## Lecture 1: Introduction to bifurcation analysis

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University of Oxford

May 29

## Can you conduct an experiment twice ...

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Axial displacement test of an Embraer aircraft stiffener.

## Can you conduct an experiment twice

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Two different, stable configurations.

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A PDE with two unknown solutions

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Start from some initial guess

When a problem has multiple solutions, it is usually crucial.


We converge to one solution, our prediction

When a problem has multiple solutions, it is usually crucial.


But nature has chosen another (unknown) solution!

When a problem has multiple solutions, it is usually crucial.

We have encountered unexpected multiple solutions in both simple and complex configurations in computational fluid dynamics (CFD); this phenomenon is both extremely important and not well understood. It has serious implications for the use of CFD as a predictive tool.

- Venkat Venkatakrishnan

Computational Aerodynamic Optimization
Boeing Research \& Technology

## Section 2

## Scope

## Mathematical formulation

Compute the multiple solutions $u^{\star}$ of a stationary nonlinear equation

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\begin{gathered}
F\left(u^{\star}, \lambda\right)=0 \\
F \in C^{1}(X \times \mathbb{R}, Y)
\end{gathered}
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as a function of a parameter $\lambda \in \mathbb{R}$.

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## Case \#1: aircraft stiffener

$u^{\star}$ displacement, $\lambda$ loading, $F$ hyperelasticity

## Case \#2: aircraft wing

$u^{\star}$ velocity and pressure, $\lambda$ angle of attack, $F$ Navier-Stokes

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with the approximation property

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\lim _{v \rightarrow 0} \frac{\left\|F(u+v, \lambda)-F(u, \lambda)-F_{u}(u, \lambda) v\right\|}{\|v\|}=0 \text { for all } v \in X
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Moreover $F_{u}$ is continuous. The same holds for $F_{\lambda}$.

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## Warning

We (usually) can't guarantee to find all solutions. But finding many is better than finding one.

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## Lecture 3

Deflation techniques for computing disconnected bifurcation diagrams.

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## Goal for the course

Develop practical numerical methods for computing multiple solutions of fine discretisations of nonlinear BVPs.

## Example: Liouville-Bratu-Gelfand problem

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Solutions of the Bratu problem


## Example: Carrier's problem

$$
\lambda^{2} u^{\prime \prime}+2\left(1-x^{2}\right) u+u^{2}-1=0, \quad u(-1)=0=u(1)
$$





Solutions of $\lambda^{2} u^{\prime \prime}+2\left(1-x^{2}\right) u+u^{2}-1=0$


## Section 3

## Great Theorems of Nonlinear Functional Analysis

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Primary references.

## Subsection 1

## Newton-Kantorovich

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First, let's recall Newton's method in $\mathbb{R}$ and $\mathbb{R}^{N}$.

Core idea: solve succession of linearised rootfinding problems.


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solve $f^{\prime}\left(x_{k}\right) \delta x_{k}=-f\left(x_{k}\right) ;$ update $x_{k+1}=x_{k}+\delta x_{k}$.

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## Good local convergence

If $f$ is smooth, the solution is isolated, and the guess close, Newton converges quadratically.

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Consider the Taylor expansion of $f$ around $x_{k}$ :

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f\left(x_{k}+\delta x_{k}\right)=f\left(x_{k}\right)+f^{\prime}\left(x_{k}\right) \delta x_{k}+\mathcal{O}\left(\delta x_{k}^{2}\right) .
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and find $\delta x$ such that $f\left(x_{k}+\delta x\right) \approx 0$ :

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This yields

$$
\delta x_{k}=\left[f^{\prime}\left(x_{k}\right)\right]^{-1} f\left(x_{k}\right)
$$

This naturally extends to $F \in C^{1}\left(\mathbb{R}^{N} ; \mathbb{R}^{N}\right)$. Newton's method is to solve $F_{x}\left(x_{k}\right) \delta x_{k}=-F\left(x_{k}\right) ;$ update $x_{k+1}=x_{k}+\delta x_{k}$, where $F_{x}$ is the Jacobian of $F$.

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For the iteration to be well-defined, we need $F_{x}\left(x_{k}\right)$ to be invertible.

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Now imagine that we change units or coordinate systems for our outputs $F$. Instead of solving $F(x)=0$, we want to solve $\tilde{F}(x)=A F(x)=0$, where $A \in \mathbb{R}^{N \times N}$ is constant and nonsingular. Of course, this doesn't change the roots $x^{\star}$.

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## Theorem (Affine covariance)

Premultiplying $F$ by a constant nonsingular $A \in \mathbb{R}^{N \times N}$ does not change the Newton sequence.

Let $\tilde{F}(x):=A F(x)$. Newton's method applied to $\tilde{F}$ from $x_{0}=\tilde{x}_{0}$ generates a sequence

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\tilde{x}_{0}, \tilde{x}_{1}, \tilde{x}_{2}, \ldots
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## Proof.

For $i=0$, we have $x_{i}=\tilde{x}_{i}$ by assumption.

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Assume the claim is true at iteration $i$. Then the Newton update for $\tilde{F}$ satisfies

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\left[\tilde{F}_{x}\left(\tilde{x}_{i}\right)\right]^{-1} \tilde{F}\left(\tilde{x}_{i}\right)=\left[A F_{x}\left(x_{i}\right)\right]^{-1} A F\left(x_{i}\right)
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Hence $x_{i+1}=\tilde{x}_{i+1}$, and the result follows by induction.
We get exactly the same iterates $x_{0}, x_{1}, \ldots$, whether we apply Newton to $F(x)=0$ or $A F(x)=0$.

## Why does this matter?

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Moreover, any sensible strategy for globalising the convergence of Newton's method from poor initial guesses $x_{0}$ must also preserve this property. This insight leads to the current state of the art for globalising Newton's method.


Peter Deuflhard, 1944-2019

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We could also think of this as a problem in $\mathbb{R}^{2}$.
We know this has three solutions,

$$
z=1, \quad z=-1 / 2+i \sqrt{3} / 2, \text { and } z=-1 / 2-i \sqrt{3} / 2
$$

Let's take a subset of the complex plane and colour each point as follows. For a given $z_{0} \in \mathbb{C}$, we

1. run Newton's method with that initial guess,
2. and colour the point according to which root it converges to.


The Newton fractal for $z^{3}-1=0$.


The Newton fractal for $z^{3}-2 z+2=0$.

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With a good initial guess, and great cleverness, it is possible to devise computer-assisted proofs of the existence of solutions to infinite-dimensional nonlinear problems.

- Invented linear programming (via industrial consultancy!).
- Instrumental in saving over a million lives during the siege of Leningrad.
- Involved in the Soviet nuclear bomb project.
- Nearly sent to the gulag for "shadow prices".
- Pseudo-Nobel prize in Economics (1975).


Leonid Kantorovich, 1912-1986

## Theorem (Kantorovich (1948))

Let $F \in C^{1}(\Omega, Y)$ for open convex $\Omega \subset X$. Given $u_{0} \in \Omega$, assume 1. $F_{u}\left(u_{0}\right)^{-1}$ exists and set $\alpha:=\left\|F_{u}\left(u_{0}\right)^{-1} F\left(u_{0}\right)\right\|$;

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Then the Newton sequence defined by

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u_{k+1}=u_{k}-F_{u}\left(u_{k}\right)^{-1} F\left(u_{k}\right)
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is well defined and remains within $\overline{B\left(u_{0}, \rho_{0}\right)}$.

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There exists $u^{\star} \in \overline{B\left(u_{0}, \rho_{0}\right)}$ which solves $F\left(u^{\star}\right)=0$, and $\left(u_{k}\right) \rightarrow u^{\star}$.

## Theorem (Kantorovich (1948))

Let $F \in C^{1}(\Omega, Y)$ for open convex $\Omega \subset X$. Given $u_{0} \in \Omega$, assume

1. $F_{u}\left(u_{0}\right)^{-1}$ exists and set $\alpha:=\left\|F_{u}\left(u_{0}\right)^{-1} F\left(u_{0}\right)\right\|$;
2. $\left\|F_{u}\left(u_{0}\right)^{-1}\left(F_{u}(v)-F_{u}(w)\right)\right\| \leq \omega_{0}\|v-w\|$ for all $v, w \in \Omega$;
3. $h_{0}:=\alpha \omega_{0} \leq \frac{1}{2}$;
4. $\overline{B\left(u_{0}, \rho_{0}\right)} \subset \Omega$ for $\rho_{0}:=\left(1-\sqrt{1-2 h_{0}}\right) / \omega_{0}$.

Then the Newton sequence defined by

$$
u_{k+1}=u_{k}-F_{u}\left(u_{k}\right)^{-1} F\left(u_{k}\right)
$$

is well defined and remains within $\overline{B\left(u_{0}, \rho_{0}\right)}$.
There exists $u^{\star} \in \overline{B\left(u_{0}, \rho_{0}\right)}$ which solves $F\left(u^{\star}\right)=0$, and $\left(u_{k}\right) \rightarrow u^{\star}$.
The solution $u^{\star}$ is unique in $\Omega \cap B\left(u_{0}, \rho^{+}\right)$for a $\rho^{+}>\rho_{0}$.

## Subsection 2

## Rall-Rheinboldt

The Newton-Kantorovich theorem is very powerful because you only need to check conditions on the initial guess (and a ball around it).

The Newton-Kantorovich theorem is very powerful because you only need to check conditions on the initial guess (and a ball around it).

If you assume the existence of roots, one gets a slightly different theory that is also useful. This allows us to place balls around the roots, such that if the Newton sequence starts within a ball, Newton's method converges to the associated root.

## Theorem (Rall-Rheinboldt (1974))

Let $F \in C^{1}(\Omega, Y)$ for open convex $\Omega \subset X$. Let $u^{\star} \in \Omega$ such that $F\left(u^{\star}\right)=0$. Assume that

1. $F_{u}\left(u^{\star}\right)^{-1}$ exists;


Louis B. Rall, 1930-


Werner C. Rheinboldt, ?-?

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Then for any $u_{0} \in B\left(u^{\star}, 2 /\left(3 \omega^{\star}\right)\right)$, the Newton sequence is well-defined and remains within the ball.

The Newton sequence converges to $u^{\star}$.


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The solution $u^{\star}$ is unique within $\Omega \cap B\left(u^{\star}, 1 / \omega^{\star}\right)$.


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Let's examine the conditions of the theorems for a simple case:

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f: \mathbb{C} \rightarrow \mathbb{C}, \quad f(z)=(z-1)(z+1)
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For $z_{0} \neq 0$, we calculate

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\alpha:=\left|f^{\prime}\left(z_{0}\right)^{-1} f\left(z_{0}\right)\right|=\left|z_{0}^{2}-1\right| / 2\left|z_{0}\right|
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we find $\omega_{0}=1 /\left|z_{0}\right|$.
We need $\alpha \omega_{0} \leq 1 / 2$, so Newton-Kantorovich guarantees convergence for

$$
\frac{\left|z_{0}^{2}-1\right|}{2\left|z_{0}\right|^{2}} \leq \frac{1}{2} \Longrightarrow\left|1-z_{0}^{-2}\right| \leq 1
$$

## Rall-Rheinboldt

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## Subsection 3

## The Implicit Function Theorem

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When does the existence of $\left(u_{0}, \lambda_{0}\right)$ such that $F\left(u_{0}, \lambda_{0}\right)=0$ imply that we can solve $F$ for nearby values of $\lambda$ ?

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## An answer

... is given by the Implicit Function Theorem.

Basically, if $F_{u}\left(u_{0}, \lambda_{0}\right)$ is invertible, then you can continue $u=H(\lambda)$ for some interval $\left(\lambda_{0}-\delta, \lambda_{0}+\delta\right)$.

## Theorem (Implicit Function Theorem)

Assume that $\Omega \subset X \times \mathbb{R}$ is open. Let $F \in C^{1}(\Omega, Y)$.
Let $\left(u_{0}, \lambda_{0}\right) \in \Omega$ such that $F\left(u_{0}, \lambda_{0}\right)=0$ with $F_{u}\left(u_{0}, \lambda_{0}\right)$ invertible.

## Then



Ulisse Dini, 1845-1918

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## Then



1. there exist $\varepsilon, \delta>0$ and $H \in C\left(B\left(\lambda_{0}, \delta\right), B\left(u_{0}, \varepsilon\right)\right)$ such that $(H(\lambda), \lambda)$ is the unique solution of $F(u, \lambda)=0$ in $B\left(\lambda_{0}, \delta\right) \times B\left(u_{0}, \varepsilon\right)$;

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$F(u, \lambda)=0$ in $B\left(\lambda_{0}, \delta\right) \times B\left(u_{0}, \varepsilon\right)$;
2. if $F \in C^{k}(\Omega, Y)$, then $H \in C^{k}\left(B\left(\lambda_{0}, \delta\right), X\right)$;
3. if $F$ is analytic, $H$ is analytic.

The history is reviewed in

# A Historical Outline of the Theorem of Implicit Functions 

Un Bosquejo Histórico del Teorema de las Funciones Implícitas
Giovanni Mingari Scarpello (giovannimingari@libero.it)
Daniele Ritelli dritelli@economia.unibo.it
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which complains
Anglo-Saxon scientific and historic literature ignores the Italian mathematician U. Dini.

## Main message

If we want to find where local uniqueness breaks down, look for $(u, \lambda)$ such that $F_{u}(u, \lambda)$ not invertible.

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If we want to find where local uniqueness breaks down, look for $(u, \lambda)$ such that $F_{u}(u, \lambda)$ not invertible.

## Note

$F_{u}(u, \lambda)$ invertible is sufficient for the existence of a local resolution $u=u(\lambda)$, but not necessary.

## Consider $F(u, \lambda)=u^{3}-\lambda$.

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$F_{u}(0,0)=0$, but the resolution $u=H(\lambda)=\sqrt[3]{\lambda}$ is unique regardless.

## Section 4

## Examples

Let's see more examples of what can happen when the IFT does not apply.

## Fold bifurcation

$$
F(u, \lambda)=\lambda-u^{2}=0
$$

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$$
F(u, \lambda)=\lambda-u^{2}=0
$$

This has solutions

$$
u= \pm \sqrt{\lambda}, \quad \lambda \geq 0
$$

and no solutions otherwise.

## Fold bifurcation

$$
F(u, \lambda)=\lambda-u^{2}=0
$$


$F_{u}(0,0)=0$. A branch of solutions is born at a fold bifurcation.

## Transcritical bifurcation

$$
F(u, \lambda)=\lambda u+u^{2}=0
$$

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$$
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$$

This has solutions

$$
u=0, \quad u=-\lambda
$$

for all values of $\lambda$.

## Transcritical bifurcation

$$
F(u, \lambda)=\lambda u+u^{2}=0
$$



Two branches cross at a transcritical bifurcation.

## Pitchfork bifurcation

$$
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$$

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$$
F(u, \lambda)=\lambda u-u^{3}=0
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This has solutions

$$
\begin{array}{ll}
u=0, & \lambda \in \mathbb{R} \\
u= \pm \sqrt{\lambda}, & \lambda \geq 0
\end{array}
$$

## Pitchfork bifurcation

$$
F(u, \lambda)=\lambda u-u^{3}=0
$$



Two branches emerge from the base branch at a pitchfork bifurcation.

## Structural stability of folds

## Fold bifurcations are structurally stable.

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## Structural stability of transcritical and pitchfork bifurcations

Transcritical and pitchfork bifurcations are not.

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Fold bifurcations are structurally stable.

## Structural stability of transcritical and pitchfork bifurcations

Transcritical and pitchfork bifurcations are not.

## Numerical implications

This will have major consequences for our algorithms.

## Perturbing a fold + transcritical bifurcation

$$
F(u, \lambda)=u^{2}-\lambda^{2}(\lambda+1)+\delta=0
$$

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Salutions of $\lambda u-u^{3}=0$

$\delta=0$


$$
\delta<0
$$

## Perturbing a pitchfork bifurcation

$$
F(u, \lambda)=\lambda u-u^{3}+\delta=0
$$



These examples motivate the following definition.

## Bifurcation point

A bifurcation point $P=\left(u^{\star}, \lambda^{\star}\right)$ is one where, for all neighbourhoods $N$ containing $P$, there exists a $\lambda \in \mathbb{R}$ such that $F(u, \lambda)=0$ has nonunique solutions within $N$.

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## Bifurcation point

A bifurcation point $P=\left(u^{\star}, \lambda^{\star}\right)$ is one where, for all neighbourhoods $N$ containing $P$, there exists a $\lambda \in \mathbb{R}$ such that $F(u, \lambda)=0$ has nonunique solutions within $N$.

In the next lectures, we will study the key question:

How do we compute these bifurcation diagrams?

Lecture 2: Classical algorithms of bifurcation analysis

Patrick E. Farrell



University of Oxford

May 30

## Challenge

How do we continue branches? How do we detect and pursue bifurcations?


## Primary references.

Basic idea of numerical bifurcation analysis:

## procedure ANALYSE $\left(u_{0}, \lambda_{0}\right)$

## end procedure



Basic idea of numerical bifurcation analysis:
procedure ANALYSE $\left(u_{0}, \lambda_{0}\right)$ continue branch of solutions;

## end procedure



## Continuation

Extending our knowledge of the branch to other values of $\lambda$.

Basic idea of numerical bifurcation analysis:
procedure ANALYSE $\left(u_{0}, \lambda_{0}\right)$ continue branch of solutions; detect bifurcations on the branch;

## end procedure



## Bifurcation detection

Discovering when a bifurcation has occurred on the branch.

Basic idea of numerical bifurcation analysis:
procedure ANALYSE $\left(u_{0}, \lambda_{0}\right)$ continue branch of solutions; detect bifurcations on the branch; localise bifurcations;

## end procedure



## Bifurcation localisation

Identifying precisely the bifurcation point.

Basic idea of numerical bifurcation analysis:
procedure ANALYSE $\left(u_{0}, \lambda_{0}\right)$ continue branch of solutions; detect bifurcations on the branch; localise bifurcations; switch branches at bifurcations, and recurse. end procedure


Herbert Keller, 1925-2008

## Branch switching

Constructing the emanating branches, and analysing them recursively.


Start with $\left(u_{0}, \lambda_{0}\right)$.


Perform a continuation step.


Detect we have passed a bifurcation.


Localise bifurcation point.


Switch branches.


Apply recursively.

## Section 1

## Continuation algorithms

Suppose we know $\left(u_{0}, \lambda_{0}\right)$, with $F_{u}\left(u_{0}, \lambda_{0}\right)$ invertible. By the IFT we know we can continue the branch for other values of $\lambda$.

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How should we do so? We will meet five algorithms:

- natural (or naïve, or first-order) continuation;

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How should we do so? We will meet five algorithms:

- natural (or naïve, or first-order) continuation;
- tangent (or second-order) continuation, and secant continuation;
- arclength continuation, and pseudo-arclength continuation.


## Subsection 1

## Natural continuation



Start with $\left(u_{0}, \lambda_{0}\right)$.


Set $u_{0}$ as our guess for $\lambda_{0}+\delta \lambda$.


Use Newton-Kantorovich to find the solution for $\lambda_{1}=\lambda_{0}+\delta \lambda$.


Piecewise-constant guess.


Newton-Kantorovich.


Guess and solve.


Guess ...

... but there are no solutions to be found for this value of $\lambda$.

## Good news

## This is cheap and easy.

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## Bad news

We can probably construct better guesses.

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## Bad news

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## Worse news

The algorithm has no hope of continuing around the fold.

## Subsection 2

## Tangent and secant continuation

## Natural continuation estimates

$$
u\left(\lambda_{i+1}\right) \approx u\left(\lambda_{i}\right)
$$

which is the first-order Taylor expansion.

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the second-order Taylor expansion.
How do we compute $u_{\lambda}\left(\lambda_{i}\right)$ ?

Since $F(u, \lambda)=0$, taking the total derivative of both sides with respect to $\lambda$ in the direction $\delta \lambda$ yields

$$
\frac{\mathrm{d}}{\mathrm{~d} \lambda} F(u, \lambda)=F_{u}(u, \lambda) u_{\lambda}+F_{\lambda}(u, \lambda)=0 .
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If $\lambda \in \mathbb{R}, u \in \mathbb{R}^{N}$, then $u_{\lambda} \in \mathbb{R}^{N}, F_{\lambda} \in \mathbb{R}^{N}$, and $F_{u} \in \mathbb{R}^{N \times N}$.

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Since the dependence of $F$ on $\lambda$ is explicit, we can calculate $F_{\lambda}$, and solve

$$
F_{u}(u, \lambda) u_{\lambda}=-F_{\lambda}(u, \lambda)
$$

at the cost of one Newton step. This is the tangent linearisation.


Start with $\left(u_{0}, \lambda_{0}\right)$.


Solve tangent linearisation to construct next guess.


Solve nonlinear problem with Newton-Kantorovich.


Solve tangent linearisation to construct next guess.

This constructs much better initial guesses, but is more expensive. We have to save at least two Newton iterations to make this worth it.

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A natural alternative is to approximate the tangent with a secant: build the line joining two previous points on the branch, and extrapolate to the next value of $\lambda$.

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Secant continuation constructs almost as good initial guesses, for almost no increase in cost over natural continuation (only memory).

## Subsection 3

## Arclength continuation

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The fundamental problem is one of parameterisation: we are thinking of our solution curve as

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## A better way

Parameterise the solution curve as

$$
(u(s), \lambda(s))
$$

where $s$ is the arclength on the curve, measured from $\left(u_{0}, \lambda_{0}\right)$.

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where $s$ is the arclength on the curve, measured from $\left(u_{0}, \lambda_{0}\right)$.

In other words, at each continuation step we will also solve for the next value of $\lambda$. This allows $\lambda$ to decrease as well as increase, to successfully traverse folds.

Since we are now solving for both $u$ and $\lambda$, we need to augment our system of equations with one more real-valued equation:

$$
A(u(s), \lambda(s)):=\left[\begin{array}{c}
F(u(s), \lambda(s)) \\
p(u(s), \lambda(s), s)
\end{array}\right]=\left[\begin{array}{l}
0 \\
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The choice arclength continuation makes is to have $p$ encode a desired change in distance:

$$
p(u, \lambda, s):=\left\|u-u_{i}\right\|^{2}+\left|\lambda-\lambda_{i}\right|^{2}-\left(s-s_{i}\right)^{2} .
$$



Start with $\left(u_{i}, \lambda_{i}\right)$.


Seek points on the curve that intersect $p(u, \lambda)=0$.


Solve nonlinear problem with Newton-Kantorovich.


## Repeat.



## Repeat.



## Repeat.



## Repeat.

## Good news

This allows us to robustly continue around folds.

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We now have to solve augmented systems with extra nonlinearity.

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## Worse news

The augmented system generically has two solutions!

## We attempt to guide Newton-Kantorovich to the solution we want by building a good initial guess.

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We compute $\left(u_{s}\left(s_{i}\right), \lambda_{s}\left(s_{i}\right)\right)$ by solving

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\frac{\mathrm{d}}{\mathrm{~d} s} A(u(s), \lambda(s))=0
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the tangent linearisation of the augmented system.
We then set the initial guess to be $\left(u\left(s_{i}\right)+u_{s}\left(s_{i}\right) \delta s, \lambda\left(s_{i}\right)+\lambda_{s}\left(s_{i}\right) \delta s\right)$.

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However, this doesn't always work: even with this good initial guess, Newton-Kantorovich can sometimes find the wrong (old) solution.

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We are free to choose the extra equation. So let's linearise it!

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## Pseudo-arclength continuation

Assuming that $X \subset L^{2}(\Omega)$, we can choose
$p(u, \lambda):=\left(u-u_{i}, u_{s}\left(s_{i}\right)\right)_{L^{2}(\Omega)}+\left(\lambda-\lambda_{i}\right) \lambda_{s}\left(s_{i}\right)-\left(s-s_{i}\right)$

This looks for points on the branch that are orthogonal (in the $L^{2}(\Omega) \times \mathbb{R}$ inner product) to the tangent, at a distance $s-s_{i}$ away.


Start with $\left(u_{i}, \lambda_{i}\right)$.


Construct the tangent to the curve.


Impose the orthogonality constraint.


Solve.


## Repeat.



Repeat.


Repeat.

## Section 2

## Bifurcation detection

## Challenge

We need some way to detect that we have passed through a bifurcation.

## Consider the problem

$$
F(u, \lambda)=-u^{\prime \prime}-\lambda u+u^{3}=0, \quad u(0)=0=u(\pi) .
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By the IFT, we know that bifurcations can only happen where its Fréchet derivative is singular. Its Fréchet derivative on the branch is

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$$

which has nonzero solutions for $v$ whenever $\lambda$ is an eigenvalue of the Dirichlet Laplacian:

$$
\lambda_{n}=n^{2}, \quad n \in \mathbb{N}
$$



The bifurcation diagram we aim to compute.


Start our continuation at $(u, \lambda)=(0,0)$.


Examine the eigenvalues of $F_{u}$ at this point.


Take a continuation step.


By chance we land on the bifurcation-Fréchet derivative is singular.


Take another continuation step.


Take another continuation step, stepping over the next bifurcation.

So how do we detect when we've continued past a bifurcation?

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## Idea A

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Recall that the determinant of a matrix is the product of its eigenvalues. So when one eigenvalue changes sign, the determinant changes sign.

## Good news

The determinant is easy to compute from an $L U$ factorisation:

$$
A=L U \Longrightarrow \operatorname{det}(A)=\operatorname{det}(L) \operatorname{det}(U)
$$

## Bad news

We usually can't afford to compute an $L U$ factorisation ...

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So we need another idea.

## Idea B

At each continuation step, compute a few (e.g. 10) eigenvalues.

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## Good news

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## Comment

This is the main choice in PDE-oriented codes (e.g. pde2path and BifurcationKit.jl).

## Section 3

## Bifurcation localisation

Our ultimate goal is to switch branches at bifurcation points. In order to do this, we'll need to locate the bifurcation points precisely.

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Apply bisection to the detection algorithm.
In other words, you know two points on the branch that straddle the bifurcation. At each iteration, cut the interval between them in half and keep the subinterval that contains the bifurcation.

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This is simple to implement (given a detector).

## Bad news

This only converges linearly, so finding many digits will take forever.

Here is an idea that will let us quickly localise (some) bifurcations to high precision.

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By the IFT, we know that a necessary condition for a bifurcation is that

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Rüdiger Seydel, 1947-

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F_{u}(u, \lambda) \text { is singular. }
$$

## Idea B: Seydel-Moore-Spence

Find $(u, v, \lambda) \in X \times X \times \mathbb{R}$ such that

$$
\begin{aligned}
F(u, \lambda) & =0 \\
F_{u}(u, \lambda) v & =0 \\
\|v\|^{2} & =1 .
\end{aligned}
$$



Gerald Moore, 1951-


## Comment

The Seydel-Moore-Spence system is highly nonlinear. However, it is easy to construct good initial guesses.

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The Fréchet derivative of the Seydel-Moore-Spence system is singular at other bifurcation points, so Newton-Kantorovich converges linearly.

## Good news

It's possible to construct other augmented systems for other kinds of bifurcations. You have to know what you're looking for, though ...

One last comment: if you want to find out how a bifurcation point varies as you vary another parameter $\mu \in \mathbb{R}$,

One last comment: if you want to find out how a bifurcation point varies as you vary another parameter $\mu \in \mathbb{R}$,
do pseudo-arclength continuation on the Seydel-Moore-Spence system!

## Section 4

## Branch switching

To learn how to switch branches at a bifurcation point, we need another Great Theorem of Nonlinear Functional Analysis.

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## Great Theorem of Nonlinear Functional Analysis.

## Lyapunov-Schmidt reduction $(1906,1908)$

Let $F\left(u_{0}, \lambda_{0}\right)=0$ with $F_{u}$ singular. Let

$$
d=\operatorname{dim} \operatorname{ker} F_{u}\left(u_{0}, \lambda_{0}\right) .
$$



Aleksandr Lyapunov, 1857-1918


Erhard Schmidt, 1876-1959

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$$

Near the bifurcation point, we can relate
Aleksandr Lyapunov, 1857-1918 solutions of $F \Longleftrightarrow$ solutions of $R$ where $R$ is a $d \times d$ algebraic system!


Erhard Schmidt, 1876-1959

For this section, we will make the following assumptions:

## Essential assumptions

- $F\left(u_{0}, \lambda_{0}\right)=0$;
- $A:=F_{u}\left(u_{0}, \lambda_{0}\right) \in L(X, Y)$ is Fredholm:
$\operatorname{dim} \operatorname{ker}(A)<\infty, \quad$ codim range $(A)<\infty ;$
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## Non-essential assumptions

- $X$ and $Y$ are Hilbert spaces;
- $\operatorname{ind}(A):=\operatorname{dim} \operatorname{ker}(A)-\operatorname{codim} \operatorname{range}(A)=0$.

Let $A^{*}: Y \rightarrow X$ be the associated adjoint operator. Construct

$$
\operatorname{ker}(A)=\operatorname{span}\left\{\phi_{1}, \ldots, \phi_{d}\right\}, \quad \operatorname{ker}\left(A^{*}\right)=\operatorname{span}\left\{\psi_{1}, \ldots, \psi_{d}\right\}
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where $\left\{\phi_{i}\right\}_{i}$ and $\left\{\psi_{i}\right\}_{i}$ are orthonormal bases.

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Then construct

$$
P x:=\sum_{i=1}^{d}\left(\phi_{i}, x\right)_{X} \phi_{i},
$$

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By construction,

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\operatorname{range}(P)=\operatorname{ker}(A), \quad \operatorname{range}(Q)=\operatorname{ker}\left(A^{*}\right)=\operatorname{range}(A)^{\perp}
$$

Then we can decompose

$$
\begin{aligned}
& X=\operatorname{range}(P) \oplus \operatorname{range}(I-P)=: X_{1} \oplus X_{2} \\
& Y=\operatorname{range}(Q) \oplus \operatorname{range}(I-Q)=: Y_{1} \oplus Y_{2} .
\end{aligned}
$$

Write

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u=P u+(I-P) u=: v+w, \quad v \in X_{1}, w \in X_{2} .
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Then the system $F(u, \lambda)=0$ is equivalent to

$$
\begin{array}{ll}
\hat{F}(v, w, \lambda):=Q F(v+w, \lambda) & =0 \in Y_{1} \\
\bar{F}(v, w, \lambda):=(I-Q) F(v+w, \lambda) & =0 \in Y_{2}
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and is thus invertible. So by the IFT we can locally write

$$
w=H(v, \lambda)
$$

We can thus write our reduced system

## Reduced system

$$
\begin{array}{r}
R(v, \lambda):=Q F(v+H(v, \lambda), \lambda)=0 \\
R: \operatorname{ker}(A) \times \mathbb{R} \rightarrow \operatorname{range}(A)^{\perp}
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This reduced system has the same symmetries and same bifurcations as the original problem, near $\left(u_{0}, \lambda_{0}\right)$.

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This reduced system has the same symmetries and same bifurcations as the original problem, near $\left(u_{0}, \lambda_{0}\right)$.

This is an extremely useful theoretical result. It forms the basis of most analytical calculations of bifurcation structures.

## Using our bases for $\operatorname{ker}(A)$ and range $(A)^{\perp}$, let's explicitly write:

## Reduced system (algebraic)

$$
\begin{gathered}
r_{j}(x, \lambda):=\left(\psi_{j}, R\left(x_{1} \phi_{1}+\cdots+x_{d} \phi_{d}, \lambda\right)\right)_{Y} \\
r: \mathbb{R}^{d} \times \mathbb{R} \rightarrow \mathbb{R}^{d}
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In practice we can never get our hands on $r$, because we don't know $H$. Instead, we compute a Taylor expansion (usually to third derivatives) of $r$.

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In practice we can never get our hands on $r$, because we don't know $H$. Instead, we compute a Taylor expansion (usually to third derivatives) of $r$.

The derivatives of $r$ can be computed from derivatives of $F$, and require solving linear systems involving $A\left(d^{2}+1\right.$ solves for third derivatives).

## Challenge

For large $d$, the Taylor expansion of the reduced equations are not easy to solve. There are techniques from numerical algebraic geometry that can provably yield all solutions, but they are too slow to use in practice.

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For large $d$, the Taylor expansion of the reduced equations are not easy to solve. There are techniques from numerical algebraic geometry that can provably yield all solutions, but they are too slow to use in practice.

The pragmatic response taken is to brute-force the system with many, many initial guesses (e.g. as in pde2path and BifurcationKit.jl).

## Example: Liouville-Bratu-Gelfand problem in 2D

$$
\nabla^{2} u-10\left(u-\lambda e^{u}\right)=0 \text { on } \Omega:=(0,1)^{2}, \quad \nabla u \cdot n=0 \text { on } \partial \Omega .
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$$

This is a famously intricate problem. I calculated the bifurcation diagram using BifurcationKit.jl. It was first computed successfully by Michiel Wouters.


Romain Veltz, 1982-


Michiel Wouters, ?-

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Solutions of the Bratu-Gelfand problem


## Example: Liouville-Bratu-Gelfand problem in 2D

$\nabla^{2} u-10\left(u-\lambda e^{u}\right)=0$ on $\Omega:=(0,1)^{2}, \quad \nabla u \cdot n=0$ on $\partial \Omega$.

Solutions of the Bratu-Gelfand problem


May 30

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Solutions of the Bratu-Gelfand problem


Lecture 3: Deflation algorithms for bifurcation analysis

## Patrick E. Farrell



University of Oxford

$$
\text { June } 1
$$

## Good news

The combination of continuation and branch switching is very powerful.

## Good news

The combination of continuation and branch switching is very powerful.

## Bad news

However, it has some disadvantages and weaknesses, too.

## Downside A

## You have to solve a lot of different problems.

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We work for years to develop a good solver for

$$
F(u, \lambda)=0 \ldots
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but now we need to solve

$$
\left[\begin{array}{c}
F(u, \lambda) \\
p(u, \lambda, s)
\end{array}\right]=\left[\begin{array}{l}
0 \\
0
\end{array}\right] \quad F_{u}(u, \lambda) v=\lambda v \quad\left[\begin{array}{c}
F(u, \lambda) \\
F_{u}(u, \lambda) v \\
\|v\|^{2}-1
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\end{array}\right]
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0 \\
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0
\end{array}\right]
$$

## Large-scale

This is OK when you can afford direct solvers, but it's hard at large scale.

## Downside B

We can only find branches connected to our initial data.

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This works fine ...


## Downside B

We can only find branches connected to our initial data.

This works fine ...

Solutions of $\lambda u-u^{3}=0$

... but this does not.


The standard approach to deal with this is to
(a) modify the problem to restore connectedness;

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## Problem A

You have to know to look for the missing branches.

## Problem B

Executing this is manual and tedious.

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## Problem A

You have to know to look for the missing branches.

## Problem B

Executing this is manual and tedious.

## Problem C

Restoring connectedness is not always possible!


The connectedness is broken by non-symmetry of the domain.

Deflation offers a complementary approach.

## Disconnected diagrams

An algorithm that can compute disconnected bifurcation diagrams.

Deflation offers a complementary approach.

## Disconnected diagrams

An algorithm that can compute disconnected bifurcation diagrams.

## Simplicity \& scaling

The computational kernel is exactly the same as Newton's method: solve

$$
F_{u}(u, \lambda) \delta u=-F(u, \lambda)
$$

## Section 2

## Deflation

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Fix parameter $\lambda$. Given

- a Fréchet differentiable residual $F: X \rightarrow Y$
- a solution $u \in X, F(u)=0, F_{u}(u)$ nonsingular


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- (Deflation property) Newton-Kantorovich applied to $G$ will never converge to $u$ again, starting from any initial guess.

Find more solutions, starting from the same initial guess.


$$
X
$$


$\star$

Newton from initial guess.

$$
X
$$

$u_{?}$
$\star$

## Deflate solution found.


$\star$

Newton from initial guess.

$$
X
$$

$u_{?}$

## Deflate solution found.

$$
X
$$



Newton from initial guess.

$$
X
$$

$u_{?}$

## Deflate solution found.



Terminate on nonconvergence.

$\star$

Terminate on nonconvergence.

Big if true. How can you do it?

## Big if true. How can you do it?

Numer. Math. 16, 334-342 (1971)
(C) by Springer-Verlag 1971

# Deflation Techniques for the Calculation of Further Solutions of a Nonlinear System 

Kenneth M. Brown and William B. Gearhart
Received March 10, 1970

Summary. This paper defines several classes of methods which can be used to find additional solutions of a nonlinear system of equations. A theory which embraces these classes is presented and the theory is extended to the multiple root problem. The techniques developed can also be used in avoiding previously found extreme points when performing function minimization. Results of computer experiments are presented.


Kenneth Brown, ?-?


Bill Gearhart, ?-

## Brown \& Gearhart's criterion

We say that $M(u ; r)$ is a deflation operator if

$$
\liminf _{u \rightarrow r}\|G(u)\|:=\liminf _{u \rightarrow r}\|M(u ; r) F(u)\|>0
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## Brown \& Gearhart's proposal

Choose

$$
M(u ; r):=\frac{1}{\|u-r\|}
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Note that $M(u, r)>0$ always, so $G(u)=0 \Longleftrightarrow F(u)=0$.

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## Brown \& Gearhart's proposal

Choose

$$
M(u ; r):=\frac{1}{\|u-r\|}
$$

Note that $M(u, r)>0$ always, so $G(u)=0 \Longleftrightarrow F(u)=0$.

Since $\|F(u)\|=O(\|u-r\|)$ as $u \rightarrow r$, this works.
... albeit sometimes not robustly.

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Why?
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(Allgower \& Georg, 1990)
[Deflation is] not . . . very reliable for larger problems.
(Kanzow, 2000)
Why?
One problem: assuming $F$ does not blow up as $\|u-r\| \rightarrow \infty$, then Newton discovers that it can achieve

$$
\|G(u)\|_{Y}<\mathrm{tol}
$$

for any tol, by taking $\|u-r\|$ large enough.

## Our proposal

$$
M_{p}(u ; r):=\left(\frac{1}{\|u-r\|^{p}}+1\right), \quad p \geq 1
$$



Ásgeir Birkisson, 1985-


Simon Funke, 1983-

## Our proposal

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## Our proposal

$$
M_{p}(u ; r):=\left(\frac{1}{\|u-r\|^{p}}+1\right), \quad p \geq 1
$$

Ásgeir Birkisson, 1985-
This has the right behaviour both as

$$
\begin{aligned}
& \|u-r\| \rightarrow 0 \\
& \|u-r\| \rightarrow \infty
\end{aligned}
$$

This makes the procedure much more reliable.


Simon Funke, 1983-


Start with $\left(u_{0}, \lambda_{0}\right)$.


Perform a continuation step.


Perform another continuation step.


Deflate the solution found.


Solve again.


Deflate the solution found.


Solve again.


Deflate the solution found.


Search again, unsuccessfully.


Repeat.

## Good news

## Deflation lets us discover disconnected branches!

## Section 3

## Solving the deflated problem

We assume we have a good solver for our discretised Newton step

$$
F_{u}(u, \lambda) \delta u_{F}=-F(u, \lambda), \quad F \in C^{1}\left(\mathbb{R}^{N} \times \mathbb{R}, \mathbb{R}^{N}\right)
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F_{u}(u, \lambda) \delta u_{F}=-F(u, \lambda), \quad F \in C^{1}\left(\mathbb{R}^{N} \times \mathbb{R}, \mathbb{R}^{N}\right)
$$

We now want to solve

$$
G_{u}(u, \lambda) \delta u_{G}=-G(u, \lambda)
$$

where

$$
G(u, \lambda)=M\left(u ; u_{1}\right) M\left(u ; u_{2}\right) \cdots M\left(u ; u_{n}\right) F(u, \lambda)=: M(u) F(u, \lambda) .
$$

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$$

## Good news

You can compute $\delta u_{G}$ easily from $\delta u_{F}$ !

## By the product rule,

$$
G_{u}(u, \lambda)=M(u) F_{u}(u, \lambda)+F(u, \lambda) M_{u}^{\top} .
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## Sherman-Morrison-Woodbury formula

$$
\left(A+u v^{\top}\right)^{-1}=A^{-1}-\left(\frac{A^{-1} u v^{\top} A^{-1}}{1+v^{\top} A^{-1} u}\right)
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Maurice Bartlett, 1910-2002

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$$

At first it looks like applying this to a vector $w$ requires two solves with $A$ : $A^{-1} u$ and $A^{-1} w$. But something magical happens ...

Applying the Sherman-Morrison-Woodbury formula, we have

$$
\delta u_{G}=-\left[G_{u}\right]^{-1} G=-\left(M F_{u}+F M_{u}^{\top}\right)^{-1}(M F)
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\begin{align*}
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& =-F_{u}^{-1} F+\frac{F_{u}^{-1} F M_{u}^{\top} M^{-1} F_{u}^{-1} F}{1+M_{u}^{\top} M^{-1} F_{u}^{-1} F}
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\end{aligned}
$$

So we only need to solve one system with $F_{u}$ !

## Solving the deflated problem

To solve

$$
G_{u} \delta u_{G}=-G
$$

do the following:

## Solving the deflated problem

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2. Evaluate

$$
p=M_{u}^{\top} \delta u_{F} .
$$

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G_{u} \delta u_{G}=-G,
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do the following:

1. Solve

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2. Evaluate

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$$

3. Evaluate

$$
\tau=1+\frac{M^{-1} p}{1-M^{-1} p}
$$

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G_{u} \delta u_{G}=-G
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do the following:

1. Solve

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2. Evaluate

$$
p=M_{u}^{\top} \delta u_{F} .
$$

3. Evaluate

$$
\tau=1+\frac{M^{-1} p}{1-M^{-1} p}
$$

4. Return

$$
\delta u_{G}=\tau \delta u_{F} .
$$

## Good news

You can apply deflation to massive discretisations.

## Section 4

## Convergence of deflation

It is possible to give sufficient conditions for deflation to find two roots.

It is possible to give sufficient conditions for deflation to find two roots.


Two solutions, with Rall-Rheinboldt balls.

It is possible to give sufficient conditions for deflation to find two roots.


Start with an initial guess within a ball.

It is possible to give sufficient conditions for deflation to find two roots.


Converge to that solution.

It is possible to give sufficient conditions for deflation to find two roots.


Deflate that solution; the other Rall-Rheinboldt ball expands.

## Section 5

## Examples

## Allen-Cahn equation

$$
F(u, \lambda)=-\lambda^{2} \nabla^{2} u+u^{3}-u=0, \quad u=g \text { on } \partial \Omega .
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Solutions found starting from $u=0$ for $\lambda=0.04$.

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Solutions found starting from $u=0$ for $\lambda=0.04$.

## Carrier's equation

$$
F(u, \lambda)=\lambda^{2} u^{\prime \prime}+2\left(1-x^{2}\right) u+u^{2}-1=0, \quad u(-1)=0=u(1)
$$





Solutions of $\lambda^{2} u^{\prime \prime}+2\left(1-x^{2}\right) u+u^{2}-1=0$


## Oseen-Frank

$$
\min J=\int_{\Omega} K_{1}(\nabla \cdot u)^{2}+K_{2}\left(u \cdot \nabla \times u+q_{0}\right)^{2}+K_{3}|u \times \nabla \times u|, \quad u \cdot u=1
$$

## Oseen-Frank

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## Section 6

## Symmetries

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## What if the equation has a continuous symmetry group?

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## Philosophy

The fundamental structures are the distinct orbits of solutions.

## Symmetries

## What if the equation has a continuous symmetry group?

## Philosophy

The fundamental structures are the distinct orbits of solutions.

## Key idea

Construct a deflation operator invariant under the action of the Lie group.


Four solutions, not related by the symmetry group.


Each solution induces a group orbit of solutions, related by symmetry.


Not enough to deflate the solution-must deflate the entire orbit.


Design a deflation operator that deflates the entire orbit.


Design a deflation operator that deflates the entire orbit.

## Gross-Pitaevskii equation

$$
-\frac{1}{2} \Delta u+\frac{x^{2}+y^{2}+z^{2}}{2} u-\mu u+|u|^{2} u=0,\left.\quad u\right|_{\partial \Omega}=0 .
$$

## Gross-Pitaevskii equation

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-\frac{1}{2} \Delta u+\frac{x^{2}+y^{2}+z^{2}}{2} u-\mu u+|u|^{2} u=0,\left.\quad u\right|_{\partial \Omega}=0 .
$$

## First symmetry group $\mathrm{SO}(2)$ : phase shifts

$$
u(\vec{x}) \mapsto e^{i \theta} u(\vec{x}), \quad \theta \in \mathbb{R}
$$

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First symmetry group $\mathrm{SO}(2)$ : phase shifts

$$
u(\vec{x}) \mapsto e^{i \theta} u(\vec{x}), \quad \theta \in \mathbb{R}
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## Invariant deflation operator

$$
M(u ; r)=\left\||u|^{2}-|r|^{2}\right\|^{-2}+1
$$

## Gross-Pitaevskii equation

$$
-\frac{1}{2} \Delta u+\frac{x^{2}+y^{2}+z^{2}}{2} u-\mu u+|u|^{2} u=0,\left.\quad u\right|_{\partial \Omega}=0 .
$$

Second symmetry group $\mathrm{SO}(3)$ : spatial rotations

$$
u(\vec{x}) \mapsto u(R \vec{x}), \quad R^{-1}=R^{T}, \quad \operatorname{det}(R)=1
$$

## Gross-Pitaevskii equation

$$
-\frac{1}{2} \Delta u+\frac{x^{2}+y^{2}+z^{2}}{2} u-\mu u+|u|^{2} u=0,\left.\quad u\right|_{\partial \Omega}=0
$$

Second symmetry group $\mathrm{SO}(3)$ : spatial rotations

$$
u(\vec{x}) \mapsto u(R \vec{x}), \quad R^{-1}=R^{T}, \quad \operatorname{det}(R)=1
$$

Invariant deflation operator

$$
\begin{gathered}
M(u ; r)=\|\bar{u}-\bar{r}\|^{-2}+1 \\
\text { where }
\end{gathered}
$$

$\bar{u}(r, \theta, \psi)$ averages $u$ over the sphere of radius $r$.

## Gross-Pitaevskii equation

$$
-\frac{1}{2} \Delta u+\frac{x^{2}+y^{2}+z^{2}}{2} u-\mu u+|u|^{2} u=0,\left.\quad u\right|_{\partial \Omega}=0
$$

Solutions for $\mu=6$.


A vortex line and a planar dark soliton.

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A pair of vortex lines.

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A vortex star.

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Four vortex lines of alternating charge.

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A vortex ring with two "handles".

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Two bent vortex rings?

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Two vortex rings and five lines?

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A vortex ring cage?

## Section 7

## Semismooth problems

Many problems feature inequality constraints.

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The natural language for formulating these is as a variational inequality.

## $\mathrm{VI}(Q, K)$

Let $X$ be a real reflexive Banach space, $K \subset X$ a closed convex subset, and $Q: K \rightarrow X^{*}$. The task is to find $u^{\star} \in K$ such that $\left\langle Q\left(u^{\star}\right), v-u^{\star}\right\rangle \geq 0$ for all $v \in K$.

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For example, if you want to minimise $f \in C^{1}(\mathbb{R}, \mathbb{R})$ over a closed interval $I \subset \mathbb{R}$, the necessary optimality condition is

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\mathrm{VI}\left(f^{\prime}, I\right)
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## The price we pay

$\ldots$ is that $S$ is not smooth.

## Good news

$S$ is just smooth enough to define a Newton-type method with superlinear convergence.


Michael Hintermüller, 1970-


Michael Ulbrich, 1967-

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## Semismoothness

Let $X$ and $Y$ be Banach spaces. Let $S: \Omega \subset X \rightarrow Y$, where $\Omega$ is an open subset of $X . S$ is semismooth at $u \in \Omega$ if it is locally Lipschitz continuous at $u$ and there exists an open neighbourhood $N \subset \Omega$ containing $u$ with a Newton derivative, i.e. a mapping $H: \Omega \rightarrow L(X, Y)$ with the property that

$$
S(u+h)-S(u)-H(u+h) h=o(h)
$$

for all $u$ in $N$.


Michael Hintermüller, 1970-


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## Good news

$S$ is just smooth enough to define a Newton-type method with superlinear convergence.

Semismooth Newton works just like normal:

$$
u_{i+1}=u_{i}-\left[H\left(u_{i}\right)\right]^{-1} S\left(u_{i}\right),
$$

Michael Hintermüller, 1970-
where $H$ is the Newton derivative.
This algorithm usually converges superlinearly.


Michael Ulbrich, 1967-

## Good news

## Deflation works for semismooth problems.

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## Theorem (F., Croci, Surowiec, 2020)

Under the same assumptions that are required for superlinear convergence of semismooth Newton, deflation works the same.


Matteo Croci, 1992-


Thomas Surowiec, 1982-

Gould gives an example where the central path is ill-behaved:

## Nonconvex quadratic programming problem

$$
\begin{aligned}
& \underset{\operatorname{minimise}}{ }-2\left(x_{1}-0.25\right)^{2}+2\left(x_{2}-0.5\right)^{2} \\
& x \in \mathbb{R}^{2} \\
& \text { subject to } \quad x_{1}+x_{2} \leq 1 \\
& 3 x_{1}+x_{2} \leq 1.5 \\
& x_{1} \geq 0 \\
& x_{2} \geq 0
\end{aligned}
$$



Nick Gould, 1957-

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## Buckling of a hyperelastic beam with contact constraints

$$
\begin{array}{ll}
\underset{u \in H^{1}\left(\Omega ; \mathbb{R}^{2}\right)}{\operatorname{minimise}} & \Pi(u)=\int_{\Omega} \psi(u) \mathrm{d} x-\int_{\Omega} B \cdot u \mathrm{~d} x \\
\text { subject to } & \left.u\right|_{\text {left }}=(0,0),\left.u\right|_{\text {right }}=(-\varepsilon, 0) \\
& \operatorname{tr}\left(u_{y}\right) \in[a, b] \text { a.e. in } \Gamma_{\text {top }}, \Gamma_{\text {bottom }}
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## Neo-Hookean strain energy density

$$
\psi(u)=\frac{\mu}{2}(\operatorname{tr}(C)-2)-\mu \log (\operatorname{det}(C))+\frac{\lambda}{2} \log (\operatorname{det}(C))^{2},
$$

where

$$
C=(I+\nabla u)^{\top}(I+\nabla u) .
$$



Multiple solutions of the beam with contact constraints.

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## Conclusions!

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## Main message

When solving nonlinear problems, think about multiple solutions!

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When solving nonlinear problems, think about multiple solutions!

## Algorithms

We now have very powerful algorithms for numerical bifurcation analysis.

## Open questions!

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## Thank you

to Josef, the organisers, and all the participants!

