# Numerical methods for large-scale dissipative dynamical systems. 

Juan Sánchez Umbría
Departament de Física, Universitat Politècnica de Catalunya.

## Overview

- Equilibria and periodic orbits of PDEs
- Newton-Krylov continuation methods
- Inexact Newton methods
- Iterative linear solvers and GMRES
- Stability
- Subspace iteration and Arnoldi methods
- Examples and exercises
- Continuation of codimension-one bifurcation points of equilibria and periodic orbits
- Continuation of invariant tori
- Periodic orbits by multiple shooting.


## Equilibria and periodic orbits of PDEs

## Equilibria and periodic orbits of ODEs

Suppose

$$
\dot{y}=f(y, p), \quad(y, p) \in \mathcal{U} \subset \mathbb{R}^{n} \times \mathbb{R}
$$

is a large-scale $(n \gg 1)$ autonomous system of ODEs obtained after the spatial discretization of a system of parabolic PDEs and that

$$
\varphi(t, x, p)
$$

is its solution with initial condition $x$ at $t=0$ for a fixed value of $p$, that is, $\varphi(0, x, p)=x$.
We will assume that this system has been obtained as the discretization of a systems of evolutionary parabolic PDEs (reaction-diffusion or Navier-Stokes equations, for instance).

We are interested in the computation of its equilibria $x$ satisfying

$$
f(x, p)=0,
$$

their dependence on the parameter $p$ and their stability.
We are also interested in the periodic regimes of the system given by the equations

$$
\begin{array}{r}
x-\varphi(T, x, p)=0, \\
g(x, p)=0,
\end{array}
$$

$x$ being a point of the periodic orbit selected by the phase condition $g(x, p)=0$ and $T>0$ its period.
In both cases one has to solve large-scale nonlinear systems of equations and to study the stability of the resulting equilibria or periodic orbits.

## Newton-Krylov continuation methods

## Continuation of zeros of a nonlinear system of equations

Consider a system of nonlinear equations depending on a parameter $p$

$$
H(x, p)=0, \quad(x, p) \in \mathcal{U} \subset \mathbb{R}^{m} \times \mathbb{R}
$$

with $m \gg 1$. We are interested in its solutions and their dependence on $p$.
Parameter and pseudo-arclength-like continuation methods are used to obtain the curves $(x(s), p(s))$ of solutions. They admit an unified formulation by adding an equation

$$
h(x, p)=0 .
$$

If $h(x, p)=p-p_{0}$ the equation fixes the parameter $p$.
If $h(x, p)=h_{x}^{\top}\left(x-x_{0}\right)+h_{p}\left(p-p_{0}\right)$, with $\left(x_{0}, p_{0}\right)$ and ( $h_{x}, h_{p}$ ) being the predicted point and the tangent to the curve of solutions, the hyperplane is transverse to the curve of solutions if the prediction is not far away from the previous point, and the algorithm allows passing turning points.

The system that determines a unique solution, $(x, p) \in \mathbb{R}^{m+1}$, is then


$$
\widetilde{H}(x, p)=\binom{H(x, p)}{h(x, p)}=0 \in \mathbb{R}^{m+1} .
$$

The system $\widetilde{H}(x, p)=0$ is solved by an inexact Newton's method:
starting from the initial $\left(x_{0}, p_{0}\right)$,

$$
\left(x_{i+1}, p_{i+1}\right)=\left(x_{i}, p_{i}\right)+\left(\Delta x_{i}, \Delta p_{i}\right)
$$

where $\left(\Delta x_{i}, \Delta p_{i}\right)$ satisfies the linear system

$$
\left(\begin{array}{cc}
D_{x} H\left(x_{i}, p_{i}\right) & D_{p} H\left(x_{i}, p_{i}\right) \\
h_{x}^{\top} & h_{p}
\end{array}\right)\binom{\Delta x_{i}}{\Delta p_{i}}=\binom{-H\left(x_{i}, p_{i}\right)}{-h\left(x_{i}, p_{i}\right)}
$$

which is solved iteratively by matrix-free methods (GMRES(m), BiCGStab, FOM, TFQMR, etc.) which only require the computation of matrix products, i.e., products of the form

$$
\left(\begin{array}{cc}
D_{x} H\left(x_{i}, p_{i}\right) & D_{p} H\left(x_{i}, p_{i}\right) \\
h_{x}^{\top} & h_{p}
\end{array}\right)\binom{\delta x}{\delta p}
$$

and, eventually, the use of preconditioners.

| GMRES $(\mathrm{m})$ | $=$ Generalized Minimal Residual (with restarting dimension m) |
| :--- | :--- |
| BiCGStab | $=$ Biconjugate Gradient Stabilized |
| FOM | $=$ Full Orthogonalization Method |
| TFQMR | $=$ Transpose-Free Quasi-Minimal Residual |

## An example of a matrix-free product

Consider the system of PDEs

$$
\begin{aligned}
& \partial_{\tau} c=\left(1 / P e_{m}\right) \partial_{s s}^{2} c-\partial_{s} c-D c \exp (\gamma(1-1 / \theta)) \\
& \partial_{\tau} \theta=\left(1 / P e_{h}\right) \partial_{s s}^{2} \theta-\partial_{s} \theta-\beta\left(\theta-\theta_{r}\right)+B D c \exp (\gamma(1-1 / \theta)),
\end{aligned}
$$

modelling a tubular exotermic chemical reactor (Heinemann and Poore 1981), with $s \in[0,1]$, and where $c, \theta$ and $\tau$ are the non-dimensional concentration of a reactant, temperature and time, respectivelly. $P e_{m}, P e_{h}, D, \beta, B, \theta_{r}$ and $\gamma$ are non-dimensional parameters of the problem.

Suppose that all of them are fixed except $D$ that will be our control parameter ( $p$ in the previous slides), and that our state variable is $x=(c, \theta)$.

Let

$$
H(x, p)=\binom{\left(1 / P e_{m}\right) \partial_{s s}^{2} c-\partial_{s} c-D c \exp (\gamma(1-1 / \theta))}{\left(1 / P e_{h}\right) \partial_{s s}^{2} \theta-\partial_{s} \theta-\beta\left(\theta-\theta_{r}\right)+B D c \exp (\gamma(1-1 / \theta))} .
$$

Then, if $\delta x=(\delta c, \delta \theta)$ and $\delta p=\delta D$,

$$
\begin{aligned}
& D_{x} H(x, p) \delta x+D_{p} H(x, p) \delta p= \\
& \binom{\left(1 / P e_{m}\right) \partial_{s s}^{2} \delta c-\partial_{s} \delta c-\exp (\gamma(1-1 / \theta))\left(D \delta c+D c\left(\gamma / \theta^{2}\right) \delta \theta+\delta D c\right)}{\left(1 / P e_{h}\right) \partial_{s s}^{2} \delta \theta-\partial_{s} \delta \theta-\beta \delta \theta+B \exp (\gamma(1-1 / \theta))\left(D \delta c+D c\left(\gamma / \theta^{2}\right) \delta \theta+\delta D c\right)} .
\end{aligned}
$$

## Inexact Newton's methods

## Types of convergence

Iterative methods can be classified by their rate of convergence.
Definition. Let $\left\{x_{k}\right\} \subset \mathbb{R}^{n}$ and $x^{*} \in \mathbb{R}^{n}$. Then

- $x_{k} \rightarrow x^{*} q$-quadratically if $x_{k} \rightarrow x^{*}$ and there is $K>0$ such that

$$
\left\|x_{k+1}-x^{*}\right\| \leq K\left\|x_{k}-x^{*}\right\|^{2} .
$$

- $x_{k} \rightarrow x^{*} q$-superlinearly with $q$-order $\alpha>1$ if $x_{k} \rightarrow x^{*}$ and there is $K>0$ such that

$$
\left\|x_{k+1}-x^{*}\right\| \leq K\left\|x_{k}-x^{*}\right\|^{\alpha} .
$$

- $x_{k} \rightarrow x^{*} q$-superlinearly if

$$
\lim _{n \rightarrow \infty}\left\|x_{k+1}-x^{*}\right\| /\left\|x_{k}-x^{*}\right\|=0
$$

- $x_{k} \rightarrow x^{*}$ q-linearly with $q$-factor $\sigma \in(0,1)$ if

$$
\left\|x_{k+1}-x^{*}\right\| \leq \sigma\left\|x_{k}-x^{*}\right\| .
$$

Definition. Let $\left\{x_{k}\right\} \subset \mathbb{R}^{n}$ and $x^{*} \in \mathbb{R}^{n}$. Then $x_{k} \rightarrow x^{*} r$-(quadratically, superlinearly, linearly) if there is a sequence $\left\{\xi_{k}\right\} \subset \mathbb{R}$ converging $q$-(quadratically, superlinearly, linearly) to zero such that

$$
\left\|x_{k}-x^{*}\right\| \leq \xi_{k},
$$

and $x_{k} \rightarrow x^{*} r$-superlinearly with $r$-order $\alpha>1$ if the sequence $\xi_{k} \rightarrow 0 q$-superlinearly with $q$-order $\alpha$.

## Newton's method

Suppose we seek to solve the system

$$
F(x)=0
$$

with $x, F(x) \in \mathbb{R}^{n}$, and assume the following standard conditions hold

- $F(x)=0$ has a solution $x^{*}$,
- There is a neighborhood of $x^{*}, \Omega \subset \mathbb{R}^{N}$, such that $D F: \Omega \rightarrow \mathbb{R}^{N \times N}$ is Lipschitz continuous with Lipschitz constant $\gamma>0$, i.e.,

$$
\|D F(x)-D F(y)\| \leq \gamma\|x-y\|
$$

for all $x, y \in \Omega$,

- $D F\left(x^{*}\right)$ is nonsingular.

Theorem. Under the above assumptions there is a $\delta>0$ such that if $\left\|x_{0}-x^{*}\right\|<\delta$ the Newton iteration

$$
x_{k+1}=x_{k}+s_{k}, \quad \text { with } \quad D F\left(x_{k}\right) s_{k}=-F\left(x_{k}\right)
$$

converges $q$-quadratically to $x^{*}$, i.e., there is a $K>0$ such that

$$
\left\|x_{k+1}-x^{*}\right\| \leq K\left\|x_{k}-x^{*}\right\|^{2} .
$$

## Inexact Newton methods

Suppose now that instead of solving

$$
D F\left(x_{k}\right) s_{k}=-F\left(x_{k}\right)
$$

exactly, the linear system is solved by an iterative method with stopping criteria

$$
\left\|D F\left(x_{k}\right) s_{k}+F\left(x_{k}\right)\right\| \leq \eta_{k}\left\|F\left(x_{k}\right)\right\|
$$

Theorem. Let the standard conditions hold. Then there exists $\delta>0$ such that if $\left\|x_{0}-x^{*}\right\|<\delta$, and $\left\{\eta_{k}\right\} \subset[0, \eta]$ with $\eta<\bar{\eta}<1$, then the inexact Newton iteration

$$
x_{k+1}=x_{k}+s_{k}, \quad \text { with } \quad\left\|D F\left(x_{k}\right) s_{k}+F\left(x_{k}\right)\right\| \leq \eta_{k}\left\|F\left(x_{k}\right)\right\|,
$$

converges $q$-linearly to $x^{*}$ with respect to the norm $\|\cdot\|_{*}=\left\|D F\left(x^{*}\right) \cdot\right\|$. Moreover

- if $\eta_{k} \rightarrow 0$ the convergence is $q$-superlinear, and
- if $\eta_{k} \leq K_{\eta}\left\|F\left(x_{k}\right)\right\|^{p}$ for some $K_{\eta}>0$ the convergence is $q$-superlinear with $q$-order $1+p$.

Proposition. Under the standard conditions, and if $x_{k} \rightarrow x^{*},\left\|x_{k}-x^{*}\right\|_{*} \rightarrow 0$-linearly if and only if $\left\|F\left(x_{k}\right)\right\|$ does.

## Iterative linear solvers and GMRES

## Krylov methods for linear systems

Large-scale non-symmetric linear systems $A x=b$ of dimension $n \gg 1$ are usually solved by iterative Krylov methods. The class of projection methods produce, from an initial guess $x_{0}$, a sequence of approximations, $x_{k}$, to the solution $x^{*}=A^{-1} b$, in the affine subspace $x_{k} \in x_{0}+\mathcal{K}_{k}$, which satisfy the Petrov-Galerkin condition

$$
b-A x_{k} \perp \mathcal{L}_{k},
$$

where $\mathcal{K}_{k}$ and $\mathcal{L}_{k}$ are two $k$-dimensional linear subspaces. If $\mathcal{L}_{k}=A \mathcal{K}_{k}$, then $x_{k}$ satisfies

$$
\left\|b-A x_{k}\right\|_{2}=\inf _{x \in x_{0}+\mathcal{K}_{k}}\|b-A x\|_{2} .
$$

It minimizes the norm $\|b-A x\|_{2}$ over $x \in x_{0}+\mathcal{K}_{k}$.
In the particular case of GMRES, $\mathcal{L}_{k}=A \mathcal{K}_{k}$, and $\mathcal{K}_{k}$ is the Krylov subspace

$$
\mathcal{K}_{k}=\left\{r_{0}, A r_{0}, A^{2} r_{0}, \ldots, A^{k-1} r_{0}\right\}, \quad \text { with } \quad r_{0}=b-A x_{0} .
$$

It follows that

$$
\begin{aligned}
r_{k}= & b-A x_{k}=b-A\left(x_{0}+z_{k}\right)=r_{0}-A z_{k}= \\
& =I r_{0}-A\left(\alpha_{1} r_{0}+\alpha_{2} A r_{0}+\cdots+\alpha_{k} A^{k-1} r_{0}\right) \\
& =\left(I-\alpha_{1} A-\alpha_{2} A^{2}-\cdots-\alpha_{k} A^{k}\right) r_{0}=p_{k}(A) r_{0}
\end{aligned}
$$

$p_{k}$ being a polynomial of degree $k$, with $p_{k}(0)=1$.

Now, by using that

- If $A=V \Lambda V^{-1}$ then $A^{l}=V \Lambda^{l} V^{-1}$ and $p(A)=V p(\Lambda) V^{-1}$
- If $\Lambda=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{n}\right)$ then $p(\Lambda)=\operatorname{diag}\left(p\left(\lambda_{1}\right), \ldots, p\left(\lambda_{n}\right)\right)$
- $\|p(A)\|_{2} \leq\|p(\Lambda)\|_{2}\|V\|_{2}\left\|V^{-1}\right\|_{2}=\kappa_{2}(V)\|p(\Lambda)\|_{2}$, with $\kappa_{2}(V)=\|V\|_{2}\left\|V^{-1}\right\|_{2}$ the norm-2 condition number of $V$.
- If $\Lambda=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{n}\right)$ then $\|p(\Lambda)\|_{2}=\max _{i=1, \ldots, n}\left|p\left(\lambda_{i}\right)\right|$
the following result is obtained.

Theorem. (Saad and Schultz 1986) Assume that $A$ is diagonalizable with $A=V \Lambda V^{-1}$, where $\Lambda=\operatorname{diag}\left(\lambda_{1}, \cdots, \lambda_{n}\right)$ is the diagonal matrix of eigenvalues, $P_{k}$ is the set of polynomials of degree at most $k$, and $\kappa_{2}(V)=\left\|V^{-1}\right\|_{2}\|V\|_{2}$ is the norm-2 condition number of $V$. Then at the $k$-th step of GMRES

$$
\frac{\left\|b-A x_{k}\right\|_{2}}{\left\|b-A x_{0}\right\|_{2}} \leq \kappa_{2}(V) \inf _{\substack{p \in P_{k} \\ p(0)=1}} \max _{i=1, \ldots, n}\left|p\left(\lambda_{i}\right)\right|
$$

Proof:

$$
\left\|b-A x_{k}\right\|_{2}=\inf _{\substack{p \in P_{k} \\ p(0)=1}}\left\|p(A) r_{0}\right\| \leq \kappa_{2}(V) \inf _{\substack{p \in P_{k} \\ p(0)=1}} \max _{i=1, \ldots, n}\left|p\left(\lambda_{i}\right)\right|\left\|b-A x_{0}\right\|_{2}
$$

## The Arnoldi's factorization

It is used in the solution of linear systems, eigenvalue problems (Arnoldi's method), and to find reduced order models by the dynamic mode decomposition (DMD).

Let $\langle.,$.$\rangle be the euclidean dot product. Given a matrix A$ and an initial unitary vector $v_{1}$
Iterate: for $j=1,2, \ldots, k$ compute

1. $h_{i, j}=\left\langle A v_{j}, v_{i}\right\rangle$, for $i=1,2, \ldots, j$
2. $w_{j}=A v_{j}-\sum_{i=1}^{j} h_{i, j} v_{i}$ (this is classical Gram-Schmidt orthogonalization)
3. $h_{j+1, j}=\left\|w_{j}\right\|_{2}$, if $h_{j+1, j}=0$ stop
4. $v_{j+1}=w_{j} / h_{j+1, j}$

If $V_{k}=\left[v_{1}, \ldots, v_{k}\right]$ is the matrix with columns $v_{1}, \ldots, v_{k}$ then

- The columns of $V_{k}$ form an orthonormal basis of $\mathcal{K}_{k}=\left\{v_{1}, A v_{1}, A^{2} v_{1}, \ldots, A^{k-1} v_{1}\right\}$.
- If $H_{k}$ is the $k \times k$ upper Hessenberg matrix whose nonzero entries are the $h_{i, j}$ then

$$
A V_{k}=V_{k} H_{k}+w_{k} e_{k}^{\top}, \quad \text { and } \quad V_{k}^{\top} A V_{k}=H_{k},
$$

with $w_{k}=h_{k+1, k} v_{k+1}$, and $e_{k}^{\top}=(0, \ldots, 0,1) \in \mathbb{R}^{k}$.

- If $\tilde{H}_{k}$ is the $(k+1) \times k$ matrix whose nonzero entries are the $h_{i, j}$ then

$$
A V_{k}=V_{k+1} \tilde{H}_{k} .
$$

The matrix $H_{k}$ is $\tilde{H}_{k}$ without its last row.

$$
\begin{aligned}
& A V_{k}=V_{k} H_{k}+h_{k+1, k} v_{k+1} e_{k}^{\top}: \\
& A\left[v_{1}, \ldots, v_{k}\right]=\left[v_{1}, \ldots, v_{k}\right]\left[\begin{array}{ccccc}
h_{1,1} & h_{1,2} & h_{1,3} & \ldots & h_{1, k} \\
h_{2,1} & h_{2,2} & h_{2,3} & \ldots & h_{2, k} \\
0 & h_{3,2} & h_{3,3} & \ldots & h_{3, k} \\
\vdots & \ddots & \ddots & \ddots & \vdots \\
0 & \ldots & 0 & h_{k, k-1} & h_{k, k}
\end{array}\right]+h_{k+1, k}\left[0, \ldots, 0, v_{k+1}\right] \\
& A V_{k}=V_{k+1} \tilde{H}_{k}: \\
& A\left[v_{1}, \ldots, v_{k}\right]=\left[v_{1}, \ldots, v_{k}, v_{k+1}\right]\left[\begin{array}{ccccc} 
\\
h_{1,1} & h_{1,2} & h_{1,3} & \ldots & h_{1, k} \\
h_{2,1} & h_{2,2} & h_{2,3} & \ldots & h_{2, k} \\
0 & h_{3,2} & h_{3,3} & \ldots & h_{3, k} \\
\vdots & \ddots & \ddots & \ddots & \vdots \\
0 & \ldots & 0 & h_{k, k-1} & h_{k, k} \\
0 & \ldots & \ldots & 0 & h_{k+1, k}
\end{array}\right]
\end{aligned}
$$

## Practical implementation of the Arnoldi's factorization

Given a matrix $A$ and an initial unitary vector $v_{1}$ Iterate: for $j=1,2, \ldots, k$ compute

1. $w=A v_{j}$
2. for $i=1,2, \ldots, j$ do (this is modified Gram-Schmidt orthogonalization)
(a) $h_{i, j}=\left\langle w, v_{i}\right\rangle$
(b) $w \leftarrow w-h_{i, j} v_{i}$
3. $h_{j+1, j}=\|w\|_{2}$, if $h_{j+1, j}=0$ stop
4. $v_{j+1}=w / h_{j+1, j}$

To complete GMRES, it remains to solve the minimizing problem

$$
\inf _{x \in x_{0}+\mathcal{K}_{k}}\|b-A x\|_{2}
$$

Suppose that $v_{1}, \ldots, v_{k}$ form an orthonormal basis of $\mathcal{K}_{k}$ and let $V_{k}=\left[v_{1}, \ldots, v_{k}\right]$. It is found by means of the Arnoldi factorization with $v_{1}=r_{0} /\left\|r_{0}\right\|$. Then

$$
x \in x_{0}+\mathcal{K}_{k} \Longrightarrow x=x_{0}+V_{k} y \quad \text { with } \quad y \in \mathbb{R}^{k}
$$

and

$$
\inf _{x \in x_{0}+\mathcal{K}_{k}}\|b-A x\|_{2}=\inf _{y \in \mathbb{R}^{k}}\left\|b-A\left(x_{0}+V_{k} y\right)\right\|=\inf _{y \in \mathbb{R}^{k}}\left\|r_{0}-A V_{k} y\right\|
$$

Now, since $A V_{k}=V_{k+1} \tilde{H}_{k}$,

$$
\left\|r_{0}-A V_{k} y\right\|_{2}=\left\|r_{0}-V_{k+1} \tilde{H}_{k} y\right\|_{2}=\left\|V_{k+1}^{\top}\left(\beta v_{1}-V_{k+1} \tilde{H}_{k} y\right)\right\|_{2}=\left\|\beta e_{1}-\tilde{H}_{k} y\right\|_{2}
$$

with $\beta=\left\|r_{0}\right\|, e_{1}=(1,0, \ldots, 0)^{\top} \in \mathbb{R}^{k+1}$ and $y \in \mathbb{R}^{k}$.
Therefore

$$
\inf _{x \in x_{0}+\mathcal{K}_{k}}\|b-A x\|_{2}=\inf _{y \in \mathbb{R}^{k}}\left\|\beta e_{1}-\tilde{H}_{k} y\right\|_{2}
$$

The latter is a least-squares problem without restrictions in $\mathbb{R}^{k}$.
One wants to solve the system in a small number of iterations to avoid the matrix $V_{k}$ to grow too much and then in practice $k \ll n$.

## The restarted GMRES(m) algorithm

The complete restarted version of the algorithm GMRES(m) is, given $x_{0}$, and $r_{0}=b-A x_{0}$,

1. Set $l=0$.
2. Start: Choose as initial unitary vector $v_{1}=r_{0} /\left\|r_{0}\right\|$, set $\rho=\beta=\left\|r_{0}\right\|, k=0$.
3. do while $\rho>\varepsilon, k<m$, and $l<l_{\max }$ :
(a) set $k=k+1$ and $l=l+1$
(b) compute the Arnoldi's factorization $A V_{k}=V_{k+1} \tilde{H}_{k}$
(c) find $y_{k}$ the minimizer of $\left\|\beta e_{1}-\tilde{H}_{k} y\right\|_{2}$
(d) set $\rho=\left\|\beta e_{1}-\tilde{H}_{k} y_{k}\right\|_{2}\left(\right.$ remember that $\left.\inf _{x \in x_{0}+\mathcal{K}_{k}}\|b-A x\|_{2}=\inf _{y \in \mathbb{R}^{k}}\left\|\beta e_{1}-\tilde{H}_{k} y\right\|_{2}\right)$
4. if $\rho<\varepsilon$ then set $x_{k}=x_{0}+V_{k} y_{k}$ as approximate solution and exit
5. if $l>l_{\max }$ (too many iterations without convergence) exit
6. if $k=m$ set $x_{0} \leftarrow x_{0}+V_{k} y_{k}$ and restart the algorithm (go to 2).

## Preconditioning

If the spectrum of $A$ is not clustered it is necessary to use preconditioners to accelerate the convergence of the iterative solvers for the linear system $A x=b$.

Suppose $M$ is a matrix which approximates $A(M \approx A)$ and is easy to invert (easy to solve systems with matrix $M$ ).

- Left preconditioning. Solve the system

$$
M^{-1} A x=M^{-1} b
$$

Its solution is that of $A x=b$.

- Right preconditioning. Solve the system

$$
A M^{-1} y=b .
$$

Then the solution of $A x=b$ is $x=M^{-1} y$.
This means that when applying a matrix-free method (GMRES, for instance) each matrix product by $A$ is substituted by a matrix product by $A$ followed by a matrix solve with matrix $M$ in the case of left preconditioning, or by a matrix solve with matrix $M$ followed by a matrix product by $A$ for right preconditioning.

## Spatial discretization of the HP problem

Consider the system of PDEs

$$
\begin{aligned}
& \partial_{\tau} c=\left(1 / P e_{m}\right) \partial_{s s}^{2} c-\partial_{s} c-D c \exp (\gamma(1-1 / \theta)) \\
& \partial_{\tau} \theta=\left(1 / P e_{h}\right) \partial_{s s}^{2} \theta-\partial_{s} \theta-\beta\left(\theta-\theta_{r}\right)+B D c \exp (\gamma(1-1 / \theta)),
\end{aligned}
$$

in the interval $s \in[0,1]$, with boundary conditions

$$
\begin{array}{lllll}
\partial_{s} c=P e_{m}(c-1) & \text { at } & s=0, & \partial_{s} c=0 & \text { at } \\
\partial_{s} \theta=P e_{h}(\theta-1) & \text { at } & s=0, & \partial_{s} \theta=0 & \text { at } \\
s=1 .
\end{array}
$$

To implement the boundary conditions easily we substitute $c=\bar{c}+1, \theta=\bar{\theta}+1$ in the equations and boundary conditions to obtain, after removing the overbars the equations

$$
\begin{aligned}
& \left.\partial_{\tau} c=\left(1 / P e_{m}\right) \partial_{s s}^{2} c-\partial_{s} c-D(c+1) \exp (\gamma \theta /(\theta+1))\right) \\
& \left.\partial_{\tau} \theta=\left(1 / P e_{h}\right) \partial_{s s}^{2} \theta-\partial_{s} \theta-\beta\left(\theta-\theta_{r}+1\right)+B D(c+1) \exp (\gamma \theta /(\theta+1))\right)
\end{aligned}
$$

with boundary conditions

$$
\begin{array}{lllll}
\partial_{s} c=P e_{m} c & \text { at } & s=0, & \partial_{s} c=0 & \text { at } \\
\partial_{s} \theta=P e_{h} \theta & \text { at } & s=0, & \partial_{s} \theta=0 & \text { at } \\
s=1 .
\end{array}
$$

If $D=0$ and $\theta_{r}=1$ then $c=0$ and $\theta=0$ is a solution of the problem.

We use collocation methods in a Gauss-Lobatto mesh. Let $s_{i}=0.5\left(1-\cos \left(\pi i / n_{d}\right)\right), i=0, \ldots, n_{d}$ and $D^{(l)}=\left\{d_{i, j}^{(l)}\right\}$ the $\left(n_{d}+1\right) \times\left(n_{d}+1\right)$ matrices which approximate the derivatives on the mesh, i.e,

$$
f^{(l)}\left(s_{i}\right) \approx \sum_{j=0}^{n_{d}} d_{i, j}^{(l)} f\left(s_{j}\right), \quad i=0, \ldots, n_{d} .
$$

Let $c_{i}=c\left(s_{i}\right), \theta_{i}=\theta\left(s_{i}\right)$ and approximate the boundary conditions (of $c$, for instance) by

$$
\sum_{j=0}^{n_{d}} d_{0, j}^{(1)} c_{j}=P e_{m} c_{0}, \quad \sum_{j=0}^{n_{d}} d_{n_{d}, j}^{(1)} c_{j}=0
$$

From these two equations the values at the end points can be obtained as a linear combination of the values at the inner points,

$$
c_{0}=\sum_{j=1}^{n_{d}-1} \alpha_{0, j} c_{j}, \quad c_{n_{d}}=\sum_{j=1}^{n_{d}-1} \alpha_{n_{d}, j} c_{j} .
$$

And then, for instance,

$$
\partial_{s s}^{2} c\left(s_{i}\right) \approx \sum_{j=1}^{n_{d}-1}\left(d_{i, j}^{(2)}+d_{i, 0}^{(2)} \alpha_{0, j}+d_{i, n_{d}}^{(2)} \alpha_{n_{d}, j}\right) c_{j}=\sum_{j=1}^{n_{d}-1} \tilde{d}_{i, j}^{(2)} c_{j}, \quad i=1 \ldots, n_{d}-1
$$

and $\tilde{D}^{(2)}=\left\{\tilde{d}_{i, j}^{(2)}\right\}$ is the $\left(n_{d}-1\right) \times\left(n_{d}-1\right)$ matrix which approximates $\partial_{s s}^{2}$ incorporating the boundary conditions and acting only on the values at the inner points.

After the spatial discretization of

$$
\begin{aligned}
& \left.\partial_{\tau} c=\left(1 / P e_{m}\right) \partial_{s s}^{2} c-\partial_{s} c-D(c+1) \exp (\gamma \theta /(\theta+1))\right) \\
& \left.\partial_{\tau} \theta=\left(1 / P e_{h}\right) \partial_{s s}^{2} \theta-\partial_{s} \theta-\beta\left(\theta-\theta_{r}+1\right)+B D(c+1) \exp (\gamma \theta /(\theta+1))\right)
\end{aligned}
$$

the following stiff system of ODEs of dimension $2\left(n_{d}-1\right)$ is obtained

$$
\begin{aligned}
\dot{c}_{i} & \left.=\sum_{j=1}^{n_{d}-1}\left(\left(1 / P e_{m}\right) \tilde{d}_{i, j}^{(2)}-\tilde{d}_{i, j}^{(1)}\right) c_{j}-D\left(c_{i}+1\right) \exp \left(\gamma \theta_{i} /\left(\theta_{i}+1\right)\right)\right) \\
\dot{\theta}_{i} & \left.=\sum_{j=1}^{n_{d}-1}\left(\left(1 / P e_{h}\right) \tilde{d}_{i, j}^{(2)}-\tilde{d}_{i, j}^{(1)}-\beta \delta_{i, j}\right) \theta_{j}-\beta\left(1-\theta_{r}\right)+B D\left(c_{i}+1\right) \exp \left(\gamma \theta_{i} /\left(\theta_{i}+1\right)\right)\right), \\
i & =1, \ldots, n_{d}-1
\end{aligned}
$$

which is integrated with the subroutine DLSODPK from the ODEPACK library.

## Effect of the preconditioner in the HP problem

In all the following examples we have taken $n_{d}=30$ and therefore the dimension of the dynamical system is $n=2\left(n_{d}-1\right)=58$.

For the next examples $P e_{m}=P e_{h}=5, B=0.5, \gamma=25, \beta=3.5, \theta_{r}=1$, and $D$ will be the free parameter.

If $\delta x=(\delta c, \delta \theta)$ then,

$$
D_{x} H(x, p) \delta x=\binom{\left(1 / P e_{m}\right) \partial_{s s}^{2} \delta c-\partial_{s} \delta c-N}{\left(1 / P e_{h}\right) \partial_{s s}^{2} \delta \theta-\partial_{s} \delta \theta-\beta \delta \theta+B N}
$$

with $N=\exp (\gamma \theta /(\theta+1)) D\left(\delta y+(y+1)\left(\gamma /(\theta+1)^{2}\right) \delta \theta\right)$
Two possible preconditioners are

$$
M_{1}=\left(\begin{array}{cc}
\left(1 / P e_{m}\right) \partial_{s s}^{2} & 0 \\
0 & \left(1 / P e_{h}\right) \partial_{s s}^{2}
\end{array}\right)
$$

and

$$
M_{2}=\left(\begin{array}{cc}
\left(1 / P e_{m}\right) \partial_{s s}^{2}-\partial_{s} & 0 \\
0 & \left(1 / P e_{h}\right) \partial_{s s}^{2}-\partial_{s}-\beta I
\end{array}\right)
$$

The following figure shows the convergence of GMRES at the first Newton iteration for $D=0.1$ starting with $c=0$ and $\theta=0$. The size of the linear system is $n=58$, and the dimension of the Krylov subspace was $m=10$ or $m=58$. Norm of the residual $=\left\|b-A x_{k}\right\|_{2}$.

Convergence of GMRES(m)


## Stability

## Stability of fixed points (equilibria)

Given an autonomous system of ODE $\dot{x}=f(x)$, with $f: \mathcal{U} \subset \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ (which we will assume to be at least $C^{1}$ ), let $\varphi(t, x)$ its solution with initial condition $x$.

Let $x_{*}$ be a fixed point (or equilibrium) of the system of EDOs, i.e., $f\left(x_{*}\right)=0$.
Definition. The fixed point is said to be Lyapunov stable if for every neighborhood $N$ of $x_{*}$ there is a neighborhood $M \subset N$ of $x_{*}$ such that if $x \in M$, then $\varphi(t, x) \in N$ for all $t \geq 0$.

An equilibrium that is not stable is called unstable.
Definition. The fixed point is said to be asymptotically stable if it is Lyapunov stable and there is a neighborhood $N$ of $x_{*}$ such if $x \in N$ then $\lim _{t \rightarrow \infty}\left\|\varphi(t, x)-x_{*}\right\|=0$.
Definition. The fixed point said to be exponentially stable if it is asymptotically stable and there exist $\alpha>0$, and $\beta>0$, and a neighborhood $N$ of $x_{*}$ such that if $x \in N$, then $\left\|\varphi(t, x)-x_{*}\right\| \leq \alpha\left\|x-x_{*}\right\| e^{-\beta t}$, for $t \geq 0$.
Theorem. If $f$ is of class $C^{1}$ and $x_{*}$ is a fixed point such that all the eigenvalues of $D f\left(x_{*}\right)$ have strictly negative real parts, then $x_{*}$ is exponentially stable (and hence asymptotically stable). If at least one eigenvalue has strictly positive real part, then $x_{*}$ is unstable.

The eigenvalues of $D f\left(x_{*}\right)$ close the imaginary axis have to be computed to detect bifurcations of fixed points.

## Stability of periodic orbits

Definition. $A$ set $\Lambda$ is said to be invariant under the flow $\varphi(t, x)$ if $\varphi(t, \Lambda)=\Lambda$ for all $t$; that is, for each $x \in \Lambda, \varphi(t, x) \in \Lambda$ for any $t$.
Definition. The invariant set is said to be stable if for every neighborhood $N$ of $\Lambda$ there is a subset $M \subset N$ of $\Lambda$ such that if $x \in M$, then $\varphi(t, x) \in N$ for all $t \geq 0$.

An set that is not stable is called unstable.
Definition. The invariant set is said to be asymptotically stable if it is stable and there is a neighborhood $N$ of $\Lambda$ such if $x \in N$ then, then $\lim _{t \rightarrow \infty} \rho(\varphi(t, x), \Lambda)=0$, with
$\rho(x, \Lambda)=\inf _{y \in \Lambda}(\|x-y\|)$.
A trajectory $x(t)=\varphi(t, x)$ is a periodic orbit if there is a minimal $T>0$ such that $\varphi(T, x)=x$.
Consider the first variational equation $\dot{M}=D f(x(t)) M$ about the periodic orbit $x(t)$, with initial condition $M(0)=I$. The solution at time $T$ is called the monodromy matrix $M(T)$. Its eigenvalues are called the Floquet multipliers of the periodic orbit.
Theorem. The monodromy matrix $M(T)$ always has a unit eigenvalue with eigenvector $\dot{x}(0)=\dot{x}(T)=f(x(0))$.

This unit eigenvalue is named the trivial eigenvalue of the periodic orbit.
Theorem. If $x(t)$ is a periodic orbit of a $C^{2}$ flow $\varphi(t, x)$ that is linearly asymptotically stable (its monodromy matrix has all the eigenvalue inside the unit circle except the trivial one), then it is asymptotically stable.

The eigenvalues of $M(T)$ of largest magnitude have to be computed to detect bifurcations of the periodic orbits.

## Subspace iteration and Arnoldi methods

## Eigenvalue problems

## Rayleigh-Ritz method.

Assume $A$ is a real matrix of dimension $n \times n(n \gg 1)$, and that $V_{m}$ is of dimension $n \times m$ ( $m \ll n$ ) such that $V_{m}^{\top} V_{m}=I_{m}\left(I_{m}\right.$ identity of dimension $m$ ) and

$$
A V_{m}=V_{m} B_{m}
$$

with $B_{m}$ of dimension $m \times m$. This expressions says that the columns of $V_{m}$ span an invariant subspace of $A$ of dimension $m$. Moreover,

$$
\begin{gathered}
\text { if }(\lambda, u) \text { eigenpair of } B_{m} \quad\left(B_{m} u=\lambda u\right) \Longrightarrow\left(\lambda, V_{m} u\right) \text { eigenpair of } A \text { : } \\
A V_{m} u=V_{m} B_{m} u=\lambda V_{m} u .
\end{gathered}
$$

If the identity $A V_{m}=V_{m} B_{m}$ is not exact, the pairs $\left(\lambda, V_{m} u\right)$ are called Ritz values and vectors of $A$, respectively.
It is clear that since $V_{m}^{\top} V_{m}=I_{m}$ then $B_{m}=V_{m}^{\top} A V_{m}$.
There are two main methods to obtain approximate bases, $V_{m}$, of the subspace corresponding to the leading (largest magnitude) eigenvalues of a large-scale matrix $A$, subspace iteration and Arnoldi's method.

## Subspace iteration

Subspace iteration is implemented, for instance, in LOPSI and SRRIT.

1. Start: Choose an initial system of orthonormal vectors $V_{m}=\left[v_{1}, \ldots, v_{m}\right],(m \ll n)$
2. Iterate: Until convergence do:
(a) For $\mathrm{I}=1, \ldots, \mathrm{k}$ do
i. Compute $Z_{m}=A V_{m}$
ii. Orthonormalize $Z_{m}$ by computing $Z_{m}=Q_{m} R_{m}$, with $Q_{m}^{\top} Q_{m}=I$ and $R_{m}$ upper-triangular, and set $V_{m}=Q_{m}$
(b) Form $B_{m}=V_{m}^{\top} A V_{m}$ and compute the eigenpairs $\left(\lambda_{i}, z_{i}\right), i=1, \cdots, m$ of $B_{m}$ by the $Q R$ method (LAPACK)
(c) Test for convergence of eigenvalues and/or eigenvectors
3. Stop: When satisfied, compute the approximate eigenvectors of $A$ as $x_{i}=V_{m} z_{i}, i=1, \cdots, m$. The $\lambda_{i}, i=1, \cdots, m$ are the approximate eigenvalues.
Theorem. Suppose that the $n$ eigenvalues of $A$ are ordered by decreasing modulus as follows: $\left|\lambda_{1}\right| \geq\left|\lambda_{2}\right| \geq \cdots\left|\lambda_{m}\right|>\left|\lambda_{m+1}\right| \geq \cdots \geq\left|\lambda_{n}\right|$. If the initial set of vectors $V_{m}$ is not deficient in the eigenvectors corresponding to $\lambda_{1}, \cdots, \lambda_{m}$, and if the total number of iterations taken, $k$, is large enough, then the previous algorithm computes approximations $\hat{\lambda}_{i, k}$ to $\lambda_{i}(i=1, \ldots, m)$ with

$$
\left|\hat{\lambda}_{i, k}-\lambda_{i}\right|=O\left(\left|\frac{\lambda_{m+1}}{\lambda_{i}}\right|+\varepsilon_{i, k}\right)^{k}, \quad \lim _{k \rightarrow \infty} \varepsilon_{i, k}=0
$$

Moreover, if $\lambda_{i}$ is simple, then $\varepsilon_{i, k}=0$.

## Arnoldi's method

Arnoldi's method is implemented, for instance, in ARPACK.

1. Start: Choose an initial unitary vector $v_{1}$.
2. Iterate: Until convergence do:
(a) Compute the Arnoldi factorization $A V_{m}=V_{m} H_{m}+w_{m} e_{m}^{\top}$ of length $m$. The columns of $V_{m}$ form an orthonormal basis of $\mathcal{K}_{m}=\left\{v_{1}, A v_{1}, A^{2} v_{1}, \ldots, A^{m-1} v_{1}\right\}$.
(b) Compute the eigenpairs $\left(\lambda_{i}, z_{i}\right), i=1, \cdots, m$ of $H_{m}=V_{m}^{\top} A V_{m}$ by the $Q R$ method (LAPACK).
(c) Test for convergence of eigenvalues and/or eigenvectors. If not converged select a new initial vector $v_{1}$ from the Arnoldi factorization.
3. Stop: When satisfied, compute approximate eigenvectors of $A$ as $x_{i}=V_{m} z_{i}, i=1, \cdots, m$. The $\lambda_{i}, i=1, \cdots, m$ are the approximate eigenvalues.
Theorem. Suppose that the $n$ eigenvalues of $A$ are simple and that $\lambda_{2}, \ldots, \lambda_{n}$ are enclosed by a circle centered at $\xi$ and passing through $\lambda_{2}$, and that $\hat{\lambda}_{1}$ is the approximation to $\lambda_{1}$ obtained by Arnoldi's method, then

$$
\left|\hat{\lambda}_{1}-\lambda_{1}\right| \leq c\left|\frac{\lambda_{2}-\xi}{\lambda_{1}-\xi}\right|^{m-1}
$$

with c a constant. This gives the same error bound as $m-1$ steps of the power method applied to $A-\xi I$.

From the Arnoldi's decomposition

$$
A V_{m}=V_{m} H_{m}+w_{m} e_{m}^{\top}, \quad \text { with } \quad w_{m}=h_{m+1, m} v_{m+1}
$$

it is easy to find the residual $\left\|\left(A-\lambda_{i} I\right) x_{i}\right\|_{2}$.
Proposition. Let $z_{i}$ be an eigenvector of $H_{m}$ associated with the eigenvalue $\lambda_{i}$, and $x_{i}=V_{m} z_{i}$ the Ritz approximate eigenvector of $A$. Then,

$$
\left(A-\lambda_{i} I\right) x_{i}=h_{m+1, m}\left(e_{m}^{\top} z_{i}\right) v_{m+1}
$$

and, therefore

$$
\left\|\left(A-\lambda_{i} I\right) x_{i}\right\|_{2}=h_{m+1, m}\left|e_{m}^{\top} z_{i}\right| .
$$

Proof:

$$
\begin{gathered}
A V_{m} z_{i}=V_{m} H_{m} z_{i}+h_{m+1, m} v_{m+1} e_{m}^{\top} z_{i} \\
A x_{i}
\end{gathered}-\lambda_{i} x_{i}=h_{m+1, m}\left(e_{m}^{\top} z_{i}\right) v_{m+1}-1 .
$$

and $v_{m+1}$ is unitary.

## Eigenvalue Transformations

To find the leading (maximal real part) eigenvalues of $A v=\lambda v$ the following transformations can be used:

Shift-invert with real or complex shift:

$$
A v=\lambda v \Longrightarrow(A-\sigma I)^{-1} v=\mu v \quad \text { with } \quad \mu=1 /(\lambda-\sigma) .
$$

The circle $C(\sigma,|\lambda-\sigma|)$ in the $\lambda$-plane is mapped to the circle $C\left(0,|\lambda-\sigma|^{-1}\right)$ in the $\mu$-plane.


Generalized Cayley transformation:

$$
A v=\lambda v \Longrightarrow(A-\sigma I)^{-1}(A-\tau I) v=\mu v \quad \text { with } \quad \mu=(\lambda-\tau) /(\lambda-\sigma)
$$

The line $\operatorname{Re}(\lambda)=(\sigma+\tau) / 2$ is mapped to the unit circle and $\operatorname{Re}(\lambda)<(\sigma+\tau) / 2$ $(\operatorname{Re}(\lambda)>(\sigma+\tau) / 2)$ is mapped to the interior (exterior) of the unit circle.


Double complex shift: If $\sigma=\rho+i \theta$,

$$
A v=\lambda v \Longrightarrow(A-\sigma I)^{-1}(A-\bar{\sigma} I)^{-1} v=\mu v \quad \text { with } \quad \mu=1 /\left((\lambda-\rho)^{2}+\theta^{2}\right) .
$$

Systems with matrix $(A-\sigma I)(A-\bar{\sigma} I)=(A-\rho I)^{2}+\theta^{2} I$ have to be solved.



Exponential:

$$
A v=\lambda v \Longrightarrow \exp (T A) v=\mu v \quad \text { with } \quad \mu=\exp (\lambda T)
$$

The line $\operatorname{Re}(\lambda)=0$ is mapped to the unit circle and $\operatorname{Re}(\lambda)<0(\operatorname{Re}(\lambda)>0)$ is mapped to the interior (exterior) of the unit circle.



By taking logarithms, and using that $\lambda=\langle v, A v\rangle /\langle v, v\rangle$ to recover the undetermined imaginary part, $\lambda$ can be obtained from $\mu$.

The previous methods (subspace or Arnoldi iterations) can be used to find the eigenvalues $\mu$ with maximal modulus of the transformed problems.

## Continuation of fixed points of ODEs

Summarizing, it is possible to find the equilibria of the system of ODEs

$$
\dot{y}=f(y, p), \quad(y, p) \in \mathcal{U} \subset \mathbb{R}^{n} \times \mathbb{R}
$$

by Newton-Krylov methods by a generic continuation code if one can provide three subroutines:

- fun $(\mathrm{X}, \mathrm{H})$ computing the function $H=f(x, p)$ from $X=(x, p)$,
- $\operatorname{dfun}(\mathrm{X}, \delta \mathbf{X}, \delta \mathrm{H})$ which computes the action by the Jacobian $\delta H=D_{y} f(x, p) \delta x+D_{p} f(x, p) \delta p$ from $X=(x, p)$ and $\delta X=(\delta x, \delta p)$, and
- $\operatorname{prec}(\mathrm{X}, \mathrm{h}, \delta \mathrm{X}, \delta \mathrm{Z})$ which solves $\mathcal{M} \delta Z=\delta X$ from $X=(x, p), h=\left(h_{x}, h_{p}\right)$, and $\delta X=(\delta x, \delta p)$, $\mathcal{M}$ being an approximation of

$$
\left(\begin{array}{cc}
D_{x} f\left(x_{i}, p_{i}\right) & D_{p} f\left(x_{i}, p_{i}\right) \\
h_{x}^{\top} & h_{p}
\end{array}\right) .
$$

In the previous example (HP problem) we used an approximation of the form

$$
\mathcal{M}=\left(\begin{array}{cc}
M & 0 \\
0 & 1
\end{array}\right)
$$

## Continuation of periodic orbits of ODEs

To compute periodic orbits of

$$
\dot{y}=f(y, p), \quad(y, p) \in \mathcal{U} \subset \mathbb{R}^{n} \times \mathbb{R}
$$

by Newton-Krylov methods two subroutines are needed:

- fun $(\mathrm{X}, \mathrm{H})$ computing the function

$$
H(x, T, p)=\binom{x-\varphi(T, x, p)}{g(x, p)}
$$

from $X=(x, T, p), g(x, p)$ being a phase condition. This involves integrating

$$
\dot{y}=f(y, p) \quad \text { with initial condition } \quad y(0)=x \quad \text { during a time } T .
$$

- $\operatorname{dfun}(\mathrm{X}, \delta \mathrm{X}, \delta \mathrm{H})$ which computes the action by the Jacobian of the system

$$
\delta H=D H(x, T, p)(\delta x, \delta T, \delta p)=\binom{\delta x-D_{x} \varphi(T, x, p) \delta x-D_{p} \varphi(T, x, p) \delta p-f(y(T), p) \delta T}{D_{x} g(x, p) \delta x+D_{p} g(x, p) \delta p}
$$

from $X=(x, T, p)$ and $\delta X=(\delta x, \delta T, \delta p)$, where $y(T)=\varphi(T, x, p)$.

The matrix product

$$
D_{x} \varphi(T, x, p) \delta x+D_{p} \varphi(T, x, p) \delta p
$$

can be computed by integrating a first variational equation. If

$$
\begin{aligned}
y(t) & =\varphi(t, x, p) \\
y_{1}(t) & =D_{x} \varphi(t, x, p) \delta x+D_{p} \varphi(t, x, p) \delta p
\end{aligned}
$$

then $y_{1}$ satisfies

$$
\dot{y}_{1}=D_{y} f(y, p) y_{1}+D_{p} f(y, p) \delta p \quad \text { and } \quad y_{1}(0)=\delta x,
$$

because $\varphi(0, x, p)=x$.
This equation must be solved coupled with that for $y$,

$$
\begin{aligned}
\dot{y} & =f(y, p) & \text { with initial conditions } & y(0)
\end{aligned}=x
$$

Finally

$$
D_{x} \varphi(T, x, p) \delta x+D_{p} \varphi(T, x, p) \delta p=y_{1}(T) .
$$

## Continuation of fixed points of ODEs by time evolution

The equilibria of the system of ODEs

$$
\dot{y}=f(y, p), \quad(y, p) \in \mathcal{U} \subset \mathbb{R}^{n} \times \mathbb{R}
$$

can also be obtained as fixed points of the map

$$
x \rightarrow \varphi(T, x, p)
$$

$\varphi(t, x, p)$ being the solution of the ODEs with initial condition $x$, because

$$
f(x, p)=0 \Rightarrow x-\varphi(T, x, p)=0 .
$$

The arbitrary time $T$ must be large enough to have most of the spectrum of $D_{x} \varphi(T, x, p)$ clustered at the origin, but as short as possible to save computing time.
The matrix products required can be computed by integrating the first variational equation

$$
\begin{array}{lll}
\dot{x} & =f(x, p) & \text { with initial conditions } \\
\dot{y} & =D_{x} f(x, p) y+D_{p} f(x, p) \delta p & \\
& & y(0)=x \\
& & =\delta x .
\end{array}
$$

Then

$$
D_{x} \varphi(T, x, p) \delta x+D_{p} \varphi(T, x, p) \delta p=y(T) .
$$

This method can be used as a second option, for instance, when the continuation method applied to $f(x, p)$ fails due to the lack a good preconditioner for the linear systems or the eigenvalue problems.

## Examples and exercises

## Results for the HP problem

In this example $P e_{m}=P e_{h}=5, B=0.5, \gamma=25, \beta=3.50$, and $\theta_{r}=1$.

Equilibria and periodic orbits


Real part of eigenvalues versus D


Equilibria and periodic orbits


Real part of eigenvalues versus D


Recent Trends in Nonlinear Science 2016, January 25-29 2016, Sevilla - p. 43

Equilibria and periodic orbits


Modulus of the leading multiplier


Period of the periodic orbits


In this example $P e_{m}=P e_{h}=5, B=0.5, \gamma=25, \beta=3.00$, and $\theta_{r}=1$.

Equilibria and periodic orbits


Real part of eigenvalues versus D


Equilibria and periodic orbits


Real part of eigenvalues versus D


Equilibria and periodic orbits


Period of the periodic orbits


Modulus of the leading multiplier


## The one-dimensional Brusselator

Consider the system of PDEs

$$
\begin{aligned}
\partial_{t} X & =\left(D_{X} / L^{2}\right) \partial_{z z}^{2} X+X^{2} Y-(B+1) X+A \\
\partial_{t} Y & =\left(D_{Y} / L^{2}\right) \partial_{z z}^{2} Y-X^{2} Y+B X
\end{aligned}
$$

in the interval $z \in[0,1]$, with boundary conditions

$$
\begin{aligned}
& X(0)=X(1)=A \\
& Y(0)=Y(1)=B / A
\end{aligned}
$$

For any value of the parameters the problem has the trivial constant solution $X=A$ and $Y=B / A$. It undergoes Hopf bifurcations for

$$
L_{k}=k \pi \sqrt{\frac{D_{X}+D_{Y}}{B-A^{2}-1}}, \quad k=1,2 \ldots
$$

If $D_{X}=0.008, D_{Y}=0.004, A=2, B=5.45$ then $L_{k} \approx 0.5130 k$.
To implement the boundary conditions we substitute $X=\bar{X}+A, Y=\bar{Y}+B / A$ in the equations and boundary conditions to obtain, after removing the overbars the equations

$$
\begin{aligned}
\partial_{t} X & =\left(D_{X} / L^{2}\right) \partial_{z z}^{2} X+X^{2} Y+(B / A) X^{2}+2 A X Y+A^{2} Y+(B-1) X, & X(0)=X(1)=0, \\
\partial_{t} Y & =\left(D_{Y} / L^{2}\right) \partial_{z z}^{2} Y-X^{2} Y-(B / A) X^{2}-2 A X Y-A^{2} Y-B X, & Y(0)=Y(1)=0 .
\end{aligned}
$$

## Results for the Brusselator problem

In this example $D_{X}=0.008, D_{Y}=0.004, A=2, B=5.45$ and $L$ is the continuation parameter.


A period doubling cascade or tori in the Brusselator problem


## Continuation of codimension-one bifurcation points

 of equilibria and periodic orbits
## Continuation of bifurcation curves

Consider an autonomous system of ODEs

$$
\dot{y}=f(y, p), \quad(y, p) \in \mathcal{U} \subset \mathbb{R}^{n} \times \mathbb{R}^{2},
$$

depending on two parameters $p=\left(p_{1}, p_{2}\right)$ obtained after spatial discretization of a system of parabolic PDEs ( $n \gg 1$ ).
Let $y(t)=\varphi(t, x, p)$ be its solution with initial condition $y(0)=x$ at $t=0$ and for a fixed $p$.
We are interested in tracking curves of codimension-one bifurcations of periodic orbits in system with or without symmetries.

Let assume a matrix-free continuation code based on Newton-Krylov methods is available to follow the curves of solutions of

$$
H(X)=0
$$

with $X \in \mathcal{U} \subset \mathbb{R}^{m+1}$ and $H(X) \in \mathbb{R}^{m}$, which requires the user to provide an initial solution $X_{0}$, and two subroutines:

- $\operatorname{fun}(\mathrm{X}, \mathrm{h})$ which computes $h=H(X)$ from $X$, and
- $\operatorname{dfun}(\mathrm{X}, \delta \mathrm{x}, \delta \mathrm{h})$ which computes $\delta h=D_{X} H(X) \delta X$ from $X$, and $\delta X$.


## Saddle-node and period doubling bifurcations

The saddle-node $(\lambda=1)$ and period doubling $(\lambda=-1)$ bifurcations of periodic orbits are solutions of the system $H(x, u, T, p)=0$ given by

$$
\begin{aligned}
x-\varphi(T, x, p) & =0 \\
g(x) & =0 \\
\lambda u-\left(D_{x} \varphi(T, x, p) u-\frac{1}{2}(1+\lambda) \frac{\langle f, u\rangle}{\langle f, f\rangle} f\right) & =0 \\
\left\langle u_{r}, u\right\rangle & =1
\end{aligned}
$$

- $g(x)=0$ is a phase condition to select a single point on the periodic orbit. We use $g(x)=\left\langle v_{\pi}, x-x^{(\pi)}\right\rangle=0$.
- $f=f(x, p)$ is the vector field evaluated at $(x, p)$.
- $\left\langle u_{r}, u\right\rangle=1$ fixes the indetermined constant of the eigenvalue problem, $u_{r}$ being a reference vector. We use $u_{r}=u$.
- The last term of the third equation is Wieland's deflation, which guarantees the regularity of the system by shifting the +1 multiplier associated with $f(x, p)$ to zero.
$X=(x, u, T, p)$ has dimension $2 n+3$, and the $2 n+2$ equations define the curve of solutions.

In order to compute $H(x, u, T, p)$, we define

$$
\begin{aligned}
y(t) & =\varphi(t, x, p) \\
y_{1}(t) & =D_{x} \varphi(t, x, p) u
\end{aligned}
$$

and, taking into account that

$$
D_{t} D_{x} \varphi(t, x, p)=D_{y} f(\varphi(t, x, p), p) D_{x} \varphi(t, x, p), \text { and } D_{x} \varphi(0, x, p)=I
$$

the following system has to be integrated during a time $T$

$$
\begin{aligned}
\dot{y} & =f(y, p), & y(0) & =x \\
\dot{y}_{1} & =D_{y} f(y, p) y_{1}, & y_{1}(0) & =u
\end{aligned}
$$

Then

$$
\begin{aligned}
\varphi(T, x, p) & =y(T) \\
D_{x} \varphi(T, x, p) u & =y_{1}(T)
\end{aligned}
$$

The action of $D_{X} H(x, u, T, p)$ on $(\delta x, \delta u, \delta T, \delta p)$ is

$$
\delta x-D_{t} \varphi(T, x, p) \delta T-D_{x} \varphi(T, x, p) \delta x-D_{p} \varphi(T, x, p) \delta p
$$

$D g(x) \delta x$,

$$
\lambda \delta u-D_{t x}^{2} \varphi(T, x, p)(u, \delta T)-D_{x x}^{2} \varphi(T, x, p)(u, \delta x)-D_{x p}^{2} \varphi(T, x, p)(u, \delta p)
$$

$$
-D_{x} \varphi(T, x, p) \delta u
$$

$$
+\frac{1+\lambda}{2\langle w, w\rangle}\left(\langle w, u\rangle z+\left(\langle z, u\rangle+\langle w, \delta u\rangle-\frac{2\langle w, z\rangle}{\langle w, w\rangle}\langle w, u\rangle\right) w\right)
$$

$$
\left\langle u_{r}, \delta u\right\rangle,
$$

where $w=f(x, p)$ and $z=D_{y} f(x, p) \delta x+D_{p} f(x, p) \delta p$. Lets define

$$
\begin{aligned}
y(t) & =\varphi(t, x, p) \\
y_{1}(t) & =D_{x} \varphi(t, x, p) u \\
y_{2}(t) & =D_{x} \varphi(t, x, p) \delta x+D_{p} \varphi(t, x, p) \delta p \\
y_{3}(t) & =D_{x x}^{2} \varphi(t, x, p)(u, \delta x)+D_{x p}^{2} \varphi(t, x, p)(u, \delta p) \\
y_{4}(t) & =D_{x} \varphi(t, x, p) \delta u \\
D_{t} \varphi(T, x, p) \delta T & =f(y(T), p) \delta T \\
D_{t x}^{2} \varphi(T, x, p)(u, \delta T) & =\delta T D_{y} f(\varphi(T, x, p), p) D_{x} \varphi(T, x, p) u=\delta T D_{y} f(y(T), p) y_{1}(T)
\end{aligned}
$$

$$
\begin{aligned}
y(t) & =\varphi(t, x, p) \\
y_{1}(t) & =D_{x} \varphi(t, x, p) u \\
y_{2}(t) & =D_{x} \varphi(t, x, p) \delta x+D_{p} \varphi(t, x, p) \delta p \\
y_{3}(t) & =D_{x x}^{2} \varphi(t, x, p)(u, \delta x)+D_{x p}^{2} \varphi(t, x, p)(u, \delta p) \\
y_{4}(t) & =D_{x} \varphi(t, x, p) \delta u
\end{aligned}
$$

the system which must be integrated to obtain $y(T), y_{i}(T), i=1, \ldots, 4$ is

$$
\begin{aligned}
\dot{y} & =f(y, p), & y(0) & =x \\
\dot{y}_{1} & =D_{y} f(y, p) y_{1}, & y_{1}(0) & =u \\
\dot{y}_{2} & =D_{y} f(y, p) y_{2}+D_{p} f(y, p) \delta p, & y_{2}(0) & =\delta x \\
\dot{y}_{3} & =D_{y} f(y, p) y_{3}+D_{y y}^{2} f(y, p)\left(y_{1}, y_{2}\right)+D_{y p}^{2} f(y, p)\left(y_{1}, \delta p\right), & y_{3}(0) & =0 \\
\dot{y}_{4} & =D_{y} f(y, p) y_{4}, & y_{4}(0) & =\delta u .
\end{aligned}
$$

## Neimark-Sacker bifurcations

The Hopf bifurcations of periodic orbits with multiplier $e^{i \theta}$ and eigenvector $u+i v$ are solutions of the system $H(x, u, v, T, \theta, p)=0$ given by

$$
\begin{aligned}
x-\varphi(T, x, p) & =0 \\
g(x) & =0 \\
u \cos \theta-v \sin \theta-D_{x} \varphi(T, x, p) u & =0 \\
u \sin \theta+v \cos \theta-D_{x} \varphi(T, x, p) v & =0 \\
\langle u, u\rangle+\langle v, v\rangle & =1 \\
\langle u, v\rangle & =0
\end{aligned}
$$

- $g(x)=0$ is the phase condition $g(x)=\left\langle v_{\pi}, x-x^{(\pi)}\right\rangle=0$.
- The third and fourth equations are the real and imaginary parts of

$$
e^{i \theta}(u+i v)-D_{x} \varphi(T, x, p)(u+i v)=0
$$

- The two last equations uniquely determine the eigenvector $u+i v$.

Now $X=(x, u, v, T, \theta, p)$ has dimension $3 n+4$, and the $3 n+3$ equations define the curve of solutions.

## Pitchfork bifurcations

If the initial system is $\mathcal{T}$-invariant, $f(\mathcal{T} x, p)=\mathcal{T} f(x, p)$ with $\mathcal{T}^{2}=I$, and $\mathcal{T} x=x$, the pitchfork bifurcation points of periodic orbits are solutions of the system $H(x, u, T, \xi, p)=0$ are given by

$$
\begin{aligned}
x-\varphi(T, x, p)+\xi \phi & =0 \\
g(x) & =0 \\
\langle x, \phi\rangle & =0 \\
u-\left(D_{x} \varphi(T, x, p) u-\frac{\langle f, u\rangle}{\langle f, f\rangle} f\right) & =0 \\
\left\langle u_{r}, u\right\rangle & =1
\end{aligned}
$$

- The slack variable $\xi$ and the third equation are introduced to make the system regular. Moreover $\xi=0$ at the solution.
- $g(x)=0$ is the phase condition $g(x)=\left\langle v_{\pi}, x-x^{(\pi)}\right\rangle=0$.
- $\phi$ is a given antisymmetric vector, $\mathcal{T} \phi=-\phi$.
- The last equation uniquely determines the eigenvector $u$.

Now $X=(x, u, T, \xi, p)$ has dimension $2 n+4$, and the $2 n+3$ equations define the curve of solutions.

Thermal convection in binary fluid mixtures
The equations in $\Omega=[0, \Gamma] \times[0,1]$ for the perturbation of the basic state $\left(\mathbf{v}_{c}=0, T_{c}=T_{c}(0)-z\right.$, and $\left.C_{c}=C_{c}(0)-z\right)$ in non-dimensional form are

$$
\begin{aligned}
& \partial_{t} \mathbf{v}+(\mathbf{v} \cdot \nabla) \mathbf{v}=\sigma \Delta \mathbf{v}-\nabla p+\sigma R a(\Theta+S C) \hat{e}_{z}, \\
& \partial_{t} \Theta+(\mathbf{v} \cdot \nabla) \Theta=\Delta \Theta+v_{z}, \\
& \partial_{t} C+(\mathbf{v} \cdot \nabla) C=L(\Delta C-\Delta \Theta)+v_{z}, \\
& \nabla \cdot \mathbf{v}=0 .
\end{aligned}
$$

The boundary conditions are non-slip for $\mathbf{v}$, constant temperatures at top and bottom and insulating lateral walls for $\Theta=T-T_{c}$, and impermeable boundaries for $C$.

The parameters are
$\Gamma \quad$ Aspect ratio (4)
$S \quad$ Separation ratio (-0.1)
$L \quad$ Lewis number (0.03)
$\sigma \quad$ Prandtl number (control)
Ra Rayleigh number (control)


To simplify the system, a streamfunction $\mathbf{v}=\left(-\partial_{z} \psi, \partial_{x} \psi\right)$, and an auxiliary function $\eta=C-\Theta$ are used. Then

$$
\begin{aligned}
\partial_{t} \Delta \psi+J(\psi, \Delta \psi) & =\sigma \Delta^{2} \psi+\sigma R a\left[(S+1) \partial_{x} \Theta+S \partial_{x} \eta\right], \\
\partial_{t} \Theta+J(\psi, \Theta) & =\Delta \Theta+\partial_{x} \psi, \\
\partial_{t} \eta+J(\psi, \eta) & =L \Delta \eta-\Delta \Theta,
\end{aligned}
$$

with $J(f, g)=\partial_{x} f \partial_{z} g-\partial_{z} f \partial_{x} g$. The boundary conditions are now

$$
\begin{array}{rll}
\psi=\partial_{n} \psi=\partial_{n} \eta=0 & \text { at } & \partial \Omega, \\
\Theta=0 & \text { at } & z=0,1, \\
\partial_{x} \Theta=0 & \text { at } & x=0, \Gamma .
\end{array}
$$

The symmetry group of the equations is $\mathbb{Z}_{2} \times \mathbb{Z}_{2}$ generated by the two reflections:

$$
\begin{aligned}
& R_{x}:(t, x, z, \psi, \Theta, \eta) \rightarrow(t, \Gamma-x, z,-\psi, \Theta, \eta) \\
& R_{z}:(t, x, z, \psi, \Theta, \eta) \rightarrow(t, x, 1-z,-\psi,-\Theta,-\eta)
\end{aligned}
$$

## Variational equations

$$
\begin{aligned}
\partial_{t} \Delta \psi_{1}+J\left(\psi, \Delta \psi_{1}\right)+J\left(\psi_{1}, \Delta \psi\right)= & \sigma \Delta^{2} \psi_{1}+\sigma R a\left[(S+1) \partial_{x} \Theta_{1}+S \partial_{x} \eta_{1}\right], \\
\partial_{t} \Theta_{1}+J\left(\psi, \Theta_{1}\right)+J\left(\psi_{1}, \Theta\right)= & \Delta \Theta_{1}+\partial_{x} \psi_{1}, \\
\partial_{t} \eta_{1}+J\left(\psi, \eta_{1}\right)+J\left(\psi_{1}, \eta\right)= & L \Delta \eta_{1}-\Delta \Theta_{1}, \\
\partial_{t} \Delta \psi_{2}+J\left(\psi, \Delta \psi_{2}\right)+J\left(\psi_{2}, \Delta \psi\right)= & \sigma \Delta^{2} \psi_{2}+\sigma R a\left[(S+1) \partial_{x} \Theta_{2}+S \partial_{x} \eta_{2}\right]+\delta \sigma \Delta^{2} \psi \\
& +(\sigma \delta R a+\delta \sigma R a)\left[(S+1) \partial_{x} \Theta+S \partial_{x} \eta\right], \\
\partial_{t} \Theta_{2}+J\left(\psi, \Theta_{2}\right)+J\left(\psi_{2}, \Theta\right)= & \Delta \Theta_{2}+\partial_{x} \psi_{2}, \\
\partial_{t} \eta_{2}+J\left(\psi, \eta_{2}\right)+J\left(\psi_{2}, \eta\right)= & L \Delta \eta_{2}-\Delta \Theta_{2}, \\
\partial_{t} \Delta \psi_{3}+J\left(\psi, \Delta \psi_{3}\right)+J\left(\psi_{3}, \Delta \psi\right)= & \sigma \Delta^{2} \psi_{3}+\sigma R a\left[(S+1) \partial_{x} \Theta_{3}+S \partial_{x} \eta_{3}\right]+\delta \sigma \Delta^{2} \psi_{1} \\
& +(\sigma \delta R a+\delta \sigma R a)\left[(S+1) \partial_{x} \Theta_{1}+S \partial_{x} \eta_{1}\right] \\
& -J\left(\psi_{1}, \Delta \psi_{2}\right)-J\left(\psi_{2}, \Delta \psi_{1}\right), \\
\partial_{t} \Theta_{3}+J\left(\psi, \Theta_{3}\right)+J\left(\psi_{3}, \Theta\right)= & \Delta \Theta_{3}+\partial_{x} \psi_{3}-J\left(\psi_{1}, \Theta_{2}\right)-J\left(\psi_{2}, \Theta_{1}\right), \\
\partial_{t} \eta_{3}+J\left(\psi, \eta_{3}\right)+J\left(\psi_{3}, \eta\right)= & L \Delta \eta_{3}-\Delta \Theta_{3}-J\left(\psi_{1}, \eta_{2}\right)-J\left(\psi_{2}, \eta_{1}\right) .
\end{aligned}
$$

## Discretization

The functions $\psi, \Theta$, and $\eta$ are approximated by a pseudo-spectral method. Collocation on a mesh of $n_{x} \times n_{z}=64 \times 16(n=3072)$ Gauss-Lobatto points is used.

Higher resolutions have been used to check the results.
The stiff system of ODEs obtained can be written as

$$
B \dot{u}=L u+N(u)
$$

with $u=\left(\psi_{i j}, \Theta_{i j}, \eta_{i j}\right)$.
It is integrated by using fifth-order BDF-extrapolation formulas:

$$
\frac{1}{\Delta t} B\left(\gamma_{0} u^{n+1}-\sum_{i=0}^{k-1} \alpha_{i} u^{n-i}\right)=\sum_{i=0}^{k-1} \beta_{i} N\left(u^{n-i}\right)+L u^{n+1}
$$

The initial points are obtained by a fully implicit BDF method.

## Some results for $\sigma=0.6$



## Some results for $\sigma=0.6$



Fixed points


## Periodic orbits



Curves of bifurcations


Curves of bifurcations


Curves of bifurcations


## Period



## Codimension-two points






## Performance



Relative distance between Newton iterates and number of GMRES iterations for the pitchfork and one of the Neimark-Sacker curves.

## Continuation of invariant tori

## Poincaré maps and its derivatives

Given the system of ODEs

$$
\dot{y}=f(y, p), \quad(y, p) \in \mathcal{U} \subset \mathbb{R}^{n} \times \mathbb{R}
$$

let $\varphi(t, x, p)$ is its solution with initial condition $x$ at $t=0$ for a fixed value of $p(\varphi(0, x, p)=x)$, and suppose there is a periodic orbit which cuts transversely a hyperplane $\Sigma_{1}$ given by

$$
\Sigma_{1}=\left\{y \in \mathbb{R}^{n} /\left\langle v_{1}, y-x_{1}^{\Sigma}\right\rangle=0\right\} .
$$

Let $\mathcal{V} \subset \Sigma_{1}$ be a neighborhood of the intersection. The Poincaré map $P: \mathcal{V} \subset \Sigma_{1} \rightarrow \Sigma_{1}$ is defined as


$$
P(x, p)=\varphi(t(x, p), x, p),
$$

$t(x, p)$ being the first positive time for which

$$
\varphi(t(x, p), x, p) \in \Sigma_{1}
$$

with

$$
\operatorname{sign}\left\langle v_{1}, f(x, p)\right\rangle=\operatorname{sign}\left\langle v_{1}, f(P(x, p), p)\right\rangle .
$$

The action of the Jacobian of $P(x, p)=\varphi(t(x, p), x, p)$ on $(\delta x, \delta p)$ with $\left\langle v_{1}, \delta x\right\rangle=0$ is

$$
w=D_{x} P(x, p) \delta x+D_{p} P(x, p) \delta p=y_{1}-\frac{\left\langle v_{1}, y_{1}\right\rangle}{\left\langle v_{1}, z\right\rangle} z
$$

where $z=f(P(x, p), p), y_{1}$ is the solution, at the arrival time, $t(x, p)$, of the first variational equation


$$
\begin{aligned}
\dot{y} & =f(y, p) \\
\dot{y_{1}} & =D_{y} f(y, p) y_{1}+D_{p} f(y, p) \delta p
\end{aligned}
$$

with initial conditions

$$
\begin{aligned}
y(0) & =x \\
y_{1}(0) & =\delta x
\end{aligned}
$$

Each matrix product requires, the time integration of a system of $2 n$ equations.

To see that

$$
D_{x} P(x, p) \delta x+D_{p} P(x, p) \delta p=y_{1}-\frac{\left\langle v_{1}, y_{1}\right\rangle}{\left\langle v_{1}, z\right\rangle} z
$$

first differentiate $P(x, p)=\varphi(t(x, p), x, p)$ to obtain

$$
\begin{aligned}
& D_{x} P(x, p) \delta x+D_{p} P(x, p) \delta p=D_{x} \varphi(t(x, p), x, p) \delta x+D_{p} \varphi(t(x, p), x, p) \delta p \\
& +f(\varphi(t(x, p), x, p), p)\left(D_{x} t(x, p) \delta x+D_{p} t(x, p) \delta p\right)=y_{1}(t(x, p))+c z
\end{aligned}
$$

if

$$
\begin{aligned}
y_{1}(t) & =D_{x} \varphi(t, x, p) \delta x+D_{p} \varphi(t, x, p) \delta p \\
c & =D_{x} t(x, p) \delta x+D_{p} t(x, p) \delta p \in \mathbb{R} \\
z & =f(\varphi(t(x, p), x, p), p)
\end{aligned}
$$

Then

- $y_{1}$ satisfies $\dot{y}_{1}=D_{y} f(y, p) y+D_{p} f(y, p) \delta p$ and $y_{1}(0)=\delta x$ (because $\varphi(0, x, p)=x$ ), $y(t)=\varphi(t, x, p)$ being the solution of $\dot{y}=f(y, p)$ with $y(0)=x$.
- Since $h_{1}(x, p) \equiv\left\langle v_{1}, P(x, p)-x_{1}^{\Sigma}\right\rangle=0 \forall x \in \mathcal{V}$ and $p$, and then

$$
D_{x} h_{1}(x, p) \delta x+D_{p} h_{1}(x, p) \delta p=\left\langle v_{1}, y_{1}\right\rangle+c\left\langle v_{1}, z\right\rangle=0
$$

## Continuation of periodic orbits of ODEs

The periodic orbits $\dot{y}=f(y, p)$ can be computed as fixed points of a parameterized version of the Poincaré map.
If $v_{1_{k}}$ is the largest component of $v_{1}$, lets define $R_{k}$ as the orthogonal projection from $\Sigma_{1}$ onto the hyperplane $y_{k}=0$.


The parameterized map

$$
\bar{P}(\bar{x}, p)=R_{k}\left(P\left(R_{k}^{-1}(\bar{x}), p\right)\right),
$$

and its fixed points verifying

$$
\bar{x}-\bar{P}(\bar{x}, p)=0, \quad \bar{x} \in \mathbb{R}^{n-1}
$$

are in one-to-one correspondence with those of $P$ by the map $x=R_{k}^{-1}(\bar{x})$.

By applying the chain rule to $\bar{P}(\bar{x}, p)=R_{k}\left(P\left(R_{k}^{-1}(\bar{x}), p\right)\right)$

$$
D_{\bar{x}} \bar{P}(\bar{x}, p) \delta \bar{x}+D_{p} \bar{P}(\bar{x}, p) \delta p=R_{k}\left(D_{x} P(x, p) D R_{k}^{-1}(\bar{x}) \delta \bar{x}+D_{p} P(x, p) \delta p\right)=R_{k} w
$$

whith

$$
w=D_{x} P(x, p) \delta x+D_{p} P(x, p) \delta p=y-\frac{\left\langle v_{1}, y\right\rangle}{\left\langle v_{1}, z\right\rangle} z,
$$

where $\delta x=D R_{k}^{-1}(\bar{x}) \delta \bar{x}, z=f(P(x, p), p)$, and $y_{1}$ is the solution, at the arrival time, $t(x, p)$, of the first variational equation


$$
\begin{aligned}
\dot{x} & =f(x, p) \\
\dot{y}_{1} & =D_{x} f(x, p) y_{1}+D_{p} f(x, p) \delta p
\end{aligned}
$$

with initial conditions

$$
\begin{aligned}
x(0) & =x=R_{k}^{-1}(\bar{x}) \\
y_{1}(0) & =\delta x=D R_{k}^{-1}(\bar{x}) \delta \bar{x}
\end{aligned}
$$

## 'Continuation' of invariant 2-tori of ODEs (first method)

Let $P: \mathcal{V} \subset \Sigma_{1} \rightarrow \Sigma_{1}$ be the Poincaré map defined on a hyperplane $\Sigma_{1}$, and $\Sigma_{2}$ another hyperplane, given by $\left\langle v_{2}, x-x_{2}^{\Sigma}\right\rangle=0$, transversal to both $\Sigma_{1}$ and the invariant 2-tori. Fix $\varepsilon$, and define the map $G(\bar{x}, p): \mathcal{U} \subset \mathbb{R}^{n-2} \times \mathbb{R} \rightarrow \mathbb{R}^{n-2}$ as follows. If $R: \Sigma_{1} \cap \Sigma_{2} \rightarrow \mathbb{R}^{n-2}$ is a parameterization of $\Sigma_{1} \cap \Sigma_{2}$ and $\bar{x} \in \mathcal{U}$ let $x=R^{-1}(\bar{x})$ and $z_{j}=P^{k_{j}}(x, p), j=1, \ldots, q+1$ be the first $q+1$ powers of $P$ on $x$ such that $\left\|P^{k_{j}}(x, p)-x\right\|<\varepsilon$. Then


The fixed points of the map $G(\bar{x}, p)$ are in one-to-one correspondence, by the map $x=R^{-1}(\bar{x})$, with approximations of the points of the invariant 2-tori in $\Sigma_{1} \cap \Sigma_{2}$.

The radius $\varepsilon$ defining $G$ must be varied adaptively during the 'continuation' process.

The action of the Jacobian of $G(\bar{x}, p)=R\left(\sum_{j=1}^{q+1} l_{j}(0) P^{k_{j}}\left(R^{-1}(\bar{x}), p\right)\right)$, on a $(\delta \bar{x}, \delta p)$ reduces to the case of the differential of the Poincaré map

$$
\begin{aligned}
D_{\bar{x}} G(\bar{x}, p) \delta \bar{x}+D_{p} G(\bar{x}, p) \delta p & =R \sum_{j=1}^{q+1}\left[l_{j}(0)\left(D_{x} P^{k_{j}}(x, p) \delta x+D_{p} P^{k_{j}}(x, p) \delta p\right)\right. \\
& \left.+P^{k_{j}}(x, p) \sum_{i=1}^{q+1} \partial_{\mu_{i}} l_{j}(0)\left\langle v_{2}, D_{x} P^{k_{i}}(x, p) \delta x+D_{p} P^{k_{i}}(x, p) \delta p\right\rangle\right]
\end{aligned}
$$

with $x=R^{-1}(\bar{x})$ and $\delta x=D_{\bar{x}} R^{-1}(\bar{x}) \delta \bar{x}$, and where

$$
\left\langle v_{2}, D_{x} P^{k_{i}}(x, p) \delta x+D_{p} P^{k_{i}}(x, p) \delta p\right\rangle=D_{x} \mu_{i} \delta x+D_{p} \mu_{i} \delta p
$$

The derivatives $\partial_{\mu_{i}} l_{j}(0)$ of

$$
l_{j}(0)=\prod_{\substack{i=1 \\ i \neq j}}^{q+1} \frac{\mu_{i}}{\mu_{i}-\mu_{j}}
$$

are trivial.

## 'Continuation' of invariant 2-tori of ODEs (second method)

Let $P: \mathcal{V} \subset \Sigma_{1} \rightarrow \Sigma_{1}$ be the Poincaré map defined on a hyperplane $\Sigma_{1}$, and $\Sigma_{2}$ another hyperplane, given by $\left\langle v_{2}, x-x_{2}^{\Sigma}\right\rangle=0$, transversal to both $\Sigma_{1}$ and the invariant 2-tori. Let $\mu_{1}, \cdots, \mu_{q+1}$ be $q+1$ fixed coordinates along the line $x=x_{2}^{\Sigma}+\mu v_{2}$. Fix $\varepsilon$, and define the map $G(X, p): \mathcal{U} \subset \mathbb{R}^{(n-1)(q+1)} \times \mathbb{R} \rightarrow \mathbb{R}^{(n-1)(q+1)}$ as follows. If $X=\left(x_{1}, \cdots, x_{q+1}\right) \in \mathcal{U}$ let $z_{j}=P^{k_{j}^{\prime}}\left(x_{j}, p\right)$ be the first power of $P$ on $x_{j}$ such that $\left\|P^{k_{j}^{\prime}}\left(x_{j}, p\right)-x_{j}\right\|<\varepsilon$ for $j=1, \ldots, q+1$. Then


$$
\begin{gathered}
G(X, p)=X^{\prime}=Z(X, p) \tilde{V}(X, p)^{-1} V, \quad \text { with } \\
X^{\prime}=\left(x_{1}^{\prime}, \cdots, x_{q+1}^{\prime}\right), Z(X, p)=\left(z_{1}, \cdots, z_{q+1}\right), \\
\text { and } V \text { and } \tilde{V} \text { the Vandermonde matrices } \\
\qquad V=\left(\begin{array}{ccc}
1 & \cdots & 1 \\
\mu_{1} & \cdots & \mu_{q+1} \\
\cdots & \cdots & \cdots \\
\mu_{1}^{q} & \cdots & \mu_{q+1}^{q}
\end{array}\right)
\end{gathered}
$$

associated with the $\mu_{j}$ and the $\tilde{\mu}_{j}=\left\langle v_{2}, P^{k_{j}^{\prime}}\left(x_{j}, p\right)-x_{2}^{\Sigma}\right\rangle, \quad j=1, \cdots, q+1$ respectively. The fixed points of the map $G(X, p)$ approximate an arc of the invariant curve in $\Sigma_{1}$.

The action by the Jacobian of $G=Z \tilde{V}^{-1} V$ also reduces to that of the Poincaré map. If $\delta X=\left(\delta x_{1}, \cdots, \delta x_{q+1}\right)$ then
$D G(X, p)(\delta X, \delta p)=$

$$
\left[D Z(X, p)(\delta X, \delta p)-Z(X, p) \tilde{V}(X, p)^{-1} D \tilde{V}(X, p)(\delta X, \delta p)\right] \tilde{V}(X, p)^{-1} V
$$

where

$$
\begin{gathered}
D Z(X, p)(\delta X, \delta p)=\left(D P^{k_{1}^{\prime}}\left(x_{1}, p\right)\left(\delta x_{1}, \delta p\right), \cdots, D P^{k_{q+1}^{\prime}}\left(x_{q+1}, p\right)\left(\delta x_{q+1}, \delta p\right)\right), \\
D \tilde{V}(X, p)(\delta X, \delta p)=\left(\begin{array}{ccc}
0 & \cdots & 0 \\
1 & \cdots & 1 \\
2 \tilde{\mu}_{1} & \cdots & 2 \tilde{\mu}_{q+1} \\
\cdots \cdots & \cdots & \cdots \\
q \tilde{\mu}_{1}^{q-1} & \cdots & q \tilde{\mu}_{q+1}^{q-1}
\end{array}\right)\left(\begin{array}{ccc}
\eta_{1} & \cdots & 0 \\
\vdots & \ddots & \vdots \\
0 & \cdots & \eta_{q+1}
\end{array}\right),
\end{gathered}
$$

and $\eta_{j}=\left\langle v_{2}, D P^{k_{j}^{\prime}}\left(x_{j}, p\right)\left(\delta x_{j}, \delta p\right)\right\rangle$. In short, $D G=\left[D Z-Z \tilde{V}^{-1} D \tilde{V}\right] \tilde{V}^{-1} V$.
The radius $\varepsilon$ and the position of the $\mu_{j}$ defining $G$ must be varied adaptively during the continuation process.

## Comparison of the two methods to compute invariant tori

Assuming the same degree $q$ of interpolation is used to approximate the invariant curve:

- The first method requires computing increasing powers $k_{1}<k_{2}<\cdots<k_{q+1}$ of the Poincaré map of a single initial condition $x$. This is a sequential process.
- The second requires the first returning powers $k_{j}^{\prime}$ for a collection of initial conditions $x_{j}$. They can be computed in parallel. If the $x_{j}$ are close enough all the $k_{j}^{\prime}$ will be the same.
- If the two methods are applied to the same problem with the same $\varepsilon$, and close initial conditions, one can expect $k_{1}^{\prime}=\cdots=k_{q+1}^{\prime}=k_{1}$. Therefore if parallelism is used in the second method the wall-clock time to compute the map or the action by the Jacobian is reduced by a factor $k_{q+1} / k_{1}$.
- If the number of iterations of Newton's method and the linear solver are the same for both methods, the time (wall-clock) to do the continuation will be reduced essentially by a factor $k_{q+1} / k_{1}$.
- If $k_{2} \approx 2 k_{1}, \cdots, k_{q+1} \approx(q+1) k_{1}$ then $k_{q+1} / k_{1} \approx q+1$.
- If $k_{1}^{\prime}=\cdots=k_{q+1}^{\prime}$ the speed-up of the second method (time sequential/time parallel) is essentially $q+1$.
- Since the second method requires, in principle, lower powers of the Poincaré map, it should be best suited for the computation of weakly unstable invariant tori.

Thermal convection in binary fluid mixtures
The equations in $\Omega=[0, \Gamma] \times[0,1]$ for the perturbation of the basic state $\left(\mathbf{v}_{c}=0, T_{c}=T_{c}(0)-z\right.$, and $\left.C_{c}=C_{c}(0)-z\right)$ in non-dimensional form are

$$
\begin{aligned}
& \partial_{t} \mathbf{v}+(\mathbf{v} \cdot \nabla) \mathbf{v}=\sigma \Delta \mathbf{v}-\nabla p+\sigma R a(\Theta+S C) \hat{e}_{z}, \\
& \partial_{t} \Theta+(\mathbf{v} \cdot \nabla) \Theta=\Delta \Theta+v_{z}, \\
& \partial_{t} C+(\mathbf{v} \cdot \nabla) C=L(\Delta C-\Delta \Theta)+v_{z}, \\
& \nabla \cdot \mathbf{v}=0 .
\end{aligned}
$$

The boundary conditions are non-slip for $\mathbf{v}$, constant temperatures at top and bottom and insulating lateral walls for $\Theta=T-T_{c}$, and impermeable boundaries for $C$.

The parameters are
$\Gamma \quad$ Aspect ratio (4)
$S \quad$ Separation ratio (-0.1)
$L \quad$ Lewis number (0.03)
$\sigma \quad$ Prandtl number (control)
Ra Rayleigh number (control)


## Some results for $\sigma=0.6$



## Invariant tori for $\sigma=0.6$




- Beginning of the branch: $R a=2066.74$
- $1 / 7$-resonance interval $2102.79<R a<2102.80$
- Pitchfork bifurcation $R a \approx 2115.92$
- $1 / 8$-resonance interval $2116.18 \leq R a \leq 2116.20$.
- First period doubling $R a \approx 2118.40$
- Second period doubling $R a \approx 2118.55$
- Breakdown of the torus $R a \approx 2118.60$

Comparison of the two methods to compute invariant tori

| Meth. | Version | $q$ | Time | N. sol. | Time/N. sol. | Ratio $_{1}$ | Ratio $_{2}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | serial | 2 | 3725 | 19 | 196 | 4.97 | 4.71 |
| 1 | serial | 3 | 3960 | 19 | 208 | 5.27 | 4.99 |
| 1 | serial | 4 | 5019 | 19 | 264 | 7.34 | 7.34 |
| 1 | serial | 5 | 5664 | 19 | 298 | 7.97 | 7.97 |
| 2 | serial | 2 | 2255 | 18 | 125 | 3.01 | 3.01 |
| 2 | serial | 3 | 2830 | 18 | 157 | 3.76 | 3.76 |
| 2 | serial | 4 | 3563 | 19 | 188 | 5.21 | 5.21 |
| 2 | serial | 5 | 4333 | 19 | 228 | 6.09 | 6.09 |
| 2 | parallel | 2 | 749 | 18 | 42 | 1.00 | 1.00 |
| 2 | parallel | 3 | 752 | 18 | 42 | 1.00 | 1.00 |
| 2 | parallel | 4 | 684 | 19 | 36 | 0.91 | 0.87 |
| 2 | parallel | 5 | 711 | 19 | 37 | 0.95 | 0.90 |

Comparison of the wall-clock times for the different algorithms, implementations, and interpolation degrees. The continuation corresponds to the interval $R a \in[210,2115]$.

The Arnold's tongue of $\rho=1 / 8(\sigma=0.6)$




## Computation of the limits of the $1 / 8$ tongue



Recent Trends in Nonlinear Science 2016, January 25-29 2016, Sevilla - p. 88

## Computation of the limits of the $1 / 8$ tongue



Recent Trends in Nonlinear Science 2016, January 25-29 2016, Sevilla - p. 89

## Continuation of periodic orbits by multiple shooting

## Computation of periodic orbits by parallel multiple shooting

Consider a system of autonomous differential equations

$$
\dot{x}=f(x, p), \quad(x, p) \in \mathcal{U} \times I \subset \mathbb{R}^{n+1}
$$

$m$ hyperplanes $\Pi_{i}$ intersecting transversally an initial periodic orbit, the partial Poincaré maps $P_{i}(i=1 \cdots, m)$, and the Poincaré map on $\Pi_{1}, P$. They satisfy

$$
P\left(x_{1}, p\right)=\left(P_{m} \circ P_{m-1} \circ \cdots \circ P_{1}\right)\left(x_{1}, p\right),
$$


if $x_{i+1}=P_{i}\left(x_{i}, p\right), i=1, \cdots, m-1$,

$$
D_{x} P\left(x_{1}, p\right)=D_{x} P_{m}\left(x_{m}, p\right) D_{x} P_{m-1}\left(x_{m-1}, p\right) \cdots D_{x} P_{1}\left(x_{1}, p\right) .
$$

We also define $X=\left(x_{1}, \cdots, x_{m}\right) \in \mathbb{R}^{m n}$, and the map $\mathcal{G}$ as

$$
\mathcal{G}(X, p)=\left(x_{1}-P_{m}\left(x_{m}, p\right), x_{2}-P_{1}\left(x_{1}, p\right), \cdots, x_{m}-P_{m-1}\left(x_{m-1}, p\right)\right) .
$$

The points $x_{1}, \cdots, x_{m}$ are on a periodic orbit if $\mathcal{G}(X, p)=0$.

## Continuation and the Newton-Krylov method

The equation $\mathcal{G}(X, p)=0$ is completed with

$$
U^{\top}\left(X-X_{0}\right)+u_{p}\left(p-p_{0}\right)=0,
$$

( $X_{0}, p_{0}$ ) being the prediction of a new point on the curve of solutions, and ( $U, u_{p}$ ) an approximation of the tangent with $U=\left(u_{1}, \cdots, u_{m}\right)$. Setting $U=0$ and $u_{p}=1$ leads to parameter continuation. The system is solved by Newton's method and the linear systems by Krylov matrix-free methods (GMRES in our case), which only require computing the action by the Jacobian of the system

$$
\left(\begin{array}{cc}
D_{X} \mathcal{G} & D_{p} \mathcal{G} \\
U^{\top} & u_{p}
\end{array}\right)=\left(\begin{array}{cccccc}
I & \ldots & \ldots & 0 & -A_{m} & -b_{m} \\
-A_{1} & I & & \vdots & 0 & -b_{1} \\
0 & \ddots & \ddots & & \vdots & \vdots \\
\vdots & \ddots & \ddots & \ddots & \vdots & \vdots \\
0 & \ldots & 0 & -A_{m-1} & I & -b_{m-1} \\
u_{1}^{\top} & \ldots & u_{m-2}^{\top} & u_{m-1}^{\top} & u_{m}^{\top} & u_{p}
\end{array}\right)
$$

with $A_{i}=D_{x} P_{i}\left(x_{i}, p\right)$, and $b_{i}=D_{p} P_{i}\left(x_{i}, p\right)$.

The computation of

$$
\mathcal{G}(X, p)=\left(x_{1}-P_{m}\left(x_{m}, p\right), x_{2}-P_{1}\left(x_{1}, p\right), \cdots, x_{m}-P_{m-1}\left(x_{m-1}, p\right)\right)
$$

involves the time integration of $\dot{x}=f(x, p)$ with $m$ initial conditions $x(0)=x_{i}$.
The computation of the action of $D \mathcal{G}(X, p)$ on $(V, \mu)=\left(v_{1}, \cdots, v_{m}, \mu\right)$,

$$
\begin{aligned}
& D \mathcal{G}(X, p)(V, \mu)= \\
& \quad\left(v_{1}-D P_{m}\left(x_{m}, p\right)\left(v_{m}, \mu\right), v_{2}-D P_{1}\left(x_{1}, p\right)\left(v_{1}, \mu\right), \cdots, v_{m}-D P_{m-1}\left(x_{m-1}, p\right)\left(v_{m-1}, \mu\right)\right)
\end{aligned}
$$

is also trivially parallelizable, and involves the time integration of the system of first variational equations (for a vector)

$$
\dot{v}=D_{x} f(x, p) v+D_{p} f(x, p) \mu
$$

with initial conditions $v(0)=v_{i}$, together with the original system $\dot{x}=f(x, p)$, with initial conditions $x(0)=x_{i}$.

Only systems of at most $2 n$ equations must be integrated.

## Speedup and efficiency

Let $T(m, k)$ be the wall-clock time required to do a continuation with $m$ sections, with each partial shoot computed on a different processor, and using a preconditioner (to be defined later) based on an invariant subspace of dimension $m k$ ( $k=0$ means no preconditioning).

The (absolute) speedup is defined as

$$
S(m, k)=T_{\text {ref }} / T(m, k), \quad(\text { the goal is } \quad S(m, k)=m)
$$

$T_{\text {ref }}$ being the wall-clock time of the best algorithm without using parallelism. We have taken $T_{\text {ref }}=T(1,0)$.
The efficiency is

$$
E(m, k)=S(m, k) / m \quad \text { (the goal is } \quad E(m, k)=1) .
$$

## Bifurcation diagram (main branch of periodic orbits)

The full circles indicate bifurcation points, and the empty circles the point at which the continuation is started in the numerical experiments $(R a=2320)$, and that at which the preconditioner has been computed ( $R a=2305$ ).


Solid and dashed lines mean stable and unstable branches, respectively. The labels beside the branch of periodic orbits indicate the number of multipliers outside the unit circle.

## Results for the unpreconditioned computations


(a) Speedup, and (b) efficiency versus the number of sections $m$ corresponding to the computation of five points without preconditioning, and starting at $R a=2320$.

## Block cyclic matrices I

## Consider the matrices

$$
\mathcal{A}=\left(\begin{array}{ccccc}
0 & \cdots & \cdots & 0 & A_{m} \\
A_{1} & \ddots & & & 0 \\
0 & \ddots & \ddots & & \vdots \\
\vdots & \ddots & \ddots & \ddots & \vdots \\
0 & \cdots & 0 & A_{m-1} & 0
\end{array}\right)
$$

$A=A_{m} A_{m-1} \cdots A_{1}$, with $\mathcal{A} \in \mathbb{R}^{m n \times m n}$, and $A_{i} \in \mathbb{R}^{n \times n}$.

Proposition. If $A$ and $\mathcal{A}$ are defined as above then
a) if $(\mu, V)$ is an eigenpair of $\mathcal{A}$, with $\mu \in \mathbb{C}$, and $V=\left(v_{1}, \cdots, v_{m}\right)^{\top}$ with $v_{i} \in \mathbb{C}^{n}$ $(i=1, \cdots, m)$, then $\left(\mu^{m}, v_{1}\right)$ is an eigenpair of $A$.
b) if $\left(\lambda, u_{1}\right)$ is an eigenpair of $A$, with $\lambda \in \mathbb{C}$ and $u_{1} \in \mathbb{C}^{n}$, and if $\mu^{m}=\lambda$, then $(\mu, V)$, with $V=\left(\mu^{m-1} u_{1}, \mu^{m-2} u_{2}, \cdots, \mu u_{m-1}, u_{m}\right)^{\top}$, and $u_{i+1}=A_{i} u_{i}(i=1, \cdots, m-1)$, is an eigenpair of $\mathcal{A}$.

This result states that the eigenvalues of $\mathcal{A}$ are placed on circles centered at the origin, and then, those of $\mathcal{I}-\mathcal{A}$ on circles centered at +1 .


Leading eigenvalues of $\mathcal{I}-\mathcal{A}$ for $m=1, m=5$, and $m=10$ at $R a=2305$. The closed curve is the unit circle centered at +1 .

## Preconditioners based on leading invariant subspaces (deflation)

Let $C x=b$ be a linear system with non-singular matrix $C \in \mathbb{R}^{n \times n}$. We look for a right preconditioner, $M$ such that the convergence of the iterative methods applied to

$$
C M^{-1} y=b, \quad x=M^{-1} y,
$$

be faster. We assume that most of the spectrum of $C$ is clustered around +1 .
Let the columns of $Q=\left[q_{1}, \cdots, q_{k}\right] \in \mathbb{R}^{n \times k}$ form an orthonormal basis of an invariant subspace of $C$ corresponding to the first $k$ leading (maximal distance to +1 ) eigenvalues of $C$, with $Q^{\top} Q=I_{k}$, and $k \ll n$. The matrix $Q$ verifies

$$
C Q=Q R, \quad \text { and } \quad C^{-1} Q=Q R^{-1},
$$

with $R$ an invertible $k \times k$ matrix. Then, we define $M=Q R Q^{\top}+\left(I-Q Q^{\top}\right)$.

- $M^{-1}=Q R^{-1} Q^{\top}+\left(I-Q Q^{\top}\right)$
- $M^{\top}=Q R^{\top} Q^{\top}+\left(I-Q Q^{\top}\right)$, and $\left(M^{\top}\right)^{-1}=Q\left(R^{\top}\right)^{-1} Q^{\top}+\left(I-Q Q^{\top}\right)$
- $C M^{-1}=Q Q^{\top}+C\left(I-Q Q^{\top}\right)$
- If $z=z_{1}+z_{2}$ with $z_{1} \in \operatorname{Span}\left\{q_{1}, \cdots, q_{k}\right\}$, and $z_{2} \in \operatorname{Span}\left\{q_{1}, \cdots, q_{k}\right\}^{\perp}$,

$$
C M^{-1} z=z_{1}+C z_{2} .
$$

In the case of multiple shooting we need a preconditioner for $(\mathcal{I}-\mathcal{A}) X=B$.

## Block cyclic matrices II

Proposition. Let the matrix $Q_{1} \in \mathbb{R}^{n \times k}$, with $Q_{1}^{\top} Q_{1}=I_{k}$, be such that $A Q_{1}=Q_{1} R$ is a partial real Schur decomposition of $A=A_{m} A_{m-1} \cdots A_{1}$. Let $Q_{i+1} \in \mathbb{R}^{n \times k}$, and $R_{i} \in \mathbb{R}^{k \times k}$, $i=1, \cdots, m-1$, form a partial periodic real Schur decomposition of $A$, defined by

$$
\begin{aligned}
A_{i} Q_{i} & =Q_{i+1} R_{i}, \quad \text { and } \\
R_{m} & =Q_{1}^{\top} A_{m} Q_{m} .
\end{aligned}
$$

Then it follows that $R=R_{m} R_{m-1} \cdots R_{2} R_{1}$, and if

$$
\mathcal{Q}=\left(\begin{array}{ccccc}
Q_{1} & \ldots & \ldots & \ldots & 0 \\
\vdots & \ddots & & & \vdots \\
\vdots & & \ddots & & \vdots \\
\vdots & & & \ddots & \vdots \\
0 & \ldots & \ldots & \ldots & Q_{m}
\end{array}\right), \quad \text { and } \quad \mathcal{R}=\left(\begin{array}{ccccc}
0 & \ldots & \ldots & 0 & R_{m} \\
R_{1} & \ddots & & & 0 \\
0 & \ddots & \ddots & & \vdots \\
\vdots & \ddots & \ddots & \ddots & \vdots \\
0 & \ldots & 0 & R_{m-1} & 0
\end{array}\right)
$$

with $\mathcal{Q} \in \mathbb{R}^{m n \times m k}$ and $\mathcal{R} \in \mathbb{R}^{m k \times m k}$, then $\mathcal{A Q}=\mathcal{Q R}$ and $(\mathcal{I}-\mathcal{A}) \mathcal{Q}=\mathcal{Q}(\mathcal{I}-\mathcal{R})$, i.e., the columns of $\mathcal{Q}$ form an orthonormal basis of an invariant subspace, of dimension $m k$, of $\mathcal{A}$ and of $\mathcal{I}-\mathcal{A}$.

## Results for the preconditioned computations



Parameter continuation of an interval of 15 units in $R a$ without limitation of the step size. (a) Speedup, and (b) efficiency.

## Preconditioners for the multiple shooting

The matrices of the linear systems corresponding to pseudo-arclength and parameter continuation are

$$
\left(\begin{array}{cc}
D_{X} \mathcal{G} & D_{p} \mathcal{G} \\
U^{\top} & u_{p}
\end{array}\right), \quad \text { and } \quad\left(\begin{array}{cc}
D_{X} \mathcal{G} & 0 \\
0 & 1
\end{array}\right)
$$

with $D_{X} \mathcal{G}=\mathcal{I}-\mathcal{A}$. If $M$ is a preconditioner for $D_{X} \mathcal{G}$ the systems can be preconditioned with

$$
\left(\begin{array}{cc}
M & D_{p} \mathcal{G} \\
U^{\top} & u_{p}
\end{array}\right), \quad \text { or } \quad\left(\begin{array}{cc}
M & 0 \\
0 & 1
\end{array}\right) .
$$

The former is a bordered system, which can be solved in a stable way by the BEM algorithm. It requires solving systems with $M$ and $M^{\top}$, and to know explicitly the column $D_{p} \mathcal{G}=-D_{p} \mathcal{P}$. It can be obtained by integrating the first variational equations with initial conditions $v(0)=0$, and $\mu(0)=1$.

The continuation and preconditioning process

There are two possible straightforward ways of linking the continuation and the updating of the preconditioner processes.
a) $m$ processors are used to compute the periodic orbits and another one is updating the preconditioner.
a)

b) The $m$ processors are used to do both tasks.
b)


