# THE APPROXIMATE INVERSE IN ACTION WITH AN APPLICATION TO COMPUTERIZED TOMOGRAPHY* 

ANDREAS RIEDER ${ }^{\dagger}$ AND THOMAS SCHUSTER ${ }^{\ddagger}$


#### Abstract

The approximate inverse is a scheme to obtain stable numerical inversion formulæ for linear operator equations of the first kind. Yet, in some applications the computation of a crucial ingredient, the reconstruction kernel, is time-consuming and instable. It may even happen that the kernel does not exist for a particular semidiscrete system. To cure this dilemma we propose and analyze a technique that is based on a singular value decomposition of the underlying operator. The results are applied to the reconstruction problem in 2D-computerized tomography where they enable the design of reconstruction filters and lead to a novel error analysis of the filtered backprojection algorithm.


Key words. approximate inverse, mollification, Radon transform, computerized tomography, filtered backprojection

AMS subject classifications. 65J10, 65R10

PII. S0036142998347619

1. Introduction. The approximate inverse is a regularization scheme which applies especially to underdetermined (semidiscrete) systems. Yet in some applications the numerical computation of the necessary reconstruction kernel $v_{\text {discrete }}$ is timeconsuming and instable. It may even happen that $v_{\text {discrete }}$ does not exist for a particular semidiscrete system. However, the reconstruction kernel $v$ of the underlying infinite dimensional (continuous) problem may be at hand. In this paper we propose a procedure to find a substitute for $v_{\text {discrete }}$ from $v$ and we show that this procedure is sound.

In the following we recall the concept of the approximate inverse which belongs to the class of mollifier methods as considered, for instance, by Murio [19]. In a systematic way the approximate inverse generalizes a technique used by Grünbaum [5] and Davison and Grünbaum [3] for tomographic inversion.

Let $A: X \rightarrow Y$ be a continuous and injective operator between the real or complex infinite dimensional Hilbert spaces $X$ and $Y$. We want to find a $f \in X$ such that

$$
\begin{equation*}
A_{n} f=g_{n}, \tag{1.1}
\end{equation*}
$$

where $A_{n}: X \rightarrow \mathbb{C}^{n}$ and $g_{n} \in \mathbb{C}^{n}$ are defined via a mapping $\Psi_{n}: Y \rightarrow \mathbb{C}^{n}$ by $A_{n}=\Psi_{n} A$ and $g_{n}=\Psi_{n} g$ with $g \in \mathrm{R}(A)$, the range of $A$. Let us assume-for the time being - that $A_{n}$ is continuous. The above setting describes most practical situations where the data can be recorded only in finitely many observation points.

Problem (1.1) is underdetermined and we can only search for its minimum norm solution $f_{n}^{\dagger}$, that is,

$$
\begin{equation*}
A_{n}^{*} A_{n} f_{n}^{\dagger}=A_{n}^{*} g_{n} \quad \text { and } f_{n}^{\dagger} \in \mathrm{N}\left(A_{n}\right)^{\perp} \tag{1.2}
\end{equation*}
$$

*Received by the editors November 12, 1998; accepted for publication (in revised form) October 14, 1999; published electronically May 23, 2000.
http://www.siam.org/journals/sinum/37-6/34761.html
${ }^{\dagger}$ Institut für Wissenschaftliches Rechnen und Mathematische Modellbildung (IWRMM), Universität Karlsruhe, 76128 Karlsruhe, Germany (andreas.rieder@math.uni-karlsruhe.de).
$\ddagger$ Fachbereich Mathematik, Geb. 36, Universität des Saarlandes, 66041 Saarbrücken, Germany (thomas@num.uni-sb.de). This author was supported by Deutsche Forschungsgemeinschaft under grant Lo310/4-1.

Here $\mathrm{N}\left(A_{n}\right)^{\perp}$ is the orthogonal complement of the null space of $A_{n}$. If the range of $A$ is nonclosed in $Y$, that is, the generalized inverse of $A$ is unbounded, instabilities appear very likely in computing $f_{n}^{\dagger}$ directly from (1.2) under erroneous data $g_{n}$.

This reasoning led Louis and Maass [13] to the approximate inverse where one tries to reconstruct moments of $f:\left\langle f, e_{n}^{i}\right\rangle_{X}, i=1, \ldots, m$, with suitable mollifiers $e_{n}^{i}$. In case $X=L^{2}(\Omega), \Omega$ a domain in $\mathbb{R}^{d}$, one can think of the $e_{n}^{i}$ 's as smooth approximations to $\delta$-distributions located at points $x_{i} \in \Omega$.

The computations of the moments is achieved by approximating $e_{n}^{i}$ in the range of $A_{n}^{*}$. To any $e_{n}^{i}$ we associate a reconstruction kernel $v_{n}^{i} \in \mathbb{C}^{n}$ by minimizing the defect $\left\|A_{n}^{*} v_{n}^{i}-e_{n}^{i}\right\|_{X}$, that is, $v_{n}^{i}$ solves the normal equation

$$
\begin{equation*}
A_{n} A_{n}^{*} v_{n}^{i}=A_{n} e_{n}^{i} \tag{1.3}
\end{equation*}
$$

The above equation for $v_{n}^{i}$ is independent of the data $g_{n}$, therefore free of noise from measurement errors. We call $\left(e_{n}^{i}, v_{n}^{i}\right)$ a mollifier/reconstruction kernel pair for $A_{n}$.

The operator $S_{n}: \mathbb{C}^{n} \rightarrow \mathbb{C}^{m}$,

$$
\begin{equation*}
\left(S_{n} h\right)_{i}=\left\langle h, v_{n}^{i}\right\rangle_{\mathbb{C}^{n}}, \quad i=1, \ldots, m \tag{1.4}
\end{equation*}
$$

is called approximate inverse of $A_{n}$. Hence, $S_{n} g_{n}$ is an approximate solution of (1.1).
Lemma 1.1. If $g_{n}=A_{n} f$ then

$$
\begin{equation*}
\left(S_{n} g_{n}\right)_{i}=\left\langle f_{n}^{\dagger}, e_{n}^{i}\right\rangle_{X}, \quad i=1, \ldots, m \tag{1.5}
\end{equation*}
$$

Proof. The reconstruction kernels satisfy $A_{n}^{*} v_{n}^{i}=\mathcal{P}_{n} e_{n}^{i}$ where $\mathcal{P}_{n}: X \rightarrow X$ is the orthogonal projector onto $\mathrm{R}\left(A_{n}^{*}\right)=\mathrm{N}\left(A_{n}\right)^{\perp}$. Hence,

$$
\left(S_{n} g_{n}\right)_{i}=\left\langle f, A_{n}^{*} v_{n}^{i}\right\rangle_{X}=\left\langle\mathcal{P}_{n} f, e_{n}^{i}\right\rangle_{X}
$$

Since $\mathcal{P}_{n} f=f_{n}^{\dagger}($ see (1.2)), we are finished with the proof.
An interpretation of the approximate inverse as regularization scheme and further details are given by Louis [12]. He also shows how invariances of $A$ improve the efficiency; see Remark 5.2 below.

For several reasons we wish to avoid solving (1.3): $A_{n} A_{n}^{*}$ may be densely populated and ill-conditioned, increasing $n$ calls for a complete new computation of the kernels; invariances of $A$ do not show in $A_{n} A_{n}^{*}$ in general.

We propose the following technique to approximate $v_{n}$ (we will drop the superscript $i$ whenever considering a single pair $\left(e_{n}, v_{n}\right)$ ). Suppose $(e, v)$ is a mollifier/reconstruction kernel pair for $A$, i.e., $A^{*} v=e\left(A\right.$ is injective!). Then we expect $\Psi_{n} v$ to be an approximate solution of (1.3) where $e_{n}$ is equal or close to $e$. In section 3.1 we show convergence of $\Psi_{n} v$ to a solution of (1.3). We also analyze the situation when the mollifier $e$ is not in $\mathrm{R}\left(A^{*}\right)$ (section 3.2). Here we approximate $v_{n}$ by $\Psi_{n} v$ where $A^{*} v$ is close to $e_{n}$. We further discuss a technique to construct $v$ from $e$ which can be implemented.

In some applications, for instance, if $A$ is the Radon transform, $A_{n}: \mathrm{D}\left(A_{n}\right) \subset$ $X \rightarrow \mathbb{C}^{n}$ is unbounded and $A_{n}^{*}$ does not exist; see section 5 . Consequently, the concept of approximate inverse cannot be applied to (1.1). Louis and Schuster [16] replaced $A$ by a truncated singular value decomposition, thus circumventing the problem. We favor another cure which is closely related to our findings for a bounded $A_{n}$ (section 4).

In section 5 we apply the results from the previous sections to the reconstruction problem in 2D-computerized tomography, mainly to illustrate our rather abstract
results by a concrete application. As a byproduct we achieve a novel error estimate for the filtered backprojection algorithm as well as an alternative to design reconstruction filters.

To start this paper we introduce our technical set-up in the next section. Especially the operator $A_{n}$ is defined precisely. In the appendix we prove an auxiliary mapping property of the Radon transform.

Hegland and Anderssen [6] investigated a mollification method being akin to our approximate inverse approach. However, the details are completely different and they require stronger conditions on $A$; for instance, $A^{-1}$ has to be densely defined. Further, an implementation of their method requires an explicit knowledge of the preimages (under $A$ ) of the chosen basis functions. On the other hand, Hegland and Anderssen relate the regularization parameter (support width of the mollifier) to the discretization step size to bound the noise amplification error. This is an issue we do not address here.
2. Preliminaries. We specify our technical assumptions that are required to hold throughout the paper if not indicated otherwise.

The operator $A$ is supposed to have the mapping property (2.1). Let there be Banach spaces $X_{1}$ and $Y_{1}$ such that the embeddings $X_{1} \hookrightarrow X$ as well as $Y_{1} \hookrightarrow Y$ are continuous, injective, and dense. Moreover,

$$
\begin{equation*}
A: X_{1} \rightarrow Y_{1} \text { is continuous. } \tag{2.1}
\end{equation*}
$$

Let $Y_{1}^{\prime}$ be the dual to $Y_{1}$. One may consider the spaces $X_{1}$ and $Y_{1}$ as abstract smoothness classes in $X$ and $Y$, respectively.

We are now able to define the observation operator $\Psi_{n}: Y_{1} \rightarrow \mathbb{C}^{n}$ precisely: given $n$ functionals $\psi_{n, k} \in Y_{1}^{\prime}, k=1, \ldots, n$, let

$$
\begin{equation*}
\left(\Psi_{n} v\right)_{k}:=\left\langle\psi_{n, k}, v\right\rangle_{Y_{1}^{\prime} \times Y_{1}}, \quad k=1, \ldots, n \tag{2.2}
\end{equation*}
$$

where $\langle\cdot, \cdot\rangle_{Y_{1}^{\prime} \times Y_{1}}$ is the duality pairing on $Y_{1}^{\prime} \times Y_{1}$.
In applications we have in mind, typically, $Y_{1}$ will be a Sobolev space of sufficient order such that point evaluations are continuous.

It will prove useful to transform equation (1.1) into an equivalent equation where $\mathbb{C}^{n}$ is replaced by a suitable subspace of $Y$; see (2.6) below. To this end we introduce a family $\left\{V_{n}\right\}_{n \in \mathbb{N}}$ of finite dimensional subspaces of $Y$ being nested: $V_{n} \subset V_{n+1}$. Furthermore, each $V_{n}$ is spanned by basis elements $\varphi_{n, k}, k=1, \ldots, n$, which build a Riesz system with respect to $Y$, that is,

$$
\begin{equation*}
\sum_{k=1}^{n}\left|a_{k}\right|^{2} \preceq\left\|\sum_{k=1}^{n} a_{k} \varphi_{n, k}\right\|_{Y}^{2} \preceq \sum_{k=1}^{n}\left|a_{k}\right|^{2} \quad \text { for all } n \in \mathbb{N} . \tag{2.3}
\end{equation*}
$$

Our notation $A \preceq B$ indicates the existence of a generic constant $c>0$ such that $A \leq c B$. The constant $c$ will not depend on the arguments of $A$ and $B$. This means that the constants involved in (2.3) do not depend on $n$.

The spaces $\mathbb{C}^{n}$ and $V_{n}$ are related one-to-one by the operator $Q_{n}: \mathbb{C}^{n} \rightarrow V_{n}$, $Q_{n} a:=\sum_{k=1}^{n} a_{k} \varphi_{n, k}$. The composition of $\Psi_{n}$ and $Q_{n}$ creates a new operator $\Pi_{n}$ : $Y_{1} \rightarrow V_{n}$ as follows:

$$
\Pi_{n} v:=Q_{n} \Psi_{n} v=\sum_{k=1}^{n}\left\langle\psi_{n, k}, v\right\rangle_{Y_{1}^{\prime} \times Y_{1}} \varphi_{n, k}
$$

The operator $\Pi_{n}$ relates the observation operator $\Psi_{n}$ to $V_{n}$. Considered as an operator mapping $Y_{1}$ into $Y, \Pi_{n}$ is assumed to be uniformly bounded in $n$ :

$$
\begin{equation*}
\left\|\Pi_{n}\right\|_{Y_{1} \rightarrow Y} \preceq 1 \quad \text { as } n \rightarrow \infty \tag{2.4}
\end{equation*}
$$

Our last ingredient is the approximation property (2.5): let there be a sequence $\left\{\rho_{n}\right\} \subset$ $[0,1]$ converging monotonically to zero such that

$$
\begin{equation*}
\left\|v-\Pi_{n} v\right\|_{Y} \preceq \rho_{n}\|v\|_{Y_{1}} \quad \text { for all } v \in Y_{1} \text { as } n \rightarrow \infty . \tag{2.5}
\end{equation*}
$$

We understand $\left\{\rho_{n}\right\}$ as optimal, that is, $\left\{\rho_{n}\right\}$ is the fastest converging admissible sequence in (2.5).

Now we apply $Q_{n}$ from the left to both sides of (1.1) yielding

$$
\begin{equation*}
\widetilde{A}_{n} f=\widetilde{g}_{n} \tag{2.6}
\end{equation*}
$$

where $\widetilde{A}_{n}=Q_{n} A_{n}: X_{1} \rightarrow V_{n}$ and $\widetilde{g}_{n}=Q_{n} g_{n}$.
For the solution of (1.1) and (2.6), respectively, by the approximate inverse we distinguish two scenarios.

First, we assume that $A_{n}: \mathrm{D}\left(A_{n}\right) \subset X \rightarrow \mathbb{C}^{n}$ is bounded where $\mathrm{D}\left(A_{n}\right):=X_{1}$ is the domain of definition of $A_{n}$. Thus, $A_{n} \in \mathcal{L}\left(X, \mathbb{C}^{n}\right)$. Typical examples are integral operators which are sufficiently smoothing.

Example 2.1. Let $A: L^{2}(0,1) \rightarrow L^{2}(0,1), A f(x):=\int_{0}^{1} \mathrm{k}(x, y) f(y) d y$, where the kernel k is such that $A: L^{2}(0,1) \rightarrow H^{1 / 2+\varepsilon}(0,1)$ is bounded for an $\varepsilon>0$. On the Sobolev space $H^{1 / 2+\varepsilon}$ point evaluations are continuous functionals, so $\Psi_{n} g=$ $\left.n^{-1 / 2}\left(g\left(x_{1}\right), \ldots, g\left(x_{n}\right)\right)^{t}, x_{i} \in\right] 0,1[$, is the right choice if we are able to observe $A f$ at $x_{i}$. Thus, $A_{n}^{*} w(y)=n^{-1 / 2} \sum_{i} \overline{\mathrm{k}\left(x_{i}, y\right)} w_{i}$ and $\left(A_{n} A_{n}^{*}\right)_{i, j}=n^{-1} \int_{0}^{1} \overline{\mathrm{k}\left(x_{i}, y\right)} \mathrm{k}\left(x_{j}, y\right) d y$.

Second, we consider $A_{n}: \mathrm{D}\left(A_{n}\right) \subset X \rightarrow \mathbb{C}^{n}$ unbounded. Hence, the Hilbert space adjoint of $A_{n}$ cannot be defined on all of $\mathbb{C}^{n}$ (otherwise $A_{n}$ would have been continuous already). Here the worst case is $\mathrm{D}\left(A_{n}^{*}\right)=\{0\}$, so that the approximate inverse is not defined meaningful for (1.1). This happens for the Radon transform; see section 5.
3. Bounded semidiscrete operators $A_{n}$ : Approximating the discrete reconstruction kernel. Let (2.1) hold true with $X_{1}=X$ (topologically):

$$
\begin{equation*}
A: X \rightarrow Y_{1} \text { is continuous, } \tag{3.1}
\end{equation*}
$$

that is, $A_{n} \in \mathcal{L}\left(X, \mathbb{C}^{n}\right)$ and $\widetilde{A}_{n} \in \mathcal{L}(X, Y)$. In what follows we will denote the adjoint of $A: X \rightarrow Y$ by $A^{*}$.

Now we study convergence of the minimum norm solution $f_{n}^{\dagger}$ of (1.1) as $n \rightarrow \infty$. From this we derive a kind of pointwise convergence of the approximate inverse $S_{n}$.

Lemma 3.1. If (3.1) then $\left\|A-\widetilde{A}_{n}\right\|_{X \rightarrow Y} \leq \rho_{n}\|A\|_{X \rightarrow Y_{1}}$.
Proof. Since $\widetilde{A}_{n} x=\Pi_{n} A x$ for $x \in X$ one needs to apply (3.1) only.
THEOREM 3.2. Let $f_{n}^{\dagger}$ (1.2) be the minimum norm solution of (1.1) with $g_{n}=$ $A_{n} f$ for $f \in X$. Then

$$
\lim _{n \rightarrow \infty}\left\|f-f_{n}^{\dagger}\right\|_{X}=0
$$

Moreover, if the sequence of mollifiers $\left\{e_{n}^{i}\right\}_{n \in \mathbb{N}}$ converges to $e^{i} \in X, i=1, \ldots, m$, we have that

$$
\begin{equation*}
\lim _{n \rightarrow \infty} S_{n} A_{n} f=E f \tag{3.2}
\end{equation*}
$$

where

$$
\begin{equation*}
E: X \rightarrow \mathbb{C}^{m} \text { is defined by }(E f)_{i}:=\left\langle f, e^{i}\right\rangle_{X}, \quad i=1, \ldots, m \tag{3.3}
\end{equation*}
$$

Proof. Recall that $f_{n}^{\dagger}=\mathcal{P}_{n} f$ where $\mathcal{P}_{n}: X \rightarrow X$ is the orthogonal projection onto $\mathrm{N}\left(A_{n}\right)^{\perp}=\mathrm{N}\left(\widetilde{A}_{n}\right)^{\perp}$. Due to Lemma 3.1 and the injectivity of $A$ we have that $\bigcap_{n \in \mathbb{N}} \mathrm{~N}\left(\widetilde{A}_{n}\right)=\{0\}$. This yields the pointwise convergence of $\mathcal{P}_{n}$ to the identity operator in $X$ as $n \rightarrow \infty$, thereby proving the first assertion. The second assertion follows readily from (1.5).

Choosing special mollifiers $e_{n}^{i}$ we will show below that $\left\|S_{n} A_{n} f-E f\right\|_{\infty} \preceq \rho_{n}\|f\|_{X}$ as $n \rightarrow \infty$; see Corollary 3.8.

For an $e \in X$ we have either $e \in \mathrm{R}\left(A^{*}\right)$ or $e \in \partial \mathrm{R}\left(A^{*}\right)$ due to the injectivity of $A\left(\partial \mathrm{R}\left(A^{*}\right)\right.$ is the topological boundary of $\left.\mathrm{R}\left(A^{*}\right)\right)$. The first situation leads to reconstruction kernels $v$ satisfying $A^{*} v=e$. In section 3.1 below we shall show that $\Psi_{n} v$ is an approximate solution of (1.3) for suitable $e_{n}$.

If we cannot find a mollifier $e$ in the range of $A^{*}$, the equation $A^{*} y=e$ has no least squares solution. Thus, no reconstruction kernel is associated with $e$. We investigate the latter situation in section 3.2.
3.1. The special case $\boldsymbol{e} \in \mathbf{R}\left(\boldsymbol{A}^{*}\right)$. In Lemma 3.3 we derive a relation between the reconstruction kernels for $A_{n}$ and $\widetilde{A}_{n}$.

LEMMA 3.3. Let $\left(e, \widetilde{v}_{n}\right)$ be a mollifier/reconstruction kernel pair for $\widetilde{A}_{n}$ where $e \in X$ is arbitrary. Then, $\left(e, Q_{n}^{*} \widetilde{v}_{n}\right)$ is a mollifier/reconstruction kernel pair for $A_{n}$.

Proof. The assertion follows from $A_{n}^{*} Q_{n}^{*} \widetilde{v}_{n}=\widetilde{A}_{n}^{*} \widetilde{v}_{n}=\mathcal{P}_{n} e$ where $\mathcal{P}_{n}$ is as in the proof of Theorem 3.2.

Below we will need the Gramian matrix $G_{n} \in \mathbb{C}^{n \times n}$ relative to $\left\{\varphi_{n, 1}, \ldots, \varphi_{n, n}\right\}$. This matrix has entries $\left(G_{n}\right)_{i, j}=\left\langle\varphi_{n, i}, \varphi_{n, j}\right\rangle_{Y}$. A quick calculation validates the equality $G_{n} \Psi_{n} z=Q_{n}^{*} \Pi_{n} z$ for all $z \in Y_{1}$.

THEOREM 3.4. Adopt all assumptions specified in section 2 and assume (3.1). Let $\left(e_{n}, \widetilde{v}_{n}\right)$ be a mollifier/reconstruction kernel pair for $\widetilde{A}_{n}$ where $e_{n}=\widetilde{A}_{n}^{*} v, v \in Y_{1}$, and $\widetilde{v}_{n} \in \mathrm{~N}\left(\widetilde{A}_{n}^{*}\right)^{\perp}$. Then,

$$
\begin{equation*}
\left\|G_{n} \Psi_{n} v-Q_{n}^{*} \widetilde{v}_{n}\right\|_{\mathbb{C}^{n}} \preceq \rho_{n}\|v\|_{Y_{1}}+\inf _{y \in \mathrm{R}\left(\widetilde{A}_{n}\right)}\|v-y\|_{Y} \tag{3.4}
\end{equation*}
$$

as $n \rightarrow \infty$. Note that $\left(e_{n}, Q_{n}^{*} \widetilde{v}_{n}\right)$ is a mollifier/reconstruction kernel pair for $A_{n}$.
Proof. Since $\left\|Q_{n}^{*}\right\|_{Y \rightarrow \mathbb{C}^{n}} \preceq 1$ by (2.3) we may estimate

$$
\begin{aligned}
\left\|G_{n} \Psi_{n} v-Q_{n}^{*} \widetilde{v}_{n}\right\|_{\mathbb{C}^{n}} & \leq\left\|Q_{n}^{*} \Pi_{n} v-Q_{n}^{*} v\right\|_{\mathbb{C}^{n}}+\left\|Q_{n}^{*} v-Q_{n}^{*} \widetilde{v}_{n}\right\|_{\mathbb{C}^{n}} \\
& \preceq\left\|\Pi_{n} v-v\right\|_{Y}+\left\|v-\widetilde{v}_{n}\right\|_{Y} \\
& \preceq \rho_{n}\|v\|_{Y_{1}}+\left\|v-\widetilde{v}_{n}\right\|_{Y}
\end{aligned}
$$

where we used (2.5) in the final step. The assertion will be proved if we bound $\left\|v-\widetilde{v}_{n}\right\|_{Y}$ by a multiple of $\inf \left\{\|v-y\|_{Y} \mid y \in \underset{\sim}{\mathrm{R}}\left(\widetilde{A}_{n}\right)\right\}$.

Recall that $\widetilde{v}_{n}$ is the unique solution in $\mathrm{N}\left(\widetilde{A}_{n}^{*}\right)^{\perp}$ of the normal equation

$$
\begin{equation*}
\widetilde{A}_{n} \widetilde{A}_{n}^{*} \widetilde{v}_{n}=\widetilde{A}_{n} e_{n}=\widetilde{A}_{n} \widetilde{A}_{n}^{*} v \tag{3.5}
\end{equation*}
$$

Let $P_{n}: Y \rightarrow Y$ be the orthogonal projector onto $\mathrm{N}\left(\widetilde{A}_{n}^{*}\right)^{\perp}$. Since $P_{n} v$ solves (3.5) as well, we obtain $\widetilde{v}_{n}=P_{n} v$. As $\mathrm{N}\left(\widetilde{A}_{n}^{*}\right)^{\perp}=\mathrm{R}\left(\widetilde{A}_{n}\right)$ we proceed with

$$
\left\|v-\widetilde{v}_{n}\right\|_{Y}=\left\|v-P_{n} v\right\|_{Y}=\inf _{y \in \mathrm{R}\left(\widetilde{A}_{n}\right)}\|v-y\|_{Y}
$$

which completes the proof.
Corollary 3.5. The assumptions are those from Theorem 3.4. If either $v \in$ $\mathrm{R}(A)$ or all $A_{n}$ 's are onto then

$$
\left\|G_{n} \Psi_{n} v-Q_{n}^{*} \widetilde{v}_{n}\right\|_{\mathbb{C}^{n}} \preceq \rho_{n}\|v\|_{Y_{1}} \quad \text { as } n \rightarrow \infty
$$

Proof. First we consider $v \in \mathrm{R}(A)$. Let $v=A z$ for $z \in X$. Now

$$
\inf _{y \in \mathrm{R}\left(\widetilde{A}_{n}\right)}\|v-y\|_{Y} \leq\left\|A z-\Pi_{n} A z\right\|_{Y} \preceq \rho_{n}\|A z\|_{Y_{1}}
$$

by the approximation property (2.5). Second, if $A_{n}: X \rightarrow \mathbb{C}^{n}$ is onto we have that $\mathrm{R}\left(\widetilde{A}_{n}\right)=V_{n}$ which gives

$$
\inf _{y \in \mathrm{R}\left(\widetilde{A}_{n}\right)}\|v-y\|_{Y}=\inf _{y \in V_{n}}\|v-y\|_{Y} \leq\left\|v-\Pi_{n} v\right\|_{Y} \preceq \rho_{n}\|v\|_{Y_{1}}
$$

In both cases the assertion follows from (3.4).
Even so $e_{n}=\widetilde{A}_{n}^{*} v$ converges to $e=A^{*} v$ due to Lemma 3.1, $e_{n}$ may be an unsuitable mollifier for fixed (possibly small) $n$. It seems natural to work with $e$ in the semidiscrete setting also. This more general situation is considered in the following lemma where we, however, allow a weighted norm in $\mathbb{C}^{n}$. Under the assumptions of Lemma 3.6 below, $\left\|A_{n} A_{n}^{*} \cdot\right\|_{\mathbb{C}^{n}}$ is a norm on $\mathbb{C}^{n}$ being, in general, weaker than the Euclidean norm in the following sense. There exist positive constants $\gamma_{n}$ and $\Gamma$ such that $\gamma_{n}\|z\|_{\mathbb{C}^{n}} \leq\left\|A_{n} A_{n}^{*} z\right\|_{\mathbb{C}^{n}} \leq \Gamma\|z\|_{\mathbb{C}^{n}}$ for all $z \in \mathbb{C}^{n}$ where $\Gamma$ does not depend on $n$ and where $\gamma_{n}$ tends to zero as $n$ grows.

Lemma 3.6. Let $e=A^{*} v$ for $v \in Y_{1}$. Further, let $\left(e, \bar{v}_{n}\right)$ be a mollifier/reconstruction kernel pair for $\widetilde{A}_{n}$ where $\bar{v}_{n} \in \mathrm{~N}\left(\widetilde{A}_{n}^{*}\right)^{\perp}$. Under the assumptions of Theorem 3.4 and provided all $A_{n}$ 's are onto we have that

$$
\left\|A_{n} A_{n}^{*}\left(G_{n} \Psi_{n} v-Q_{n}^{*} \bar{v}_{n}\right)\right\|_{\mathbb{C}^{n}} \preceq \rho_{n}\|v\|_{Y_{1}} \quad \text { as } n \rightarrow \infty
$$

Proof. Let $\widetilde{v}_{n}$ be as in Theorem 3.4. Hence,

$$
\begin{aligned}
\left\|A_{n} A_{n}^{*}\left(G_{n} \Psi_{n} v-Q_{n}^{*} \bar{v}_{n}\right)\right\|_{\mathbb{C}^{n}} \leq\|A\|_{X \rightarrow Y_{1}}^{2} \| G_{n} & \Psi_{n} v-Q_{n}^{*} \widetilde{v}_{n} \|_{\mathbb{C}^{n}} \\
& +\left\|A_{n} \widetilde{A}_{n}^{*} \widetilde{v}_{n}-A_{n} \widetilde{A}_{n}^{*} \bar{v}_{n}\right\|_{\mathbb{C}^{n}} \\
& \preceq \rho_{n}\|v\|_{Y_{1}}+\left\|\widetilde{A}_{n} \widetilde{A}_{n}^{*} \widetilde{v}_{n}-\widetilde{A}_{n} \widetilde{A}_{n}^{*} \bar{v}_{n}\right\|_{Y}
\end{aligned}
$$

where we used Corollary 3.5, (2.3), and the estimate

$$
\begin{equation*}
\left\|A_{n} A_{n}^{*}\right\|_{\mathbb{C}^{n} \rightarrow \mathbb{C}^{n}}=\left\|A_{n}\right\|_{X \rightarrow \mathbb{C}^{n}}^{2} \preceq\left\|\Pi_{n} A\right\|_{X \rightarrow Y}^{2} \preceq\|A\|_{X \rightarrow Y_{1}}^{2} \tag{3.6}
\end{equation*}
$$

by (2.3) and (2.4). Since $\widetilde{A}_{n} \widetilde{A}_{n}^{*} \widetilde{v}_{n}=\widetilde{A}_{n} \widetilde{A}_{n}^{*} v$ and $\widetilde{A}_{n} \widetilde{A}_{n}^{*} \bar{v}_{n}=\widetilde{A}_{n} A^{*} v$ we obtain that

$$
\left\|\widetilde{A}_{n} \widetilde{A}_{n}^{*} \widetilde{v}_{n}-\widetilde{A}_{n} \widetilde{A}_{n}^{*} \bar{v}_{n}\right\|_{Y} \preceq\left\|\widetilde{A}_{n}^{*}-A^{*}\right\|_{Y \rightarrow X}\|v\|_{Y}
$$

The assertion of Lemma 3.6 is now due to Lemma 3.1.
We discuss the implications of Corollary 3.5 on the approximate inverse $S_{n}$ of $A_{n}$ (1.4). Here one has mollifier/reconstruction kernel pairs $\left(e_{n}^{i}, v_{n}^{i}\right), i=$ $1, \ldots, m$; see (1.3). Now let $e_{n}^{i}=\widetilde{A}_{n}^{*} v^{i}$ where $v^{i} \in Y_{1}, i=1, \ldots, m$. Our investigations from above suggest to replace the (unknown) approximate inverse $S_{n}$ by the (computable) operator $\Sigma_{n}$ defined by

$$
\begin{equation*}
\left(\Sigma_{n} b\right)_{i}=\left\langle b, G_{n} \Psi_{n} v^{i}\right\rangle_{\mathbb{C}^{n}}, \quad i=1, \ldots, m \tag{3.7}
\end{equation*}
$$

As a direct consequence of Corollary 3.5, we can show that $\Sigma_{n}$ is a reasonable substitute for $S_{n}$.

ThEOREM 3.7. The assumptions are those from Theorem 3.4. Further, let $\left(e_{n}^{i}, v_{n}^{i}\right), i=1, \ldots, m$, be mollifier/reconstruction kernel pairs for $A_{n}$ where $e_{n}^{i}=$ $\widetilde{A}_{n}^{*} v^{i}$. Assume that all $v^{i}$ 's are in $Y_{1}$. If either all $v^{i}$ 's are in $\mathrm{R}(A)$ or all $A_{n}$ 's are onto, then

$$
\begin{equation*}
\left\|S_{n} A_{n} f-\Sigma_{n} A_{n} f\right\|_{\infty} \preceq \rho_{n} \max _{1 \leq i \leq m}\left\|v^{i}\right\|_{Y_{1}}\|f\|_{X} \quad \text { as } n \rightarrow \infty \tag{3.8}
\end{equation*}
$$

Proof. Let $\left(e_{n}^{i}, \widetilde{v}_{n}^{i}\right)$ be the mollifier/reconstruction kernel pair for $\widetilde{A}_{n}$ where $\widetilde{v}_{n}^{i} \in$ $\mathrm{N}\left(\widetilde{A}_{n}^{*}\right)^{\perp}$. From Lemma 3.3 we know that $\left(e_{n}^{i}, Q_{n}^{*} \widetilde{v}_{n}^{i}\right)$ is a mollifier/reconstruction kernel pair for $A_{n}$. Note that $Q_{n}^{*} \widetilde{v}_{n}^{i}$ may be different from the kernel $v_{n}^{i}$ used in $S_{n}$; however, $A_{n}^{*} v_{n}^{i}=A_{n}^{*} Q_{n}^{*} \widetilde{v}_{n}^{i}$. Thus,

$$
\left(S_{n} A_{n} f\right)_{i}=\left\langle f, A_{n}^{*} v_{n}^{i}\right\rangle_{X}=\left\langle f, A_{n}^{*} Q_{n}^{*} \widetilde{v}_{n}^{i}\right\rangle_{X}=\left\langle A_{n} f, Q_{n}^{*} \widetilde{v}_{n}^{i}\right\rangle_{\mathbb{C}^{n}}
$$

which implies that
$\left|\left(S_{n} A_{n} f\right)_{i}-\left(\Sigma_{n} A_{n} f\right)_{i}\right|=\left|\left\langle A_{n} f, Q_{n}^{*} \widetilde{v}_{n}^{i}-G_{n} \Psi_{n} v^{i}\right\rangle_{\mathbb{C}^{n}}\right| \leq\left\|A_{n} f\right\|_{\mathbb{C}^{n}}\left\|Q_{n}^{*} \widetilde{v}_{n}^{i}-G_{n} \Psi_{n} v^{i}\right\|_{\mathbb{C}^{n}}$.
The estimate (3.8) follows now from (3.6) and from Corollary 3.5.
The following fact on the convergence speed of the approximate inverse is worthwhile to mention; compare (3.2).

Corollary 3.8. We have that

$$
\left\|S_{n} A_{n} f-E f\right\|_{\infty} \preceq \rho_{n}\|f\|_{X} \max _{1 \leq i \leq m}\left\|v^{i}\right\|_{Y_{1}} \quad \text { as } n \rightarrow \infty
$$

Proof. By the triangle inequality and by (3.8) it suffices to show that $\| \Sigma_{n} A_{n} f-$ $E f\left\|_{\infty} \preceq \rho_{n}\right\| f\left\|_{X} \max _{1 \leq i \leq m}\right\| v^{i} \|_{Y_{1}}$. This is obtained from

$$
\begin{aligned}
\left|\left(\Sigma_{n} A_{n} f\right)_{i}-\left\langle f, e^{i}\right\rangle_{X}\right| & =\left|\left\langle\Psi_{n} A f, G_{n} \Psi_{n} v^{i}\right\rangle_{\mathbb{C}^{n}}-\left\langle f, A^{*} v^{i}\right\rangle_{X}\right| \\
& =\left|\left\langle\Pi_{n} g, \Pi_{n} v^{i}\right\rangle_{Y}-\left\langle g, v^{i}\right\rangle_{Y}\right|
\end{aligned}
$$

where $g=A f$. The difference on the right-hand side may now be estimated as follows:

$$
\begin{aligned}
\left|\left\langle\Pi_{n} g, \Pi_{n} v^{i}\right\rangle_{Y}-\left\langle g, v^{i}\right\rangle_{Y}\right| & \leq\left\|\Pi_{n} g-g\right\|_{Y}\left\|\Pi_{n} v^{i}\right\|_{Y}+\left\|\Pi_{n} v^{i}-v^{i}\right\|_{Y}\|g\|_{Y} \\
& \preceq \rho_{n}\|g\|_{Y_{1}}\left\|v^{i}\right\|_{Y_{1}} \preceq \rho_{n}\|f\|_{X}\left\|v^{i}\right\|_{Y_{1}}
\end{aligned}
$$

where we used the uniform boundedness (2.4), the approximation property (2.5), and the continuity (3.1).
3.2. The general case $\boldsymbol{e} \in \boldsymbol{X}$. The range of $A^{*}$ is dense in $X$ due to the injectivity of $A$. Therefore, we will assume only that the mollifier can be approximated arbitrarily close by an element in $\mathrm{R}\left(A^{*}\right)$.

Let $e^{i} \in X$ be mollifiers for $i=1, \ldots, m$. To any $\varepsilon_{i}>0$ we can find a $v^{i} \in Y_{1}$ so that

$$
\begin{equation*}
\left\|e^{i}-A^{*} v^{i}\right\|_{X} \leq \varepsilon_{i}, \quad i=1, \ldots, m \tag{3.9}
\end{equation*}
$$

Below we will demonstrate how to get $v^{i}$ from $e^{i}$ knowing a singular value decomposition of $A$.

Since, in general, no reconstruction kernel is associated with $e^{i}$ there will be no counterparts of Theorems 3.4 and 3.7, respectively. Instead, we are directly heading towards an estimate of $\Sigma_{n} A_{n} f-E f$. Based on the $e^{i}$ 's and the $v^{i}$,s from above the operators $E$ (3.3) and $\Sigma_{n}(3.7)$ are well defined.

Theorem 3.9. Adopt the assumptions specified in section 2 and assume (3.1). Let the operators $E$ and $\Sigma_{n}$ be defined as in (3.3) and (3.7), respectively, where $e^{i} \in X$ and $v^{i} \in Y_{1}$ are related by (3.9). Then

$$
\begin{equation*}
\left\|\Sigma_{n} A_{n} f-E f\right\|_{\infty} \preceq\left(\rho_{n} \max _{1 \leq i \leq m}\left\|v^{i}\right\|_{Y_{1}}+\max _{1 \leq i \leq m} \varepsilon_{i}\right)\|f\|_{X} \quad \text { as } n \rightarrow \infty . \tag{3.10}
\end{equation*}
$$

Proof. By the triangle inequality and by (3.9) we get

$$
\left|\left(\Sigma_{n} A_{n} f\right)_{i}-\left\langle f, e^{i}\right\rangle_{X}\right| \leq\left|\left\langle\Psi_{n} A f, G_{n} \Psi_{n} v^{i}\right\rangle_{\mathbb{C}^{n}}-\left\langle f, A^{*} v^{i}\right\rangle_{X}\right|+\|f\|_{X} \varepsilon_{i}
$$

We may now proceed as in the proof of Corollary 3.8.
We will now discuss the vital issue of constructing $v^{i} \in Y_{1}$ from $e^{i} \in X$ which satisfy (3.9) for $\varepsilon_{i}$ arbitrarily small. For convenience let us suppress the superscript $i$.

The tool we employ is a singular value decomposition (SVD) of the operator $A$. In medical imaging SVDs are explicitly known; see, e.g., $[9,10,15,17,18,21]$.

Let $A: X \rightarrow Y$ be a compact operator and let $\left\{\mathrm{v}_{k}, u_{k} ; \sigma_{k} \mid k \in \mathbb{N}_{0}\right\}$ be its singular system, that is,

$$
A x=\sum_{k=0}^{\infty} \sigma_{k}\left\langle x, \mathrm{v}_{k}\right\rangle_{X} u_{k} .
$$

The sets of singular functions $\left\{\mathrm{v}_{k}\right\}$ and $\left\{u_{k}\right\}$ are orthonormal bases in $X(A$ is injective) and $\overline{\mathrm{R}(A)}$, respectively. The positive numbers $\sigma_{k}$ are the singular values of $A$ satisfying $\lim _{k \rightarrow \infty} \sigma_{k}=0$ (monotonically). The singular functions and the singular values are related via

$$
A \mathrm{v}_{k}=\sigma_{k} u_{k} \quad \text { and } \quad A^{*} u_{k}=\sigma_{k} \mathrm{v}_{k} .
$$

We assume that all $u_{k}$ 's are in $Y_{1}$. For an arbitrary $e \in X$ we follow the approach of Dietz [4] and define

$$
\begin{equation*}
v_{M}:=\sum_{k=0}^{M-1} \sigma_{k}^{-1}\left\langle e, \mathrm{v}_{k}\right\rangle_{X} u_{k} \tag{3.11}
\end{equation*}
$$

which is an element of $Y_{1}$. Dietz [4] implemented (3.11) to solve the cone beam reconstruction problem in 3D utilizing the formula of Grangeat.

Obviously,

$$
\begin{equation*}
\left\|e-A^{*} v_{M}\right\|_{X}^{2}=\sum_{k=M}^{\infty}\left|\left\langle e, \mathrm{v}_{k}\right\rangle_{X}\right|^{2} \rightarrow 0 \quad \text { as } M \rightarrow \infty \tag{3.12}
\end{equation*}
$$

Incorporating an abstract smoothness assumption on $e$, we are able to give convergence rates of $\left\|e-A^{*} v_{M}\right\|_{X}$ as $M \rightarrow \infty$.

Lemma 3.10. Suppose that $e \in \mathrm{R}\left(\left(A^{*} A\right)^{\alpha}\right)=\mathrm{D}\left(\left(A^{*} A\right)^{-\alpha}\right)$ for $a \alpha \geq 0$. Then

$$
\lim _{M \rightarrow \infty} \sigma_{M}^{-\alpha}\left\|e-A^{*} v_{M}\right\|_{X}=0
$$

Moreover, the following error estimate holds:

$$
\left\|e-A^{*} v_{M}\right\|_{X}<\sigma_{M}^{\alpha} \sqrt{\|e\|_{X}\left\|\left(A^{*} A\right)^{-\alpha} e\right\|_{X}}
$$

Proof. We have that

$$
\begin{aligned}
\left\|e-A^{*} v_{M}\right\|_{X}^{2} & =\sum_{k=M}^{\infty} \sigma_{k}^{-2 \alpha}\left|\left\langle e, \mathrm{v}_{k}\right\rangle_{X}\right| \sigma_{k}^{2 \alpha}\left|\left\langle e, \mathrm{v}_{k}\right\rangle_{X}\right| \\
& \leq\left(\sum_{k=M}^{\infty} \sigma_{k}^{-4 \alpha}\left|\left\langle e, \mathrm{v}_{k}\right\rangle_{X}\right|^{2}\right)^{1 / 2} \sigma_{M}^{2 \alpha}\left(\sum_{k=M}^{\infty}\left|\left\langle e, \mathrm{v}_{k}\right\rangle_{X}\right|^{2}\right)^{1 / 2}
\end{aligned}
$$

and both assertions follow readily.
In view of (3.10) we realize that controlling the $\varepsilon_{i}$ 's tells only half of the story. To learn the whole story we look at $\left\|v_{M}\right\|_{Y_{1}}$.

Lemma 3.11. Suppose that $e \in \mathrm{R}\left(\left(A^{*} A\right)^{\alpha}\right)=\mathrm{D}\left(\left(A^{*} A\right)^{-\alpha}\right)$ for $a \geq 0$. Further, let there exist a $\beta \geq 0$ such that $\left\|u_{k}\right\|_{Y_{1}} \preceq \sigma_{k}^{-\beta}$ for all $k$. Then

$$
\left\|v_{M}\right\|_{Y_{1}} \preceq\left\|\left(A^{*} A\right)^{-\alpha} e\right\|_{X}\left(\sum_{k=0}^{M-1} \sigma_{k}^{4 \alpha-2(1+\beta)}\right)^{1 / 2}
$$

Proof. The straightforward estimates

$$
\begin{aligned}
\left\|v_{M}\right\|_{Y_{1}} & \preceq \sum_{k=0}^{M-1} \sigma_{k}^{-2 \alpha}\left|\left\langle e, \mathrm{v}_{k}\right\rangle_{X}\right| \sigma_{k}^{2 \alpha-(1+\beta)} \\
& \preceq\left(\sum_{k=0}^{M-1} \sigma_{k}^{-4 \alpha}\left|\left\langle e, \mathrm{v}_{k}\right\rangle_{X}\right|^{2}\right)^{1 / 2}\left(\sum_{k=0}^{M-1} \sigma_{k}^{4 \alpha-2(1+\beta)}\right)^{1 / 2}
\end{aligned}
$$

verify the claim.
TheOrem 3.12. Let $A: X \rightarrow Y$ be compact with singular system $\left\{\mathrm{v}_{k}, u_{k} ; \sigma_{k} \mid k \in\right.$ $\left.\mathbb{N}_{0}\right\}$. Assume that $\sigma_{k} \asymp(k+1)^{-\gamma}$ for $a \gamma>0$ as $k \rightarrow \infty(a \asymp b$ abbreviates $a \preceq b \preceq a)$ and that $\left\|u_{k}\right\|_{Y_{1}} \preceq \sigma_{k}^{-\beta}$ for $\beta \geq 0$.

Assume the hypotheses of Theorem 3.9; in particular, let the operators $E$ and $\Sigma_{n}$ be defined as in (3.3) and (3.7), respectively, where $e^{i} \in \mathrm{D}\left(\left(A^{*} A\right)^{-\alpha}\right)$ and $v_{M_{i}}^{i}$ are related by (3.11).

If $\alpha>(1+\beta) / 2+1 /(4 \gamma)$ and $M_{i}=M_{i}(n) \succeq \rho_{n}^{-1 /(\alpha \gamma)}$ as $n \rightarrow \infty\left(\rho_{n}\right.$ from (2.5)), then

$$
\begin{equation*}
\left\|\Sigma_{n} A_{n} f-E f\right\|_{\infty} \preceq \rho_{n}\|f\|_{X} \max _{1 \leq i \leq m}\left\|\left(A^{*} A\right)^{-\alpha} e^{i}\right\|_{X} \quad \text { as } n \rightarrow \infty \tag{3.13}
\end{equation*}
$$

Proof. Since $\left\|e^{i}\right\|_{X} \preceq\left\|\left(A^{*} A\right)^{-\alpha} e^{i}\right\|_{X}$ we have that

$$
\begin{aligned}
\varepsilon_{i}=\left\|e^{i}-A^{*} v_{M_{i}}^{i}\right\|_{X} & \preceq \sigma_{M_{i}}^{\alpha}\left\|\left(A^{*} A\right)^{-\alpha} e^{i}\right\|_{X} \\
& \preceq\left(M_{i}+1\right)^{-\alpha \gamma}\left\|\left(A^{*} A\right)^{-\alpha} e^{i}\right\|_{X} \preceq \rho_{n}\left\|\left(A^{*} A\right)^{-\alpha} e^{i}\right\|_{X}
\end{aligned}
$$

by Lemma 3.10 and our assumption on $M_{i}=M_{i}(n)$ as $n \rightarrow \infty$. Further, by Lemma 3.11,

$$
\left\|v_{M_{i}}^{i}\right\|_{Y_{1}} \preceq\left\|\left(A^{*} A\right)^{-\alpha} e^{i}\right\|_{X}\left(\sum_{k=0}^{\infty}(k+1)^{-\gamma(4 \alpha-2(1+\beta))}\right)^{1 / 2}
$$

where the series converges due to $\gamma(4 \alpha-2(1+\beta))>1$. Recalling Theorem 3.9 we are finished with the proof of (3.13).
4. Unbounded semidiscrete operators $\boldsymbol{A}_{\boldsymbol{n}}$. Here we consider (2.1) where $X_{1}$ is a proper subspace of $X$ with a stronger topology.

As we will see in the next section it may happen that $A_{n}: X_{1} \subset X \rightarrow \mathbb{C}^{n}$ is unbounded. In the extremest case we even have to deal with $\mathrm{D}\left(A_{n}^{*}\right)=\{0\}$, that is, the approximate inverse with respect to the topology in $X$ is not defined for (1.1).

Basically, this leaves us with the situation already investigated in section 3.2. Indeed, if $\left(e^{i}, v^{i}\right) \in X \times Y_{1}, i=1, \ldots, m$, are mollifier/reconstruction kernel pairs satisfying (3.9) then $E$ (3.3) as well as $\Sigma_{n}$ (3.7) are well defined. Even for unbounded operators $A_{n}$ both Theorems 3.9 and 3.12 remain valid with a slight modification: we have to assume that $f \in X_{1}$. In (3.10) as well as in (3.13) we have to replace $\|f\|_{X}$ by $\|f\|_{X_{1}}$.
5. Application to the reconstruction problem in 2D-computerized tomography. We apply our abstract results of the former sections to the reconstruction problem in 2D-computerized tomography, that is, the reconstruction of a function from its line integrals. For further applications of our results in vector and local tomography we refer to [24] and [22], respectively.

The underlying operator is the Radon transform $\mathbf{R}$ mapping a function $f \in L^{2}(\Omega)$ to its line integrals. Here, $\Omega$ is the unit ball in $\mathbb{R}^{2}$ centered at the origin. More precisely,

$$
\begin{equation*}
\mathbf{R} f(s, \vartheta):=\int_{L(s, \vartheta) \cap \Omega} f(x) d \sigma(x) \tag{5.1}
\end{equation*}
$$

The lines are parameterized by $L(s, \vartheta)=\left\{\tau \omega^{\perp}(\vartheta)+s \omega(\vartheta) \mid \tau \in \mathbb{R}\right\}$ where $\left.s \in\right]-1,1[$, $\omega(\vartheta)=(\cos \vartheta, \sin \vartheta)^{t}$ and $\omega^{\perp}(\vartheta)=(-\sin \vartheta, \cos \vartheta)^{t}$ for $\left.\vartheta \in\right] 0, \pi[$. By this parameterization of lines we are dealing with the parallel scanning geometry.

The Radon transform maps $X=L^{2}(\Omega)$ continuously to $Y=L^{2}(Z)$ where $Z:=$ $]-1,1[\times] 0, \pi[$; see, e.g., Natterer [20, Chap. II.1]. In the appendix we will verify the following mapping property (see Theorem A. 2 below):

$$
\mathbf{R}: H_{0}^{\alpha}(\Omega) \rightarrow H^{\alpha+1 / 2}(Z) \quad \text { is continuous for any } \alpha \geq 0
$$

The involved Sobolev spaces are defined as follows. By $H_{0}^{\alpha}(\Omega)$ we denote the closure of $\mathcal{C}_{0}^{\infty}(\Omega)$, the space of infinitely differentiable functions with compact support in $\Omega$, with respect to the norm $\|f\|_{\alpha}^{2}=\int_{\mathbb{R}^{2}}\left(1+\|\xi\|^{2}\right)^{\alpha}|\widehat{f}(\xi)|^{2} d \xi$. Here, $\widehat{f}$ is the Fourier transform of $f$.

The space $H^{\beta}(Z)=W_{2}^{\beta}(Z)$ is an $L^{2}$-Sobolev space defined on the rectangular domain $Z$; see, e.g., Wloka [25].

Since point evaluations are continuous linear functionals on $H^{\beta}(Z)$ for $\beta>1$ we set $X_{1}=H_{0}^{1 / 2+\kappa}(\Omega)$ and $Y_{1}=H^{1+\kappa}(Z)$ for a $\kappa>0$; cf. (2.1).

For $q, p \in \mathbb{N}$ let $h_{s}=1 / q$ and $h_{\vartheta}=\pi / p$ be the discretization step sizes and set $s_{i}=i h_{s}, i=-q, \ldots, q$, and $\vartheta_{j}=j h_{\vartheta}, j=0, \ldots, p$. Let $\ell \in\{1,2\}$. With this index $\ell$ we will be able to distinguish between two different settings using a compact notation.

To the pairs $\left(s_{i}, \vartheta_{j}\right)$ we associate the Dirac-distributions $\psi_{i, j}^{(\ell)}$ given by

$$
\left\langle\psi_{i, j}^{(\ell)}, g\right\rangle_{Y_{1}^{\prime} \times Y_{1}}:=\varsigma_{i, j}^{(\ell)} g\left(s_{i}, \vartheta_{j}\right), \quad i=-q, \ldots, q_{\ell}, j=0, \ldots, p_{\ell}
$$

for any $g \in H^{1+\kappa}(Z)$ where $q_{1}=q-1, q_{2}=q, p_{1}=p-1$ and $p_{2}=p$. The $\varsigma_{i, j}^{(\ell)}$,s are normalization factors to be defined below in (5.3). We define the mapping $\Psi_{q, p}^{(\ell)}: H^{1+\kappa}(Z) \rightarrow \mathbb{R}^{n_{\ell}}$ according to (2.2) using the $\psi_{i, j}^{(\ell)}$ 's. The respective dimensions are $n_{1}=2 q p$ and $n_{2}=(2 q+1)(p+1)$.

THEOREM 5.1. The operator $\mathbf{R}_{q, p}^{(\ell)}=: \Psi_{q, p}^{(\ell)} \mathbf{R}: H_{0}^{1 / 2+\kappa}(\Omega) \subset L^{2}(\Omega) \rightarrow \mathbb{R}^{n_{\ell}}$ is unbounded for any $\kappa>0$. Moreover, $\mathrm{D}\left(\left(\mathbf{R}_{q, p}^{(\ell)}\right)^{*}\right)=\{0\}$.

Proof. We construct a sequence $\left\{f_{r}\right\}_{r \in \mathbb{N}} \subset H_{0}^{1 / 2+\kappa}(\Omega)$ with $\left\|f_{r}\right\|_{L^{2}(\Omega)} \leq 1$ and $\left\|\mathbf{R}_{q, p}^{(\ell)} f_{r}\right\|_{\mathbb{R}^{n_{\ell}}} \rightarrow \infty$ as $r \rightarrow \infty$.

We will define $f_{r}$ as the tensor product of two univariate functions $\chi_{r}$ and $\mu_{r}$. Let $\left\{\alpha_{r}\right\}$ and $\left\{\beta_{r}\right\}$ be monotonically decreasing zero sequences with $0<\alpha_{r}<1$, $0<\beta_{r}<1 / 2$, and $\alpha_{r}^{2}+\left(1-\beta_{r}\right)^{2}<1$.

Let $\chi_{r} \in \mathcal{C}_{0}^{\infty}\left(-\alpha_{r}, \alpha_{r}\right)$ with values in $\left[0,1 / \sqrt{2 \alpha_{r}}\right]$ such that $\chi_{r}(t)=1 / \sqrt{2 \alpha_{r}}$ for $|t| \leq \alpha_{r} / 2$. Similarly, let $\mu_{r} \in \mathcal{C}_{0}^{\infty}\left(-1+\beta_{r}, 1-\beta_{r}\right)$ with values in $\left[0,1 / \sqrt{2\left(1-\beta_{r}\right)}\right]$ such that $\mu_{r}(t)=1 / \sqrt{2\left(1-\beta_{r}\right)}$ for $|t| \leq 1-2 \beta_{r}$. Both functions can be constructed explicitly using a partition of unity; see, e.g., Wloka [25, Chap. 1.2].

For $f_{r}(x):=\chi_{r}\left(x_{1}\right) \mu_{r}\left(x_{2}\right)$ we have $0<\left\|f_{r}\right\|_{L^{2}(\Omega)}=\left\|\chi_{r}\right\|_{L^{2}(\mathbb{R})}\left\|\mu_{r}\right\|_{L^{2}(\mathbb{R})} \leq 1$ and $\operatorname{supp} f_{r} \subset\left[-\alpha_{r}, \alpha_{r}\right] \times\left[-1+\beta_{r}, 1-\beta_{r}\right] \subset \Omega$ because $\alpha_{r}^{2}+\left(1-\beta_{r}\right)^{2}<1$. Thus, $f_{r} \in H_{0}^{1 / 2+\kappa}(\Omega)$ for any $\kappa>0$.

Now consider $\mathbf{R} f_{r}$ at $s_{0}=0$ and $\vartheta_{0}=0$ :

$$
\begin{aligned}
\left|\mathbf{R} f_{r}\left(s_{0}, \vartheta_{0}\right)\right| & =\int_{\mathbb{R}} f_{r}(0, t) d t=\chi_{r}(0) \int_{\mathbb{R}} \mu_{r}(t) d t \\
& \geq \chi_{r}(0) \int_{-1+2 \beta_{r}}^{1-2 \beta_{r}} \mu_{r}(t) d t=\frac{1-2 \beta_{r}}{\sqrt{\alpha_{r}\left(1-\beta_{r}\right)}} \xrightarrow{r \rightarrow \infty} \infty .
\end{aligned}
$$

Hence, $\left\|\mathbf{R}_{q, p}^{(\ell)} f_{r}\right\|_{\mathbb{R}^{n} \ell} \rightarrow \infty$ as $r \rightarrow \infty$.
We are now going to verify the second statement of Theorem 5.1. To this end observe that $\lim _{r \rightarrow \infty} \mathbf{R} f_{r}\left(s_{i}, \vartheta_{j}\right)=0$ if $(i, j) \neq(0,0)$. This limit holds since supp $f_{r}$ "converges" to the line segment $L(0,0) \cap \Omega$.

The construction principle from above can be repeated for any pair $\left(s_{i}, \vartheta_{j}\right),\left|s_{i}\right|<$ 1, leading to a sequence of functions $\left\{f_{r}^{i, j}\right\}_{r \in \mathbb{N}} \subset H_{0}^{1 / 2+\kappa}(\Omega)$ with $\left\|f_{r}^{i, j}\right\|_{L^{2}(\Omega)} \leq 1$ and

$$
\lim _{r \rightarrow \infty} \mathbf{R} f_{r}^{i, j}\left(s_{k}, \vartheta_{l}\right)=\left\{\begin{array}{ccc}
\infty & : & (k, l)=(i, j) \\
0 & : & \text { otherwise }
\end{array}\right.
$$

Assume that $0 \neq w \in \mathrm{D}\left(\left(\mathbf{R}_{q, p}^{(\ell)}\right)^{*}\right)$. Then the linear functional $f \mapsto\left\langle\mathbf{R}_{q, p}^{(\ell)} f, w\right\rangle_{\mathbb{R}^{n} \ell}$ is continuous on $\mathrm{D}\left(\mathbf{R}_{q, p}^{(\ell)}\right)$ with respect to the $L^{2}(\Omega)$-topology. With $w_{i, j} \neq 0$ we obtain

$$
\left\langle\mathbf{R}_{q, p}^{(\ell)} f_{r}^{i, j}, w\right\rangle_{\mathbb{R}^{n} \ell}=w_{i, j} \varsigma_{i, j}^{(\ell)} \mathbf{R} f_{r}^{i, j}\left(s_{i}, \vartheta_{j}\right)+\sum_{\substack{(k, l) \\(k, l) \neq(i, j)}} w_{k, l} \varsigma_{k, l}^{(\ell)} \mathbf{R} f_{r}^{i, j}\left(s_{k}, \vartheta_{l}\right)
$$

which implies that $\lim _{r \rightarrow \infty}\left\langle\mathbf{R}_{q, p}^{(\ell)} f_{r}^{i, j}, w\right\rangle_{\mathbb{R}^{n_{\ell}}}=\operatorname{sgn}\left(w_{i, j}\right) \infty$. However, this unboundedness contradicts $w \in \mathrm{D}\left(\left(\mathbf{R}_{q, p}^{(\ell)}\right)^{*}\right)$. $\quad \square$

Due to Theorem 5.1 the approximate inverse cannot be applied to the 2D-reconstruction problem

$$
\text { given } g_{q, p} \in \mathbb{R}^{n_{\ell}} \text { find } f \in L^{2}(\Omega) \text { such that } \mathbf{R}_{q, p}^{(\ell)} f=g_{q, p}
$$

Here we are facing the situation from section 4, that is, we have to replace the "nonexisting" $S_{q, p}$ by $\Sigma_{q, p}$; compare (1.3), (1.4), and (3.7), respectively.

Canonical candidates for approximation spaces related to $\Psi_{q, p}^{(\ell)}$ are the tensor product spline spaces $V_{q, p}^{(\ell)}=S_{s}^{(\ell)} \otimes S_{\vartheta}^{(\ell)}, \ell=1,2$. Here, $S_{s}^{(\ell)}$ and $S_{\vartheta}^{(\ell)}$ are either the piecewise constant $(\ell=1)$ or linear $(\ell=2)$ spline spaces with respect to the knot sequences $\left\{s_{i}\right\}$ and $\left\{\vartheta_{j}\right\}$, respectively. As basis in $V_{q, p}^{(\ell)}$ we choose the tensor product B-spline basis

$$
\begin{equation*}
\left\{B_{q, i}^{(\ell)} \otimes B_{p, j}^{(\ell)} / \varsigma_{i, j}^{(\ell)} \mid-q \leq i \leq q_{\ell}, 0 \leq j \leq p_{\ell}\right\} . \tag{5.2}
\end{equation*}
$$

The B-splines $B_{q, i}^{(\ell)} \in S_{s}^{(\ell)}$ and $B_{p, j}^{(\ell)} \in S_{\vartheta}^{(\ell)}$ are uniquely determined by ( $\chi_{D}$ is the indicator function of the set $D$ )

$$
B_{q, i}^{(1)}=\chi_{\left[s_{i}, s_{i+1}[ \right.}, \quad B_{p, j}^{(1)}=\chi_{\left[\vartheta_{j}, \vartheta_{j+1}[ \right.}
$$

and

$$
B_{q, i}^{(2)}\left(s_{k}\right)=\left\{\begin{array}{l}
1: \quad i=k, \\
0: \text { otherwise },
\end{array} \quad B_{p, j}^{(2)}\left(\vartheta_{l}\right)=\left\{\begin{array}{l}
1: \quad j=l \\
0: \text { otherwise }
\end{array}\right.\right.
$$

respectively. The normalization factors $\varsigma_{i, j}^{(\ell)}$ are just the $L^{2}$-norms of the B-splines:

$$
\begin{equation*}
\varsigma_{i, j}^{(\ell)}:=\left\|B_{q, i}^{(\ell)} \otimes B_{p, j}^{(\ell)}\right\|_{L^{2}(Z)}, \quad i=-q, \ldots, q_{\ell}, j=0, \ldots, p_{\ell} \tag{5.3}
\end{equation*}
$$

Thus, the normalized tensor product B-spline basis (5.2) is an $L^{2}(Z)$-Riesz system where the constants in the corresponding norm equivalence do not depend on $h_{s}$ or $h_{\vartheta}$; compare (2.3).

We next define the interpolation operator $\Pi_{q, p}^{(\ell)}: H^{1+\kappa}(Z) \rightarrow V_{q, p}^{(\ell)}$ which links $V_{q, p}^{(\ell)}$ to $\Psi_{q, p}^{(\ell)}$ :

$$
\Pi_{q, p}^{(\ell)} v:=\sum_{i=-q}^{q_{\ell}} \sum_{j=0}^{p_{\ell}}\left(\Psi_{q, p}^{(\ell)} v\right)_{i, j} B_{q, i}^{(\ell)} \otimes B_{p, j}^{(\ell)} / \varsigma_{i, j}^{(\ell)}=\sum_{i=-q}^{q_{\ell}} \sum_{j=0}^{p_{\ell}} v\left(s_{i}, \vartheta_{j}\right) B_{q, i}^{(\ell)} \otimes B_{p, j}^{(\ell)} .
$$

Let $h=\max \left\{h_{s}, h_{\vartheta}\right\}$. Then the uniform boundedness

$$
\left\|\Pi_{q, p}^{(\ell)} v\right\|_{L^{2}(Z)} \preceq\|v\|_{H^{\alpha}(Z)}, \quad \alpha>1
$$

as well as the approximation property

$$
\left\|v-\Pi_{q, p}^{(\ell)} v\right\|_{L^{2}(Z)} \preceq h^{\min \{\alpha, \ell\}}\|v\|_{H^{\alpha}(Z)}, \quad \alpha>1
$$

hold true whenever the right-hand sides are finite. Both estimates are standard results from spline approximation theory; see, e.g., Schumaker [23, Chap. 12].

In the following we apply our results of section 3.2 to the 2D-reconstruction problem. In a first step we therefore construct reconstruction kernels from mollifiers using a SVD of the Radon transform. Unfortunately, a SVD of $\mathbf{R}: L^{2}(\Omega) \rightarrow L^{2}(Z)$ is not known explicitly. However, it can be shown that the Radon transform maps $L^{2}(\Omega)$ compactly to $L^{2}\left(\widetilde{Z}, w^{-1}\right)$ where $\left.\widetilde{Z}=\right]-1,1[\times] 0,2 \pi[$; see, e.g., Natterer [20, Chap. IV.3]. The weight function is given by $w(s):=\sqrt{1-s^{2}}$ and acts on the first variable only.

Let $\left\{\mathrm{v}_{m, l}, u_{m, l} ; \sigma_{m}\left|m \in \mathbb{N}_{0}, l \in \mathbb{Z},|l| \leq m, m+l \in 2 \mathbb{Z}\right\}\right.$ be the singular system of $\mathbf{R}: L^{2}(\Omega) \rightarrow L^{2}\left(\widetilde{Z}, w^{-1}\right)$, that is,

$$
\begin{equation*}
\mathbf{R} f=\sum_{m=0}^{\infty} \sum_{l=-m}^{m} \sigma_{m}\left\langle f, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)} u_{m, l} \tag{5.4}
\end{equation*}
$$

where $\star$ restricts the summation over those $l$ 's with $m+l \in 2 \mathbb{Z}$.
Later on we will need an explicit representation of the $\sigma_{m}$ 's and the $u_{m, l}$ 's only. We therefore give analytic expressions

$$
\begin{equation*}
\sigma_{m}=2 \sqrt{\frac{\pi}{m+1}} \quad \text { and } \quad u_{m, l}(s, \varphi)=\frac{1}{\pi} w(s) U_{m}(s) \mathrm{e}^{\imath l \varphi} \tag{5.5}
\end{equation*}
$$

where $U_{m}(s)=\sin ((m+1) \arccos s) / \sin (\arccos s), m \in \mathbb{N}_{0}$, are the Chebyshev polynomials of the second kind. For the $\mathrm{v}_{m, l}$ 's see Louis [11] or Natterer [20].

Denoting by $\mathbf{R}^{*}$ and $\mathbf{R}^{\#}$ the adjoints of $\mathbf{R}: L^{2}(\Omega) \rightarrow L^{2}(Z)$ and $\mathbf{R}: L^{2}(\Omega) \rightarrow$ $L^{2}\left(\widetilde{Z}, w^{-1}\right)$, respectively, we have that

$$
2 \mathbf{R}^{*} w^{-1} u_{m, l}=\mathbf{R}^{\#} u_{m, l}=\sigma_{m} \mathrm{v}_{m, l}
$$

The first equality can be checked by straightforward calculations. Given a mollifier $e \in L^{2}(\Omega)$ normalized by $\int_{\Omega} e(x) d x=1$ and centered about the origin we define

$$
v_{M}:=2 \sum_{m=0}^{M-1} \sum_{l=-m}^{m} \sigma_{m}^{-1}\left\langle e, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)} w^{-1} u_{m, l}
$$

which then gives

$$
\begin{equation*}
\left\|\mathbf{R}^{*} v_{M}-e\right\|_{L^{2}(\Omega)}^{2}=\sum_{m=M}^{\infty} \sum_{l=-m}^{m}{ }^{\star}\left|\left\langle e, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)}\right|^{2} \tag{5.6}
\end{equation*}
$$

compare (3.11) and (3.12).
Let us assume from now on that the mollifier $e$ is a radial function, that is, $e(x)=\mathbf{e}\left(\|x\|_{\mathbb{R}^{2}}\right)$. Since $\left\langle e, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)}=0$ for $l \neq 0$ the representation of $v_{M}$ simplifies to

$$
\begin{equation*}
v_{M}=2 \sum_{k=0}^{(M-1) / 2} \sigma_{2 k}^{-1}\left\langle e, \mathrm{v}_{2 k, 0}\right\rangle_{L^{2}(\Omega)} w^{-1} u_{2 k, 0} \tag{5.7}
\end{equation*}
$$

Hence, the reconstruction kernel does not depend on the angle $\vartheta$. Moreover $v_{M}$ is an even function in $s$ as so are the Chebyshev polynomials of even degree. See Figure 5.1 for an example.

Let $x_{i} \in \Omega, i=1, \ldots, m$, be points in which we would like to reconstruct moments $\left\langle f, e^{i}\right\rangle_{L^{2}(\Omega)}$ from the data $g_{q, p}$. The mollifiers $e^{i}$ are derived from $e$ by translation and dilation:

$$
e^{i}(\cdot)=T_{1}^{x_{i}} e(\cdot):=\frac{1}{4} e\left(\frac{\cdot-x_{i}}{2}\right)
$$

At the present time the choice of the dilation factor 2 seems to be artificial; however, it will become clear in the proof of Lemma 5.3 below.



Fig. 5.1. Reconstruction kernels (left) and radial parts of the related mollifiers (right). Solid curves: reconstruction kernel $v_{501}(5.7)$ corresponding to $e(x)=5 \gamma^{-2} p\left(\|x\|_{\mathbb{R}^{2}} / \gamma\right) / \pi$ with $\gamma=0.05$ where $p(t)=\left(1-t^{2}\right)^{4}$ for $t \leq 1$ and $p(t)=0$ otherwise. From a numerical point of view we have $\mathbf{R}^{*} v_{501}=e . \quad$ Dashed curves: reconstruction kernel $v$ (5.14) corresponding to the Gaußian $e(x)=(2 \pi)^{-1} \gamma^{-2} \exp \left(-\gamma^{-2}\|x\|_{\mathbb{R}^{2}}^{2} / 2\right)$ with $\gamma=0.013$. Please note that both kernels are negative in $[0.2,1]$ and monotonically increasing.

The invariance property

$$
\begin{equation*}
\mathbf{R}^{*} T_{2}^{x_{i}}=T_{1}^{x_{i}} \mathbf{R}^{*}, \quad \text { where } \quad T_{2}^{x_{i}} v(s, \vartheta):=\frac{1}{4} v\left(\frac{s-x_{i}^{t} \omega(\vartheta)}{2}, \vartheta\right) \tag{5.8}
\end{equation*}
$$

suggests to define the reconstruction kernel $v_{M}^{i}$ associated with $e^{i}$ by

$$
v_{M}^{i}(s, \vartheta):=T_{2}^{x_{i}} v_{M}(s, \vartheta)=\frac{1}{4} v_{M}\left(\frac{s-x_{i}^{t} \omega(\vartheta)}{2}\right), \quad i=1, \ldots, m
$$

Thus,

$$
\mathbf{R}^{*} v_{M}^{i}=T_{1}^{x_{i}} \mathbf{R}^{*} v_{M} \longrightarrow T_{1}^{x_{i}} e=e^{i} \quad \text { as } \quad M \rightarrow \infty .
$$

Remark 5.2. Thanks to the invariance property (5.8) only the kernel $v_{M}$ has to be computed and stored. The kernels for the reconstruction points $x_{i}$ are simply found by the action of $T_{2}^{x_{i}}$ on $v_{M}$.

Lemma 5.3. Let $v \in H^{r}(Z), r \geq 0$. Then

$$
\begin{equation*}
\left\|T_{2}^{x} v\right\|_{H^{r}(Z)} \preceq\|v\|_{H^{r}(Z)} \quad \text { uniformly in } x \in \Omega \tag{5.9}
\end{equation*}
$$

Proof. The transformation $\Phi(s, \vartheta):=\left(\left(s-x^{t} \omega(\vartheta)\right) / 2, \vartheta\right)$ maps $Z$ to $Z^{\prime}$ bijectively where

$$
Z^{\prime}=\left\{(\sigma, \varphi) \in Z \left\lvert\, \sigma \in\left[\frac{-1-x^{t} \omega(\varphi)}{2}, \frac{1-x^{t} \omega(\varphi)}{2}\right]\right.\right\} .
$$

Moreover, $\Phi$ is a $\mathcal{C}^{\infty}$-diffeomorphism with $\operatorname{det} J \Phi(s, \vartheta)=1 / 2$ where $J \Phi$ is the Jacobian of $\Phi$. Since $T_{2}^{x} v=v \circ \Phi / 4$ the assertion follows from transformation results for Sobolev norms, see; e.g., Wloka [25].

Define $E: L^{2}(\Omega) \rightarrow \mathbb{R}^{m}$ by $(E f)_{i}:=\left\langle f, e^{i}\right\rangle_{L^{2}(\Omega)}, i=1, \ldots, m$; see (3.3), and $\Sigma_{q, p}^{(\ell)}: \mathbb{R}^{n_{\ell}} \rightarrow \mathbb{R}^{m}$ by

$$
\left(\Sigma_{q, p}^{(\ell)} b\right)_{i}:=\left\langle b, G_{q, p}^{(\ell)} \Psi_{q, p}^{(\ell)} v_{M}^{i}\right\rangle_{\mathbb{R}^{n} \ell}, \quad \ell=1,2
$$

see (3.7). Here, $G_{q, p}^{(\ell)}$ is the Gramian matrix with respect to the B-spline basis in $V_{q, p}^{(\ell)}$. Especially $G_{q, p}^{(1)}$ is the identity matrix.

The process of evaluating $\Sigma_{q, p}^{(\ell)} g_{q, p}$ coincides with (and may be implemented exactly as) the filtered backprojection algorithm in computerized tomography with filter function $v_{M}$; see, e.g., Natterer [20, Chap. V.1]. Indeed, for $\ell=1$,

$$
\left(\Sigma_{q, p}^{(1)} g_{q, p}\right)_{i}=\frac{\pi}{4 q p} \sum_{l=-q}^{q-1} \sum_{j=0}^{p-1} g\left(s_{l}, \vartheta_{j}\right) v_{M}\left(\frac{s_{l}-x_{i}^{t} \omega\left(\vartheta_{j}\right)}{2}\right)
$$

A reformulation of Theorem 3.12 in the present context results therefore in a novel error analysis of the filtered backprojection algorithm (Theorem 5.4 below). Compared to already known error estimates, see, e.g., Natterer [20, Th. V.1.1], we allow mild smoothness assumptions on the density distribution $f$. Further, the kernel $v$ needs only to be known approximately. The error bound reflects clearly the influence of the smoothness of $f$ and $e$ on the convergence rate.

THEOREM 5.4. Let $f \in H_{0}^{1 / 2+\kappa}(\Omega)$ for $0<\kappa \leq 1$. Assume that the radial mollifier $e$ is in $H_{0}^{\alpha}(\Omega)$ for $\alpha>4+2 \kappa$. Let $\lambda_{\ell}=\min \{\overline{1}+\kappa, \ell\}$ for $\ell=1,2$.

If $M=M(h) \succeq h^{-2 \lambda_{\ell} / \alpha}$, then

$$
\begin{equation*}
\left\|\Sigma_{q, p}^{(\ell)} \Psi_{q, p}^{(\ell)} \mathbf{R} f-E f\right\|_{\infty} \preceq h^{\lambda_{\ell}}\|f\|_{1 / 2+\kappa}\|e\|_{\alpha} \quad \text { as } h \rightarrow 0 \tag{5.10}
\end{equation*}
$$

Proof. We follow the line of proof of Theorem 3.9 to obtain

$$
\begin{aligned}
\mid\left(\Sigma_{q, p}^{(\ell)} \Psi_{q, p}^{(\ell)}\right. & \mathbf{R} f)_{i}-\left\langle f, e^{i}\right\rangle_{L^{2}(\Omega)} \mid \\
& \leq\left|\left(\Sigma_{q, p}^{(\ell)} \Psi_{q, p}^{(\ell)} \mathbf{R} f\right)_{i}-\left\langle f, \mathbf{R}^{*} v_{M}^{i}\right\rangle_{L^{2}(\Omega)}\right|+\left|\left\langle f, \mathbf{R}^{*} v_{M}^{i}-e^{i}\right\rangle_{L^{2}(\Omega)}\right| \\
& =\left|\left(\Sigma_{q, p}^{(\ell)} \Psi_{q, p}^{(\ell)} \mathbf{R} f\right)_{i}-\left\langle\mathbf{R} f, T_{2}^{x_{i}} v_{M}\right\rangle_{L^{2}(Z)}\right|+\left|\left\langle f, T_{1}^{x_{i}}\left(\mathbf{R}^{*} v_{M}-e\right)\right\rangle_{L^{2}(\Omega)}\right| \\
& \preceq\|f\|_{1 / 2+\kappa}\left(h^{\lambda_{\ell}}\left\|v_{M}\right\|_{H^{1+\kappa}(Z)}+\left\|\mathbf{R}^{*} v_{M}-e\right\|_{L^{2}(\Omega)}\right)
\end{aligned}
$$

where we used the invariance property (5.8) as well as (5.9). Since

$$
\left\|\left(\mathbf{R}^{\#} \mathbf{R}\right)^{-\alpha} e\right\|_{L^{2}(\Omega)} \preceq\|e\|_{\alpha}
$$

(see Lemma A. 3 below), we immediately infer from (5.6) and from the proof of Lemma 3.10 that

$$
\left\|\mathbf{R}^{*} v_{M}-e\right\|_{L^{2}(\Omega)} \preceq \sigma_{M}^{\alpha}\|e\|_{\alpha} \preceq(M+1)^{-\alpha / 2}\|e\|_{\alpha} \preceq h^{\lambda_{\ell}}\|e\|_{\alpha} .
$$

It remains to bound $\left\|v_{M}\right\|_{H^{1+\kappa}(Z)}$; see (5.7). We will be guided by the proof of Lemma 3.11. Using the interpolation inequality for Sobolev norms; see, e.g., Lions and Magenes [7, Chap. 2.5], we may estimate as follows:

$$
\begin{aligned}
\left\|w^{-1} u_{2 k, 0}\right\|_{H^{1+\kappa}(Z)} & \preceq\left\|w^{-1} u_{2 k, 0}\right\|_{H^{1}(Z)}^{1-\kappa}\left\|w^{-1} u_{2 k, 0}\right\|_{H^{2}(Z)}^{\kappa} \\
& =\left\|U_{2 k}\right\|_{H^{1}(-1,1)}^{1-\kappa}\left\|U_{2 k}\right\|_{H^{2}(-1,1)}^{\kappa} .
\end{aligned}
$$

A bound on the Sobolev norms of the Chebyshev polynomials may be obtained by Markov's inequality (5.11); see, e.g., Lorentz [8, Chap. 3.3]: let $P_{r}$ be a polynomial of degree $r$ then

$$
\begin{equation*}
\left|P_{r}^{\prime}(s)\right| \leq r^{2} \max _{-1 \leq t \leq 1}\left|P_{r}(t)\right|, \quad|s| \leq 1 \tag{5.11}
\end{equation*}
$$

With $\max _{-1 \leq t \leq 1}\left|U_{r}(t)\right|=r+1$ we easily find that

$$
\left\|U_{r}\right\|_{H^{1}(-1,1)} \preceq(r+1)^{3} \preceq \sigma_{r}^{-6} \quad \text { and } \quad\left\|U_{r}\right\|_{H^{2}(-1,1)} \preceq(r+1)^{5} \preceq \sigma_{r}^{-10}
$$

which result in

$$
\left\|w^{-1} u_{2 k, 0}\right\|_{H^{1+\kappa}(Z)} \preceq \sigma_{2 k}^{-2(3+2 \kappa)}
$$

Recalling the representation (5.7) of $v_{M}$ we get

$$
\begin{aligned}
\left\|v_{M}\right\|_{H^{1+\kappa}(Z)} & \preceq\left(\sum_{k=0}^{(M-1) / 2} \sigma_{2 k}^{-4 \alpha}\left|\left\langle e, \mathrm{v}_{2 k, 0}\right\rangle_{L^{2}(\Omega)}\right|^{2}\right)^{1 / 2}\left(\sum_{k=0}^{(M-1) / 2}(2 k+1)^{-2 \alpha+7+4 \kappa}\right)^{1 / 2} \\
& \preceq\left\|\left(\mathbf{R}^{\#} \mathbf{R}\right)^{-\alpha} e\right\|_{L^{2}(\Omega)} \preceq\|e\|_{\alpha} .
\end{aligned}
$$

We used the fact that the second sum is bounded in $M$ since $2 \alpha-7-4 \kappa>1$. The proof of Theorem 5.4 is now complete.

Because $h=\max \left\{h_{s}, h_{\vartheta}\right\}$ it is most efficient-in view of (5.10)-to work with discretization step sizes $h_{s}$ and $h_{\vartheta}$ which coincide: $h_{s}=h_{\vartheta}$, that is, $p=\pi q$. So we recovered the optimal sampling relation for the parallel scanning geometry; see, e.g., Natterer [20, Chap. III].

In the remainder of this section we comment briefly on another way to design reconstruction kernels for the Radon transform; see (5.13) below. This approach is based on the inversion formula (5.12) of the Radon transform

$$
\begin{equation*}
e=(2 \pi)^{-1} \mathbf{R}^{*} \mathrm{I}^{-1} \mathbf{R} e \quad \text { for } \quad e \in H_{0}^{\alpha}(\Omega), \quad \alpha \geq 1 / 2 \tag{5.12}
\end{equation*}
$$

see, e.g., Natterer [20, Chap. II.2]. The operator $\mathrm{I}^{-1}: H_{0}^{1}(-1,1) \rightarrow L^{2}(\mathbb{R})$ is the Riesz potential: $\widehat{\left(\mathrm{I}^{-1} f\right)}(\xi)=|\xi| \widehat{f}(\xi)$. In (5.12), the Riesz potential acts on the first variable of $\mathbf{R} e$. Motivated by (5.12) we make the ansatz $v:=\mathrm{I}^{-1} \mathbf{R} e /(2 \pi)$. Assuming radial symmetry of $e$ the latter formula may be expressed as

$$
\begin{equation*}
v(s)=\frac{1}{\pi} \int_{0}^{\infty} \sigma \widehat{e}(\sigma \omega(0)) \cos (s \sigma) d \sigma \tag{5.13}
\end{equation*}
$$

compare Natterer [20, (1.5), p. 103].
For instance, let $e$ be the Gaußian $e(x)=(2 \pi)^{-1} \gamma^{-2} \exp \left(-\gamma^{-2}\|x\|_{\mathbb{R}^{2}}^{2} / 2\right), \gamma>0$. Clearly, these mollifiers are not supported in $\Omega$. However, for $\gamma$ small, they decay fast enough to consider them elements of $H_{0}^{1 / 2}(\Omega)$. Thus,

$$
\begin{aligned}
v(s) & =\frac{1}{2 \pi^{2}} \int_{0}^{\infty} \sigma \exp \left(-\gamma^{2} \sigma^{2} / 2\right) \cos (s \sigma) d \sigma \\
& =\frac{-1}{2 \pi^{2} \gamma^{2}} \int_{0}^{\infty} \frac{\mathrm{d}}{\mathrm{~d} \sigma}\left(\exp \left(-\gamma^{2} \sigma^{2} / 2\right)\right) \cos (s \sigma) d \sigma
\end{aligned}
$$

Now applying integration by parts and using formulæ (7.4.7) and (7.1.3) from [1] yields

$$
\begin{equation*}
v(s)=\frac{1}{2 \pi^{2} \gamma^{2}}\left(1+\sqrt{\frac{\pi}{2}} \frac{s}{\gamma} \exp \left(-s^{2} /\left(2 \gamma^{2}\right)\right) \imath \operatorname{erf}(\imath s /(\sqrt{2} \gamma))\right) \tag{5.14}
\end{equation*}
$$

where $\operatorname{erf}(t)=(2 / \sqrt{\pi}) \int_{0}^{t} \exp \left(-z^{2}\right) d z$ is the error function. Figure 5.1 displays $v$ (5.14) for $\gamma=0.013$ (dashed curves).

Remark 5.5. We recommend the filter design methods from above whenever one wants to impose certain conditions on the mollifier, e.g., nonnegativity and compact support; see Figure 5.1. The widely used Shepp-Logan filter and its noncompactly supported mollifier have frequent sign changes. To avoid artifacts in the reconstructions these oscillations require a certain fine-tuning: the dilation parameter $\gamma$ (compare (5.14)) needs to be selected carefully. In contrast, the reconstructions based on the filters from Figure 5.1 are more robust with respect to the support width of the mollifier.
A. Appendix: A Sobolev space estimate of the 2D-radon transform. In this appendix we will show that the Radon transform (5.1) maps $H_{0}^{\alpha}(\Omega)$ boundedly to $H_{\mathrm{p}}^{\alpha+1 / 2}(\widetilde{Z}), \alpha \geq 0$, where $\left.\widetilde{Z}=\right]-1,1[\times] 0,2 \pi\left[\right.$. The space $H_{\mathrm{p}}^{\beta}(\widetilde{Z})$ is a Sobolev space of periodic functions. Let $g \in L^{2}(\widetilde{Z})$ be expressed in its Fourier series, that is,

$$
g(s, \varphi)=\sum_{k \in \mathbb{Z}} \sum_{n \in \mathbb{Z}} \widehat{g}_{k, n} \mathrm{e}^{\imath(\pi k s+n \varphi)}, \quad \widehat{g}_{k, n}=\frac{1}{4 \pi} \int_{\tilde{Z}} g(s, \varphi) \mathrm{e}^{-\imath(\pi k s+n \varphi)} d \varphi d s
$$

Then, $g \in H_{\mathrm{p}}^{\beta}(\widetilde{Z}), \beta \geq 0$, iff the norm

$$
\|g\|_{\mathrm{p}, \beta}^{2}=\sum_{k \in \mathbb{Z}} \sum_{n \in \mathbb{Z}}\left(1+k^{2}+n^{2}\right)^{\beta}\left|\widehat{g}_{k, n}\right|^{2}
$$

is finite.
Remark A.1. Interpreting periodic functions in $L^{2}(\widetilde{Z})$ as functions defined on the torus $\mathcal{T} \subset \mathbb{R}^{3}$ we may identify $H_{\mathrm{p}}^{\beta}(\widetilde{Z})$ with the Sobolev space $H^{\beta}(\mathcal{T})$ defined on the smooth compact manifold $\mathcal{T}$ by means of local coordinates; see, e.g., Wloka [25].

In proving our main result in Theorem A. 2 below we will benefit from a known Sobolev space estimate (A.1) for the Radon transform due to Louis and Natterer [14]; see also [20, Chap. II.5].

Let $H^{(\beta, 0)}(\widetilde{Z}), \beta>0$, be the tensor product $H_{0}^{\beta}(-1,1) \widehat{\otimes} L^{2}(0,2 \pi)$ (for the tensor product of Sobolev spaces, see, e.g., Aubin [2]); then, for $\alpha \geq 0$,

$$
\begin{equation*}
\|\mathbf{R} f\|_{H^{(\alpha+1 / 2,0)}(\widetilde{Z})} \preceq\|f\|_{\alpha} \quad \text { for all } f \in H_{0}^{\alpha}(\Omega) \tag{A.1}
\end{equation*}
$$

In view of Remark A. 1 our estimate (A.2) below is intrinsically different from a result of Natterer which looks similar at first glance; see [20, Chap. II, Thrm. 5.3].

Theorem A.2. The Radon transform maps $H_{0}^{\alpha}(\Omega)$ continuously to $H_{\mathrm{p}}^{\alpha+1 / 2}(\widetilde{Z})$, $\alpha \geq 0$, that is,

$$
\begin{equation*}
\|\mathbf{R} f\|_{\mathrm{p}, \alpha+1 / 2} \preceq\|f\|_{\alpha} \quad \text { for all } f \in H_{0}^{\alpha}(\Omega) . \tag{A.2}
\end{equation*}
$$

Proof. Let $g=\mathbf{R} f$. Since $\left(1+k^{2}+n^{2}\right)^{\beta} \leq 2^{\beta}\left(\left(1+k^{2}\right)^{\beta}+\left(1+n^{2}\right)^{\beta}\right)$ for $\beta \geq 0$ and for all $k, n \in \mathbb{Z}$ we have that

$$
\|g\|_{\mathrm{p}, \alpha+1 / 2} \preceq A(g)+B(g)
$$

with

$$
A(g)^{2}=\sum_{k, n \in \mathbb{Z}}\left(1+k^{2}\right)^{\alpha+1 / 2}\left|\widehat{g}_{k, n}\right|^{2} \quad \text { and } \quad B(g)^{2}=\sum_{k, n \in \mathbb{Z}}\left(1+n^{2}\right)^{\alpha+1 / 2}\left|\widehat{g}_{k, n}\right|^{2} .
$$

We will bound $A(g)$ as well as $B(g)$ by a multiple of $\|f\|_{\alpha}$. Both relations

$$
\sum_{n \in \mathbb{Z}}\left|\widehat{g}_{k, n}\right|^{2}=\int_{0}^{2 \pi}\left|\frac{1}{2} \int_{-1}^{1} g(s, \varphi) \mathrm{e}^{-\imath \pi k s} d s\right|^{2} d \varphi
$$

and

$$
\sum_{k \in \mathbb{Z}}\left|\widehat{g}_{k, n}\right|^{2}=\int_{-1}^{1}\left|\frac{1}{2 \pi} \int_{0}^{2 \pi} g(s, \varphi) \mathrm{e}^{-\imath n \varphi} d \varphi\right|^{2} d s
$$

follow from Parseval's identity. Hence,

$$
\begin{aligned}
A(g)^{2} & =\frac{1}{4} \int_{0}^{2 \pi} \sum_{k \in \mathbb{Z}}\left(1+k^{2}\right)^{\alpha+1 / 2}\left|\int_{-1}^{1} g(s, \varphi) \mathrm{e}^{-\imath \pi k s} d s\right|^{2} d \varphi \\
& \preceq\|g\|_{H^{(\alpha+1 / 2,0)}(\widetilde{Z})}^{2} \preceq\|f\|_{\alpha}^{2}
\end{aligned}
$$

where the first inequality follows from a Sobolev norm equivalence given by Natterer in [20, Chap. VII, Lem. 4.4]. The second inequality comes from (A.1).

Estimating $B(g)$ is a little bit more involved. From the singular value expansion (5.4) of $g=\mathbf{R} f$ we deduce that

$$
\begin{aligned}
\frac{1}{2 \pi} \int_{0}^{2 \pi} g(s, \varphi) \mathrm{e}^{-\imath n \varphi} d \varphi & =\frac{1}{\pi} \sum_{m=0}^{\infty} \sum_{l=-m}^{m} g_{m, l} w(s) U_{m}(s) \underbrace{\frac{1}{2 \pi} \int_{0}^{2 \pi} \mathrm{e}^{-\imath(n-l) \varphi} d \varphi}_{=\delta_{n, l}} \\
& =\frac{1}{\pi} \sum_{\mu=0}^{\infty} g_{|n|+2 \mu, n} w(s) U_{|n|+2 \mu}(s)
\end{aligned}
$$

with $g_{m, l}=\sigma_{m}\left\langle f, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)}$. Thus,

$$
\begin{aligned}
B(g)^{2} & =\frac{1}{\pi^{2}} \sum_{n \in \mathbb{Z}}\left(1+n^{2}\right)^{\alpha+1 / 2} \int_{-1}^{1}\left|\sum_{\mu=0}^{\infty} g_{|n|+2 \mu, n} w(s) U_{|n|+2 \mu}(s)\right|^{2} d s \\
& \leq \frac{1}{2 \pi} \sum_{n \in \mathbb{Z}}\left(1+n^{2}\right)^{\alpha+1 / 2} \int_{-1}^{1}\left|\sum_{\mu=0}^{\infty} g_{|n|+2 \mu, n} \sqrt{\frac{2}{\pi}} w(s) U_{|n|+2 \mu}(s)\right|^{2} w^{-1}(s) d s \\
& =\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}}\left(1+n^{2}\right)^{\alpha+1 / 2} \sum_{\mu=0}^{\infty}\left|g_{|n|+2 \mu, n}\right|^{2}
\end{aligned}
$$

because $\left\{\sqrt{2 / \pi} w(\cdot) U_{m}(\cdot) \mid m \in \mathbb{N}\right\}$ is an orthonormal basis in $L^{2}(]-1,1\left[, w^{-1}\right)$. Further,

$$
\begin{aligned}
B(g)^{2} & \preceq \sum_{n \in \mathbb{Z}}(1+|n|)^{2 \alpha+1} \sum_{\mu=0}^{\infty} \sigma_{|n|+2 \mu}^{2}\left|\left\langle f, \mathrm{v}_{|n|+2 \mu, n}\right\rangle_{L^{2}(\Omega)}\right|^{2} \\
& \preceq \sum_{n \in \mathbb{Z}} \sum_{\mu=0}^{\infty} \sigma_{|n|+2 \mu}^{-4 \alpha}\left|\left\langle f, \mathrm{v}_{|n|+2 \mu, n}\right\rangle_{L^{2}(\Omega)}\right|^{2} \\
& \leq \sum_{m=0}^{\infty} \sum_{l=-m}^{m} \sigma_{m}^{-4 \alpha}\left|\left\langle f, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)}\right|^{2} .
\end{aligned}
$$

In Lemma A. 3 below we bound the latter double sum by a multiple of $\|f\|_{\alpha}^{2}$ which finally proves (A.2).

LEmmA A.3. The operator $\left(\mathbf{R}^{\#} \mathbf{R}\right)^{-\alpha}: H_{0}^{\alpha}(\Omega) \rightarrow L^{2}(\Omega), \alpha \geq 0$, is continuous, that is,

$$
\begin{equation*}
\left\|\left(\mathbf{R}^{\#} \mathbf{R}\right)^{-\alpha} f\right\|_{0}^{2}=\sum_{m=0}^{\infty} \sum_{l=-m}^{m} \sigma_{m}^{-4 \alpha}\left|\left\langle f, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)}\right|^{2} \preceq\|f\|_{\alpha}^{2} \tag{A.3}
\end{equation*}
$$

for all $f \in H_{0}^{\alpha}(\Omega)$.
Proof. Since $\mathcal{C}_{0}^{\infty}(\Omega)$ is dense $H_{0}^{\alpha}(\Omega)$ it suffices to consider $f \in \mathcal{C}_{0}^{\infty}(\Omega)$. Let us start with $\alpha=2 r+1 / 2, r \in \mathbb{N}_{0}$. We have

$$
\begin{equation*}
\left\langle f, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)}=\sigma_{m}^{-1}\left\langle f, \mathbf{R}^{\#} u_{m, l}\right\rangle_{L^{2}(\Omega)}=\sigma_{m}^{-1}\left\langle\mathbf{R} f, u_{m, l}\right\rangle_{L^{2}\left(\tilde{Z}, w^{-1}\right)} \tag{A.4}
\end{equation*}
$$

Further, $\mathbf{R} f(\cdot, \varphi) \in \mathcal{C}_{0}^{\infty}(-1,1)$ for any $\varphi \in[0,2 \pi]$; see, e.g., Natterer [20]. We estimate the inner product on the right-hand side of (A.4). Let $g_{\varphi}(s)=\mathbf{R} f(s, \varphi)$. Integration by parts yields

$$
\begin{aligned}
\int_{-1}^{1} g_{\varphi}(s) U_{m}(s) d s & =\int_{0}^{\pi} g_{\varphi}(\cos \vartheta) \sin ((m+1) \vartheta) d \vartheta \\
& =\frac{-1}{m+1} \int_{0}^{\pi} g_{\varphi}^{\prime}(\cos \vartheta) \sin \vartheta \cos ((m+1) \vartheta) d \vartheta
\end{aligned}
$$

where we used that $g_{\varphi}(-1)=g_{\varphi}(1)=0$. Repeating integration by parts $2 r+1$-times we obtain

$$
\int_{-1}^{1} g_{\varphi}(s) U_{m}(s) d s=(m+1)^{-2 r-1} \int_{0}^{\pi} \rho_{r}(\vartheta, \varphi) \cos ((m+1) \vartheta) d \vartheta
$$

with $\rho_{r}(\vartheta, \varphi)=\sin \vartheta \sum_{i=1}^{2 r+1}\left(\frac{\partial^{i}}{\partial s^{i}} \mathbf{R} f\right)(\cos \vartheta, \varphi) p_{i-1}(\vartheta)$ where $p_{i-1}$ is a real trigonometric polynomial of degree $i-1$ at most. Thus,

$$
\left\langle g, u_{m, l}\right\rangle_{L^{2}\left(\widetilde{Z}, w^{-1}\right)}=\pi^{-1}(m+1)^{-2 r-1} \underbrace{\int_{0}^{2 \pi} \int_{0}^{\pi} \rho_{r}(\vartheta, \varphi) \cos ((m+1) \vartheta) d \vartheta \mathrm{e}^{-\imath l \varphi} d \varphi}_{=: c_{m, l}(\mathbf{R} f)}
$$

implying

$$
\left|\left\langle f, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)}\right|^{2} \preceq \sigma_{m}^{8 r+2}\left|c_{m, l}(\mathbf{R} f)\right|^{2}=\sigma_{m}^{4 \alpha}\left|c_{m, l}(\mathbf{R} f)\right|^{2}
$$

by (5.5) and (A.4). Summing up results in

$$
\begin{equation*}
\sum_{m=0}^{\infty} \sum_{l=-m}^{m} \sigma_{m}^{-4 \alpha}\left|\left\langle f, \mathrm{v}_{m, l}\right\rangle_{L^{2}(\Omega)}\right|^{2} \preceq \sum_{m=0}^{\infty} \sum_{l=-m}^{m}\left|c_{m, l}(\mathbf{R} f)\right|^{2} \tag{A.5}
\end{equation*}
$$

Since $\left\{\cos (n \vartheta) \mathrm{e}^{\imath l \varphi} / \pi \mid n \in \mathbb{N}, l \in \mathbb{Z}\right\}$ is an orthonormal system in $L^{2}([0, \pi] \times[0,2 \pi])$ we get from the Bessel inequality that

$$
\sum_{m=0}^{\infty} \sum_{l=-m}^{m}\left|c_{m, l}(\mathbf{R} f)\right|^{2} \leq \pi^{-2} \int_{0}^{2 \pi} \int_{0}^{\pi}\left|\rho_{r}(\vartheta, \varphi)\right|^{2} d \vartheta d \varphi
$$

$$
\begin{aligned}
& =\pi^{-2} \int_{0}^{2 \pi} \int_{0}^{\pi}\left|\sin \vartheta \sum_{i=1}^{2 r+1}\left(\frac{\partial^{i}}{\partial s^{i}} \mathbf{R} f\right)(\cos \vartheta, \varphi) p_{i-1}(\vartheta)\right|^{2} d \vartheta d \varphi \\
& \preceq \int_{0}^{2 \pi} \int_{-1}^{1}\left(\sum_{i=1}^{2 r+1}\left|\left(\frac{\partial^{i}}{\partial s^{i}} \mathbf{R} f\right)(s, \varphi)\right|\right)^{2} d s d \varphi \\
& \preceq\|\mathbf{R} f\|_{H^{(\alpha+1 / 2,0)}(\widetilde{Z})}^{2} \preceq\|f\|_{\alpha}^{2}
\end{aligned}
$$

In the last step we used (A.1). The latter estimate together with (A.5) verifies (A.3) for $\alpha=2 r+1 / 2$. For arbitrary $\alpha \geq 0$ one can use arguments from interpolation theory of Sobolev spaces; see, e.g., Lions and Magenes [7, Chap. 5.1].

Acknowledgments. We thank Rainer Dietz for many fruitful discussions. We are further indebted to Alfred K. Louis for his suggestions which helped to improve the presentation of our results. Our paper benefits from improvements suggested by one of the referees.

## REFERENCES

[1] M. Abramowitz and I. Stegun, eds., Handbook of Mathematical Functions, Dover, New York, 1972.
[2] J.-P. Aubin, Applied Functional Analysis, Wiley, Chichester, UK, 1979.
[3] M. E. Davison and F. A. Grünbaum, Tomographic reconstructions with arbitrary directions, Comm. Pure Appl. Math., 34 (1981), pp. 77-120.
[4] R. Dietz, Die Approximative Inverse als Rekonstruktionsmethode in der Röntgen-Computertomographie (Approximate Inverse as Reconstruction Method in Computerized Tomography), Ph.D. thesis, Universität des Saarlandes, Fachbereich Mathematik, 66041 Saarbrücken, Germany, 1999.
[5] F. A. Grünbaum, Reconstruction with arbitrary directions: Dimensions two and three, in Mathematical Aspects of Computerized Tomography, G. T. Hermann and F. Natterer, eds., Springer-Verlag, Heidelberg, 1980, pp. 112-126.
[6] M. Hegland and R. S. Anderssen, A mollification framework for improperly posed problems, Numer. Math., 78 (1998), pp. 549-575.
[7] J. L. Lions and E. Magenes, Non-Homogeneous Boundary Value Problems and Applications, Vol. 1, Springer-Verlag, New York, 1972.
[8] G. G. Lorentz, Approximation of Functions, Holt, Rinehart and Winston, New York, 1966.
[9] A. K. Louis, Orthogonal function series expansions and the null space of the Radon transform, SIAM J. Math. Anal., 15 (1984), pp. 621-633.
[10] A. K. Louis, Incomplete data problems in X-ray computerized tomography I: Singular value decomposition of the limited angle transform, Numer. Math., 48 (1986), pp. 251-262.
[11] A. K. Louis, Inverse und schlecht gestellte Probleme, Studienbücher Mathematik, B. G. Teubner, Stuttgart, Germany, 1989.
[12] A. K. Louis, Approximate inverse for linear and some nonlinear problems, Inverse Problems, 12 (1996), pp. 175-190.
[13] A. K. Louis and P. MaAss, A mollifier method for linear operator equations of the first kind, Inverse Problems, 6 (1990), pp. 427-440.
[14] A. K. Louis and F. Natterer, Mathematical problems in computerized tomography, Proc. IEEE, 71 (1983), pp. 379-389.
[15] A. K. Louis and A. Rieder, Incomplete data problems in X-ray computerized tomography II: Truncated projections and region-of-interest tomography, Numer. Math., 56 (1989), pp. 371-383.
[16] A. K. Louis and T. Schuster, A novel filter design technique in $2 D$ computerized tomography, Inverse Problems, 12 (1996), pp. 685-696.
[17] P. MaASs, The x-ray transform: Singular value decomposition and resolution, Inverse Problems, 3 (1987), pp. 729-741.
[18] P. MaAss, The interior Radon transform, SIAM J. Appl. Math., 52 (1992), pp. 710-724.
[19] D. A. Murio, The Mollification Method and the Numerical Solution of Ill-Posed Problems, Wiley, New York, 1993.
[20] F. Natterer, The Mathematics of Computerized Tomography, Wiley, Chichester, UK, 1986.
[21] E. T. Quinto, Singular value decomposition and inversion methods for the exterior Radon transform and a spherical transform, J. Math. Anal. Appl., 95 (1985), pp. 437-448.
[22] A. Rieder, R. Dietz, and T. Schuster, Approximate Inverse Meets Local Tomography, Preprint 99/6, Institut für Wissenschaftliches Rechnen und Mathematische Modellbildung, Universität Karlsruhe, Karlsruhe, Germany; Math. Methods Appl. Sci, to appear.
[23] L. L. Schumaker, Spline Functions: Basic Theory, Wiley, New York, 1981.
[24] T. Schuster, Schnelle Rekonstruktion von Geschwindigkeitsfeldern und Theorie der Approximativen Inversen (Fast reconstruction of velocity fields and theory of approximate inverse), Ph.D. Thesis, Universität des Saarlandes, Fachbereich Mathematik, Saarbrücken, Germany, 1999.
[25] J. Wloka, Partial Differential Equations, Cambridge University Press, Cambridge, UK, 1987.

