of type 1 and 2 subintervals. Type 3 subintervals are

$$
\begin{equation*}
\left[x_{i+\beta}, x_{i+1}\right]=\left[x_{i}, x_{i+1}\right]-\left[x_{i}, x_{i+\beta}\right] \tag{4.20}
\end{equation*}
$$

and type 4 subintervals are

$$
\begin{equation*}
\left[x_{i+\beta}, x_{i+\alpha}\right]=\left[x_{i}, x_{i+\alpha}\right]-\left[x_{i}, x_{i+\beta}\right] . \tag{4.21}
\end{equation*}
$$

This gives a way of breaking the integral over any subinterval into a sum of integrals over type 1 and type 2 subintervals. Since the integral over a type 1 subinterval is $\Delta x \bar{u}_{i}$, it suffices to demonstrate how to compute a high order reconstruction of the integral over a type 2 subinterval.

To get the high order reconstruction of the integral over a type 2 subinterval a WENO scheme is used. Define $m+1$ polynomials, $p_{0}, \ldots, p_{m}$, of degree $m$ such that

$$
\begin{equation*}
\int_{I_{i}} p_{k} d x=\Delta x \bar{u}_{i}, \quad i=-m+k, \ldots, k, \tag{4.22}
\end{equation*}
$$

as shown in Figure 2.1. The value of $m$ is 1 or 2 for the $3^{\text {rd }}$ and $5^{\text {th }}$ order schemes respectively. We denote $P_{k}$ as the integral of $p_{k}$ over a type 2 subinterval for the $3^{\text {rd }}$ order scheme these polynomials are

$$
\begin{align*}
& P_{0}=\frac{\alpha \Delta x}{2}\left(\bar{u}_{i-1}(1-\alpha)+\bar{u}_{i}(1+\alpha)\right),  \tag{4.23}\\
& P_{1}=\frac{-\alpha \Delta x}{2}\left(\bar{u}_{i}(-3+\alpha)+\bar{u}_{i+1}(1-\alpha)\right), \tag{4.24}
\end{align*}
$$

and for the $5^{\text {th }}$ order scheme they are

$$
\begin{align*}
& P_{0}=\frac{\alpha \Delta x}{6}\left(\bar{u}_{i-2}\left(-1+\alpha^{2}\right)+\bar{u}_{i-1}\left(5-3 \alpha-2 \alpha^{2}\right)\right. \\
&\left.+\bar{u}_{i}\left(2+3 \alpha+\alpha^{2}\right)\right),  \tag{4.25}\\
& P_{1}=\frac{\alpha \Delta x}{6}( \bar{u}_{i-1}\left(2-3 \alpha+\alpha^{2}\right)+\bar{u}_{i}\left(5+\alpha-2 \alpha^{2}\right) \\
&\left.\quad+\bar{u}_{i+1}\left(-1+\alpha^{2}\right)\right),  \tag{4.26}\\
& P_{2}=\frac{\alpha \Delta x}{6}\left(\bar{u}_{i}\left(11-6 \alpha+\alpha^{2}\right)+\bar{u}_{i+1}\left(-7+9 \alpha-2 \alpha^{2}\right)\right. \\
&\left.+\bar{u}_{i+2}\left(2-3 \alpha+\alpha^{2}\right)\right) . \tag{4.27}
\end{align*}
$$

Linear weighting coefficients are chosen such that

$$
\begin{equation*}
\sum_{k=0}^{m} c_{k} P_{k}-\int_{x_{i}}^{x_{i+\alpha}} u(x) d x=O\left(\Delta x^{2 m+1}\right) \tag{4.28}
\end{equation*}
$$

For the $3^{\text {rd }}$ order scheme the linear coefficients are

$$
\begin{equation*}
c_{0}=(2-\alpha) / 3 \quad \text { and } \quad c_{1}=(1+\alpha) / 3 . \tag{4.29}
\end{equation*}
$$

For the fifth order scheme the linear coefficients are

$$
\begin{align*}
& c_{0}=(-3+\alpha)(-2+\alpha) / 20,  \tag{4.30}\\
& c_{1}=(3-\alpha)(2+\alpha) / 10,  \tag{4.31}\\
& c_{2}=(1+\alpha)(2+\alpha) / 20 . \tag{4.32}
\end{align*}
$$

The nonlinear weighting polynomials $w_{k}(x)$ are defined as a function of the linear coefficients and of a set of smoothness indicators $I S$ as

$$
\begin{equation*}
w_{k}=\frac{\sigma_{k}}{\sum_{i=0}^{m} \sigma_{i}} \quad \text { with } \quad \sigma_{i}=\frac{c_{i}}{\left(\epsilon+\mathrm{IS}_{i}\right)^{p}} \quad \text { for } \quad k=0, \ldots, m \tag{4.33}
\end{equation*}
$$

Here $p \geq 1$ is an integer and $0<\epsilon \ll 1$ is to avoid dividing by zero. The standard smoothness indicator for a WENO scheme are chosen, these are defined as

$$
\begin{equation*}
\sum_{l=1}^{m} \int_{I} h^{2 l-1}\left(p_{k}^{(l)}\right)^{2} d x \tag{4.34}
\end{equation*}
$$

The smoothness indicators for the third order scheme are

$$
\begin{align*}
& I S_{0}=\left(\bar{u}_{i-1}-\bar{u}_{i}\right)^{2},  \tag{4.35}\\
& I S_{1}=\left(\bar{u}_{i}-\bar{u}_{i+1}\right)^{2}, \tag{4.36}
\end{align*}
$$

and for the fifth order scheme they are

$$
\begin{align*}
& I S_{0}=\frac{13}{12}\left(\bar{u}_{i-2}-2 \bar{u}_{i-1}+\bar{u}_{i}\right)^{2}+\frac{1}{4}\left(\bar{u}_{i-2}-4 \bar{u}_{i-1}+3 \bar{u}_{i}\right)^{2},  \tag{4.37}\\
& I S_{1}=\frac{13}{12}\left(\bar{u}_{i-1}-2 \bar{u}_{i}+\bar{u}_{i+1}\right)^{2}+\frac{1}{4}\left(\bar{u}_{i-1}-\bar{u}_{i+1}\right)^{2}  \tag{4.38}\\
& I S_{2}=\frac{13}{12}\left(\bar{u}_{i}-2 \bar{u}_{i+1}+\bar{u}_{i+2}\right)^{2}+\frac{1}{4}\left(\bar{u}_{i}-4 \bar{u}_{i+1}+3 \bar{u}_{i+2}\right)^{2} . \tag{4.39}
\end{align*}
$$

The high order approximation of the mass in a type 2 subinterval is

$$
\begin{equation*}
\int_{x_{i}}^{x_{i}+\alpha} u(x) d x \approx \sum_{k=0}^{m} w_{k} P_{k}, \tag{4.40}
\end{equation*}
$$

which is computed using equations (4.23)-(4.39). This completes the details for the high order mass integration over the trace region $\check{E}$.

### 4.2.2 Flux Correction

The flux correction is computed by solving equation (4.15). The value of $u$ changes along the side $S$ in time according to the derivative

$$
\begin{equation*}
\frac{d u(x(t), t)}{d t}=\frac{\partial u}{\partial x} \frac{d x}{d t}+\frac{\partial u}{\partial t} . \tag{4.41}
\end{equation*}
$$

Where along $S$ the relation between $x$ and $t$ is given by equation (4.16) and $\frac{d x}{d t}=v$. Also $u_{t}=-f(u)_{x}$ substituting these gives

$$
\begin{equation*}
\frac{d u(x(t), t)}{d t}=u_{x} v-f(u)_{x} \tag{4.42}
\end{equation*}
$$

This equation is solved numerically using a NCE Runge-Kutta scheme. The value of $u$ is used to evolved along $S$ according to

$$
\begin{align*}
& u\left(\left(\check{x}\left(t^{n}\right), t^{n}\right)+\theta \Delta t(v, 1)\right)=u\left(\check{x}\left(t^{n}\right), t^{n}\right)+\Delta t \sum_{i=1}^{v} b_{i}(\theta) g^{(i)}  \tag{4.43}\\
& g^{(i)}=F\left(\left(\check{x}\left(t^{n}\right), t^{n}\right)+c_{i} \Delta t(v, 1), u\left(\check{x}\left(t^{n}\right), t^{n}\right)+\Delta t \sum_{j=1}^{v} a_{i j} g^{(j)}\right)  \tag{4.44}\\
& F((x, t), u(x, t))=\left(u_{x} v-f_{x}(u)\right) . \tag{4.45}
\end{align*}
$$

The matrix $a$ and the vector $b(\theta)$ for the third order scheme are

$$
\begin{equation*}
b_{1}(\theta)=-1 / 2 \theta^{2}+\theta, \quad b_{2}(\theta)=1 / 2 \theta^{2} \tag{4.46}
\end{equation*}
$$

and

$$
a=\left(\begin{array}{ll}
0 & 0  \tag{4.47}\\
1 & 0
\end{array}\right)
$$

For the fifth order scheme they are

$$
\begin{array}{ll}
b_{1}(\theta)=2 / 3 \theta^{3}-3 / 2 \theta^{2}+\theta, & b_{2}(\theta)=-2 / 3 \theta^{3}+\theta^{2} \\
b_{3}(\theta)=b_{2}(\theta), & b_{4}(\theta)=2 / 3 \theta^{3}-1 / 2 \theta^{2} \tag{4.49}
\end{array}
$$

and

$$
a=\left(\begin{array}{cccc}
0 & 0 & 0 & 0  \tag{4.50}\\
1 / 2 & 0 & 0 & 0 \\
0 & 1 / 2 & 0 & 0 \\
0 & 0 & 1 & 0
\end{array}\right)
$$

To use this Runge-Kutta scheme the values of $u$ are first reconstructed on a local stencil. For the trace point $\check{x}$ the stencil is $\check{x}+\gamma \Delta x$ where for the $3^{\text {rd }}$ order scheme

$$
\begin{equation*}
\gamma=-1 / 2,0,1 / 2 \tag{4.51}
\end{equation*}
$$

and for the $5^{\text {th }}$ order scheme

$$
\begin{equation*}
\gamma=-1 / 2,-1 / 4,0,1 / 4,1 / 2 \tag{4.52}
\end{equation*}
$$

At each point on the stencil the concentration $u\left(\check{x}+\gamma \Delta x, t^{n}\right)$ is evaluated using an IWENO reconstruction over the elements $E_{k-m}, \ldots, E_{k+m}$ where $\check{x} \in E_{k}$, here $m$ is the degree of the lower order polynomial used in the reconstruction.

In order to compute $u_{x}$ and $f(u)_{x}$ from these points a nodal variation of the IWENO scheme is used. The following explains the computation of $u_{x}$ and it is straightforward to use the points $f(u)$ in the same processes to get $f(u)_{x}$. Polynomials $P_{k}(x)$ are constructed such that

$$
\begin{equation*}
P_{k}\left(x_{i}\right)=u_{i}, \quad i=-m+k, \ldots, k+1, \tag{4.53}
\end{equation*}
$$

and linear coefficients $c_{k}(x)$ are found optimizing the value of the reconstruction at the point $x \in E_{k}$

$$
\begin{equation*}
\sum_{k=0}^{m} c_{k}(x) P_{k}(x)=u(x)+O\left(\Delta x^{2 m+1}\right) \tag{4.54}
\end{equation*}
$$

The results in chapter 2 guarantee the existence of these polynomials and coefficients and provides an algorithm for computing them. Since this is a nodal based interpolation instead of a cell average based interpolation these
polynomials are different than those presented in Chapter 2, for completeness they are listed at the end of this section. Following the IWENO strategy we take the derivative of the reconstruction

$$
\begin{equation*}
\sum_{k=0}^{m} c_{k}^{\prime}(x) P_{k}(x)+c_{k}(x) p_{k}(x)=u^{\prime}(x)+O\left(\Delta x^{2 m}\right) \tag{4.55}
\end{equation*}
$$

where $p_{k}(x)=P_{k}^{\prime}(x)$. The weighting coefficients are modified as in (4.33) with smoothness indicators

$$
\begin{equation*}
\sum_{l=1}^{m} \int_{I} h^{2 l-1}\left(p_{k}^{(l)}\right)^{2} d x \tag{4.56}
\end{equation*}
$$

This gives the reconstruction

$$
\begin{equation*}
u^{\prime}(x) \approx \sum_{k=0}^{m} w_{k}^{\prime}(x) P_{k}(x)+w_{k}(x) p_{k}(x) \tag{4.57}
\end{equation*}
$$

which is evaluated at all the points on the stencil $\check{x}+\gamma \Delta x$. These values are used to advance the Runge-Kutta Scheme to the next stage where the same procedure is applied.

The polynomials, linear coefficients, and smoothness indicators for the $3^{\text {rd }}$ order scheme are

$$
\begin{array}{cl}
P_{0}=(1+\alpha) u_{i}-\alpha u_{i-1}, & p_{0}=\left(u_{i}-u_{i-1}\right) / h, \\
P_{1}=(1-\alpha) u_{i}+\alpha u_{i+1}, & p_{1}=\left(u_{i+1}-u_{i}\right) / h, \\
c_{0}=(1-\alpha) / 2, & c_{0}^{\prime}=-1 / 2, \\
c_{1}=(\alpha+1) / 2, & c_{1}^{\prime}=1 / 2, \\
I S_{0}=\left(u_{i}-u_{i-1}\right)^{2}, & \text { and }  \tag{4.62}\\
I S_{1}=\left(u_{i+1}-u_{i}\right)^{2} .
\end{array}
$$

For the fifth order scheme they are

$$
\begin{gather*}
P_{0}=\frac{1}{2}\left(\alpha(\alpha+1) u_{i-2}-2 \alpha(\alpha+2) u_{i-1}+(\alpha+1)(\alpha+2) u_{i}\right),  \tag{4.63}\\
P_{1}=\frac{1}{2}\left(-2\left(\alpha^{2}-1\right) u_{i}+(\alpha-1) \alpha u_{i-1}+\alpha(\alpha+1) u_{i+1}\right),  \tag{4.64}\\
P_{2}=\frac{1}{2}\left((\alpha-2)\left((\alpha-1) u_{i}-2 \alpha u_{i+1}\right)+(\alpha-1) \alpha u_{i+2}\right),  \tag{4.65}\\
p_{0}=\frac{1}{2 h}\left((2 \alpha+1) u_{i-2}-4(\alpha+1) u_{i-1}+(2 \alpha+3) u_{i}\right),  \tag{4.66}\\
p_{1}=\frac{1}{2 h}\left((2 \alpha-1) u_{i-1}+-4 \alpha u_{i}+(2 \alpha-1) u_{i+1}\right),  \tag{4.67}\\
p_{2}=\frac{1}{2 h}\left((2 \alpha-3) u_{i}-4(\alpha-1) u_{i+1}+(2 \alpha-1) u_{i+2}\right),  \tag{4.68}\\
c_{0}=\frac{1}{12}(\alpha-2)(\alpha-1), \quad c_{0}^{\prime}=(2 \alpha-3) /(12 h),  \tag{4.69}\\
\quad c_{1}=-\frac{1}{6}(\alpha-2)(\alpha+2), \quad c_{1}^{\prime}=-\alpha /(3 h),  \tag{4.70}\\
c_{2}=\frac{1}{12}(\alpha+1)(\alpha+2), \quad c_{2}^{\prime}=(2 \alpha+3) /(12 h),  \tag{4.71}\\
I S_{0}=\frac{13}{12}\left(\bar{u}_{i-2}-2 \bar{u}_{i-1}+\bar{u}_{i}\right)^{2}+\frac{1}{4}\left(\bar{u}_{i-2}-4 \bar{u}_{i-1}+3 \bar{u}_{i}\right)^{2}  \tag{4.72}\\
I S_{1}=\frac{13}{12}\left(\bar{u}_{i-1}-2 \bar{u}_{i}+\bar{u}_{i+1}\right)^{2}+\frac{1}{4}\left(\bar{u}_{i-1}-\bar{u}_{i+1}\right)^{2}  \tag{4.73}\\
I S_{2}=\frac{13}{12}\left(\bar{u}_{i}-2 \bar{u}_{i+1}+\bar{u}_{i+2}\right)^{2}+\frac{1}{4}\left(\bar{u}_{i}-4 \bar{u}_{i+1}+3 \bar{u}_{i+2}\right)^{2} . \tag{4.74}
\end{gather*}
$$

The $3^{\text {rd }}$ order scheme uses $\alpha=-1,0,1$ and $h=\Delta x / 2$; the $5^{\text {th }}$ order scheme uses $\alpha=-2,-1,0,1,2$ and $h=\Delta x / 4$.

### 4.2.3 The Relaxed CFL Condition

For the nonlinear problem there are 2 natural choices for the approximate trace velocity $\tilde{v}$. The first is $\tilde{v}=f^{\prime}(u)$ for some fixed value of $u$. Since $u$ is constant along characteristics, this choice should give a small error in the flux correction step. The nonlinear problem can admit solutions where the characteristics collide and if no CFL constraint is applied the scheme will collapse. The standard CFL constraint for a fixed mesh method limits the time step to the time when characteristics at grid point $x_{i}$ arrives at the next grid point $x_{i+1}$ or $x_{i-1}$ this time step is given as $\Delta t \leq \frac{\Delta x}{\max \left|f^{\prime}(u)\right|}$. The related problem $\left(u_{t}+(v u)_{x}\right)+(f(u)-v u)_{x}$ has a relaxed CFL condition

$$
\begin{equation*}
\Delta t \leq \Delta t_{\mathrm{CFL}}:=\frac{h}{\max \left|f^{\prime}(u)-v\right|} \tag{4.75}
\end{equation*}
$$

as illustrated in Figure 4.4 This CFL condition was identified by Stockie Mackenzie and Russell [73] where they used a velocity perturbation to solve their Moving Mesh Method. Later used by Arbogast and Huang in [46] note that this is the same condition needed for the flux corrected EulerianLagrangian schemes. This condition states that the time step is small enough that waves from neighboring cells don't intersect across tracelines.

A second choice for $\tilde{v}$ is the particle velocity field $f(u) / u$ where again $u(x, t)$ is fixed. This minimized the number of particles flowing across the boundary, and again should lead to a small flux correction.

The choice of $\tilde{v}$ depends on the solution of the given problem, however for the problems we ran numerical tests on there was no significant difference.


Figure 4.4: Relaxed CFL Constraint An illustration of relaxed CFL condtion, orange line is fixed grid location, blue line has slope $\max \left(f^{\prime}(u)\right)^{-1}$, green line has slope $v^{-1}$. The fixed grid CFL condition is when the wave reaches the next grid point. The relaxed CFL condition describes the time when the wave reaches the Lagrangian traceline.

The numerical results for this dissertation are presented with $\tilde{v}_{i}=f^{\prime}\left(u_{i}\right)$.

### 4.3 Numerical Results

We follow the numerical examples provided by Arbogast and Huang in [46] since this method differs only in the computation of the flux integral it provides the natural comparison for studying the changes to the scheme. In their paper they demonstrate that without computing the flux term, inaccurate particle tracing leads to both a lack of accuracy and poor performance at capturing shocks. In this section we will not reproduce results of this type, we will simply analyze the performance of this algorithm in comparison with their scheme. Numerical results will first be presented for the linear problem, followed by the non-linear problem. Error results are presented with discrete $L^{1}$ and $L^{\infty}$ norms as given in equations (3.24) and (3.25).


Figure 4.5: Shu's Linear Test Numerical solutions for Shu's Linear Test at time $t=2$ with 4 time steps. On the left is the $3^{\text {rd }}$ order scheme, and on the right is the $5^{\text {th }}$ order scheme.

We first present results to modeling the linear transport problem. This is useful for the following reasons, first it illustrates the scheme in a context that is simple to interpret. The second reason is that exactly computing tracelines while taking large time steps can be an expensive procedure. As a way to minimize the computational effort one can trace through the approximated velocity field and then compute the flux integral provided in this chapter.

### 4.3.0.1 Constant Speed Transport

We first test our scheme on the linear problem $u_{t}+u_{x}=0$ with analytic solution $u(x, t)=u_{0}(x-t)$. Since trace points are simply $x-\Delta t$ we randomly perturb these points with a uniform distribution in $[-0.2 \Delta x, 0.2 \Delta x]$, so that the flux across the traceline is nonzero. We apply this to the initial condition given by Shu's Linear Test as described in Section 3.2.4.4. The test is run on a mesh with 200 elements. The results at $T=2$ using 4 time steps shown in Figure 4.5 show no visible oscillations or diffusion.

| $N$ | Error $L^{1}$ | Order | Error $L^{\infty}$ | Order |
| :---: | :---: | :---: | :---: | :---: |
| $3^{\text {rd }}$ order scheme |  |  |  |  |
| 80 | 1.74 E-1 | - | 8.33 E-1 | - |
| 160 | $2.06 \mathrm{E}-3$ | 6.40 | 6.97 E-3 | 6.90 |
| 320 | $4.36 \mathrm{E}-5$ | 5.56 | 5.43 E-4 | 3.68 |
| 640 | 6.63 E-6 | 2.71 | $1.22 \mathrm{E}-4$ | 2.15 |
| 1280 | 4.73 E-7 | 3.81 | $1.20 \mathrm{E}-5$ | 3.34 |
| 2560 | $2.48 \mathrm{E}-8$ | 4.26 | $4.45 \mathrm{E}-7$ | 4.76 |
| 5120 | 1.50 E-9 | 4.05 | $1.39 \mathrm{E}-8$ | 5.00 |
| 10240 | $9.72 \mathrm{E}-11$ | 3.95 | 1.40 E-9 | 3.30 |
| LR |  | 4.21 |  | 3.96 |
| $5^{\text {th }}$ order scheme |  |  |  |  |
| 40 | 2.25 E-3 | - | $1.77 \mathrm{E}-3$ | - |
| 80 | $4.08 \mathrm{E}-5$ | 5.79 | $3.20 \mathrm{E}-5$ | 5.79 |
| 160 | $6.80 \mathrm{E}-7$ | 5.91 | $5.34 \mathrm{E}-7$ | 5.90 |
| 320 | $1.08 \mathrm{E}-8$ | 5.96 | 8.58 E-9 | 5.96 |
| 640 | 1.73 E-10 | 5.98 | $1.36 \mathrm{E}-10$ | 5.98 |
| 1280 | $3.48 \mathrm{E}-12$ | 5.63 | $2.62 \mathrm{E}-11$ | 2.37 |

Table 4.1: Time Dependent Transport Test 1D Results for the flux corrected EL-WENO method with $v=\sin (t)$.

### 4.3.0.2 Time Dependent Velocity

This example uses the velocity $v=\sin (t)$. The initial condition is again $u_{0}(x)=\sin (\pi x)$ for $x \in[0,2]$. The analytical solution to this problem is $u(x, t)=u_{0}(x+1+\cos (t))$. The trace points are evaluated using RK4 with a single time step. The error results at $t=4$ using $\Delta t=20 \Delta x$ are listed in Table 4.1 show high convergence rates. The $5^{\text {th }}$ order results here are comparable to those of the EL-WENO method of Arbogast and Huang who also see a converge rate at around 6 for this problem.

### 4.3.0.3 Space Dependent Velocity

For this test the velocity is given as $v=\sin (x)$ with an initial condition that $u_{0}(x)=1$ for $x \in[0,2 \pi]$. The analytical solution to this problem is

$$
\begin{equation*}
u(x, t)=\frac{\sin \left(2 \arctan \left(e^{-t} \tan (x / 2)\right)\right)}{\sin (x)} \tag{4.76}
\end{equation*}
$$

These results are presented in Table 4.2 at $t=1$ with $\Delta t=5 \Delta x$. The trace points are evaluated using a single time step of RK4. Numerical results confirm formally high order convergence. The results of the $5^{\text {th }}$ order scheme are similar in magnitude to the method of Arbogast and Huang.

### 4.3.0.4 Burgers' Equation

The final numerical example we give is for Burgers' equation $u_{t}+$ $\left(1 / 2 u^{2}\right)_{x}=0$. We first solve this equation for the initial condition

$$
u_{0}(x)= \begin{cases}1.0, & x \in[0.6,1.4]  \tag{4.77}\\ 0.5, & \text { otherwise }\end{cases}
$$

which develops shock and rarefaction waves. Numerical results are shown in Figure 4.6 on a grid with 100 points at time $t=0.5$ with $\Delta t=.02$. The $5^{\text {th }}$ order scheme has better resolution near the rarefaction wave than the $3^{\text {rd }}$ order scheme.

The second initial condition we use is $u(x, 0)=u_{0}(x)=0.75+0.25 \sin (\pi x)$ for $x \in[0,2]$. The mass forms a shock at time $t=4 / \pi$. Numerical results in Table 4.3 are computed before the shock forms at $t=1$ with $\Delta t=4 \Delta x$. The analytical solution is computed numerically by solving the fixed point problem

$$
\begin{equation*}
u(x, t)-(0.75+0.25 \sin (\pi \check{x}))=0 \tag{4.78}
\end{equation*}
$$

| $N$ | Error $L^{1}$ | Order | Error $L^{\infty}$ | Order |
| :---: | :---: | :---: | :---: | :---: |
| $3^{\text {rd }}$ order scheme |  |  |  |  |
| 20 | $9.63 \mathrm{E}-3$ | - | $5.11 \mathrm{E}-2$ |  |
| 40 | $5.25 \mathrm{E}-3$ | 0.87 | $4.41 \mathrm{E}-3$ | . 212 |
| 80 | $1.08 \mathrm{E}-3$ | 2.28 | $1.86 \mathrm{E}-3$ | 1.25 |
| 160 | $1.75 \mathrm{E}-4$ | 2.62 | $5.65 \mathrm{E}-4$ | 1.72 |
| 320 | 2.61 E-5 | 2.74 | 1.52 E-4 | 1.89 |
| 640 | 3.33 E-6 | 2.97 | $3.77 \mathrm{E}-5$ | 2.01 |
| 1280 | $4.23 \mathrm{E}-7$ | 2.98 | 7.54 E-6 | 2.32 |
| 2560 | $5.28 \mathrm{E}-8$ | 3.00 | $9.25 \mathrm{E}-7$ | 3.02 |
| $5^{\text {th }}$ order scheme |  |  |  |  |
| 20 | $9.40 \mathrm{E}-4$ | - | $6.94 \mathrm{E}-4$ | - |
| 40 | $1.47 \mathrm{E}-4$ | 2.67 | 1.60 E-4 | 2.11 |
| 80 | $1.01 \mathrm{E}-5$ | 3.87 | 1.75 E-5 | 3.19 |
| 160 | $4.25 \mathrm{E}-7$ | 4.57 | 7.45 E-7 | 4.55 |
| 320 | $1.51 \mathrm{E}-8$ | 4.81 | 2.61 E-8 | 4.83 |
| 640 | 5.05 E-10 | 4.91 | 8.64 E-10 | 4.92 |
| 1280 | 1.65 E-11 | 4.94 | 2.79 E-11 | 4.95 |

Table 4.2: Spatially Dependent Transport Test 1D Results for the flux corrected EL-WENO method with $v=\sin (x)$.


Figure 4.6: Burgers' Equation: Example 1 Numerical solutions for Burgers' equation at time $t=0.5$. On the left is the $3^{\text {rd }}$ order scheme, and on the right is the $5^{\text {th }}$ order scheme. The red is the analytical solution.

| $N$ | Error $L^{1}$ | Order | Error $L^{\infty}$ | Order |
| :---: | :---: | :---: | :---: | :---: |
| 40 $8.34 \mathrm{E}-4$ - $7.43 \mathrm{E}-3$ -  <br> 80 $1.51 \mathrm{E}-4$ 2.47 $1.56 \mathrm{E}-3$ 2.25  <br> 160 $1.91 \mathrm{E}-5$ 2.95 $2.96 \mathrm{E}-4$ 2.40  <br> 320 $2.22 \mathrm{E}-6$ 3.13 $6.79 \mathrm{E}-5$ 2.13  <br> 640 $2.90 \mathrm{E}-7$ 2.93 $1.09 \mathrm{E}-5$ 2.64  <br> $5^{\text {th }}$ order scheme      <br> 40 $2.13 \mathrm{E}-4$ - $1.89 \mathrm{E}-3$ -  <br> 80 $8.34 \mathrm{E}-6$ 4.67 $1.12 \mathrm{E}-4$ 4.08  <br> 160 $6.99 \mathrm{E}-7$ 3.58 $1.35 \mathrm{E}-5$ 3.05  <br> 320 $2.87 \mathrm{E}-8$ 4.60 $7.34 \mathrm{E}-7$ 4.21  <br> 640 $1.17 \mathrm{E}-9$ 4.62 $3.21 \mathrm{E}-8$ 4.51  |  |  |  |  |$.$

Table 4.3: Burgers' Equation Results for the flux corrected EL-WENO method.
for the trace location $\check{x}$ using a Newton iteration. This is solved at quadrature points $x$ on the element which are then used to compute the average value of $u$ used in the error equation. The code implements a 5 point Gaussian quadrature rule for this computation. The code is then tested after the shock formation at $t=2$ with $\Delta t=2 \Delta x$ the results in Figure 4.7 show that the method captures a sharp front at the shock.


Figure 4.7: Burgers' Equation: Example 2 Numerical solutions for Burgers' equation at time $t=1$. On the left is a grid with 40 elements and on the right there are 80 elements. The black dots are the solution of the $3^{\text {rd }}$ order scheme and the red circles are for the $5^{\text {th }}$ order scheme.

## Chapter 5

## QEL-WENO for Nonlinear Advection in 2D

In this chapter we will combine the ideas of the previous chapters in order to solve the 2 dimensional nonlinear conservation problem

$$
\begin{array}{r}
u_{t}+\nabla \cdot f(u)=0 \\
u(x, y, t=0)=u_{0}(x, y) \tag{5.2}
\end{array}
$$

where $f(u)=\left(f_{1}(u), f_{2}(u)\right)$.

### 5.1 The Space-Time Region, $\mathcal{E}(v)$

We define the approximate velocity field by taking an approximate characteristic tracing at the boundary nodes. To solve this we define an approximated velocity field $\tilde{v}(x, y, t)$, that is specific to the flux function $f(u)$. We apply a characteristic tracing, at the nodes of the grid $x_{i}, y_{j}$ through the chosen approximate velocity field $\tilde{v}(x, y, t)$ :

$$
\begin{align*}
& \left(\xi_{1}, \xi_{2}\right)_{t}=\left(\tilde{v}_{1}, \tilde{v}_{2}\right),  \tag{5.3}\\
& \xi\left(t^{n+1}\right)=\left(x_{i}, y_{j}\right) . \tag{5.4}
\end{align*}
$$

This tracing defines the traceback points $\left(\check{x}_{i}, \check{y}_{j}\right)=\xi\left(t^{n}\right)$. We define the velocity $v$ at the nodes as the constant value $v_{i, j}=\left(\frac{x_{i}-\check{x}_{i}}{\Delta t}, \frac{y_{j}-\check{y}_{j}}{\Delta t}\right)$. These are
used to define the velocity field over the element $E_{i j}=\left[x_{i}, x_{i+1}\right] \times\left[y_{j}, y_{j+1}\right]$ as the bilinear interpolation

$$
\begin{align*}
v\left(x_{i+\alpha}, y_{j+\beta}\right)= & (1-\beta)(1-\alpha) v_{i, j}+(1-\beta) \alpha v_{i+1, j} \\
& +\beta(1-\alpha) v_{i, j+1}+\alpha \beta v_{i+1, j+1} \tag{5.5}
\end{align*}
$$

for $(\alpha, \beta) \in[0,1]^{2}$. The velocity field on $E_{i j}$ has two components $v(x, y)=$ $\left(a_{1}(x, y), a_{2}(x, y)\right)$. This defines the space-time region $\mathcal{E}=\mathcal{E}(v)$ and traceback region $\check{E}=\check{E}(v)$ as a function of the velocity field.

### 5.2 The Flux Corrected Problem

We now solve the related problem

$$
\begin{align*}
& \left(u_{t}+\nabla \cdot\left(a_{1} u, a_{2} u\right)\right)+\nabla \cdot\left(f_{1}(u)-a_{1} u, f_{2}(u)-a_{2} u\right)=0  \tag{5.6}\\
& u(x, y, t=0)=u_{0}(x, y) \tag{5.7}
\end{align*}
$$

Integrating over $\mathcal{E}$ gives

$$
\begin{equation*}
\int_{\mathcal{E}} u_{t}+\nabla \cdot\binom{a_{1} u}{a_{2} u} d V+\int_{\mathcal{E}} \nabla \cdot\binom{f_{1}(u)-a_{1} u}{f_{2}(u)-a_{2} u} d V=0 . \tag{5.8}
\end{equation*}
$$

Applying the divergence theorem over the space time region this equation becomes

$$
\int_{E} u^{n+1} d A-\int_{\check{E}} u^{n} d A+\oint_{\partial \mathcal{E}}\left(\begin{array}{c}
f_{1}(u)-a_{1} u  \tag{5.9}\\
f_{2}(u)-a_{2} u \\
0
\end{array}\right) \cdot\left(\begin{array}{l}
\nu_{x} \\
\nu_{y} \\
\nu_{t}
\end{array}\right) d S
$$

The mass on cell $E$ at $t^{n+1}$ is given by

$$
\int_{E} u^{n+1} d A=\int_{\check{E}} u^{n} d A-\oint_{\partial \mathcal{E}}\left(\begin{array}{c}
f_{1}(u)-a_{1} u  \tag{5.10}\\
f_{2}(u)-a_{2} u \\
0
\end{array}\right) \cdot\left(\begin{array}{l}
\nu_{x} \\
\nu_{y} \\
\nu_{t}
\end{array}\right) d S .
$$

In chapter 3 we gave the algorithm to solve the integral over $\check{E}$ using a quadrature scheme. We remark that since the velocity field at every point is constant we do not need to apply an ODE solver to compute the trace points as they are simply $(\check{x}, \check{y})=(x, y)-v\left(x_{i+\alpha}, y_{j+\beta}\right) \Delta t$ for $(x, y) \in E_{i, j}$. The Jacobian can be computed directly using the linear trace and

$$
\begin{align*}
\frac{\partial v}{\partial x} & =\frac{1}{\Delta x}\left[(1-\beta)\left(v_{i+1, j}-v_{i, j}\right)+\beta\left(v_{i+1, j+1}-v_{i, j+1}\right)\right]  \tag{5.11}\\
\frac{\partial v}{\partial y} & =\frac{1}{\Delta y}\left[(1-\alpha)\left(v_{i, j+1}-v_{i, j}\right)+\alpha\left(v_{i+1, j+1}-v_{i+1, j}\right)\right] \tag{5.12}
\end{align*}
$$

to solve for the components $\frac{\partial \check{x}}{\partial y}, \frac{\partial \check{x}}{\partial x}, \frac{\partial \check{y}}{\partial y}$ and $\frac{\partial \check{y}}{\partial x}$ directly.
This leaves a method to solve for

$$
\oint_{\partial \mathcal{E}}\left(\begin{array}{c}
f_{1}(u)-a_{1} u  \tag{5.13}\\
f_{2}(u)-a_{2} u \\
0
\end{array}\right) \cdot\left(\begin{array}{c}
\nu_{x} \\
\nu_{y} \\
\nu_{t}
\end{array}\right) d S=\sum_{i=0}^{3} \int_{S_{i}}\left(\begin{array}{c}
f_{1}(u)-a_{1} u \\
f_{2}(u)-a_{2} u \\
0
\end{array}\right) \cdot\left(\begin{array}{c}
\nu_{x} \\
\nu_{y} \\
\nu_{t}
\end{array}\right) d S
$$

which will use the technique detailed in Chapter 4. However we will be explicit about this extension to higher dimensions. We note that the faces $E$ and $\check{E}$ have outward normals $\pm(0,0,1)$ and do not contribute to this sum, which is over the 4 remaining surfaces.

### 5.2.1 Surface Parameterization

The goal of this section is to define a parameterization of the Surface $S_{i}$ of the form

$$
\begin{equation*}
r(\gamma, t)=x(\gamma, t) \hat{\mathbf{i}}+y(\gamma, t) \hat{\mathbf{j}}+t \hat{\mathbf{k}} \tag{5.14}
\end{equation*}
$$

where $(\gamma, t) \in[0,1] \times\left[t^{n}, t^{n+1}\right]=D$. This allows us compute (5.13) as

$$
\int_{D}\left(\begin{array}{c}
f_{1}(u)-a_{1} u  \tag{5.15}\\
f_{2}(u)-a_{2} u \\
0
\end{array}\right) \cdot\left(\begin{array}{c}
\nu_{x(\gamma, t)} \\
\nu_{y(\gamma, t)} \\
\nu_{t}
\end{array}\right)\left\|r_{\gamma} \times r_{t}\right\| d A
$$

The outward unit normal on a surface with a parameterization $r(\gamma, t)$ is given by

$$
\begin{equation*}
\nu=\frac{r_{\gamma} \times r_{t}}{\left\|r_{\gamma} \times r_{t}\right\|} \tag{5.16}
\end{equation*}
$$

which gives the flux integral over $D$ as

$$
\int_{D}\left(\begin{array}{c}
f_{1}(u)-a_{1} u  \tag{5.17}\\
f_{2}(u)-a_{2} u \\
0
\end{array}\right) \cdot\left(r_{\gamma} \times r_{t}\right) d A
$$

We will solve this using quadrature for both $\gamma$ and $t$.
This parameterization is defined on each surface independently. Let $\left(x_{p}, y_{p}\right)$ and $\left(x_{q}, y_{q}\right)$ be adjacent nodes of $E_{i, j}$ We will consider the face described by these 2 points and their traceback points $\left(\check{x}_{p}, \check{y}_{p}\right)$ and $\left(\check{x}_{q}, \check{y}_{q}\right)$. The surface $r(\gamma, t)=(x(\gamma, t), y(\gamma, t), t)$ is described by

$$
\begin{align*}
x(\gamma, t)= & x_{p}(1-\gamma)+x_{q} \gamma \\
& -\left(\frac{\left(x_{p}-\check{x}_{p}\right)}{\Delta t}(1-\gamma)+\frac{\left(x_{q}-\check{x}_{q}\right)}{\Delta t} \gamma\right)\left(t^{n+1}-t\right),  \tag{5.18}\\
y(\gamma, t)= & y_{p}(1-\gamma)+y_{q} \gamma \\
& -\left(\frac{\left(y_{p}-\check{y}_{p}\right)}{\Delta t}(1-\gamma)+\frac{\left(y_{q}-\check{y}_{q}\right)}{\Delta t} \gamma\right)\left(t^{n+1}-t\right) \tag{5.19}
\end{align*}
$$

The partial derivatives of $r$ with respect to $\gamma$ and $t$ are

$$
\begin{align*}
r_{\gamma}= & \frac{\partial x(\gamma, t)}{\partial \gamma} \hat{\mathbf{i}}+\frac{\partial y(\gamma, t)}{\partial \gamma} \hat{\mathbf{j}}+\frac{\partial t}{\partial \gamma} \hat{\mathbf{k}} \\
= & \left(x_{q}-x_{p}\right)+\left(\frac{\left(x_{p}-\check{x}_{p}\right)}{\Delta t}-\frac{\left(x_{q}-\check{x}_{q}\right)}{\Delta t}\right)\left(t^{n+1}-t\right) \hat{\mathbf{\imath}} \\
& +\left(y_{q}-y_{p}\right)+\left(\frac{\left(y_{p}-\check{y}_{p}\right)}{\Delta t}-\frac{\left(y_{q}-\check{y}_{q}\right)}{\Delta t}\right)\left(t^{n+1}-t\right) \hat{\mathbf{j}} \\
& +0 \hat{\mathbf{k}} \tag{5.20}
\end{align*}
$$

and

$$
\begin{align*}
r_{t}= & \frac{\partial x(\gamma, t)}{\partial t} \hat{\mathbf{\imath}}+\frac{\partial y(\gamma, t)}{\partial t} \hat{\mathbf{j}}+\frac{\partial t}{\partial t} \hat{\mathbf{k}} \\
= & \left(\frac{\left(x_{p}-\check{x}_{p}\right)}{\Delta t}(1-\gamma)+\frac{\left(x_{q}-\check{x}_{q}\right)}{\Delta t} \gamma\right) \hat{\mathbf{i}} \\
& +\left(\frac{\left(y_{p}-\check{y}_{p}\right)}{\Delta t}(1-\gamma)+\frac{\left(y_{q}-\check{y}_{q}\right)}{\Delta t} \gamma\right) \hat{\mathbf{j}} \\
& +1 \hat{\mathbf{k}} . \tag{5.21}
\end{align*}
$$

Using these the flux integral over the surface is the integral

$$
\begin{equation*}
\int_{D}\left(f_{1}(u)-a_{1} u\right)\left(r_{\gamma} \cdot(0,1,0)\right)+\left(f_{2}(u)-a_{2} u\right)\left(r_{\gamma} \cdot(-1,0,0)\right) d A \tag{5.22}
\end{equation*}
$$

### 5.2.2 Quadrature

We solve (5.22) with a quadrature rule in space $\gamma$ and time $t$

$$
\begin{align*}
\sum_{k=1}^{n_{q}} \tilde{w}_{k} \sum_{l=1}^{n_{q}} w_{l}( & \left(f_{1}(u)-a_{1} u\right)\left(r_{\gamma} \cdot(0,1,0)\right) \\
& \left.+\left(f_{2}(u)-a_{2} u\right)\left(r_{\gamma} \cdot(-1,0,0)\right)\right) \tag{5.23}
\end{align*}
$$

where the function $f$ and the mass $u$ and the vector $r$ are functions of the quadrature points $\gamma_{k}$ and $t_{l}$. To compute this we need to know the value of $u\left(x\left(\gamma_{k}, t_{l}\right), y\left(\gamma_{k}, t_{l}\right), t_{l}\right)$. We compute this term using the strategy in Chapter 4. We note that for a fixed $\gamma$, the mass $u$ changes in time along the line $(x(t), y(t), t)$ as

$$
\begin{align*}
\frac{d u}{d t} & =\frac{\partial u}{\partial x} \frac{\partial x}{\partial t}+\frac{\partial u}{\partial y} \frac{\partial y}{\partial t}+\frac{\partial u}{\partial t} \\
& =u_{x} a_{1}+u_{y} a_{2}-f_{1}(u)_{x}-f_{2}(u)_{y} \tag{5.24}
\end{align*}
$$

This is solved using the NCE Runge Kutta Scheme in Chapter 4. The details of the IWENO scheme are presented in Chapter 2, taking the tensor product of the 1D polynomials given in Chapter 4. The smoothness indicators are defined in the standard way.

### 5.3 Numerical Results

### 5.3.0.1 Example: Swirling Deforming Flow

Two dimensional problem with a swirling and deforming velocity field as described in section 3.2.5.3 is applied. A Runge-Kutta scheme using a single time step is used to trace the element nodes. Figure 5.1 shows the numerical results for 10 times steps on a grid with 60 elements in both directions. The error results listed in Table 5.1 show that for coarse grids we get the similar error results as the linear QEL-WENO method. The formal convergence rates are not seen here possibly due to rounding error accumulation or a subtle bug in the code.


Figure 5.1: Flux Corrected Swirling and Deforming Flow Numerical results for a radial bump function in a swirling and deforming velocity field at various times with $\Delta t=0.1$.

| $N$ | Error $L^{1}$ | Order | Error $L^{\infty}$ | Order |
| :---: | :---: | :---: | :---: | :---: |
| 10 | $3.20 \mathrm{E}-3$ | - | $2.64 \mathrm{E}-2$ | - |
| 20 | $1.03 \mathrm{E}-3$ | 1.63 | $8.44 \mathrm{E}-3$ | 1.64 |
| 40 | $1.10 \mathrm{E}-4$ | 3.22 | $9.06 \mathrm{E}-4$ | 3.21 |
| 60 | $3.47 \mathrm{E}-5$ | 2.84 | $6.41 \mathrm{E}-4$ | 0.85 |
| 80 | $3.70 \mathrm{E}-5$ | -0.22 | $1.78 \mathrm{E}-3$ | -3.54 |

Table 5.1: Flux Corrected Swirling and Deforming Flow. Error and convergence order for transport with a time and spatially dependent velocity field at time $t=1$ with $\Delta t=0.1$.

### 5.3.0.2 Example: Burgers' Equation

The final example is the two dimensional Burgers' Equation

$$
\begin{equation*}
u_{t}+\left(u^{2} / 2\right)_{x}+\left(u^{2} / 2\right)_{y}=0, \quad(x, y) \in[0,2]^{2}, t>0 \tag{5.25}
\end{equation*}
$$

with the initial condition

$$
u(x, y, 0)= \begin{cases}\sin ^{2}(\pi x) \sin ^{2}(\pi y), & \text { for }(x, y) \in(0,1)^{2}  \tag{5.26}\\ 0, & \text { otherwise }\end{cases}
$$

The problem is solved using the wave velocity $v_{i j}=u_{i j}$ for the approximate tracing. The solution with 40 grid elements in each direction and $\Delta t=\Delta x$ is shown in Figure 5.2 at times $t=0,1,2$ and 3 . The solution shows the method captures the sharp front of the shock.

$$
t=0
$$



$$
t=2
$$



$$
t=1
$$



$$
t=3
$$



Figure 5.2: Flux Corrected 2-D Burgers' Solution at times $t=0,1,2$ and 3 using 40 grid elements in each direction and $\Delta t=\Delta x$.

## Chapter 6

## Two Phase Flow

One application of the Quadrature Eulerian Lagrangian method is to simulate the transport part of the incompressible two phase flow problem. There is a broad set of literature $[24,57,66,67,69,74,80,81]$ developing the physical equations and analyzing numerical methods to solve flow and transport in porous media. This chapter will provide an overview of the twophase flow problem and include details of our implementation.

Let $\alpha=w, n$ refere to the wetting and non-wetting phases of the fluid. The two-phase flow problem for the Darcy velocity $u_{\alpha}$, pressure $p_{\alpha}$ and saturation $s_{\alpha}$ is described by the partial differential equations

$$
\begin{align*}
u_{\alpha} & =-K_{a} \lambda_{r \alpha}\left(s_{\alpha}\right)\left(\nabla p_{\alpha}-\rho_{\alpha} g\right) & & \text { in } \Omega \times J,  \tag{6.1}\\
\partial_{t}\left(\phi s_{\alpha}\right)+\nabla \cdot u_{\alpha} & =q_{a} & & \text { in } \Omega \times J,  \tag{6.2}\\
u_{\alpha} \cdot \nu & =0 & & \text { on } \partial \Omega \times J,  \tag{6.3}\\
s_{n}(x, 0) & =s_{n}^{0}(x) & & \text { in } \Omega, \tag{6.4}
\end{align*}
$$

where $\Omega$ is the spatial domain, $J$ is the time domain, $K_{a}$ is the medium absolute permeability, Let $k_{r \alpha}$ is the relative permeability, $\mu_{\alpha}$ is the viscosity and $\lambda_{r \alpha}$
is the relative mobility defined as

$$
\begin{equation*}
\lambda_{r \alpha}=\frac{k_{r \alpha}}{\mu_{\alpha}} \tag{6.5}
\end{equation*}
$$

$q_{\alpha}$ is the external source/sink, $\nu$ is the normal, $\rho_{\alpha}$ is the density of the fluid and $g$ is the gravitational acceleration.

$$
\begin{equation*}
s_{w}+s_{n}=1 \tag{6.6}
\end{equation*}
$$

$p$ is the pressure, We also consider the effects of capillary pressure, which is modeled as

$$
\begin{equation*}
p_{c}\left(s_{w}\right)=p_{n}-p_{w} . \tag{6.7}
\end{equation*}
$$

### 6.1 Global Pressure Formulation

The global pressure formulation of the above flow equation was introduced for incompressible two phase flow in [1, 21] , and extended as a method to solve three phase flow in $[22,23]$.

### 6.1.1 Pressure Equation

To minimize the coupling between the pressure and saturation, we use a global pressure formulation. This global pressure is

$$
\begin{equation*}
p=p\left(s_{w}\right)=p_{0}-\int_{0}^{s_{w}} \frac{\lambda_{r w}}{\lambda_{r t}} p_{c}^{\prime} d s \tag{6.8}
\end{equation*}
$$

We also define a total velocity, volumetric injection rate and mobility as

$$
\begin{gather*}
u_{t}=u_{n}+u_{w},  \tag{6.9}\\
q_{t}=q_{n}+q_{w},  \tag{6.10}\\
\lambda_{r t}=\lambda_{r n}+\lambda_{r w} . \tag{6.11}
\end{gather*}
$$

By summing each phase in Darcy's Law (6.1) and the Mass Conservation (6.2) we get an elliptic pressure equation in mixed form

$$
\begin{align*}
\nabla \cdot u_{t} & =q_{t} & \text { in } \Omega \times J  \tag{6.12}\\
u_{t} & =-K_{a} \lambda_{r t}\left(\nabla p-\frac{\lambda_{r w} \rho_{w}+\lambda_{r n} \rho_{n}}{\lambda_{r t}} g\right) & \text { in } \Omega \times J \tag{6.13}
\end{align*}
$$

### 6.1.2 Saturation Equation

We solve for the wetting phase $\alpha=w$ and use (6.6) to solve for the nonwetting phase. Combining the mass conservation equation (6.2) with Darcy's Law (6.1) we derive the saturation equation as

$$
\begin{equation*}
\partial_{t}\left(\phi s_{w}\right)+\nabla\left(\frac{\lambda_{r w}}{\lambda_{r t}} u_{t}\right)+\nabla \cdot\left[a \frac{\lambda_{r n} \lambda_{r w}}{\lambda_{r t}}\left(\nabla p_{c}+\left(\rho_{w}-\rho_{n}\right) g\right)\right]=q_{w} \tag{6.14}
\end{equation*}
$$

### 6.2 Discretization

To discretize the problem we choose $p, u_{t}$, and $s_{w}$ as the primary variables for solving equations (6.1)-(6.4). We use an implicit pressure explicit saturatuation (IMPES) scheme for solving the two phase problem.

### 6.2.1 Discretization of Flow

Mortar mixed finite elements are a well-known technique used to solve the flow problem. In $[58,59]$ a fully implicit formulation is used to solve a non-linear interface problem. This work was extended in [11, 37, 38, 83, 84] where preconditioning techniques are developed to aid in solving the problem. We use these types of techniques, and use code developed by Hailong Xiao; details of his method and can be found in [82].

For the examples provided we use Raviart-Thomas [64] spaces of lowest order (RT0). The space $\left(V_{h}, W_{h}\right)$ is used to approximate

$$
\begin{array}{rlrl}
\left(\left(K_{a} \lambda_{r t}\right)^{-1} u_{t}, v\right)-(p, \nabla \cdot v) & =-\left(\frac{\lambda_{r w} \rho_{w}+\lambda_{r n \rho_{n}}}{\lambda_{r t}} g, v\right) & \forall v \in V_{h} \\
\left(\nabla \cdot u_{t}, w\right) & =\left(q_{t}, w\right) & & \forall w \in W_{h} \tag{6.16}
\end{array}
$$

### 6.2.2 Discretization of Transport

As discussed in Chapter 1, the two terms in the transport problem have very different behaviors. In order to solve this equation over a time step we use an operator splitting technique to solve for the hyperbolic part (advection) and parabolic part (diffusion) separately. We solve

$$
\begin{equation*}
\partial_{t}\left(\phi s_{w}\right)+\nabla\left(\frac{\lambda_{r w}}{\lambda_{r t}} u_{t}\right)=0 \tag{6.17}
\end{equation*}
$$

using the QEL-WENO techniques provided in Chapters 5. Then we use the solution as an initial condition to solve for

$$
\begin{equation*}
\partial_{t}\left(\phi s_{w}\right)+\nabla \cdot\left[a \frac{\lambda_{r n} \lambda_{r w}}{\lambda_{r t}}\left(\nabla p_{c}+\left(\rho_{w}-\rho_{n}\right) g\right)\right]=q_{w} . \tag{6.18}
\end{equation*}
$$

In order to deal with the degenerate diffusion coefficient $\alpha(s)=K_{a} \lambda_{r n} \lambda_{r w} / \lambda_{r t}$ we use a mixed finite element method $[10,25,68]$ computed with a rectangle quadrature rule, or cell centered finite differencing. This method is given as

$$
\begin{align*}
\left(\phi s^{n+1}, w\right)+(\Delta t \nabla \cdot \gamma, w) & =\left(\phi s^{n}, w\right)  \tag{6.19}\\
(\gamma, v)_{Q} & =\left(\alpha(s)\left[p_{c}^{\prime}(s) \hat{\gamma}+\left(\rho_{w}-\rho_{n}\right) g \nabla z\right], v\right)_{Q},  \tag{6.20}\\
(\hat{\gamma}, \omega)_{Q} & =-(s, \nabla \cdot \omega)_{Q} \tag{6.21}
\end{align*}
$$

### 6.3 Numerical Example

We run a simulation for 7 days taking time steps of 0.02 [days]. The computational domain is $(x, y) \in[0,35] \times[0,35][\mathrm{m}]$ solved on a uniform mesh with 35 elements in each spatial dimension. The initial condition reflects gravitation equilibrium with initial saturation set at 0.4 at a depth 35 [m]. The boundary has a no flow condition. The inflow source term is on the elements with $(x, y) \in[0,1] \times[0,32][\mathrm{m}]$ and the outflow is on $(x, y) \in[31,32] \times[0,32][\mathrm{m}]$ with respective flow rates of $\pm 10[/$ day $]$. The rock porosity is constant at 0.2 , and the permeability is given in Figure 6.1. This figure also shows the initial pressure and velocity field for the problem. The Brooks-Corey model [17, 18] is used to generate the capillary pressure and relative permeability curves

$$
\begin{align*}
p_{c} & =p_{e} s_{w}^{-1 / \lambda}  \tag{6.22}\\
k_{r w} & =s_{w}^{(2+3 \lambda) / \lambda}  \tag{6.23}\\
k_{r n} & =\left(1-s_{w}\right)^{2}\left(1-s_{w}^{(2+\lambda) / \lambda}\right), \tag{6.24}
\end{align*}
$$



Figure 6.1: Permeability and Flow On the left is the permeability. On the right is the initial pressure and velocity.
where $p_{e}$ is the entry pressure and $\lambda$ is the pore-size distribution index. For all of the numerical examples provided we take $p_{e}=0.5[\mathrm{psi}]$ and $\lambda=2$ as shown in Figure 6.2. The water and oil viscosities are $0.5[\mathrm{cp}]$ and $2.5[\mathrm{cp}]$ respectively. The densities of water and oil are respectively $1\left[\mathrm{~g} / \mathrm{cm}^{3}\right]$ and $0.7\left[\mathrm{~g} / \mathrm{cm}^{3}\right]$.

The numerical results of this simulation are shown in 6.3. These results show the fluid moving through the channels. The solution is not monotone, suggesting that a flux splitting technique might be needed for this problem. There are two improvements we suggest for future investigation. The first is to improve the velocity model to include a cross flow. This can be done using the higher order Raviart Thomas space RT1 or the Brezzi Douglas Marini space BDM1 [16]. The higher order velocity field will increase the accuracy of the particle tracing in the QEL-WENO scheme and allow for larger time steps to


Figure 6.2: Brooks-Corey Model Capillary pressure is on left and relative mobilities are on the right.
be taken.

The second improvement is to change the implementation of the source term. Currently the source term is implemented in the diffusion equation. This is a not a physical interpretation of the injection process. The source term should be included with the advection solve and implemented in the method described in the VCCMM where they trace the fluid forward out of the injection wells. The current implementation is subject to local oscillation around the wells since the QEL-WENO method can transport the fluid a larger spatial distance than where the injected fluid can reach with the diffusion solver.


Figure 6.3: Two-Phase Transport Solution The solution to the two-phase flow problem using the Flux Corrected QEL-WENO Scheme.

## Chapter 7

## Conclusions and Future Work

### 7.1 Conclusions

Eulerian-Lagrangian algorithms are numerical methods solving advection equations. The goal in the development of these types of schemes is to utilize time steps significantly larger than the CFL limited time step of fixed mesh methods. The computational time of tracing the characteristics is mitigated since the problems naturally can be decomposed to solve different cell blocks on multiple processors. These schemes have recently been combined with the reconstruction techniques of WENO schemes which provide a high order accuracy.

In order to aid in the development of these schemes we have presented the new IWENO interpolation technique. This interpolation procedure allows for a high order reconstruction of the function at any point in the rectangular domain. The key idea in this method is to reconstruct the primitive function and the differentiate. These polynomials are weighted in the standard WENO way so that the interpolation avoids using a stencil that contains a shock in the data. These schemes achieve the theoretical order of accuracy.

This interpolation technique aided the development of two new strate-
gies for solving the advection problem. The first is to solve the linear advection problem by using a quadrature rule to integrate the traceback region. This allows for reconstruction of the mass in the traceback region without out resolving its boundary. This method fails to be either mass or volume conservative, however criteria was developed to choose a sufficient number of quadrature points to ensure that the volume error was within an acceptable tolerance. Numerical results demonstrate the formally high order convergence rate for smooth problems.

The second application of IWENO is for the nonlinear problem which is solved using a flux correction technique. The IWENO technique is used to compute the integral of the flux across the approximated boundary. This technique was implemented in one spatial dimension. Numerical results were presented demonstrating its performance. The numerical results are very similar to the previous technique of Arbogast and Huang. The compact stencil improves the ability for the scheme to be run across multiple processors.

The final addition of this thesis is to combine the quadrature scheme with the flux-correction scheme in higher spatial dimensions. This is done through approximate characteristic tracing at element nodes to define a bilinear velocity field which is used in the quadrature scheme. Numerical results were presented.

### 7.2 Future Work

### 7.2.1 Conservation

The major limitation of the QEL-WENO method is that it fails to be conservative. There is an underlying algebraic system to the scheme $\bar{u}^{n+1}=$ $M \bar{u}^{n}$, where $M$ is a matrix with coefficients that depend on the scheme. It has been shown that the scheme is volume conservative is the column sum of $M=1$, and mass conservative if the row sum of $M=1$. Future research would include an investigation of methods to adjust the resulting algebraic system in a way to improve the mass conservation of the scheme.

### 7.2.2 CWENO3

Central WENO schemes are important fixed mesh methods since they allow for mass transport without an artificial flux splitting technique. However due to technical reasons these schemes did not exist for odd orders until recently due to challenges in reconstructing the midpoint. The IWENO technique presented here gives a way around this problem. In addition the idea of using a compact stencil to compute the flux across the edge has the potential of providing a less diffusive solution.

### 7.2.3 Applications

There are multiple applications in reservoir simulation such as unstable miscible displacement that the QEL-WENO algorithm can be applied to. Another direction for future research is testing the performance on systems of
coupled conservation laws such as the Euler Equations.

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