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## CORRIGENDUM

The version of this paper that was published in Inverse Problems 1995 1211-33 was an earlier draft of the paper and did not include the many refinements that had been made. The correct version of this paper is published in full below.

# Approximate inverse for linear and some nonlinear problems 

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Received 14 March 1995, in final form 7 August 1995


#### Abstract

In this paper we present a method for solving problems such as $A f=g$ by constructing an approximate inverse which maps the data $g$ to a regularized solution of this equation of the first kind. No discretization for $f$ is needed. The solution operator can be precomputed independently of the data. This works for linear problems and for nonlinear problems with a special structure. The regularization is achieved by computing mollified versions of the (minimum-norm) solution. It is shown that this class of regularization operators contains, as special cases, the classical methods such as Tikhonov-Phillips, iteration methods and also discretization methods. In the case where the operator has some invariance properties the storage needs are dramatically reduced.


## 1. Introduction

We study operator equations $A f=g$ for operators between Hilbert spaces $X$ and $Y$. Both the cases of linear operators $A$ and of nonlinear $A$ with a special structure are treated. The approximate inverse means a solution operator which maps the data $g$ to a stable approximation of the solution of the ill-posed problem $A f=g$. This inversion operator is precomputed without using the data $g$.

The method is based on two ideas. First, the computation of moments of the solution is stable, i.e. instead of $f$ we compute the approximation $\left\langle f, e_{\gamma}\right\rangle$ with a suitable mollifier $e_{\gamma}$, thus reducing the high-frequency components in the solution which are mostly affected by the data noise. This can be reformulated as using a weaker topology in the space $X$, see [3,5]. Examples for $e_{\gamma}$ are given in section $2 ; e_{\gamma}$ can be a basis function for projection methods, it can be chosen such that $\left\langle f, e_{\gamma}\right\rangle$ approximates a derivative of $f$; in wavelet language it can be a scaling function or a wavelet. Second, in the case of linear operators the computation of $\left\langle e_{\gamma}, f\right\rangle$ is then achieved by approximating $e_{\gamma}$ in the range of the adjoint operator $A^{*}$ by the reconstruction kernel $v_{\gamma}: A^{*} v_{\gamma} \simeq e_{\gamma}$. Then

$$
\left\langle f, e_{\gamma}\right\rangle \simeq\left\langle f, A^{*} v_{\gamma}\right\rangle=\left\langle A f, v_{\gamma}\right\rangle=\left\langle g, v_{\gamma}\right\rangle
$$

This is the mollifier method presented in [7]. For nonlinear operators we combine these ideas with the results of Snieder [13] who generalized the Backus-Gilbert method to nonlinear problems. A detailed analysis of the method shows that the computational effort is much smaller than in [13].

In general, the regularization of ill-posed problems is achieved by regularization methods like Tikhonov-Phillips, truncated singular-value decomposition or iterative methods like Landweber or conjugate gradients. For references see, for example, [2, 4, 5, 10, 12]. These methods can easily be represented by the singular-value decomposition of the compact operator $A$ and a corresponding filter. On the other hand there are regularization methods like Backus-Gilbert or mollifier methods, see [1, 7, 12]. In section 2 we derive the approximate inverse and we show that in the case of linear operators we compute a mollified version of the minimum-norm solution. We show in section 3 that the abovementioned regularization methods are special cases of these mollifier methods. Also the relation to the Backus-Gilbert method is made precise. The next section presents the central difference quotient for approximating the first derivative as an approximate inverse where local averages of the solution are computed. Further on it shows an example where we determine the derivative of the solution of an integral equation of the first kind by properly choosing $e_{\gamma}$.

If the operator has some invariance properties the storage needs are reduced by applying group representations which intertwine with the operator. This is the content of section 5 and the base for an efficient implementation. Section 6 contains the nonlinear case where for the sake of simplicity we start with a quadratic problem. The approach is then generalized to a larger class of problems. The last section contains numerical experiments for a nonlinear problem.

## 2. Approximate inverse for linear problems

In the following we assume $A$ to be a linear, continuous operator between the Hilbert spaces $X$ and $Y$. Particularly, we think of $X$ as a space of functions and $Y$ as a finite-dimensional space of measurements. Hence, if necessary, we use $X=L_{2}(\Omega)$ for a suitable set $\Omega \subset \mathbb{R}^{d}$. Examples for mollifiers are

$$
e_{\gamma}(x, y)=\frac{d}{\operatorname{vol}\left(S^{d-1}\right) \gamma^{d}} \chi_{\gamma}(x-y)
$$

where $\chi_{\gamma}$ is the characteristic function of the ball around 0 with radius $\gamma$ and $\operatorname{vol}\left(S^{d-1}\right)$ is the measure of the surface of the unit ball in $\mathbb{R}^{d}$. Here local averages of the solution are computed. With the band-limiting filter

$$
e_{\gamma}(x, y)=\left(\frac{\gamma}{\pi}\right)^{d} \operatorname{sinc}(\gamma(x-y))
$$

the high-frequency components in the solution are eliminated. The kernel of the heat equation is fast decaying,

$$
e_{\gamma}(x, y)=(2 \pi)^{-d / 2} \gamma^{-d} \exp \left(-|x-y|^{2} /\left(2 \gamma^{2}\right)\right)
$$

In all cases the parameter $\gamma$ acts as a regularization parameter. The mollifier $e_{\gamma}$ is not necessarily a function with mean value 1 . When the essential information we need is discontinuities in $f$ we can use as $e_{\gamma}$ a function such that $\left\langle f, e_{\gamma}(x, \cdot)\right\rangle$ approximates a derivative of $f(x)$. This means that $e_{\gamma}$ can also be a wavelet, see e.g. [8].

For a motivation we consider the problem

$$
A f(x)=\int_{0}^{x} f(t) \mathrm{d} t=g(x) \quad 0 \leqslant x \leqslant 1
$$

where the inverse operator $A^{-1}$ is the differentiation of the first order. For a sufficiently smooth function $\psi$ and for homogeneous or periodic boundary conditions we compute, see [5],

$$
\langle\psi, f\rangle=\left\langle\psi, A^{-1} g\right\rangle=\left\langle\left(A^{-1}\right)^{*} \psi, g\right\rangle=\langle\Lambda, g\rangle
$$

where

$$
\Lambda=\left(A^{-1}\right)^{*} \psi=-\psi^{\prime}
$$

The solution depends continuously on the data $g$ as is seen by the simple calculation

$$
|\langle\psi, f\rangle|=\left|\left\langle\psi, A^{-1} g\right\rangle\right| \leqslant\|\Lambda\|\|g\| .
$$

The smoothness needed for the function $\psi$ depends on the smoothing of the operator $A$. In the example considered here, the function $g$ is once more differentiable than $f$ and the function $\psi$ has to be at least once differentiable.

After this motivation we start with the procedure described in the introduction. First we assume the equation $A^{*} v_{\gamma}=e_{\gamma}$ to be solvable. Then we put

$$
\begin{equation*}
\left\langle f, e_{\gamma}\right\rangle=\left\langle f, A^{*} v_{\gamma}\right\rangle=\left\langle A f, v_{\gamma}\right\rangle=\left\langle g, v_{\gamma}\right\rangle=: S_{\gamma} g \tag{1}
\end{equation*}
$$

This is the technique to derive inversion formulae in x-ray computer tomography resulting in the so-called filtered backprojection methods, see e.g. [9]. If the equation $A^{*} v_{\gamma}=e_{\gamma}$ is not solvable we approximate $v_{\gamma}$ by minimizing the defect $\left\|A^{*} v_{\gamma}-e_{\gamma}\right\|$ for sufficiently smooth $g$ leading to the equation

$$
\begin{equation*}
A A^{*} v_{\gamma}=A e_{\gamma} \tag{2}
\end{equation*}
$$

Then we get

$$
\left\langle f, e_{\gamma}\right\rangle \simeq\left\langle f, A^{*} v_{\gamma}\right\rangle=\left\langle A f, v_{\gamma}\right\rangle=\left\langle g, v_{\gamma}\right\rangle=: S_{\gamma} g
$$

It is important to mention that no artificial discretization of $f$ is needed as introduced by projection methods, see e.g. [5, 9, 10].
Definition 1. Let $e_{\gamma}$ be a suitable function, and let $v_{\gamma}$ be the solution of (2). Then $S_{\gamma} g:=\left\langle g, v_{\gamma}\right\rangle$ is called the approximate inverse of the operator $A$ and $v_{\gamma}$ is the reconstruction kernel.

In the situation where $X=L_{2}(\Omega), Y=\mathbb{R}^{N}$ this means

$$
S_{\gamma} g(x)=\sum_{n=1}^{N} g_{n} v_{\gamma, n}(x)
$$

The minimum-norm solution of the equation $A f=g$ is the solution with smallest norm. It is contained in the range of the adjoint operator $A^{*}$ and therefore computed as

$$
f_{M}=A^{*} u \quad \text { where } \quad A A^{*} u=g
$$

Hence

$$
\begin{equation*}
f_{M}=A^{*}\left(A A^{*}\right)^{-1} g \tag{3}
\end{equation*}
$$

The next theorem gives a simple relation between the approximate inverse and the minimumnorm solution.
Theorem 2. The approximate inverse maps the right-hand side of $A f=g$ to the mollified version of the minimum-norm solution $f_{M}$ defined in (3), i.e.

$$
\begin{equation*}
S_{\gamma} g=\left\langle f_{M}, e_{\gamma}\right\rangle \tag{4}
\end{equation*}
$$

Proof. The function $v_{\gamma}$ can be written following (2) as

$$
v_{\gamma}=\left(A A^{*}\right)^{-1} A e_{\gamma}
$$

resulting in

$$
S_{\gamma} g=\left\langle g, v_{\gamma}\right\rangle=\left\langle g,\left(A A^{*}\right)^{-1} A e_{\gamma}\right\rangle=\left\langle A^{*}\left(A A^{*}\right)^{-1} g, e_{\gamma}\right\rangle=\left\langle f_{M}, e_{\gamma}\right\rangle
$$

If $A: L_{2}(\Omega) \rightarrow \mathbb{R}^{N}$ is given as $(A f)_{n}=\left\langle f, k_{n}\right\rangle=\int_{\Omega} k_{n}(y) f(y) \mathrm{d} y$ for $n=1, \ldots N$, then the adjoint operator $A^{*}: \mathbb{R}^{N} \rightarrow L_{2}(\Omega)$ is given as

$$
A^{*} g(x)=\sum_{n=1}^{N} g_{n} k_{n}(x)
$$

The matrix representing $A A^{*}$ is the Gram matrix of the system $\left\{k_{n}: 1 \leqslant n \leqslant N\right\}$ and computed as

$$
\left(A A^{*}\right)_{m n}=\left\langle k_{n}, k_{m}\right\rangle=\int_{\Omega} k_{m}(y) k_{n}(y) \mathrm{d} y
$$

or

$$
A A^{*}=\int_{\Omega} k(y) k^{\top}(y) \mathrm{d} y
$$

where $k k^{\top}$ denotes the dyadic product $\left(k k^{\top}\right)_{i j}=k_{i} k_{j}$ for $k=\left(k_{1}, \ldots, k_{N}\right)^{\top} \in \mathbb{R}^{N}$.
The discrete equation inherits the instability of the integral operator, which means that the matrix is ill-conditioned. The $v_{\gamma}$ are then computed by using Tikhonov-Phillips regularization with a rather small regularization parameter. The fine regularization is achieved by the parameter $\gamma$ in $e_{\gamma}$, compare this with [11].

## 3. Comparison with other methods

In this section we show that most methods used for stabilizing the solution are special cases of the approximate inverse. To this end we first present some of these methods. We then compare our result with the Backus-Gilbert method.

The most prominent regularization method for ill-posed problems is Tikhonov-Phillips regularization where the defect is penalized with a priori information on the solution. For the sake of simplicity we use the $L_{2}$-norm of the solution and get

$$
J_{\gamma}(f)=\|A f-g\|^{2}+\gamma\|f\|^{2}
$$

The minimum of this functional is computed as solution of

$$
\begin{equation*}
\left(A^{*} A+\gamma I\right) f_{\gamma}=A^{*} g . \tag{5}
\end{equation*}
$$

Now let $A$ be a compact operator between the Hilbert spaces $X$ and $Y$. Then it has a singular-value decomposition

$$
\left\{v_{n}, u_{n} ; \sigma_{n}\right\}_{n}
$$

where $v_{n}, u_{n}$ are normalized and

$$
A v_{n}=\sigma_{n} u_{n} \quad \text { and } \quad A^{*} u_{n}=\sigma_{n} v_{n}
$$

For non-degenerate operators $A$, i.e. where the dimension of the range is not finite, the singular values $\sigma_{n}$ decay to 0 . The degree of the decay is a measure for the ill-posedness of the operator. With the singular-value decomposition the solution of (5) can be written as

$$
\begin{equation*}
T_{\gamma} g=\sum_{n} F_{\gamma}\left(\sigma_{n}\right) \sigma_{n}^{-1}\left\langle g, u_{n}\right\rangle v_{n} \tag{6}
\end{equation*}
$$

where the filter $F_{\gamma}$ is given as

$$
F_{\gamma}(\sigma)=\frac{\sigma^{2}}{\sigma^{2}+\gamma} .
$$

Of course, the truncated singular-value decomposition can also be written in the form (6) where

$$
F_{\gamma}(\sigma)= \begin{cases}1 & \sigma \geqslant \gamma \\ 0 & \sigma<\gamma\end{cases}
$$

For the Landweber iteration,

$$
f^{m+1}=f^{m}-\beta A^{*}\left(A f^{m}-g\right)
$$

the filter is, when identifying the regularization parameter $\gamma$ with $1 / m$ where $m$ is the iteration index,

$$
F_{m}(\sigma)=1-\left(1-\beta \sigma^{2}\right)^{m}
$$

The last example we want to mention is the method of conjugate gradients where the filter is given as

$$
F_{m}(\sigma)=P_{m-1}\left(\sigma^{2} ; g\right) \sigma^{2}
$$

and the polynomials $P_{m-1}$ describe the iterations; they depend on the data $g$, hence this is a nonlinear regularization method for a linear problem, and therefore this extremely well behaving method does not fit into the frame of the approximate inverse we study here. For more details see e.g. [5].

The next result shows that the regularization methods as discussed above with the exception of the conjugate gradient method are special cases of the mollifier methods. To describe the result we make use of function spaces $X$, the mollifier is then a function of two variables and we consider $S_{\gamma} g(x)=\left\langle f, e_{\gamma}(x, \cdot)\right\rangle$ which means that the scalar product is taken with respect to the second variable of $e_{\gamma}$.

Theorem 3. Let the regularization method $T_{\gamma}$ in (6) be given with a filter $F_{\gamma}$. Then this method can be written as an approximate inverse with mollifier

$$
\begin{equation*}
e_{\gamma}(x, y)=\sum_{n} F_{\gamma}\left(\sigma_{n}\right) v_{n}(x) v_{n}(y) \tag{7}
\end{equation*}
$$

Proof. The definition of $v_{\gamma}$ as solution of $A A^{*} v_{\gamma}=A e_{\gamma}$ in (2) leads with

$$
v_{\gamma}(x)=\sum_{n} \sigma_{n}^{-1}\left\langle e_{\gamma}(x, \cdot), u_{n}\right\rangle v_{n}
$$

to

$$
v_{\gamma}(x)=\sum_{n} F_{\gamma}\left(\sigma_{n}\right) \sigma_{n}^{-1} u_{n} v_{n}(x)
$$

Then

$$
S_{\gamma} g(x)=\left\langle g, v_{\gamma}(x)\right\rangle=T_{\gamma} g(x)
$$

Because of the versatility in the selection of the $e_{\gamma}$, which are functions of two variables, $e_{\gamma}(x, y)$, and which can differ for different $x$, it is not always possible to write the approximate inverse as a special regularization method. We can state the following result.

Theorem 4. Let $e_{\gamma}$ be decomposed into

$$
e_{\gamma}(x, y)=\sum_{m n} e_{\gamma, m n} v_{n}(x) v_{m}(y)
$$

Then the approximate inverse can be represented as a regularization method based on a filter $F_{\gamma}$ if and only if the filter satisfies

$$
e_{\gamma, m n}=F_{\gamma}\left(\sigma_{n}\right) \delta_{m n}
$$

Proof. If we expand $e_{\gamma}$ in terms of the $v_{n}$ as

$$
e_{\gamma}(x, y)=\sum_{m n} e_{\gamma, m n} v_{n}(x) v_{m}(y)
$$

then

$$
v_{\gamma}(x, \cdot)=\sum_{n} \sigma_{n}^{-1} u_{n} \sum_{m} e_{\gamma, m n} v_{m}(x)
$$

Hence

$$
\begin{aligned}
\left\langle v_{\gamma}(x, \cdot), g\right\rangle & =\sum_{n} \sigma_{n}^{-1}\left\langle u_{n}, g\right\rangle \sum_{m} e_{\gamma, m n} v_{m}(x) \\
& =\sum_{n} \sigma_{n}^{-1} F_{\gamma}\left(\sigma_{n}\right)\left\langle u_{n}, g\right\rangle v_{n}(x) \Longleftrightarrow \sum_{m} e_{\gamma, m n} v_{m}=F_{\gamma}\left(\sigma_{n}\right) v_{n}
\end{aligned}
$$

With the orthonormality of the $v_{n}$ it follows that

$$
e_{\gamma, m n}=F_{\gamma}\left(\sigma_{n}\right) \delta_{m n}
$$

This leads to the above statement and shows that this is, at least from a theoretical point of view, the much more general approach. The regularization property can be interpreted in the following sense. For many integral operators the singular functions belonging to small singular values; i.e. which produce the instability, are highly oscillating, see for example [ $5,6,14]$. The local averaging in this method thus stabilizes the problem.

In the following we compare the method of approximate inverse with the BackusGilbert method for the situation $A: L_{2}(\Omega) \rightarrow \mathbb{R}^{N}$. In the classical Backus-Gilbert method an approximation to the solution at point $x$ is computed as

$$
S g(x)=\langle g, v(x)\rangle=\sum_{n=1}^{N} g_{n} w_{n}(x)
$$

where $w_{n}(x)$ is determined by minimizing

$$
\int_{\Omega}|x-y|^{2}\left|\sum_{n=1}^{N} w_{n}(x) k_{n}(y)\right|^{2} \mathrm{~d} y=\int_{\Omega}|x-y|^{2}\left|A^{*} w(x)(y)\right|^{2} \mathrm{~d} y
$$

with the normalization

$$
\sum_{n=1}^{N} w_{n}(x) \int_{\Omega} k_{n}(y) \mathrm{d} y=\int_{\Omega} A^{*}(w(x))(y)=\left\langle A^{*} w(x), 1\right\rangle=1
$$

The first condition is interpreted as finding the $w(x)$ such that $A^{*} w(x)$ approximates the delta distribution at the point $x$, the second condition gives a normalization because otherwise the solution to the minimization would be $w=0$. The $w_{n}(x)$ are computed with Lagrangian multipliers; together with the multiplier $\mu$ they are the solution of

$$
\left(\begin{array}{cc}
C(x) & h \\
h^{\top} & 0
\end{array}\right)\binom{w(x)}{\mu}=\binom{0}{1}
$$

where the matrix elements of $C$ are

$$
C_{i j}(x)=\int_{\Omega}|x-y|^{2} k_{i}(y) k_{j}(y) \mathrm{d} y
$$

and

$$
h_{i}=\int_{\Omega} k_{i}(y) \mathrm{d} y .
$$

The matrix for determining the $w_{n}(x)$ depends here on the reconstruction point $x$, so it has to be decomposed for each $x$. Besides that, the amount of smoothing, the spread of the point spread function, depends on the kernel $|x-y|^{2}$. If this should be changed then the kernel changes and the calculations have to start from the beginning. So the main difference to the approximate inverse is the attempt to approximate the delta distribution instead of a mollifier and the extreme numerical effort.

## 4. Invariances reduce the storage needs

Due to the fact that in the applications only a finite number of data is available we have $Y=\mathbb{R}^{N}$ or $Y=\mathbb{C}^{N}$. Then

$$
S_{\gamma} g(z)=\sum_{n=1}^{N} g_{n} v_{\gamma, n}(z)
$$

which means that we have to compute the reconstruction kernel $v_{\gamma}(x) \in \mathbb{R}^{N}$ for all reconstruction points $x$. These values are precomputed independently of the data. So computing time is not a big issue, but the storage needs can be prohibitively large.

In the following we use invariance properties of the operator $A$ to reduce dramatically the storage needs for the method. To present the ideas we start with the following example.
Example 5. Let $A$ be a convolution operator

$$
A f(x)=\int_{\mathbb{R}^{d}} k(x-y) f(y) \mathrm{d} y
$$

and $X=Y=L_{2}\left(\mathbb{R}^{d}\right)$. Let $T^{z}$ be the translation defined as

$$
T^{z} f(x)=f(x-z)
$$

Then

$$
T^{z} A=A T^{z}
$$

and

$$
T^{z} A A^{*}=A A^{*} T^{z}
$$

Proof. We easily see that

$$
\begin{aligned}
A T^{z} f(x) & =\int k(x-y) f(y-z) \mathrm{d} y \\
& =\int k(x-z-u) f(u) \mathrm{d} u \\
& =A f(x-z) \\
& =T^{z} A f(x)
\end{aligned}
$$

The operator $A A^{*}$ is computed as

$$
A A^{*} g(x)=\int_{\mathbb{R}^{d}} k^{(2)}(x, y) g(y) \mathrm{d} y
$$

where

$$
k^{(2)}(x, y)=\int_{\mathbb{R}^{d}} k(x-u) k(y-u) \mathrm{d} u
$$

Then

$$
\begin{aligned}
A A^{*} T^{z} g(x) & =\int k^{(2)}(x, y) g(y-z) \mathrm{d} y \\
& =\int k^{(2)}(x, u+z) g(u) \mathrm{d} u \\
& =\int k^{(2)}(x-z, u) g(u) \mathrm{d} u \\
& =T^{z} A A^{*} g(x)
\end{aligned}
$$

In the following let $E_{\gamma}$ be a function of one variable, we can think of $E_{\gamma}(y)=e_{\gamma}(0, y)$. Instead of a translation we consider arbitrary group representations.

Theorem 6. Let $A: X \rightarrow Y$ and let $T_{1}^{x}$ be a group representation on $X$ and $T_{2}^{x}, T_{3}^{x}$ be group representations on $Y$ such that

$$
\begin{equation*}
A T_{1}^{x}=T_{2}^{x} A \tag{8}
\end{equation*}
$$

and

$$
\begin{equation*}
T_{2}^{x} A A^{*}=A A^{*} T_{3}^{x} \tag{9}
\end{equation*}
$$

Let $w_{\gamma}$ be the minimum-norm solution of

$$
\begin{equation*}
A A^{*} w_{\gamma}=A E_{\gamma} \tag{10}
\end{equation*}
$$

Then the solution of

$$
A A^{*} v_{\gamma}(x)=A T_{1}^{x} E_{\gamma}
$$

is

$$
\begin{equation*}
v_{\gamma}(x)=T_{3}^{x} w_{\gamma} . \tag{11}
\end{equation*}
$$

Proof. From the invariance properties follows:

$$
\begin{aligned}
A T_{1}^{x} E_{\gamma} & =T_{2}^{x} A E_{\gamma} \\
& =A A^{*} T_{3}^{x} w_{\gamma}
\end{aligned}
$$

which completes the proof.
This means that only the solution $w_{\gamma}$ has to be computed and stored, the kernels for other reconstruction points $x$ are found by the action of $T_{1}^{x}$ on $E_{\gamma}$ and by $T_{3}^{x}$ on $w_{\gamma}$.

For the above example of the convolution equation this can be applied in the following way. Let $T=T_{1}=T_{2}=T_{3}$ be the translation. The approximate inverse is then

$$
\begin{aligned}
S_{\gamma} g(x) & =\left\langle T^{x} w_{\gamma}, g\right\rangle \\
& =\int_{\mathbb{R}^{d}} w_{\gamma}(y-x) g(y) \mathrm{d} y
\end{aligned}
$$

which again is of convolution type. Hence instead of a function of two variables $v_{\gamma}(x, y)$ only the function $w_{\gamma}$ of one variable is needed.

In the case of a finite number of data where $(A f)_{n}=A f\left(x_{n}\right), n=1, \ldots, N$ for suitable points $x_{n}$ the reconstruction kernel $w_{\gamma}$ is a vector in $\mathbb{R}^{N}$ with $\left(w_{\gamma}\right)_{n}=w_{\gamma}\left(x_{n}\right)$. Then $w_{\gamma}\left(x_{n}-x\right)$ can be evaluated by linear interpolation between $\left(w_{\gamma}\right)_{m}$ and $\left(w_{\gamma}\right)_{m+1}$ with $x_{m} \leqslant x_{n}-x<x_{m+1}$.

For $Y=\mathbb{R}^{N}$ and $M$ reconstruction points the storage needs originally was $M \times N$ real numbers, now it is only $N$ real numbers!

## 5. Numerical examples

We study two examples in this section. First we consider $X=Y=L_{2}(0,1)$ and again the equation of the first kind

$$
A f(x)=\int_{0}^{x} f(t) \mathrm{d} t=g(x) \quad 0 \leqslant x \leqslant 1
$$

with the solution $f=g^{\prime}$. The adjoint operator is

$$
A^{*} g(y)=\int_{y}^{1} g(t) \mathrm{d} t
$$

and therefore the solution of

$$
A^{*} v_{\gamma}(x, \cdot)=e_{\gamma}(x, \cdot)
$$

is

$$
v_{\gamma}(x, y)=-\frac{\partial}{\partial y} e_{\gamma}(x, y)
$$

With $e_{\gamma}(x, y)=\chi_{[-\gamma, \gamma]}(x-y) /(2 \gamma)$ we get, using distributional derivatives,

$$
v_{\gamma}(x, y)=\frac{1}{2 \gamma}\left(\delta_{x+\gamma}-\delta_{x-\gamma}\right)
$$

and hence the approximate inverse is the central difference quotient

$$
S_{\gamma} g(x)=\frac{1}{2 \gamma}(g(x+\gamma)-g(x-\gamma))
$$

which means that the central difference quotient with stepsize $\gamma$ results in a local averaging of the searched-for solution with the characteristic function of the interval $[-\gamma, \gamma]$.

The band-limiting filter $e_{\gamma}(x, y)=\gamma / \pi \operatorname{sinc}(\gamma(x-y))$ results in the band-limiting differentiation formula with reconstruction kernel

$$
v_{\gamma}(x, y)=\frac{\gamma(x-y) \cos \gamma(x-y)-\sin \gamma(x-y)}{\pi(x-y)^{2}}
$$

Both the mollifiers $e_{\gamma}$ and the reconstruction kernels $v_{\gamma}$ are of convolution type which minimizes both the storage and the computational effort, as discussed in the last section. For the numerical experiment we consider $A: L_{2}(\Omega) \rightarrow \mathbb{R}^{N}$ and the integral equation

$$
\begin{equation*}
(A f)_{n}=\int_{0}^{1} k\left(x_{n}, y\right) f(y) \mathrm{d} y=g\left(x_{n}\right) \quad 1 \leqslant n \leqslant N \tag{12}
\end{equation*}
$$

for equally distributed points $x_{n}=(n-1 / 2) / N$ where the kernel is given as

$$
k(x, y)= \begin{cases}x(1-y) & x \leqslant y \\ y(1-x) & y \leqslant x\end{cases}
$$

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Figure 1. Solution of (12) with exact data.


Figure 2. Solution of (12) with noisy data, 5\% noise.


Figure 3. Derivative of the solution of (12) with exact data.

Reconstructions with $e_{\gamma}(x, y)=\gamma \operatorname{sinc}(\gamma(x-y)) / \pi$ are shown in [7]. Here we want to apply the approximate inverse with $e_{\gamma}(x, y)=(2 \pi)^{1 / 2} \gamma^{-1} \exp \left((x-y)^{2} /\left(2 \gamma^{2}\right)\right)$ to determine both the solution and the derivative of the solution. For computing the derivative we start from

$$
f_{\gamma}^{\prime}(x)=\left\langle f^{\prime}, e_{\gamma}(x, \cdot)\right\rangle=-\left\langle f, e_{\gamma}^{\prime}(x, \cdot)\right\rangle+\left.f e_{\gamma}(x, \cdot)\right|_{a} ^{b}
$$

which means that we determine $v_{\gamma}$ such that $\left\|A^{*} v_{\gamma}(x)+e_{\gamma}^{\prime}(x, \cdot)\right\|$ is minimized. When $f(a)$ and $f(b)$ are known we can use this formula directly, otherwise we have to approximate these two values. The critical task of computing derivatives of the searched-for function is achieved here by differentiating the analytically given function $e_{\gamma}$ instead of differentiating the numerically determined function $f$. The price to pay is that $A e_{\gamma}^{\prime}$ has to be evaluated with high accuracy.

The right-hand side in the numerical example was chosen such that the exact solution is piecewise linear as shown in the following figures. Data errors occur also by the evaluation of $A e_{\gamma}$. Figure 2 shows the result with $5 \%$ data error.


Figure 4. Derivative of the solution of (12) with noisy data, $3 \%$ noise.

Figures 3 and 4 show the computation of the derivative of the above function. Note that this function is not differentiable in the classical sense. The discontinuities and the behaviour of the derivative are clearly visible.

## 6. Nonlinear problems

We start the presentation of the approach for the nonlinear problem

$$
A: L_{2}(\Omega) \rightarrow \mathbb{R}^{N}
$$

by considering, for the sake of simplicity, only a quadratic problem. Let the operator $A$ be given as

$$
A f=A_{1} f+A_{2} f
$$

where $A_{1}$ is linear,

$$
\left(A_{1} f\right)_{n}=\int_{\Omega} k_{n}^{1}(y) f(y) \mathrm{d} y
$$

and $A_{2}$ is a quadratic operator defined as

$$
\left(A_{2} f\right)_{n}=\int_{\Omega} \int_{\Omega} k_{n}^{2}\left(y_{1}, y_{2}\right) f\left(y_{1}\right) f\left(y_{2}\right) \mathrm{d} y_{1} \mathrm{~d} y_{2}
$$

With $k^{1}$ we denote the vector of the $N$ components of $k_{n}^{1}$ and similarly with $k^{2}$ the vector of the $N$ components of $k_{n}^{2}$. For the approximate inverse we follow Snieder [13] and make the following ansatz:

$$
f_{\gamma}(x)=\left\langle g, v_{\gamma}(x)\right\rangle+\left\langle g, V_{\gamma}(x) g\right\rangle
$$

where $V_{\gamma}(x)$ is an $N \times N$ matrix. We replace $g$ by $A f$ and get

$$
\begin{equation*}
f_{\gamma}(x) \simeq\left\langle A_{1} f, v_{\gamma}(x)\right\rangle+\left\langle A_{2} f, v_{\gamma}(x)\right\rangle+\left\langle A_{1} f, V_{\gamma}(x) A_{1} f\right\rangle \tag{13}
\end{equation*}
$$

where we omitted the higher order terms.
In the following we always use $x \in \Omega$ as the reconstruction point which is arbitrarily fixed; and with $y_{1}, y_{2}$ and so on we denote the integration variables in the integral operators.

Of course, we can attempt to approximate with the right-hand side an expression like $\left\langle f, e_{\gamma}(x, \cdot)\right\rangle$. But then the approximate inverse is not independent of the data, and this does not lead to an approximate inverse as aimed for. We therefore follow [13] and consider the terms separately. We approximate the mollified solution with the linear term as well as possible which means

$$
\left\langle A_{1} f, v_{\gamma}(x)\right\rangle \simeq\left\langle f, e_{\gamma}(x, \cdot)\right\rangle
$$

hence

$$
A_{1} A_{1}^{*} v_{\gamma}=A_{1} e_{\gamma}
$$

The computation of $v_{\gamma}$ thus follows exactly the same lines as in the linear case. When this term presents the solution it only remains to make the rest as small as possible. Denoting the remainder by

$$
\begin{equation*}
R(x)\left(y_{1}, y_{2}\right)=v_{\gamma}(x)^{\top} k^{2}\left(y_{1}, y_{2}\right)+k^{1}\left(y_{1}\right)^{\top} V_{\gamma}(x) k^{1}\left(y_{2}\right) \tag{14}
\end{equation*}
$$

where $a^{\top} b$ is the scalar product in the image space $\mathbb{R}^{N}$, we can write the rest in (13) as

$$
\int_{\Omega} \int_{\Omega} R(x)\left(y_{1}, y_{2}\right) f\left(y_{1}\right) f\left(y_{2}\right) \mathrm{d} y_{1} \mathrm{~d} y_{2} .
$$

In order to make this term as small as possible we minimize with respect to $V_{\gamma}(x)$ for each reconstruction point $x$ the norm

$$
\|R(x)(\cdot, \cdot)\|_{L_{2}(\Omega \times \Omega)}^{2} \rightarrow \min
$$

We get

$$
\begin{aligned}
\|R(x)\|^{2}= & \int_{\Omega \times \Omega}\left(v_{\gamma}^{\top}(x) k^{2}\left(y_{1}, y_{2}\right)+k^{1}\left(y_{1}\right)^{\top} V_{\gamma}(x) k^{1}\left(y_{2}\right)\right)^{2} \mathrm{~d} y_{1} \mathrm{~d} y_{2} \\
= & \int_{\Omega \times \Omega}\left(\left(v_{\gamma}(x)^{\top} k^{2}\left(y_{1}, y_{2}\right)\right)^{2}+2 v_{\gamma}(x)^{\top} k^{2}\left(y_{1}, y_{2}\right) k^{1}\left(y_{1}\right)^{\top} V_{\gamma}(x) k^{1}\left(y_{2}\right)\right. \\
& \left.+\left(k^{1}\left(y_{1}\right)^{\top} V_{\gamma}(x) k^{1}\left(y_{2}\right)\right)^{2}\right) \mathrm{d} y_{1} \mathrm{~d} y_{2}
\end{aligned}
$$

Differentiating this expression with respect to the fixed matrix element $\left(V_{\gamma}\right)_{i j}(x)$ and equating the derivative to zero we find

$$
\begin{aligned}
& \int_{\Omega} k_{i}^{1}\left(y_{1}\right) k^{1}\left(y_{1}\right)^{\top} \mathrm{d} y_{1} V_{\gamma}(x) \int_{\Omega} k_{j}^{1}\left(y_{2}\right) k^{1}\left(y_{2}\right)^{\top} \mathrm{d} y_{2} \\
&=-\int_{\Omega \times \Omega} v_{\gamma}^{\top}(x) k^{2}\left(y_{1}, y_{2}\right) k_{i}^{1}\left(y_{1}\right) k_{j}^{1}\left(y_{2}\right) \mathrm{d} y_{1} \mathrm{~d} y_{2} \quad i, j=1, \ldots, N .
\end{aligned}
$$

In matrix notation this is simply

$$
\begin{equation*}
\left(A_{1} A_{1}^{*}\right) V_{\gamma}(x)\left(A_{1} A_{1}^{*}\right)=-\sum_{n=1}^{N} v_{\gamma, n}(x) B_{n} \tag{15}
\end{equation*}
$$

where the matrices $B_{n}$ on the right-hand side are

$$
B_{n}=\int_{\Omega \times \Omega} k^{1}\left(y_{1}\right) k_{n}^{2}\left(y_{1}, y_{2}\right) k^{1}\left(y_{2}\right)^{\top} \mathrm{d} y_{1} \mathrm{~d} y_{2} .
$$

This means that with $C_{n}=\left(A_{1} A_{1}^{*}\right)^{-1} B_{n}\left(A_{1} A_{1}^{*}\right)^{-1}$ the function $k^{1}\left(y_{1}\right)^{\top} C_{n} k^{1}\left(y_{2}\right)$ is the orthogonal projection of the function $k_{n}^{2}\left(y_{1}, y_{2}\right)$ on the linear space $\operatorname{span}\left\{k_{i}^{1}\left(y_{1}\right) k_{j}^{1}\left(y_{2}\right): i, j=\right.$ $1, \ldots, N\}$. Hence, if the $k_{n}^{1}$ are linearly independent, the matrix $V_{\gamma}$ is simply

$$
V_{\gamma}(x)=-\sum_{n=1}^{N} v_{\gamma, n}(x) C_{n}
$$

and we get

$$
f_{\gamma}(x)=\left\langle g, v_{\gamma}(x)\right\rangle-\sum_{n=1}^{N} v_{\gamma, n}(x)\left\langle g, C_{n} g\right\rangle .
$$

For the approximation $f_{\gamma}$ this means that

$$
\begin{aligned}
f_{\gamma}(x) & =\left\langle v_{\gamma}(x), g\right\rangle+\left\langle g, V_{\gamma} g\right\rangle \\
& =\left\langle v_{\gamma}, A_{1} f\right\rangle+\sum_{n=1}^{N} v_{\gamma, n}(x)\left(\left(A_{2} f\right)_{n}-g^{\top} C_{n} g\right) .
\end{aligned}
$$

If the kernels $k_{n}^{2}$ are already in the span of the product of the $k_{m}^{1}$ which means that $k_{n}^{2}\left(y_{1}, y_{2}\right)=k^{1}\left(y_{1}\right)^{\top} C_{n} k^{1}\left(y_{2}\right)$ and $\left(A_{2} f\right)_{n}=\left(A_{1} f\right)^{\top} C_{n} A_{1} f$, then the contribution of the square term comes only from

$$
\left\langle A_{1} f, V_{\gamma}(x) A_{2} f\right\rangle+\left\langle A_{2} f, V_{\gamma}(x) A_{1} f\right\rangle+\left\langle A_{2} f, V_{\gamma}(x) A_{2} f\right\rangle
$$

(see also the discussion in [13]).
The generalization of this approach to operators of arbitrary order in the following sense:

$$
\begin{equation*}
A f=\sum_{\ell=1}^{\infty} A_{\ell} f \tag{16}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(A_{\ell} f\right)_{n}=\int_{\Omega^{\ell}} k_{n}^{\ell}\left(y_{1}, \ldots, y_{\ell}\right) f\left(y_{1}\right) \ldots f\left(y_{\ell}\right) \mathrm{d} y_{1} \ldots \mathrm{~d} y_{\ell} \tag{17}
\end{equation*}
$$

is now straightforward, but notationally difficult. Examples for operators of this type are given in [13]. With $b^{\ell}\left(g_{1}, \ldots, g_{\ell}\right)$ we denote a multilinear form: i.e. $b^{\ell}\left(g_{1}, \ldots, g_{\ell}\right)$ is linear in each argument. This form is generated by the tensor $v^{\ell}$. For $\ell=1$ we have

$$
b^{1}(g)=\left\langle v_{\gamma}, g\right\rangle \quad v^{1}=v_{\gamma}
$$

for $\ell=2$

$$
b^{2}\left(g_{1}, g_{2}\right)=\left\langle g_{1}, V_{\gamma} g_{2}\right\rangle \quad v^{2}=V_{\gamma}
$$

If we always use the same argument we simply write $b^{\ell}(g)$. To be precise we have to indicate the dependence of $b^{\ell}$ on the reconstruction point $x$, hence $b^{\ell}(x ; g)$. The ansatz for the approximate inverse is now

$$
\begin{equation*}
f_{\gamma}^{L}(x)=\sum_{\ell=1}^{L} b^{\ell}(x ; g) \tag{18}
\end{equation*}
$$

where we determine the tensors generating the multilinear form $b^{\ell}$ as described above for $L=2$. We replace $g$ by $A f$ and collect the terms with same number of factors $f$ :

$$
\begin{aligned}
f_{\gamma}(x)=b^{1}(x ; & \left.A_{1} f\right)+b^{1}\left(x ; A_{2} f\right)+b^{2}\left(x ; A_{1} f, A_{1} f\right)+b^{1}\left(x ; A_{3} f\right)+b^{2}\left(x ; A_{2} f, A_{1} f\right) \\
& +b^{2}\left(x ; A_{1} f, A_{2} f\right)+b^{3}\left(x ; A_{1} f, A_{1} f, A_{1} f\right)+\cdots
\end{aligned}
$$

where $b^{1}$ and $b^{2}$ are determined in the way just described and we continue recursively. The remainder corresponding to (14) is denoted by $R_{\ell}$ and we minimize

$$
\left\|R_{\ell}\left(x ; y_{1}, \ldots, y_{\ell}\right)\right\|_{L_{2}\left(\Omega^{\ell}\right)}
$$

We denote the components of the tensor $v^{\ell}$ by $v_{\alpha}^{\ell}$ with the multi-index $\alpha \in \mathbb{N}^{\ell}$. Then

$$
R_{\ell}\left(x ; y_{1}, \ldots, y_{\ell}\right)=\sum_{\alpha_{1}=1}^{N} \cdots \sum_{\alpha_{\ell}=1}^{N} v_{\alpha}^{\ell} k_{\alpha_{1}}^{1}\left(y_{1}\right) \cdots k_{\alpha_{\ell}}^{1}\left(y_{\ell}\right)+r_{\ell}\left(x ; y_{1}, \ldots, y_{\ell}\right)
$$

where $r_{\ell}$ depends on terms determined before. The norm of the rest is

$$
\begin{aligned}
\left\|R_{\ell}(x)\right\|_{L_{2}\left(\Omega^{\ell}\right)}^{2}= & \int_{\Omega^{\ell}}\left\{\left(\sum_{\alpha_{1}=1}^{N} \cdots \sum_{\alpha_{\ell}=1}^{N} v_{\alpha}^{\ell} k_{\alpha_{1}}^{1}\left(y_{1}\right) \cdots k_{\alpha_{\ell}}^{1}\left(y_{\ell}\right)\right)^{2}\right. \\
& +2 \sum_{\alpha_{1}=1}^{N} \cdots \sum_{\alpha_{\ell}=1}^{N} v_{\alpha}^{\ell} k_{\alpha_{1}}^{1}\left(y_{1}\right) \cdots k_{\alpha_{\ell}}^{1}\left(y_{\ell}\right) \\
& \left.\times r_{\ell}\left(x ; y_{1}, \ldots, y_{\ell}\right)\right\} \mathrm{d} y_{1} \ldots \mathrm{~d} y_{\ell}+\left\|r_{\ell}\right\|^{2} .
\end{aligned}
$$

Differentiating with respect to a fixed $v_{\alpha}^{\ell}$ leads to the following system of equations for the $v_{\beta}^{\ell}$ :

$$
\begin{aligned}
\sum_{\beta_{1}=1}^{N} \cdots \sum_{\beta_{\ell}=1}^{N} & v_{\beta}^{\ell} \int_{\Omega} k_{\alpha_{1}}^{1}\left(y_{1}\right) k_{\beta_{1}}^{1}\left(y_{1}\right) \mathrm{d} y_{1} \cdots \int_{\Omega} k_{\alpha_{\ell}}^{1}\left(y_{\ell}\right) k_{\beta_{\ell}}^{1}\left(y_{\ell}\right) \mathrm{d} y_{\ell} \\
& =-\int_{\Omega^{\ell}} k_{\alpha_{1}}^{1}\left(y_{1}\right) \cdots k_{\alpha_{\ell}}^{1}\left(y_{\ell}\right) r_{\ell}\left(x ; y_{1}, \ldots, y_{\ell}\right) \mathrm{d} y_{1} \cdots \mathrm{~d} y_{\ell}
\end{aligned}
$$

Note that $\int_{\Omega} k_{\alpha_{m}}^{1}\left(y_{m}\right) k_{\beta_{m}}^{1}\left(y_{m}\right) \mathrm{d} y_{m}=\left(A_{1} A_{1}^{*}\right)_{\alpha_{m} \beta_{m}}$. We can collect the results in the following theorem.

Theorem 7. The matrix for determining the tensor $v^{\ell}$ does not depend on $x$, it consists of factors of $A_{1} A_{1}^{*}$. The tensor has the form

$$
\begin{equation*}
v^{\ell}(x)=\sum_{n=1}^{N} v_{n}^{1}(x) B_{n}^{\ell} \tag{19}
\end{equation*}
$$

where $B_{n}^{\ell}$ is a tensor of order $\ell$ and it is independent of $x$.
Proof. The first assertion is obvious from the above derived system of equations, the second statement follows by induction. For $\ell=2$ we have

$$
r_{2}\left(x ; y_{1}, y_{2}\right)=v_{\gamma}(x)^{\top} k^{2}\left(y_{1}, y_{2}\right)
$$

In the step from $\ell$ to $\ell+1$ the sum starts with $v_{1}^{\top}(x) k^{\ell+1}\left(y_{1}, \ldots, y_{\ell+1}\right)$ plus terms which are contained already in $r_{\ell}$.

## 7. Numerical example

We present numerical results for a problem discussed in [13] where the remainder $r^{2}$ was displayed but where numerical reconstructions are not given.

The differential equation for the vibrating string is

$$
u^{\prime \prime}(x)+\frac{\rho(x) \omega^{2}}{T} u(x)=0 \quad u(0)=u(L)=0
$$

We assume the function $\rho$ to be

$$
\rho(x)=\rho_{0}(1+f(x)) .
$$

For $f \equiv 0$ the eigenfunctions are

$$
u_{n}^{0}(x)=\left(\frac{2}{L}\right)^{1 / 2} \sin \frac{n \pi x}{L}
$$

and the eigenvalues

$$
\omega_{n}^{0}=\left(\frac{T}{\rho_{0}}\right)^{1 / 2} \frac{n \pi}{L}
$$

The data $g_{n}$ in our problem are the relative frequency shifts given as

$$
g_{n}=\frac{\omega_{n}^{2}-\left(\omega_{n}^{0}\right)^{2}}{\left(\omega_{n}^{0}\right)^{2}} \quad n=1, \ldots, N
$$

which means that we study $A: L_{2}(\Omega) \rightarrow \mathbb{R}^{N}$. For the integral operator we get the kernels

$$
k_{n}^{1}(y)=-\frac{2}{L} \sin ^{2}\left(\frac{n \pi y}{L}\right)
$$

and

$$
\begin{aligned}
k_{n}^{2}\left(y_{1}, y_{2}\right)= & \frac{4}{L^{2}} \sin ^{2} \frac{n \pi y_{1}}{L} \sin ^{2} \frac{n \pi y_{2}}{L} \\
& +\sum_{m \neq n} \frac{4}{L^{2}} \frac{n^{2}}{n^{2}-m^{2}} \sin \frac{n \pi y_{1}}{L} \sin \frac{m \pi y_{1}}{L} \sin \frac{n \pi y_{2}}{L} \sin \frac{m \pi y_{2}}{L} .
\end{aligned}
$$

The matrix $A_{1} A_{1}^{*}=\frac{1}{2} I+e e^{\top}$, where $e=(1, \ldots, 1)^{\top}$. For the reconstruction we used the piecewise linear mollifier

$$
e_{\gamma}(x-y)=\frac{1}{\gamma} \begin{cases}1-\frac{y-x}{\gamma} & x \leqslant y \leqslant x+\gamma \\ 1+\frac{y-x}{\gamma} & x-\gamma \leqslant y \leqslant x \\ 0 & \text { otherwise }\end{cases}
$$

Here $A_{1} e_{\gamma}(x, \cdot)$ can be computed exactly.
The reconstructions are presented here for $N=15$ without noise (figure 5) and with $3 \%$ noise in figure 6 . Exact data means that the first 15 eigenvalues are computed with high precision using hierarchical finite elements. The function to be reconstructed is piecewise linear, the broken curves ( --- ) are the results from the linear part of the problem and the straight lines are the reconstructions including also the quadratic part. These examples show that the additional effort in computing the quadratic approximation improves the result considerably.


Figure 5. Solution of inverse eigenvalue problem with exact data, linear and quadratic approximations of the exact solution.


Figure 6. Solution of inverse eigenvalue problem with noisy data (3\%), linear and quadratic approximations of the exact solution.

## Acknowledgments

The work of the author was supported by the Deutsche Forschungsgemeinschaft in the Sonderforschungsbereich SFB 273 and by the Bundesministerium für Bildung, Wissenschaft, Forschung und Technologie (BMBF), grant no 20M360-03LO7SAA-8.

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