# Generalized Sampling: Extension to Frames and Inverse and Ill-Posed Problems 

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#### Abstract

Generalized sampling is new framework for sampling and reconstruction in infinite-dimensional Hilbert spaces. Given measurements (inner products) of an element with respect to one basis, it allows one to reconstruct in another, arbitrary basis, in a way that is both convergent and numerically stable. However, generalized sampling is thus far only valid for sampling and reconstruction in systems that comprise bases. Thus, in the first part of this paper we extend this framework from bases to frames, and provide fundamental sampling theorems for this more general case. The second part of the paper is concerned with extending the idea of generalized sampling to the solution of inverse and ill-posed problems. In particular, we introduce two generalized sampling frameworks for such problems, based on regularized and non-regularized approaches. We furnish evidence of the usefulness of the proposed theories by providing a number of numerical experiments.


## 1 Introduction

A vital task in applied mathematics and engineering is the recovery of an object - a signal or image, for example - from a collection of its samples (or measurements). This problem, which lies at the heart of modern sampling theory, can be modelled in a Hilbert space $\mathcal{H}$, with the samples of the unknown $f \in \mathcal{H}$ being of the form

$$
\left\langle f, s_{j}\right\rangle, \quad j \in \mathbb{N},
$$

for some system $\left\{s_{j}\right\}_{j \in \mathbb{N}} \subset \mathcal{H}$. In many important cases, one encounters the situation where the system $\left\{s_{j}\right\}_{j \in \mathbb{N}}$ is fixed. Also, one usually only has access to a finite number of measurements, say

$$
\left\langle f, s_{j}\right\rangle, \quad j=1, \ldots, m
$$

Thus, the computational problem is to recover $f$ from a fixed, and finite, collection of its measurements.
Consider, for example, the Magnetic Resonance Imaging (MRI) problem of reconstructing an image from pointwise samples of its Fourier transform. Here $\mathcal{H}=\mathrm{L}_{2}([-1,1])$ (in the 1D setting) and $s_{j}(t)=$ $\mathrm{e}^{\pi i \rho(j) t}$, where $\rho: \mathbb{N} \rightarrow \mathbb{Z}$ is the standard re-indexing function. The particular design of the MRI scanner means that the sampling scheme cannot be easily altered. Moreover, it is infeasible in practice to take too many measurements, meaning that the parameter $m$ may well be small. Of course, reconstructions from Fourier samples are precisely the setting of arguably the most classical result in sampling theory, the celebrated Shannon Sampling Theorem [21, 25]. However, whilst this theorem guarantees a reconstruction of $f$ via its Fourier series, this is often an extremely poor way to proceed. Indeed, one typically requires an intolerably large number of samples $m$ to recover $f$ to any reasonably accuracy [13, 25]. Other effects, such as the Gibbs phenomenon, can also be highly problematic.

This concern aside, note that the MRI problem is fundamentally well posed: it involves the inversion of the Fourier transform. However, in applications, one may also encounter considerably more challenging situations. In particular, it may well be the case that $f$ is the solution of $A f=g$, where $A$ is some bounded linear operator, and, rather than measurements of $f$, we are actually given measurements of $g$. Thus, to recover $f$, one has to invert $A$, making the Shannon Sampling Theorem (for example) no longer directly

[^0]applicable. Another concern is that the operator $A$ may well in practice be ill-posed. This is the case for the problem of X-Ray tomography, where one is tasked with inverting the famously ill-posed Radon transform.

With this in mind, the purpose of this paper is to study novel approaches for these types of reconstruction problems, based on the framework of generalized sampling (GS) introduced by Adcock \& Hansen [3, 2, 4, 5]. Specifically, we provide two generalizations of this framework: in the first part of the paper we augment this framework to systems $\left\{s_{j}\right\}_{j \in \mathbb{N}}$ that comprise frames (as opposed to previously studied case of Riesz bases), and in the second part we present an extension of GS to the solution of inverse and ill-posed problems.

We explain GS in more detail in the forthcoming section. Before doing so, let us first mention one point that is essential to what follows. Given that we only have access to a finite number of measurements of $f$ (or $g$ ), it is quite common to treat the reconstruction problem as finite dimensional, i.e. living in a finite-dimensional vector space. However, as we explain in due course, there are a number of potential pitfalls of ignoring the true infinite dimensionality of the problem. In particular, one may well end up with reconstructions that are neither stable nor convergent (in the sense $m \rightarrow \infty$ ). Fortunately, as we explain, GS allows one to work directly with the infinite-dimensional model, whilst giving stable, convergent numerical methods.

When introducing new sampling techniques (such as those we present in this paper) in the "age of sparsity" it is natural to ask the following question: can such one exploit sparsity in order to subsample? In many cases, the answer turns out to be yes. In [1,19] (see also [6]), a new theory of compressed sensing for infinite-dimensional problems was introduced, based directly on ideas from GS. This development relates directly to the point made above: current compressed sensing theory (and techniques) are based largely on finite-dimensional models, which are not sufficient in all circumstances. Having said this, the focus of this paper will be on classical (i.e. nonsparse) sampling. In Section 5 we discuss how to extend and combine the sampling theorems introduced in this paper with infinite-dimensional compressed sensing techniques.

The outline of the remainder of this paper is as follows. In Section 2 we describe the problems we consider in more detail and recap GS. Section 3 extends GS to the setting of frames. In Section 4 we discussed inverse and ill-posed problems.

## 2 Generalized sampling

### 2.1 The Problems

The purpose of this paper is to study the following recovery problems. We are given a signal $f \in \mathcal{H}$, where $\mathcal{H}$ is an infinite-dimensional separable Hilbert space. We are also given a sampling system $\left\{s_{k}\right\}_{k \in \mathbb{N}} \subset \mathcal{H}$ and a bounded linear operator $A$ on $\mathcal{H}$, where $A$ may be compact. The problems are as follows:
(i) Recover a good approximation to $f$ from the samples $\left\{\left\langle f, s_{k}\right\rangle\right\}_{k=1}^{m}$, for some $m \in \mathbb{N}$.
(ii) Recover a good approximation to $f$ from the samples $\left\{\left\langle g, s_{k}\right\rangle\right\}_{k=1}^{m}$, for some $m \in \mathbb{N}$, where $g=A f$.

We shall also consider the case where the samples are contaminated by noise.
As mentioned, a basic example of Problem (i) is the recovery of $f \in \mathrm{~L}_{2}([-1,1])$ from its Fourier coefficients. In this case, the sampling system consists of the usual complex exponentials. Of course, one could construct the Fourier series approximation $f_{m}=\sum_{k=1}^{m}\left\langle f, s_{k}\right\rangle s_{k}$, knowing that $f_{m} \rightarrow f$ as $m \rightarrow \infty$. Yet, as discussed, such convergence is often in practice intolerably slow.

An alternative in this case is to proceed as follows. Suppose that we know that $f$ has a 'nice' expansion in another system $\left\{w_{j}\right\}_{j \in \mathbb{N}} \subset \mathcal{H}$. In other words, there exists $\alpha_{j} \in \mathbb{C}$ such that

$$
\sum_{j=1}^{n} \alpha_{j} w_{j} \rightarrow f
$$

rapidly as $n \rightarrow \infty$. In the above problem of Fourier sampling it is easy to think that this may well be the case. For example, the $w_{k}$ could be wavelets of some variety, or more exotic objects such as curvelets, shearlets or contourlets. Now, had we access to $\left\{\alpha_{j}\right\}_{j=1}^{n}$ we could have recovered $f$ to high accuracy. However, we only know $\left\{\left\langle f, s_{j}\right\rangle\right\}_{j=1}^{m}$. Nonetheless, given the additional information that $f$ is well represented in $\left\{w_{j}\right\}_{j \in \mathbb{N}}$, we can now ask the following question: how do we obtain the values $\left\{\alpha_{j}\right\}_{j=1}^{n}$ (or some approximation thereof) from the given measurements?

This question is more subtle than it appears. As we explain in the next section, the most straightforward approach is not guaranteed to succeed.

### 2.2 Consistent Reconstructions

Suppose now that $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ are orthonormal vectors in a Hilbert space $\mathcal{H}$, and that $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal basis of $\mathcal{H}$. Let $f \in \mathcal{H}$ be given by

$$
\begin{equation*}
f=\sum_{j=1}^{m} \alpha_{j} w_{j}, \quad \alpha_{j}=\left\langle f, w_{j}\right\rangle \tag{2.1}
\end{equation*}
$$

where $\langle\cdot, \cdot\rangle$ is the inner product on $\mathcal{H}$. Suppose that we can access samples $\left\langle f, s_{k}\right\rangle$ for $k \in \mathbb{N}$. It is then tempting to form

$$
U_{m}=\left(\begin{array}{ccc}
u_{11} & \ldots & u_{1 m} \\
\vdots & \ddots & \vdots \\
u_{m 1} & \ldots & u_{m m}
\end{array}\right), \quad u_{i j}=\left\langle w_{j}, s_{i}\right\rangle
$$

and approach the problem of obtaining the $\alpha_{j} \mathrm{~s}$ by solving the linear system of equations

$$
\begin{equation*}
U_{m} x=y_{m}, \quad y_{m}=\left\{\left\langle f, s_{1}\right\rangle, \ldots,\left\langle f, s_{m}\right\rangle\right\} \tag{2.2}
\end{equation*}
$$

This approach is often referred to as a consistent reconstruction technique (i.e. the first $m$ measurements of the reconstruction $\tilde{f}=\sum_{k=1}^{m} x_{k} w_{k}$ coincide with those of $f$ ). Introduced (in the context of sampling) by Unser \& Aldroubi [26, 27, 7], and later generalized significantly by Eldar et al [10, 12, 11, 14], this technique is quite widely used in applications [25]. Note in the case above that $f$ has only $m$ nonzero coefficients in the basis $\left\{w_{j}\right\}_{j \in \mathbb{N}}$. Thus, ignoring any noise or numerical issues, the solution $x$ of (2.2) will coincide with $\left\{\alpha_{1}, \ldots, \alpha_{m}\right\}$, and thus $f$ should, at least theoretically, be recovered perfectly by (2.2).

Let us now consider an example of this approach:
Example 2.1. Suppose that the $w_{k} \mathrm{~s}$ are the Haar wavelets on $[-1 / 2,1 / 2]$, and take $m=769$. Let $\left\langle f, s_{k}\right\rangle$ be Fourier samples of $f$, i.e. $s_{k}=e^{2 \pi \mathrm{i} \rho(k)}$, where $\rho: \mathbb{N} \rightarrow \mathbb{N}$ is defined by $\rho(1)=0, \rho(2)=1, \rho(3)=$ $-1, \rho(4)=2, \rho(5)=-2 \ldots$ ( $\rho$ is merely a re-indexing function).

If we form the matrix $U_{m}$, then, although $U_{m}$ is invertible, and in particular the problem (2.2) has a unique solution, one finds that its condition number is intolerably large (in this instance it is $\approx 4.62 \times 10^{16}$ ). Thus, if the observations $y_{m}$ are perturbed by $\delta_{m} \in \mathbb{C}^{m}$, and we solve $U_{m} x=y_{m}+\delta_{m}$, then we expect a catastrophic failure. Indeed, this is the case even in the simplest setting. For $f=\sum_{j=1}^{m} \alpha_{j} w_{j}$, where $\alpha_{1}=\alpha_{2}=\alpha_{6}=\alpha_{7}=1, \alpha_{769}=-1 / 20$ and $\alpha_{j}=0$ otherwise we get that when $\|\delta\|=3.85 \times 10^{-13}$, the reconstruction $\tilde{f}=\sum_{j=1}^{m} x_{j} w_{j}$ satisfies $\|f-\tilde{f}\|_{\mathrm{L}_{2}}=146.85$. The disastrous error $f-\tilde{f}$ as well as $f$ is visualized in Figure 1. Note that, even though the perturbation is small, the exceedingly large condition number leads to a complete failure in the reconstruction.

One may think that a remedy could be to increase $m$. In particular, maybe taking more Fourier samples would cure the problem. Let us increase $m$ to 1351 , but keeping $f$ as before as well as the noise level. Unfortunately, this also gives a completely unsatisfactory result (even worse than before) with $\|f-\tilde{f}\|_{\mathrm{L}_{2}}=$ 360.12 as visualized in Figure 1.

There is, however, a remedy to the problem. In particular, by defining $\tilde{f}_{n, m}$ as in (3.1), with parameters $m=1351$ and $n=769$, we obtain a recovery error that is very close to the noise level. This approach was first introduced in $[4,5]$, and we will in this paper extend this framework even further.

This failure is somewhat surprising (for similar examples, see [2, 4]). Since $U_{m}$ is nothing more than the change-of-basis matrix between two systems of $m$ orthogonal functions, it is tempting to think that $U_{m}$ is unitary (and therefore well-conditioned). However, herein lies a critical misunderstanding. There is no reason why $U_{m}$ should be unitary (or even invertible). The trained eye of an operator theorist will immediately spot the mistake: $U_{m}$ is unitary if and only if $\operatorname{span}\left\{w_{k}\right\}_{k=1}^{m}=\operatorname{span}\left\{s_{k}\right\}_{k=1}^{m}$. The problem is that this may very well not be the case since $\mathcal{H}$ is infinite-dimensional. Indeed, it is straightforward to see that the two bases in Example 2.1 do not lead to a unitary matrix: the first $m$ Haar wavelets span a space of piecewise constant functions, whereas the standard complex exponentials span a space of (smooth) trigonometric polynomials.


Figure 1: The figure displays $f$ (upper left) as well as $f-\tilde{f}$ for two different values of $m$ (upper right and lower left), and $f-\tilde{f}_{n, m}$ (lower right).

The approach considered above can be loosely described as 'finite-dimensional': one simply interpolates between two $m$-dimensional spaces. Herein lies the problem: the recovery of $f$, an element of the infinite-dimensional Hilbert space, is inherently infinite-dimensional. In particular, a typical $f \in \mathcal{H}$ will have an infinite expansion in $\sum_{k=1}^{\infty} \alpha_{k} w_{k}$ in the basis $\left\{w_{k}\right\}_{k \in \mathbb{N}}$, and a countably infinite collection of measurements $\left\{\left\langle f, s_{k}\right\rangle\right\}_{k \in \mathbb{N}}$. Thus, to tackle the reconstruction problem effectively, i.e. to avoid the pitfalls encountered above, it is necessary to pursue an infinite-dimensional approach.

### 2.3 The new approach

To compensate for the problems that consistent reconstructions cause, a new type of theory, and associated technique, was introduced, known as generalized sampling (GS) [2, 4, 5]. This theory generalizes the Shannon Sampling Theorem [21], as well as some fundamental work by Hrycak \& Gröchenig [20].

The main message of GS is the following: one can obtain samples via inner products with respect to one basis (e.g. Fourier coefficients) and reconstruct in another basis of choice, and this can always be done in a completely stable and convergent way. Hence, one can combat the problem of potentially slow convergence in the sampling basis by simply reconstructing in a basis more appropriate for the signal to be reconstructed. Stability also implies complete robustness with respect to noise.

### 2.4 Overview of the paper

GS was introduced in a series of papers [2, 3, 4, 5]. However, there are several important problems that have not yet been addressed. These form the content of this paper:

- Frames: GS is so far only valid for sampling and reconstruction systems $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ that comprise Riesz bases. Given the vast literature on frames and the wide variety of applications it is crucial that the GS is extended to frames. This is content of the first part of the paper.
- Ill-posed problems: GS is concerned with the problem of reconstructing an element $f$ from samples $\left\langle f, s_{j}\right\rangle$. However, what if we cannot sample $f$, but rather $g:=A f$, where $A$ is a compact operator. The topic of ill-posed problems is vast and of great importance. Therefore it is vital to have an extension of GS that includes such problems. In the second part of this paper, we will introduce new techniques with GS using regularization, as well as a non-regularized technique.

Let us remark in passing that the ideas leading to GS originate computational spectral theory [17, 18]. In particular, the technique of uneven sections of infinite matrices. Spectral theory has been used in sampling before (although in a very different context than what will be presented here) to great success, with the fundamental paper [22] of Landau being an important example. The problems above will also addressed through extensions and refinements of the uneven section technique (see also [8] and [16] for other results on even and uneven section techniques).

### 2.5 Background and Notation

We will let $\mathcal{H}$ denote a separable Hilbert space. Let $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ be a frame for $\mathcal{W}=\overline{\operatorname{span}}\left\{w_{k}\right\}_{k \in \mathbb{N}} \subset \mathcal{H}$ and $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ be a frame for $\mathcal{S}=\overline{\operatorname{span}}\left\{s_{k}\right\}_{k \in \mathbb{N}} \subset \mathcal{H}$. In particular, this means that for all $f \in \mathcal{W}$ there exist constants $0<A \leq B<\infty$ such that

$$
\begin{equation*}
A\|f\|^{2} \leq \sum_{k \in \mathbb{N}}\left|\left\langle f, w_{k}\right\rangle\right|^{2} \leq B\|f\|^{2} \tag{2.3}
\end{equation*}
$$

and that for all $f \in \mathcal{S}$ there exist constants $0<C \leq D<\infty$ such that

$$
\begin{equation*}
C\|g\|^{2} \leq \sum_{k \in \mathbb{N}}\left|\left\langle g, s_{k}\right\rangle\right|^{2} \leq D\|g\|^{2} . \tag{2.4}
\end{equation*}
$$

Recall the so-called synthesis operators $S, W: l^{2}(\mathbb{N}) \rightarrow \mathcal{H}$ are defined by

$$
S x=x_{1} s_{1}+x_{2} s_{2}+\ldots, \quad W y=y_{1} w_{1}+y_{2} w_{2}+\ldots,
$$

and their adjoints (the analysis operators) $S^{*}, W^{*}: \mathcal{H} \rightarrow l^{2}(\mathbb{N})$ are easily seen to be

$$
S^{*} g=\left\{\left\langle g, s_{1}\right\rangle,\left\langle g, s_{2}\right\rangle, \ldots\right\}, \quad W^{*} h=\left\{\left\langle h, w_{1}\right\rangle,\left\langle h, w_{2}\right\rangle \ldots\right\} .
$$

Define also the operators $S_{m}, W_{m}: \mathbb{C}^{m} \rightarrow \mathcal{H}$ by

$$
\begin{equation*}
S_{m} x=x_{1} s_{1}+\ldots+x_{m} s_{m}, \quad W_{m} y=y_{1} w_{1}+\ldots+y_{m} w_{m} \tag{2.5}
\end{equation*}
$$

with adjoints $S^{*}, W^{*}: \mathcal{H} \rightarrow \mathbb{C}^{m}$ given by

$$
S_{m}^{*} g=\left\{\left\langle g, s_{1}\right\rangle, \ldots,\left\langle g, s_{m}\right\rangle\right\}, \quad W_{m}^{*} h=\left\{\left\langle h, w_{1}\right\rangle, \ldots,\left\langle h, w_{m}\right\rangle\right\} .
$$

Define also the corresponding spaces

$$
\mathcal{S}_{m}=\operatorname{span}\left\{s_{1}, \ldots, s_{m}\right\}, \quad \mathcal{W}_{n}=\operatorname{span}\left\{w_{1}, \ldots, w_{n}\right\}
$$

For any closed subspace $\mathcal{V} \subset \mathcal{H}$ we denote the orthogonal projection onto $\mathcal{V}$ by $P_{\mathcal{V}}$. We shall also let $\left\{e_{j}\right\}_{j \in \mathbb{N}}$ denote the canonical basis of $l^{2}(\mathbb{N})$, and, for any $k \in \mathbb{N}$, write $P_{k}$ for the orthogonal projection onto $\operatorname{span}\left\{e_{1}, \ldots, e_{k}\right\}$. Lastly, if $A: \mathcal{H} \rightarrow \mathcal{H}$ is bounded, where $\mathcal{H}$ is a Hilbert space and $\operatorname{Ran}(A)$ is closed, we denote its pseudo-inverse by $A^{\dagger}$.

## 3 Generalized Sampling with Frames

Let $f \in \mathcal{H}$ and $\left\{w_{k}\right\}_{k \in \mathbb{N}},\left\{s_{k}\right\}_{k \in \mathbb{N}}$ be as in the previous section. Suppose also that we can, for some $m \in \mathbb{N}$, access the samples $S_{m}^{*} f$. The task is to reconstruct an approximation to $f$ from these samples. As mentioned, if $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ was an orthonormal basis (and $f \in \mathcal{S}$ ) we could have formed the approximation $f_{m}=\sum_{k=1}^{m}\left\langle f, s_{k}\right\rangle s_{k}$. However, it may well be the case that $f_{m} \rightarrow f$ intolerably slowly and therefore $f_{m}$ would be a rather poor approximation. Moreover, if $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ is merely a frame, then $f_{m}$ in general will not converge to $f$, and hence will be completely useless.

Suppose now that $P_{\mathcal{W}_{n}} f \rightarrow f$ rapidly - in particular, much faster than $f_{m} \rightarrow f$ (if valid). Given this information we are lead to the following question: can we reconstruct $P_{\mathcal{W}_{n}} f$ (or at least an approximation) from the samples $S_{m}^{*} f$ ?

### 3.1 Reconstruction

Suppose now that we are given $f \in \mathcal{W}$ and that we can access the samples $S_{m}^{*} f$. Based on this information we will construct the GS approximation $\tilde{f}_{n, m} \in \mathcal{W}_{n}$ as follows:

$$
\begin{equation*}
\tilde{f}_{n, m}=W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} f \tag{3.1}
\end{equation*}
$$

The following theorem concerns the behaviour of $\tilde{f}_{n, m}$ :
Theorem 3.1. Let $\mathcal{H}$ be a separable Hilbert space and $\mathcal{S}, \mathcal{W} \subset \mathcal{H}$ be closed subspaces such that $\mathcal{W} \cap \mathcal{S}^{\perp}=$ $\{0\}$. Suppose that $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ are frames for $\mathcal{S}$ and $\mathcal{W}$ respectively with frame bounds as in (2.3) and (2.4) respectively. Then, for each $n \in \mathbb{N}$ and any $f \in \mathcal{W}$, there is an $M \in \mathbb{N}$ such that, for all $m \geq M$, the reconstruction $\tilde{f}_{n, m}$ (defined in (3.1)) exists and is unique. Moreover,

$$
\left\|P_{\mathcal{W}_{n}}^{\perp} f\right\|_{\mathcal{H}} \leq\left\|f-\tilde{f}_{n, m}\right\|_{\mathcal{H}} \leq\left(1+K_{n, m}\right)\left\|P_{\mathcal{W}_{n}}^{\perp} f\right\|_{\mathcal{H}},
$$

where $P_{\mathcal{W}_{n}}$ is the orthogonal projection onto $\mathcal{W}_{n}, P_{\mathcal{W}_{n}}^{\perp}=I-P_{\mathcal{W}_{n}}$ and

$$
\begin{equation*}
K_{n, m}=\left\|W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} P_{\mathcal{W}_{n}}^{\perp}\right\| \tag{3.2}
\end{equation*}
$$

For each fixed $n$ we have that $K_{n, m} \rightarrow K_{n}$ as $m \rightarrow \infty$, where

$$
K_{n}=\left\|W_{n}\left(W_{n}^{*} S S^{*} W_{n}\right)^{\dagger} W_{n}^{*} S S^{*} P_{\mathcal{W}_{n}}^{\perp}\right\| \leq \frac{D}{C}
$$

Moreover, when $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal basis and $\mathcal{W} \subset \mathcal{S}$ then, for fixed $n, K_{n, m} \rightarrow 0$ as $m \rightarrow \infty$.
Note that this theorem generalizes the results in $[4,5]$ to frames (we give a proof in the next section).

### 3.2 Stability

The stability of GS is very important to analyze. One may at first glance suspect that this method could be unstable, due to the fact that $\left\|\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger}\right\|$ could blow up for large $n$, even when $m$ is arbitrarily large. Indeed, this could very well happen as suggested in the next proposition.

Proposition 3.2. There exists a Hilbert space $\mathcal{H}$ with frames $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ such that

$$
\lim _{m \rightarrow \infty}\left\|\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger}\right\| \longrightarrow \infty, \quad n \rightarrow \infty
$$

Proof. Let $\mathcal{H}=l^{2}(\mathbb{N})$ and write $\left\{e_{k}\right\}_{k \in \mathbb{N}}$ for its canonical basis. Let $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ be defined by $w_{1}=e_{1}$ and $w_{k}=e_{k-1}+\frac{e_{k}}{k}, k \geq 2$, and let $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ be given by $s_{k}=e_{k}$. Then, as shown in [9], $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ is a frame with frame bounds $1-\sqrt{\pi^{2} / 6-1}$ and 3 respectively. Note that

$$
W_{n} y=y_{1} e_{1}+\sum_{k=2}^{n} y_{k}\left(\frac{e_{k}}{k}+e_{k-1}\right), \quad y \in \mathbb{C}^{n}
$$

and therefore

$$
S_{m}^{*} W_{n}=\left(\begin{array}{cccccc}
1 & 1 & 0 & 0 & \ldots & 0 \\
0 & \frac{1}{2} & 1 & 0 & \ddots & 0 \\
0 & 0 & \frac{1}{3} & \ddots & \ddots & \vdots \\
0 & 0 & 0 & \ddots & \ddots & \vdots \\
0 & 0 & 0 & \ddots & \ddots & 1 \\
0 & 0 & 0 & \ldots & \ldots & \frac{1}{n} \\
0 & 0 & 0 & \ldots & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & 0 & \ldots & \ldots & 0
\end{array}\right) \in \mathbb{C}^{m \times n}, \quad m \geq n
$$

In particular, we have that $\left\|\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger}\right\| \geq n^{2}$, yielding the assertion.

Despite this rather pessimistic result, the method is perfectly stable. Indeed, although $\left\|\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger}\right\|$ may blow up, it is always the case that $\left\|W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*}\right\|$ stays bounded provided $m$ is sufficiently large. Specifically, we have the following theorem.

Theorem 3.3. Let $\mathcal{H}, \mathcal{S}, \mathcal{W},\left\{s_{k}\right\}_{k \in \mathbb{N}},\left\{w_{k}\right\}_{k \in \mathbb{N}}, n$ and $m$ be as in Theorem 3.1. Suppose that the samples of $f \in \mathcal{H}$ are contaminated by noise. In particular, suppose that we are given $S_{m}^{*}(f+h)$ for some $h \in \mathcal{H}$. If

$$
\tilde{f}_{n, m}^{h}=W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*}(f+h)
$$

then

$$
\left\|f-\tilde{f}_{n, m}^{h}\right\|_{\mathcal{H}} \leq\left(1+K_{n, m}\right)\left(\left\|P_{\mathcal{W}_{n}}^{\perp} f\right\|_{\mathcal{H}}+\|h\|_{\mathcal{H}}\right)
$$

for all $m \geq M$ where $M$ is as in Theorem 3.1.

### 3.3 Proofs

Proof of Theorem 3.1. Let us first observe that due to the fact that $\underset{\sim}{\operatorname{Ran}}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)$ is finite dimensional, and hence closed, then the pseudo-inverse is unique and hence $\tilde{f}_{n, m}$ exists and is unique. Moreover, we observe that

$$
W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*} W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*}=W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} P_{\operatorname{Ran}\left(\mathrm{W}_{\mathrm{n}}^{*} \mathrm{~W}_{\mathrm{n}}\right)} W_{n}^{*}=W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*}
$$

In particular, by self-adjointness, it follows that

$$
\begin{equation*}
P_{\mathcal{W}_{n}}=W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*} . \tag{3.3}
\end{equation*}
$$

Note that

$$
\begin{align*}
W_{n}^{*} S_{m} S_{m}^{*} f & =W_{n}^{*} S_{m} S_{m}^{*}\left(P_{\mathcal{W}_{n}} f+P_{\mathcal{W}_{n}}^{\perp} f\right)  \tag{3.4}\\
& =W_{n}^{*} S_{m} S_{m}^{*} W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*} f+W_{n}^{*} S_{m} S_{m}^{*} P_{\mathcal{W}_{n}}^{\perp} f
\end{align*}
$$

Equation (3.4) now gives

$$
\begin{align*}
W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} f= & W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*} f  \tag{3.5}\\
& +W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} P_{\mathcal{W}_{n}}^{\perp} f
\end{align*}
$$

We now claim that there exists an $M$ such that

$$
W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*} f=P_{\mathcal{W}_{n}} f, \quad \forall m \geq M
$$

To see this, we first observe that

$$
W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} W_{n}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*} f=W_{n} P_{\operatorname{Ran}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)}\left(W_{n}^{*} W_{n}\right)^{\dagger} W_{n}^{*} f
$$

and it therefore suffices to show that

$$
\begin{equation*}
\operatorname{Ran}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)=\operatorname{Ran}\left(W_{n}^{*} W_{n}\right) \tag{3.6}
\end{equation*}
$$

for all sufficiently large $m$. By assumption we have $\mathcal{W} \cap \mathcal{S}^{\perp}=\{0\}$. Thus, for all $n \in \mathbb{N}$ we deduce that $\mathcal{W}_{n} \cap \mathcal{S}^{\perp}=\{0\}$. It follows that

$$
\inf \left\{\left\|S^{*} x\right\|: x \in \operatorname{Ran}\left(W_{n}\right),\|x\|=1\right\}>0
$$

In particular, we have that $\operatorname{Rank}\left(S^{*} W_{n}\right)=\operatorname{Rank}\left(W_{n}^{*} W_{n}\right)$, and therefore $\operatorname{Rank}\left(W_{n}^{*} S S^{*} W_{n}\right)=\operatorname{Rank}\left(W_{n}^{*} W_{n}\right)$, which yields that $\operatorname{Ran}\left(W_{n}^{*} S S^{*} W_{n}\right)=\operatorname{Ran}\left(W_{n}^{*} W_{n}\right)$. Note that $S_{m} S_{m}^{*} \rightarrow S S^{*}$ strongly, however, by compactness of $W_{n}$ it follows that $W_{n}^{*} S_{m} S_{m}^{*} W_{n} \rightarrow W_{n}^{*} S S^{*} W_{n}$ in norm as $m \rightarrow \infty$. Hence, we deduce (3.6) for all sufficiently large $m$. Thus,

$$
\left\|f-\tilde{f}_{n, m}\right\|_{\mathcal{H}} \leq\left\|P_{\mathcal{W}_{n}}^{\perp}-W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} P_{\mathcal{W}_{n}}^{\perp}\right\|_{\mathcal{H}}\left\|P_{\mathcal{W}_{n}}^{\perp} f\right\|_{\mathcal{H}}
$$

which gives the first part of the theorem.

To prove the second assertion note that

$$
W_{n}\left(W_{n}^{*} S S^{*} W_{n}\right)^{\dagger} W_{n}^{*} S S^{*} P_{\mathcal{W}_{n}}^{\perp}=\left(S S^{*}\right)^{-1} S S^{*} W_{n}\left(W_{n}^{*} S S^{*} W_{n}\right)^{\dagger} W_{n}^{*} S S^{*} P_{\mathcal{W}_{n}}^{\perp}
$$

and therefore, since $S^{*} W_{n}\left(W_{n}^{*} S S^{*} W_{n}\right)^{\dagger} W_{n}^{*} S$ is the orthogonal projection $P_{\operatorname{Ran}\left(S^{*} W_{n}\right)}$, it follows that

$$
\begin{aligned}
\left\|W_{n}\left(W_{n}^{*} S S^{*} W_{n}\right)^{\dagger} W_{n}^{*} S S^{*} P_{\mathcal{W}_{n}}^{\perp}\right\| & \leq\left\|\left(S S^{*}\right)^{-1}\right\|_{\mathcal{S} \rightarrow \mathcal{S}}\|S\|_{l^{2}(\mathbb{N}) \rightarrow \mathcal{S}}\left\|S^{*}\right\|_{\mathcal{S} \rightarrow l^{2}(\mathbb{N})} \\
& =\left\|\left(S S^{*}\right)^{-1}\right\|_{\mathcal{S} \rightarrow \mathcal{S}}\left\|S S^{*}\right\|_{\mathcal{S} \rightarrow \mathcal{S}} \leq \frac{D}{C}
\end{aligned}
$$

Also, by a compactness argument, since $S_{m} S_{m}^{*} \rightarrow S S^{*}$ strongly, it follows that

$$
\begin{equation*}
W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} P_{\mathcal{W}_{n}}^{\perp} \longrightarrow W_{n}\left(W_{n}^{*} S S^{*} W_{n}\right)^{\dagger} W_{n}^{*} S S^{*} P_{\mathcal{W}_{n}}^{\perp}, \quad m \rightarrow \infty \tag{3.7}
\end{equation*}
$$

Therefore $K_{n, m} \rightarrow\left\|W_{n}\left(W_{n}^{*} S S^{*} W_{n}\right)^{\dagger} W_{n}^{*} S S^{*} P_{\mathcal{W}_{n}}^{\perp}\right\|$ as $m \rightarrow \infty$ and we have proved the second assertion.

To prove that $K_{n, m} \rightarrow 0$ as $m \rightarrow \infty$ when $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal basis, we start by making the following observation: since $\mathcal{W} \subset \mathcal{S}$ we certainly have that $\mathcal{W}_{n} \subset \mathcal{S}$. Therefore $P_{\mathcal{S}} P_{\mathcal{W}_{n}}=P_{\mathcal{W}_{n}}=$ $P_{\mathcal{W}_{n}} P_{\mathcal{S}}$. Also, $W_{n}^{*}=W_{n}^{*} P_{\mathcal{W}_{n}}$, and therefore

$$
W_{n}^{*} P_{\mathcal{S}} P_{\mathcal{W}_{n}}^{\perp}=W_{n}^{*} P_{\mathcal{S}}\left(I-P_{\mathcal{W}_{n}}\right)=0 .
$$

We therefore immediately obtain

$$
W_{n}\left(W_{n}^{*} S S^{*} W_{n}\right)^{\dagger} W_{n}^{*} S S^{*} P_{\mathcal{W}_{n}}^{\perp}=W_{n}\left(W_{n}^{*} P_{\mathcal{S}} W_{n}\right)^{\dagger} W_{n}^{*} P_{\mathcal{S}} P_{\mathcal{W}_{n}}^{\perp}=0
$$

So, by (3.7) the assertion follows, and we are done.
Proof of Theorem 3.3. Note that (3.5) and the reasoning in the proof of Theorem 3.3 give

$$
f-\tilde{f}_{n, m}^{h}=P_{\mathcal{W}_{n}}^{\perp} f-W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} P_{\mathcal{W}_{n}}^{\perp}(f+h)-h,
$$

which yields the assertion.

### 3.4 Determining $m$ : The Stable Sampling Rate

For Theorems 3.1 and 3.3 we require $m \geq M$ samples, where $M$ depends only on $n$ and the sampling and reconstruction systems. This leads to the question: how large must $m$ be? This was first explored in [4, 5] in the context of Riesz bases. Continuing in the same way, let us define the stable sampling rate

$$
\begin{equation*}
\Phi(n, \theta)=\min \left\{m \in \mathbb{N}: K_{n, m} \leq \theta\right\}, \quad(n, \theta) \in \mathbb{N} \times \mathbb{R}_{+}, \tag{3.8}
\end{equation*}
$$

where $K_{n, m}$ is as in (3.2). This quantity tells us how many samples one has to take, as a function of the number of coefficients to be computed, in order to secure a stable and convergent solution. It is of course very important to get analytical bounds for $\Phi$. However, these may be very difficult to derive for arbitrary sampling and reconstruction systems. Fortunately, as we now explain, we can always obtain numerical estimates for $\Phi$.

When confronted with the problem of computing $\Phi$ there are two obvious obstacles:
(i) The operator $P_{\mathcal{W}_{n}}^{\perp}$ does not have finite rank (and this suggests therefore the need for an infinite amount of information to compute $K_{n, m}$ ).
(ii) The operator $W_{n}\left(W_{n}^{*} S_{m} S_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} S_{m} S_{m}^{*} P_{\mathcal{W}_{n}}^{\perp}$ acts on the abstract Hilbert space $\mathcal{H}$. In it therefore unclear how to obtain numerical norm estimates when the information we will (at best) be able to access are the numerical values of the inner products $\left\langle s_{i}, w_{j}\right\rangle$ and $\left\langle w_{i}, w_{j}\right\rangle$ for $i, j \in \mathbb{N}$.


Figure 2: The left figure shows the error $f-f_{m}$ (the truncated Fourier series) for $m=501$. The right figure shows the error $f-\tilde{f}_{n, m}$ for $m=501, n=45$. Note that both $f_{m}$ and $\tilde{f}_{n, m}$ use exactly the same samples.

Fortunately, one can at least obtain a computable upper bound for $K_{n, m}$ (and therefore $\Phi$ ). Indeed, if $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal basis for $\mathcal{S}, \mathcal{W} \subset \mathcal{S}$ and $U$ and $V$ are the infinite matrices

$$
U=\left(\begin{array}{cccc}
u_{11} & u_{12} & u_{13} & \ldots  \tag{3.9}\\
u_{21} & u_{22} & u_{23} & \ldots \\
u_{31} & u_{32} & u_{33} & \ldots \\
\vdots & \vdots & \vdots & \ddots
\end{array}\right), \quad V=\left(\begin{array}{cccc}
v_{11} & v_{12} & v_{13} & \ldots \\
v_{21} & v_{22} & v_{23} & \ldots \\
v_{31} & v_{32} & v_{33} & \ldots \\
\vdots & \vdots & \vdots & \ddots
\end{array}\right) \quad u_{i j}=\left\langle s_{i}, w_{j}\right\rangle, v_{i j}=\left\langle w_{i}, w_{j}\right\rangle
$$

then it is a straightforward exercise (along the same lines as arguments given in [4]) to show that $K_{n, m} \leq$ $\tilde{K}_{n, m}$, where

$$
\begin{aligned}
\tilde{K}_{n, m}= & \sqrt{B}\left\|\left(P_{n} U^{*} P_{m} U P_{n}\right)^{\dagger} P_{n} U^{*} P_{m}-\left(P_{n} V P_{n}\right)^{\dagger} P_{n} U^{*} P_{m}\right\| \\
& +\left\|\left(P_{n} V P_{n}\right)^{\dagger}\left(P_{n} U^{*} P_{m} U P_{n}-P_{n} V P_{n}\right)\left(P_{n} V P_{n}\right)^{\dagger}\right\|^{1 / 2} .
\end{aligned}
$$

Thus, we have

$$
\begin{equation*}
\Phi(n, \theta) \leq \Psi(n, \theta)=\min \left\{m \in \mathbb{N}: \tilde{K}_{n, m} \leq \theta\right\}, \quad \forall(n, \theta) \in \mathbb{N} \times \mathbb{R}_{+} \tag{3.10}
\end{equation*}
$$

where the right-hand side can be computed numerically, since $\tilde{K}_{n, m}$ involves only sections of the matrices $U$ and $V$. Thus, in practice, one can always ensure, via a numerical computation, the bounds of Theorems 3.1 and 3.3, and therefore convergence and stability. In the next section we give an example of the behaviour of $\Psi(n, \theta)$ for a typical problem (further examples in the case of bases are given in $[4,5]$ ).

### 3.5 Numerical Examples

In this section we will test the framework suggested in the previous section with the following frames: let $\left\{p_{k}\right\}_{k \in \mathbb{N}}$ denote the Legendre polynomials on $[-1,1]$ and $\left\{h_{k}\right\}_{k \in \mathbb{N}}$ the Haar wavelets on $[0,1]$. Define the frames

$$
\begin{equation*}
\left\{w_{k}\right\}_{k \in \mathbb{N}}, \quad w_{2 k-1}=p_{k}, \quad w_{2 k}=h_{k}, \quad\left\{\tilde{w}_{k}\right\}_{k \in \mathbb{N}}, \quad \tilde{w}_{k}=p_{k} \tag{3.11}
\end{equation*}
$$

Note that $\left\{\tilde{w}_{k}\right\}_{k \in \mathbb{N}}$ is actually an orthonormal basis. When it comes to indexing the frame we simply make the most natural choice. However, one should note that one has complete freedom to choose the indexing, although the synthesis and analysis operators will obviously depend on such a choice. We will denote the operators and spaces associated with this basis (as defined in Section 3) by simply adding a tilde, for example $\widetilde{W}_{n}, \widetilde{W}_{n}^{*}, \widetilde{W}$ etc. Also, define the orthonormal basis for $\mathrm{L}_{2}([-1,1])$

$$
\left\{s_{k}\right\}_{k \in \mathbb{Z}}, \quad s_{k}(t)=\frac{1}{\sqrt{2}} \mathrm{e}^{\pi i k t}
$$

Thus, our samples of any function $f \in \mathrm{~L}_{2}([-1,1])$ will be its Fourier coefficients. In particular,

$$
\hat{f}(k)=\left\langle f, s_{k}\right\rangle=\frac{1}{\sqrt{2}} \int_{-1}^{1} f(t) \mathrm{e}^{\pi i k t} \mathrm{~d} t
$$



Figure 3: The figure shows the errors $f_{d}-f_{d, m}$ (left), $f_{d}-\tilde{f}_{d, n, m}$ (middle) and $f_{d}-f_{d, k, m}$ (right) for $m=501, n=45, k=150$. Note that $f_{d, m}, \tilde{f}_{d, n, m}$ and $f_{d, k, m}$ use exactly the same samples.

Thus, in this case we have $\mathcal{W}=\widetilde{\mathcal{W}}=\mathcal{S}=\mathrm{L}_{2}([-1,1])$. In the previous framework we have indexed the operators over the natural numbers, so to make it compatible with the setup in this section we will let (for odd $m$ ) $S_{m}: \mathbb{C}^{m} \rightarrow \mathrm{~L}_{2}([-1,1])$ and $S_{m}^{*}: \mathrm{L}_{2}([-1,1]) \rightarrow \mathbb{C}^{m}$ be defined by

$$
S_{m} x=x_{1} s_{-(m-1) / 2}+\ldots+x_{m} s_{(m-1) / 2}, \quad S_{m}^{*} g=\left\{\left\langle g, s_{-(m-1) / 2}\right\rangle, \ldots,\left\langle g, s_{(m-1) / 2}\right\rangle\right\}
$$

Also, for $f \in \mathrm{~L}_{2}([-1,1])$ and for odd $m \in \mathbb{N}$ we let $f_{m}$ denote the truncated Fourier series of $f$. In particular,

$$
f_{m}(t)=\left(S_{m} S_{m}^{*} f\right)(t)=\frac{1}{2} \sum_{j=-(m-1) / 2}^{(m-1) / 2} \hat{f}(j) \mathrm{e}^{\pi i j t}
$$

Example 3.4. In order to illustrate the theory developed in the previous sections we have tested the functions

$$
f(t)=t^{3} e^{t}, \quad f_{d}(t)=t^{3} e^{t}\left(1-\chi_{[0.5,0.75]}\right)(t)
$$

Note that the analyticity of $f$ strongly favors reconstruction in a polynomial basis. In particular, the fact that $f$ is non-periodic prevents rapid convergence of the Fourier series as shown in Figure 2, where we have displayed the error $f-f_{m}$ for $m=501$. As an alternative we have defined the GS reconstruction

$$
\tilde{f}_{n, m}=\widetilde{W}_{n}\left(\widetilde{W}_{n}^{*} S_{m} S_{m}^{*} \widetilde{W}_{n}\right)^{\dagger} \widetilde{W}_{n}^{*} S_{m} S_{m}^{*} f
$$

where $m=501$ and $n=45$. By using exactly the same samples as for $f_{m}$, namely $S_{m}^{*} f$, we simply reconstruct in a different basis (the Legendre basis), and dramatically reduce the error. This is illustrated in Figure 2. For more examples with reconstructions in different bases see [3, 2, 4, 5]. However, this paper is primarily about the use of frames. In particular, frames that can be used to capture specific characteristics of the function to be recovered.

To illustrate the effectiveness of the use of frames, consider the following. What happens if we introduce a discontinuity in $f$ ? In particular, what happens if we replace $f$ by the discontinuous function $f_{d}$ ? This does have an effect on the truncated Fourier series (as documented in Figure 3) $f_{d, m}=S_{m} S_{m}^{*} f_{d}$, but not near the striking impact it has on

$$
\tilde{f}_{d, n, m}=\widetilde{W}_{n}\left(\widetilde{W}_{n}^{*} S_{m} S_{m}^{*} \widetilde{W}_{n}\right)^{\dagger} \widetilde{W}_{n}^{*} S_{m} S_{m}^{*} f_{d}
$$

which uses the Legendre basis in the same way as for $\tilde{f}_{n, m}$. The quite dramatic effect of the discontinuity is visualized in Figure 3. The question is then: what to do? One idea that comes to mind is that one could try to mix two bases, one that would favor smoothness, and one that would favor discontinuities. This is exactly what the frame $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ (defined in (3.11)) would do. By introducing

$$
f_{d, k, m}=W_{k}\left(W_{k}^{*} S_{m} S_{m}^{*} W_{k}\right)^{\dagger} W_{k}^{*} S_{m} S_{m}^{*} f_{d}
$$

we can again dramatically reduce the error. This is documented in Figure 3. Note that $f_{d, m}, \tilde{f}_{d, n, m}$ and $f_{d, k, m}$ use exactly the same samples.


Figure 4: The left figure shows $\Psi(n, 1)$ for $n=1, \ldots, 100$ together with the mapping $t \mapsto 2.45 t^{2}$. The right figure shows $\Psi(n, 2)$ for $n=1, \ldots, 100$ together with the mapping $t \mapsto 1.3 t^{2}$

Example 3.5. In this example we consider numerical estimates of the function $\Phi$ defined in (3.8). However, as pointed out above, the function $\Psi$ defined in (3.10) is much simpler to compute and we have

$$
\Phi(S, W, n, \theta) \leq \Psi(S, W, n, \theta), \quad \forall(n, \theta) \in \mathbb{N} \times \mathbb{R}_{+}
$$

Sections of the graphs of $\Psi(\cdot, 1)$ and $\Psi(\cdot, 2)$ are displayed in Figure 4. Note that it is not a big surprise that

$$
\Psi(S, W, n, \theta)=\mathcal{O}\left(n^{2}\right), \quad \theta=1,2 .
$$

In fact, it is not hard to use the fact that

$$
\Phi(S, \widetilde{W}, n,)=\mathcal{O}\left(n^{2}\right), \quad \forall \theta \in \mathbb{R}_{+}
$$

(see [5]) in order to prove that

$$
\Phi(S, W, n,)=\mathcal{O}\left(n^{2}\right), \quad \forall \theta \in \mathbb{R}_{+}
$$

## 4 Generalized Sampling and Inverse Problems

Suppose now that we are given Hilbert spaces $\mathcal{X}$ and $\mathcal{Y}$ and a linear and bounded operator $A$ that maps between $\mathcal{X}$ and $\mathcal{Y}$. Assume $A$ is equipped with singular system $\left\{\sigma_{k}, v_{k}, u_{k}\right\}$, where the orthonormal system $\left\{v_{k}\right\}_{k \in \mathbb{N}}$ spans $\mathcal{V}=\mathcal{N}(A)^{\perp}$ and where the orthonormal system $\left\{u_{k}\right\}_{k \in \mathbb{N}}$ spans $\mathcal{U}=\mathcal{N}\left(A^{*}\right)^{\perp}$. We aim to solve

$$
A f=g,
$$

where we are typically faced with noisy data $g^{\delta}=g+z$ with $\|z\| \leq \delta$. We assume that we have a
 $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ spanning $\mathcal{W}=\overline{\operatorname{span}}\left\{w_{k}\right\}_{k \in \mathbb{N}}=\mathcal{N}(A)^{\perp} \subset \mathcal{X}$ at our disposal. With $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ we are able to take samples of $g$,

$$
\eta=S^{*} g=\left\{\left\langle g, s_{k}\right\rangle\right\}_{k \in \mathbb{N}}
$$

and with $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ we expand $f$,

$$
f=W \beta=\sum_{k \in \mathbb{N}} \beta_{k} w_{k}
$$

As we cannot deal with infinitely many samples of $g$ and infinite series expansions for $f$, the goal is to reconstruct the best possible approximation of the form $\tilde{f}=W_{n} \tilde{\beta}=\sum_{k=1}^{n} \tilde{\beta}_{k} w_{k}$ based on the finite subset $\left\{\left\langle g, s_{k}\right\rangle\right\}_{k=1}^{r}=S_{r}^{*} g$ of the full sampling information $\eta$. We shall do this by exploiting the ideas of generalized sampling. As we now allow the problem to be ill-posed and the data noisy, we rephrase the problem as follows: given $n$, is there some $r$ (determine if possible) such that

$$
\left\|S_{r}^{*} g^{\delta}-S_{r}^{*} A W_{n} \tilde{\beta}\right\|^{2} \rightarrow \min _{\tilde{\beta}}
$$

has a unique solution? The corresponding normal equation is given by

$$
\left(W_{n}^{*} A^{*} S_{r} S_{r}^{*} A W_{n}\right)\left(\begin{array}{c}
\tilde{\beta}_{1}  \tag{4.1}\\
\vdots \\
\tilde{\beta}_{n}
\end{array}\right)=W_{n}^{*} A^{*} S_{r}\left(\begin{array}{c}
\left\langle g^{\delta}, s_{1}\right\rangle \\
\vdots \\
\left\langle g^{\delta}, s_{r}\right\rangle
\end{array}\right)
$$

and therefore

$$
\begin{equation*}
f_{n, r}^{\delta}=W_{n} \tilde{\beta}=W_{n}\left(W_{n}^{*} A^{*} S_{r} S_{r}^{*} A W_{n}\right)^{\dagger} W_{n}^{*} A^{*} S_{r} S_{r}^{*} g \tag{4.2}
\end{equation*}
$$

As already mentioned, the problem can be ill-posed and therefore the generalized inverse in (4.2) need not exist. Hence we are first faced with regularization issues. The second task is to analyze (4.2) by both establishing existence of approximations of the form $f_{n, r}^{\delta}$ and providing error bounds for $\left\|f-f_{n, r}^{\delta}\right\|$.

In what follows we discuss two different treatments of (4.2). Both proposed techniques heavily rely on the singular value decomposition of the operator $A$. This allows a splitting into a sampling and recovery step (which is quite natural). The sampling step in both algorithms is almost the same, whereas the recovery steps are rather different. In the first approach the recovery step relies on (classical) regularization principles, whereas in the second approach we seek to stabilize the ill-posed problem by adequate discretizations (which, of course, is also a regularization technique).

### 4.1 Regularized Reconstruction

Let us consider the normal equation $A^{*} A f=A^{*} g$ and let $A^{\dagger}$ denote the generalized inverse of $A$. If $g \in \mathcal{D}\left(A^{\dagger}\right)$, we can define $f^{\dagger}:=A^{\dagger} g$. If $A$ is injective it makes sense to define $A^{\dagger}:=\left(A^{*} A\right)^{-1} A^{*}$. A stabilized version of $f^{\dagger}$ can be constructed through $f^{\alpha}=R_{\alpha} g$ with $R_{\alpha}=F_{\alpha}\left(A^{*} A\right) A^{*}$ with properly chosen $F_{\alpha}$ (for an extensive discussion on the choice of $F_{\alpha}$ see [15] or [23] and references therein).

Therefore, we can write

$$
f^{\alpha}=W \beta^{\alpha}=\sum_{k} \beta_{k}^{\alpha} w_{k}=F_{\alpha}\left(A^{*} A\right) A^{*} g=\sum_{l} F_{\alpha}\left(\sigma_{l}^{2}\right) \sigma_{l}\left\langle g, u_{l}\right\rangle v_{l} .
$$

For the singular system we associate the corresponding synthesis and analysis operators which we denote by $V, U$ and $V^{*}, U^{*}$ as well as its finite versions $V_{m}, U_{m}$ and $V_{m}^{*}, U_{m}^{*}$. Therefore, for all $j \in \mathbb{N}$ we have

$$
\left(V^{*} f^{\alpha}\right)_{j}=\sum_{k} \beta_{k}^{\alpha}\left\langle w_{k}, v_{j}\right\rangle=\sum_{l} F_{\alpha}\left(\sigma_{l}^{2}\right) \sigma_{l}\left\langle g, u_{l}\right\rangle\left\langle v_{l}, v_{j}\right\rangle,
$$

which is nothing other than

$$
V^{*} W \beta^{\alpha}=\Theta_{\alpha} \Sigma \gamma \quad \text { or equivalently, } \quad \Theta_{\alpha}^{-1} V^{*} W \beta^{\alpha}=\Sigma \gamma,
$$

where

$$
V^{*} W=\left(\begin{array}{ccc}
\left\langle w_{1}, v_{1}\right\rangle & \left\langle w_{1}, v_{2}\right\rangle & \ldots \\
\left\langle w_{2}, v_{1}\right\rangle & \left\langle w_{2}, v_{2}\right\rangle & \ldots \\
\vdots & \vdots & \ddots
\end{array}\right), \Theta_{\alpha}=\left(\begin{array}{ccc}
F_{\alpha}\left(\sigma_{1}^{2}\right) & 0 & \ldots \\
0 & F_{\alpha}\left(\sigma_{2}^{2}\right) & \ldots \\
\vdots & \vdots & \ddots
\end{array}\right), \Sigma=\left(\begin{array}{ccc}
\sigma_{1} & 0 & \ldots \\
0 & \sigma_{2} & \ldots \\
\vdots & \vdots & \ddots
\end{array}\right) .
$$

The vector $\gamma=U^{*} g=\left\{\left\langle g, u_{l}\right\rangle\right\}_{l \in \mathbb{N}}$ is not accessible in practice and must be related to the samples of $g$ with respect to our sampling frame $\left\{s_{k}\right\}_{k \in \mathbb{N}}$,

$$
\eta=S^{*} g=S^{*} U U^{*} g=S^{*} U \gamma, \text { where } \quad S^{*} U=\left(\begin{array}{ccc}
\left\langle u_{1}, s_{1}\right\rangle & \left\langle u_{1}, s_{2}\right\rangle & \ldots \\
\left\langle u_{2}, s_{1}\right\rangle & \left\langle u_{2}, s_{2}\right\rangle & \ldots \\
\vdots & \vdots & \ddots .
\end{array}\right)
$$

Consequently, in order to compute $f^{\alpha}=W \beta^{\alpha}$ we are faced with two infinite dimensional linear equations that need to be solved,

$$
\begin{align*}
S^{*} U \gamma & =\eta  \tag{4.3}\\
\Theta_{\alpha}^{-1} V^{*} W \beta^{\alpha} & =\Sigma \gamma \tag{4.4}
\end{align*}
$$

We proceed in two steps. To find an approximation to the solution of (4.4) we assume in a first step that a finite length approximation $\gamma_{m}=U_{m}^{*} g$ to the solution of (4.3) is at our disposal (with $m$ specified below). Then, in a second step, we derive an approximation $\gamma_{m, r}^{\delta}$ to $\gamma_{m}$ where we also allow the data $g$ to be noisy.

Suppose now that we have the samples $U_{m}^{*} g$ at our disposal. Based on this information we construct an approximation $f_{n, m}^{\alpha}=W_{n} \beta_{n, m}^{\alpha} \in \mathcal{W}_{n}$ as follows:

$$
\begin{equation*}
f_{n, m}^{\alpha}=W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Theta_{\alpha, m}^{-1} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Sigma_{m} U_{m}^{*} g \tag{4.5}
\end{equation*}
$$

where $\Theta_{\alpha, m}=\left.P_{m} \Theta_{\alpha}\right|_{P_{m} l^{2}(\mathbb{N})}$. As the following theorem indicates, $f_{n, m}^{\alpha}$ is uniquely defined for sufficiently large $m$ :

Theorem 4.1. Let $\mathcal{X}, \mathcal{Y}$ be separable Hilbert spaces and $A: \mathcal{X} \rightarrow \mathcal{Y}$ be a linear and bounded operator with singular system $\left\{\sigma_{k}, v_{k}, u_{k}\right\}$. Moreover, let $\mathcal{W}=\mathcal{N}(A)^{\perp} \subset \mathcal{X}$ and $\mathcal{S}=\mathcal{N}\left(A^{*}\right)^{\perp} \subset \mathcal{Y}$. Suppose $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ are frames for $\mathcal{S}$ and $\mathcal{W}$ respectively. Then, for the generalized solution $f^{\dagger} \in \mathcal{W} \subset$ $\mathcal{X}$ and $n \in \mathbb{N}$ and $\alpha>0$, there is an $M \in \mathbb{N}$ such that, for all $m \geq M$ the approximation $f_{n, m}^{\alpha}$ exists and is unique. Moreover,

$$
\left\|f^{\dagger}-f_{n, m}^{\alpha}\right\|_{\mathcal{X}} \leq\left(2+K_{n, m, \alpha}\right)\left\|f^{\dagger}-f^{\alpha}\right\|_{\mathcal{X}}+\left(1+K_{n, m, \alpha}\right)\left\|P_{\mathcal{W}_{n}}^{\perp} f^{\dagger}\right\|_{\mathcal{X}}
$$

where

$$
K_{n, m, \alpha}:=\left\|W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Theta_{\alpha, m}^{-1} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Theta_{\alpha, m}^{-1} V_{m}^{*} P_{\mathcal{W}_{n}}^{\perp}\right\|
$$

In particular, for m large enough,

$$
K_{n, m, \alpha} \leq \frac{\lambda_{\max }\left(\Theta_{\alpha, m}^{-2}\right)}{\lambda_{\min }\left(\Theta_{\alpha, m}^{-2}\right)},
$$

where $\lambda_{\max }\left(\Theta_{\alpha, m}^{-2}\right)$ and $\lambda_{\min }\left(\Theta_{\alpha, m}^{-2}\right)$ denote the largest and smallest eigenvalue of $\Theta_{\alpha, m}^{-2}$ respectively.
Proof. Solving the optimization problem

$$
\begin{equation*}
\left\|\Sigma_{m} U_{m}^{*} g-\Theta_{\alpha, m}^{-1} V_{m}^{*} W_{n} \beta_{n, m}^{\alpha}\right\|^{2} \rightarrow \min _{\beta_{n, m}^{\alpha}} \tag{4.6}
\end{equation*}
$$

results in

$$
\begin{aligned}
\beta_{n, m}^{\alpha} & =\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Sigma_{m} U_{m}^{*} g \\
& =\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*}\left(P_{\mathcal{W}_{n}} f^{\alpha}+P_{\mathcal{W}_{n}}^{\perp} f^{\alpha}\right)
\end{aligned}
$$

and therefore, as argued in the proof of Theorem 3.1, for $m \geq M$ and $M$ large enough,

$$
\begin{aligned}
f_{n, m}^{\alpha} & =W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*}\left(P_{\mathcal{W}_{n}} f^{\alpha}+P_{\mathcal{W}_{n}}^{\perp} f^{\alpha}\right) \\
& =P_{\mathcal{W}_{n}} f^{\alpha}+W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*} P_{\mathcal{W}_{n}}^{\perp} f^{\alpha}
\end{aligned}
$$

and $f_{n, m}^{\alpha}$ is unique. Consequently,

$$
\begin{aligned}
\left\|f^{\dagger}-f_{n, m}^{\alpha}\right\|_{\mathcal{X}} & \leq\left\|f^{\dagger}-f^{\alpha}\right\|_{\mathcal{X}}+\left\|f^{\alpha}-f_{n, m}^{\alpha}\right\|_{\mathcal{X}} \\
& \leq\left\|f^{\dagger}-f^{\alpha}\right\|_{\mathcal{X}}+\left\|P_{\mathcal{W}_{n}}^{\perp} f^{\alpha}-W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-2} V_{m}^{*} P_{\mathcal{W}_{n}}^{\perp} f^{\alpha}\right\|_{\mathcal{X}} \\
& \leq\left\|f^{\dagger}-f^{\alpha}\right\|_{\mathcal{X}}+\left(1+K_{n, m, \alpha}\right)\left\|P_{\mathcal{W}_{n}}^{\perp} f^{\alpha}\right\|_{\mathcal{X}} \\
& \leq\left(2+K_{n, m, \alpha}\right)\left\|f^{\dagger}-f^{\alpha}\right\|_{\mathcal{X}}+\left(1+K_{n, m, \alpha}\right)\left\|P_{\mathcal{W}_{n}}^{\perp} f^{\dagger}\right\|_{\mathcal{X}} .
\end{aligned}
$$

To prove that, for fixed $n$ and $\alpha>0$ and $m$ large enough, the constant $K_{n, m, \alpha}$ can be bounded by $\lambda_{\max }\left(\Theta_{\alpha, m}^{-2}\right) \lambda_{\min }^{-1}\left(\Theta_{\alpha, m}^{-2}\right)$, we first define $\tilde{v}_{k}:=F_{\alpha}\left(\sigma_{k}^{2}\right)^{-1} v_{k}$. This results in

$$
K_{n, m, \alpha}=\left\|W_{n}\left(W_{n}^{*} \tilde{V}_{m} \tilde{V}_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} \tilde{V}_{m} \tilde{V}_{m}^{*} P_{\mathcal{W}_{n}}^{\perp}\right\|
$$

As $\left\{v_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal basis, the system $\left\{\tilde{v}_{k}\right\}_{k \in \mathbb{N}}$ is a Riesz basis and therefore we obtain

$$
K_{n, m, \alpha} \leq\left\|\left(\tilde{V}_{m} \tilde{V}_{m}^{*}\right)^{-1}\right\| \mathcal{V}_{m} \rightarrow \mathcal{V}_{m}\left\|\tilde{V}_{m} \tilde{V}_{m}^{*}\right\|_{\mathcal{V}_{m} \rightarrow \mathcal{V}_{m}}
$$

For $f \in \mathcal{V}_{m}$, it follows that

$$
\min _{j=1, \ldots, m} F_{\alpha}\left(\sigma_{j}^{2}\right)^{-2}\|f\|^{2} \leq\left\langle\tilde{V}_{m} \tilde{V}_{m}^{*} f, f\right\rangle=\sum_{j=1}^{m} F_{\alpha}\left(\sigma_{j}^{2}\right)^{-2}\left|\left\langle f, v_{j}\right\rangle\right|^{2} \leq \max _{j=1, \ldots, m} F_{\alpha}\left(\sigma_{j}^{2}\right)^{-2}\|f\|^{2}
$$

resulting in

$$
K_{n, m, \alpha} \leq \frac{\max _{j=1, \ldots, m} F_{\alpha}\left(\sigma_{j}^{2}\right)^{-2}}{\min _{j=1, \ldots, m} F_{\alpha}\left(\sigma_{j}^{2}\right)^{-2}}
$$

as required
As already mentioned, the samples to which we have access in practice are given by $\eta^{\delta}=S^{*} g^{\delta}$. Based on this (possible noisy) information we construct an approximation $\gamma_{m, r}^{\delta}$ to the solution of (4.3) as follows:

$$
\begin{equation*}
\gamma_{m, r}^{\delta}=\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} g^{\delta} \text { and hence, } g_{m, r}^{\delta}=U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} g^{\delta} \tag{4.7}
\end{equation*}
$$

Consequently, combining (4.5) and (4.7), we obtain an approximation to the solution of (4.4) that can indeed be realized in practice,

$$
\begin{equation*}
f_{n, m, r}^{\alpha, \delta}=W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Theta_{\alpha, m}^{-1} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Sigma_{m} U_{m}^{*} g_{m, r}^{\delta} \tag{4.8}
\end{equation*}
$$

As a simple consequence of the Theorems 3.1 and 4.1 we have the following:
Theorem 4.2. If the assumptions made in Theorems 3.1 and 4.1 hold true, then

$$
\left\|f^{\dagger}-f_{n, m, r}^{\alpha, \delta}\right\| \leq\left\|f^{\dagger}-f_{n, m}^{\alpha}\right\|+K_{n, m, \alpha}^{1}\left(\sigma_{m+1}\left(1+K_{m, r}^{2}\right)\left\|P_{\mathcal{V}_{m}}^{\perp} f\right\|+K_{m, r}^{3} \delta\right)
$$

with constants defined by

$$
\begin{gathered}
K_{n, m, \alpha}^{1}:=\left\|W_{n}\left(W_{n}^{*} \tilde{V}_{m} \tilde{V}_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} \tilde{V}_{m} V_{m}^{*} A^{*}\right\|, K_{m, r}^{2}:=\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp}\right\| \\
K_{m, r}^{3}:=\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*}\right\|
\end{gathered}
$$

where $K_{n, m, \alpha}^{1} \leq \frac{\lambda_{\max }\left(\Theta_{\alpha, m}^{-1}\right)}{\lambda_{\min }\left(\Theta_{\alpha, m}^{-2}\right)}\|A\|$, and, for sufficiently large $r, K_{m, r}^{2}, K_{m, r}^{3} \leq 2 \frac{D}{C}$. Moreover, if $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal basis, then, for fixed $m, K_{m, r}^{2} \rightarrow 0, K_{m, r}^{3} \rightarrow 1$ as $r \rightarrow \infty$.

Proof. First, with $\Sigma_{m} U_{m}^{*}=V_{m}^{*} A^{*}$, we observe that

$$
\begin{aligned}
f_{n, m, r}^{\alpha, \delta} & =W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Theta_{\alpha, m}^{-1} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Sigma_{m} U_{m}^{*}\left(g-g+g_{m, r}^{\delta}\right) \\
& =f_{n, m}^{\alpha}+W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Theta_{\alpha, m}^{-1} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} V_{m}^{*} A^{*}\left(g_{m, r}^{\delta}-g\right)
\end{aligned}
$$

Moreover, with the help of Theorem 3.1,

$$
g-g_{m, r}^{\delta}=P_{\mathcal{U}_{m}}^{\perp} g+U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} P_{\mathcal{U}_{m}}^{\perp} g+U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*}\left(g^{\delta}-g\right),
$$

and since

$$
P_{\mathcal{U}_{m}}^{\perp} g=\left(I-U P_{m} U^{*}\right) g=U\left(I-P_{m}\right) U^{*} g=U P_{m}^{\perp} \gamma=U P_{m}^{\perp} \Sigma V^{*} W \beta
$$

it follows that

$$
\begin{aligned}
\left\|g-g_{m, r}^{\delta}\right\| \leq & \sigma_{m+1}\left(1+\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp}\right\|\right)\left\|P_{\mathcal{V}_{m}}^{\perp} f\right\| \\
& +\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*}\right\| \delta
\end{aligned}
$$

This implies, with $\tilde{V}_{m}=V_{m} \Theta_{\alpha, m}^{-1}$, that

$$
\begin{aligned}
\left\|f^{\dagger}-f_{n, m, r}^{\alpha, \delta}\right\| & \leq\left\|f^{\dagger}-f_{n, m}^{\alpha}\right\|+\left\|W_{n}\left(W_{n}^{*} \tilde{V}_{m} \tilde{V}_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} \tilde{V}_{m} V_{m}^{*} A^{*}\right\|\left\|g_{m, r}^{\delta}-g\right\| \\
& \leq\left\|f^{\dagger}-f_{n, m}^{\alpha}\right\|+K_{n, m, \alpha}^{1}\left(\sigma_{m+1}\left(1+K_{m, r}^{2}\right)\left\|P_{\mathcal{V}_{m}}^{\perp} f\right\|+K_{m, r}^{3} \delta\right) .
\end{aligned}
$$

Now, for $m$ sufficiently large,

$$
\begin{aligned}
K_{n, m, \alpha}^{1} & =\left\|W_{n}\left(W_{n}^{*} \tilde{V}_{m} \tilde{V}_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} \tilde{V}_{m} V_{m}^{*} A^{*}\right\| \\
& =\left\|\left(\tilde{V}_{m} \tilde{V}_{m}^{*}\right)^{-1} \tilde{V}_{m} \tilde{V}_{m}^{*} W_{n}\left(W_{n}^{*} \tilde{V}_{m} \tilde{V}_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} \tilde{V}_{m} V_{m}^{*} A^{*}\right\| \\
& \leq\left\|\left(\tilde{V}_{m} \tilde{V}_{m}^{*}\right)^{-1}\right\|\left\|\tilde{V}_{m}\right\|\left\|V_{m}^{*}\right\|\left\|A^{*}\right\| \leq \frac{\max _{j=1, \ldots, m} F_{\alpha}\left(\sigma_{j}^{2}\right)^{-1}}{\min _{j=1, \ldots, m} F_{\alpha}\left(\sigma_{j}^{2}\right)^{-2}}\|A\| .
\end{aligned}
$$

Moreover, if $r$ is also sufficiently large we obtain

$$
K_{m, r}^{2} \leq\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*}\right\| \leq 2\left\|\left(S S^{*}\right)^{-1}\right\|\left\|S S^{*}\right\|, K_{m, r}^{3} \leq 2\left\|\left(S S^{*}\right)^{-1}\right\|\left\|S S^{*}\right\|
$$

To see that $K_{m, r}^{2} \rightarrow 0$ as $r \rightarrow \infty$ we observe that, since $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{u_{k}\right\}_{k \in \mathbb{N}}$ span the same space, we have that $P_{\mathcal{S}} U=U$, where $P_{\mathcal{S}}=S S^{*}$. Thus,

$$
U_{m}\left(U_{m}^{*} S S^{*} U_{m}\right)^{\dagger} U_{m}^{*} S S^{*} U P_{m}^{\perp}=U_{m}\left(U_{m}^{*} S S^{*} U_{m}\right)^{\dagger} P_{m} U^{*} U P_{m}^{\perp}=0
$$

Also,

$$
U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp} \longrightarrow U_{m}\left(U_{m}^{*} S S^{*} U_{m}\right)^{\dagger} U_{m}^{*} S S^{*} U P_{m}^{\perp}, \quad r \rightarrow \infty
$$

by the strong convergence of $S_{r} S_{r}^{*}$ to $S S^{*}$ and compactness of $U_{m}$. The assertion about $K_{m, r}^{2}$ follows.

### 4.2 Non-Regularized Reconstruction

The approach in the previous section was essentially based on the normal equation $A^{*} A f=A^{*} g$. As an alternative, we now propose an approach based on directly utilizing the singular value decomposition of $A$. In other words, since $A=U \Sigma V^{*}$, we have

$$
\eta=S^{*} g=S^{*} U \Sigma V^{*} W \beta=S^{*} U \gamma
$$

and, as in the previous section, this results in two linear equations

$$
\begin{align*}
S^{*} U \gamma & =\eta  \tag{4.9}\\
V^{*} W \beta & =\Sigma^{-1} \gamma . \tag{4.10}
\end{align*}
$$

Based on equation (4.9) we construct an approximation to $g$,

$$
\begin{equation*}
g_{m, r}^{\delta}=U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} g^{\delta}, \tag{4.11}
\end{equation*}
$$

and with the help of equation (4.10) we construct

$$
\begin{equation*}
f_{n, m}=W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1} U_{m}^{*} g \tag{4.12}
\end{equation*}
$$

Combining (4.11) and (4.12) suggests the following approximation to $f$ :

$$
\begin{equation*}
f_{n, m, r}^{\delta}=W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1} U_{m}^{*} g_{m, r}^{\delta} \tag{4.13}
\end{equation*}
$$

Theorem 4.3. Let $\mathcal{X}, \mathcal{Y}$ be separable Hilbert spaces and $A: \mathcal{X} \rightarrow \mathcal{Y}$ be a linear and bounded operator with singular system $\left\{\sigma_{k}, v_{k}, u_{k}\right\}$ (such that $\left\{v_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal system for $\mathcal{N}(A)^{\perp}$ and $\left\{u_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal system for $\left.\mathcal{N}\left(A^{*}\right)^{\perp}\right)$. Moreover, let $\mathcal{W}=\mathcal{N}(A)^{\perp} \subset \mathcal{X}$ and $\mathcal{S}=\mathcal{N}\left(A^{*}\right)^{\perp} \subset \mathcal{Y}$. Suppose $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ are frames for $\mathcal{S}$ and $\mathcal{W}$ respectively. Then, for the solution $f \in \mathcal{W} \subset \mathcal{X}$ and $n \in \mathbb{N}$, there exist $M, R \in \mathbb{N}$ such that, for all $m \geq M$ and $r \geq R$ the approximation $f_{n, m, r}^{\delta}$ exists and is unique. Moreover,

$$
\begin{equation*}
\left\|f-f_{n, m, r}^{\delta}\right\|_{\mathcal{X}} \leq\left(1+K_{n, m}^{1}\right)\left\|P_{\mathcal{W}_{n}}^{\perp} f\right\|+K_{m, r}^{2}\left\|P_{\mathcal{V}_{m}}^{\perp} f\right\|+K_{m, r}^{3} \frac{\delta}{\sigma_{m}} \tag{4.14}
\end{equation*}
$$

with

$$
\begin{gathered}
K_{n, m}^{1}=\left\|W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} V_{m}^{*} P_{\mathcal{W}_{n}}^{\perp}\right\|, K_{m, r}^{2}=\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp}\right\|, \\
K_{m, r}^{3}=\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*}\right\|
\end{gathered}
$$

In particular, $K_{n, m}^{1} \rightarrow 0$ as $m \rightarrow \infty$, and, for sufficiently large $r, K_{m, r}^{2}, K_{m, r}^{3} \leq 2 \frac{D}{C}$. Moreover, if $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ forms an orthonormal system, then $K_{m, r}^{2} \rightarrow 0$ as $r \rightarrow \infty$ and $K_{m, r}^{3} \rightarrow 1$ as $r \rightarrow \infty$.

Proof. Following the same lines of the proof of Theorem 3.1, we obtain for (4.11),

$$
\begin{aligned}
g_{m, r}^{\delta} & =U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*}(g+z) \\
& =P_{\mathcal{U}_{m}} g+U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} P_{\mathcal{U}_{m}}^{\perp} g+U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} z
\end{aligned}
$$

Therefore, we have for (4.13) with $\Sigma_{m}^{-1} U_{m}^{*} P_{\mathcal{U}_{m}} g=V_{m}^{*} f$,

$$
\begin{aligned}
f_{n, m, r}^{\delta}= & W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} V_{m}^{*}\left(P_{\mathcal{W}_{n}} f+P_{\mathcal{W}_{n}}^{\perp} f\right) \\
& +W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1} U_{m}^{*} U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} P_{\mathcal{U}_{m}}^{\perp} g \\
& +W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1} U_{m}^{*} U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} z
\end{aligned}
$$

Note that we may argue exactly as in the proof of Theorem 3.1, by using the fact that $\left\{v_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal system for $\mathcal{N}(A)^{\perp}=\mathcal{W}$ and deduce that there is an $M \in \mathbb{N}$ such that

$$
\begin{equation*}
W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} V_{m}^{*} P_{\mathcal{W}_{n}} f=P_{\mathcal{W}_{n}} f, \quad m \geq M \tag{4.15}
\end{equation*}
$$

Thus, by simply plugging in the expression for $f_{n, m, r}^{\delta}$, it follows that

$$
\begin{align*}
\| f- & f_{n, m, r}^{\delta} \| \\
& \left.\leq \|\left(I-W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} V_{m}^{*}\right) P_{\mathcal{W}_{n}}^{\perp} f\right) \| \\
& +\left\|W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1} U_{m}^{*} U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} P_{\mathcal{U}_{m}}^{\perp} g\right\|  \tag{4.16}\\
& +\left\|W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1} U_{m}^{*} U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} z\right\| \quad m \geq M
\end{align*}
$$

To get the asserted bounds we shall bound separately the three terms on the right-hand side of (4.16). The first term is obvious, so we focus on the second and third. To bound the second term observe that

$$
\left\|W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1}\right\|=\left\|\left(V_{m} V_{m}^{*}\right)^{-1} V_{m} P_{\operatorname{Ran}\left(V_{m}^{*} W_{n}\right)} \Sigma_{m}^{-1}\right\| \leq \frac{1}{\sigma_{m}}
$$

and that

$$
P_{\mathcal{U}_{m}}^{\perp} g=\left(I-U P_{m} U^{*}\right) g=U\left(I-P_{m}\right) U^{*} g=U P_{m}^{\perp} \gamma=U P_{m}^{\perp} \Sigma V^{*} W \beta .
$$

In particular, since $\sigma_{j} \geq \sigma_{j+1}$ for all $j \in \mathbb{N}$, it follows that

$$
\begin{align*}
& \left\|W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1} U_{m}^{*} U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} P_{\mathcal{U}_{m}}^{\perp} g\right\| \\
& \leq\left\|W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1}\right\|\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp}\right\|\left\|P_{m}^{\perp} \Sigma V^{*} W \beta\right\| \\
& \leq \frac{\sigma_{m+1}}{\sigma_{m}}\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp}\right\|\left\|P_{m}^{\perp} V^{*} W \beta\right\|  \tag{4.17}\\
& \leq\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp}\right\|\left\|P_{\mathcal{V}_{m}}^{\perp} f\right\| .
\end{align*}
$$

As for the third term we use (4.15) and the assumption on $z$ and get that

$$
\begin{gather*}
\left\|W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1} U_{m}^{*} U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} z\right\| \\
\leq \frac{1}{\sigma_{m}}\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*}\right\| \delta \tag{4.18}
\end{gather*}
$$

Hence, by (4.16), (4.17) and (4.18) we obtain

$$
\begin{aligned}
\left\|f-f_{n, m, r}^{\delta}\right\| \leq & \left.\left(1+\| W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} V_{m}^{*}\right) P_{\mathcal{W}_{n}}^{\perp}\right)\left\|\left\|P_{\mathcal{W}_{n}}^{\perp} f\right\|\right. \\
& +\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp}\right\|\left\|P_{\mathcal{V}_{m}}^{\perp} f\right\| \\
& +\left\|U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*}\right\| \frac{\delta}{\sigma_{m}}
\end{aligned}
$$

which yields (4.14).

To see that $K_{n, m}^{1} \rightarrow 0$ as $m \rightarrow \infty$ we can simply use Theorem 3.1 since, by assumption, $\left\{v_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{w_{k}\right\}_{k \in \mathbb{N}}$ span the same space and $\left\{v_{k}\right\}_{k \in \mathbb{N}}$ is an orthonormal basis. Similarly, the statement about $K_{m, r}^{3}$ when $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ is orthonormal follows from Theorem 3.1. To see that $K_{m, r}^{2} \rightarrow 0$ as $r \rightarrow \infty$ we observe that, since $\left\{s_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{u_{k}\right\}_{k \in \mathbb{N}}$ span the same space, we have $P_{\mathcal{S}} U=U$. Thus,

$$
U_{m}\left(U_{m}^{*} S S^{*} U_{m}\right)^{\dagger} U_{m}^{*} S S^{*} U P_{m}^{\perp}=U_{m}\left(U_{m}^{*} S S^{*} U_{m}\right)^{\dagger} P_{m} U^{*} U P_{m}^{\perp}=0
$$

Also,

$$
U_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} U P_{m}^{\perp} \longrightarrow U_{m}\left(U_{m}^{*} S S^{*} U_{m}\right)^{\dagger} U_{m}^{*} S S^{*} U P_{m}^{\perp}, \quad r \rightarrow \infty
$$

by the strong convergence of $S_{r} S_{r}^{*}$ to $S S^{*}$ and compactness of $U_{m}$. The assertion about $K_{m, r}^{2}$ follows.

### 4.3 Numerical Examples

In this section we will test the frameworks proposed in the previous subsections. First, we discuss a onedimensional example for which we analyze the suggested regularized and non-regularized reconstruction methods. Thereafter, we consider a two-dimensional experiment. The goal is to verify that we can achieve, even in the presence of noise, a reasonable reconstruction by the proposed sampling-recovery technique.

Example 4.4. In order illustrate the proposed sampling theorems (Theorem 4.1, 4.2, and 4.3), we consider the linear operator $A: \mathrm{L}_{2}([0,1]) \rightarrow \mathrm{L}_{2}([0,1])$ defined by

$$
g(t)=A f(t)=\int_{0}^{t} f(s) d s
$$

with singular system $\left\{\sigma_{k}, v_{k}, u_{k}\right\}$ given by

$$
\sigma_{k}=\frac{1}{(k+1 / 2) \pi}, \quad v_{k}=\sqrt{2} \cos (k+1 / 2) \pi t, \quad u_{k}=\sqrt{2} \sin (k+1 / 2) \pi t
$$

Note that $\left\{v_{k}\right\}_{k \in \mathbb{N}}$ and $\left\{u_{k}\right\}_{k \in \mathbb{N}}$ form orthonormal systems for $\mathrm{L}_{2}([0,1])$. To keep technicalities at a reasonable level, we choose the Fourier basis as both the recovery system $\left\{w_{k}\right\}_{k \in \mathbb{Z}}$ and sampling system $\left\{s_{k}\right\}_{k \in \mathbb{Z}}$, i.e.

$$
w_{k}(t)=e^{2 \pi i k t} \quad \text { and } \quad s_{k}(t)=e^{2 \pi i k t}
$$

Let the signal $f$ to be reconstructed be defined by $f(t)=\cos 2 \pi t$. Consequently, $f$ can be expanded as follows,

$$
f(t)=W \beta=\sum_{k \in \mathbb{Z}} \beta_{k} e^{2 \pi i k t}=\frac{1}{2} e^{2 \pi i(-1) t}+\frac{1}{2} e^{2 \pi i(+1) t}
$$

Consequently, $\beta_{-1}=1 / 2, \beta_{1}=1 / 2$, and $\beta_{k}=0$ for $k \in \mathbb{Z} \backslash\{-1,+1\}$. Moreover, the data $g$ are given through $g(t)=A f(t)=1 /(2 \pi) \sin 2 \pi t$. In this particular example we also have explicit expression for all further required quantities,

$$
\begin{aligned}
\gamma & =U^{*} g=\left\{\gamma_{l}\right\}_{l \in \mathbb{N}}=\left\{\frac{4 \sqrt{2} \cos (l \pi)}{\pi^{2}\left(4 l^{2}+4 l-15\right)}\right\}_{l \in \mathbb{N}} \\
\eta & =S^{*} g=\left\{\eta_{k}\right\}_{k \in \mathbb{Z}} \text { with } \eta_{-1}=\frac{-i}{4 \pi}, \eta_{+1}=\frac{i}{4 \pi} \text { and } \eta_{k}=0, k \neq \pm 1 \\
V^{*} W & =\left(\frac{\sqrt{2}((l \pi+\pi / 2) \cos (l \pi)-2 \pi i k)}{(2 \pi i k)^{2}+(l \pi+\pi / 2)^{2}}\right)_{l \in \mathbb{N}, k \in \mathbb{Z}} \\
S^{*} U & =\left(\frac{(-1)^{l+1} \sqrt{2}((l \pi+\pi / 2) \cos (l \pi)+2 \pi i k)}{(2 \pi i k)^{2}+(l \pi+\pi / 2)^{2}}\right)_{k \in \mathbb{Z}, l \in \mathbb{N}}
\end{aligned}
$$

The approximations to $f$ from the $r$ samples $S_{r}^{*} g^{\delta}$ may now be derived by

$$
\begin{aligned}
f_{n, m, r}^{\alpha, \delta} & =W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Theta_{\alpha, m}^{-1} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Sigma_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} g^{\delta} \text { and } \\
f_{n, m, r}^{\delta} & =W_{n}\left(W_{n}^{*} V_{m} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Sigma_{m}^{-1}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} g^{\delta},
\end{aligned}
$$



Figure 5: Recovery results for $\varepsilon_{r e l}=0 \%$. Top (from left to right): $\eta^{\delta}=S_{r}^{*}(g+z)(\cdot)$ and $S_{r}^{*} U_{m} \gamma(\circ)$, $\beta_{20,30,40}^{\delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{\delta}=W_{n} \beta_{20,30,40}^{\delta}(-)$ and $f(--)$. Bottom (from left to right): $\beta_{20,30,40}^{0, \delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{0, \delta}=W_{n} \beta_{20,30,40}^{0, \delta}(-)$ and $f(--), \beta_{20,30,40}^{\alpha_{o p t}, \delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{\alpha_{o p t}, \delta}=W_{n} \beta_{20,30,40}^{\alpha_{o p}, \delta}(-)$ and $f(-$ -).
where we focus here on Tikhonov regularization, i.e. the entries in $\Theta_{\alpha, m}$ are given by $F_{\alpha}\left(\sigma_{k}^{2}\right)=1 /(\alpha+$ $\left.\sigma_{k}^{2}\right)$. We discuss now several different recovery scenarios.

In the first case we choose a fixed (and reasonable) setting for $n, m$ and $r$ and vary the noise level $\delta$ and compare the recovery quality of $f_{n, m, r}^{\alpha, \delta}$ and $f_{n, m, r}^{\delta}$ while experimentally tuning the regularization parameter $\alpha$ towards optimal recovery. This experiment shall show that for a fixed number of data samples and coefficients in the series expansion of the solution an optimal choice of regularization parameter induces a substantially improved recovery.

In the second case we fix the number $n$ of coefficients in series expansion of the solution and try to find for different noise levels $\delta$ reasonable integers $m$ and $r$ to derive $f_{n, m, r}^{\delta}$. For the same numbers $m$ and $r$ we then experimentally determine an optimal $\alpha$ to compute $f_{n, m, r}^{\alpha, \delta}$. This experiment shall show that a reasonable choice of $m$ and $r$ may feasibly stabilize the recovery and providing approximations that cannot be significantly improved by a fine tuning of $\alpha$.

First case: vary $z=g-g^{\delta}$ such the relative error $\varepsilon_{\text {rel }}=100 \cdot\|z\| /\|g\|$ is $0 \%, 5 \%$ and $10 \%$ and let $n=20, m=30$ and $r=40$. The numerical results are illustrated in the following table and visualized in Figures 5,6, and 7 .

| $\varepsilon_{\text {rel }}, \delta$ | $\left\\|f-f_{20,30,40}^{\delta}\right\\|$ | $\left\\|f-f_{20,30,40}^{0, \delta}\right\\|$ | $\left\\|f-f_{20,30,40}^{\alpha_{\text {opt }}, \delta}\right\\|$ | $\alpha_{\text {opt }}$ | Fig. |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $0 \%, 0$ | 0.6262 | 0.4995 | 0.0071 | 0.00017 | 5 |
| $5 \%, 0.0056$ | 1.1738 | 0.9728 | 0.1536 | 0.00037 | 6 |
| $10 \%, 0.0113$ | 1.7593 | 1.5268 | 0.2265 | 0.00061 | 7 |

Second case: we first fix $n=10$ and ask then, for different relative errors $\varepsilon_{r e l} \in\{0 \%, 5 \%, 10 \%\}$, for an adequate choice (numerically determined) of $m$ and $r$ in order to derive an optimal approximation $f_{n, m, r}^{\delta}$. Then, we try by fine tuning $\alpha$ to obtain with $f_{n, m, r}^{\alpha, \delta}$ a comparable or possibly better approximation. The results are documented in the following table. The illustrations of this experiment are given in Figure 8 (the illustrations for $\varepsilon_{r e l}=0 \%$ are not provided since there is no visual difference).

| $\varepsilon_{\text {rel }}, \delta$ | $m$ | $r$ | $\left\\|f-f_{10, m, r}^{\delta}\right\\|$ | $\left\\|f-f_{10, m, r}^{0, \delta}\right\\|$ | $\left\\|f-f_{10, m, r}^{\alpha_{\text {opt }, ~} \\|}\right\\|$ | $\alpha_{\text {opt }}$ | Fig. |
| :--- | ---: | ---: | :--- | :--- | :--- | :--- | :---: |
| $0 \%, 0.0$ | 10 | 1000 | 0.002114839173 | 0.002114839160 | 0.000112 | 0.0000035 | - |
| $5 \%, 0.0042$ | 40 | 100 | 0.0303 | 0.0433 | 0.0371 | 0.000025 | 8 |
| $10 \%, 0.0075$ | 40 | 80 | 0.1044 | 0.2990 | 0.2732 | 0.0001 | 8 |

Example 4.5. In the second example we are concerned with a two dimensional deconvolution problem. The operator under consideration is the so-called Sobolev smoothing operator $A^{s}: H_{p}^{r}\left([0,1]^{2}\right) \rightarrow$


Figure 6: Recovery results for $\varepsilon_{r e l}=5 \%$. Top (from left to right): $\eta^{\delta}=S_{r}^{*}(g+z)(\cdot)$ and $S_{r}^{*} U_{m} \gamma(\circ)$, $\beta_{20,30,40}^{\delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{\delta}=W_{n} \beta_{20,30,40}^{\delta}(-)$ and $f(--)$. Bottom (from left to right): $\beta_{20,30,40}^{0, \delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{0, \delta}=W_{n} \beta_{20,30,40}^{0, \delta}(-)$ and $f(--), \beta_{20,30,40}^{\alpha_{o p t}, \delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{\alpha_{o p t}, \delta}=W_{n} \beta_{20,30,40}^{\alpha_{o p}, \delta}(-)$ and $f(-$ -).





Figure 7: Recovery results for $\varepsilon_{r e l}=10 \%$. Top (from left to right): $\eta^{\delta}=S_{r}^{*}(g+z)(\cdot)$ and $S_{r}^{*} U_{m} \gamma(\circ)$, $\beta_{20,30,40}^{\delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{\delta}=W_{n} \beta_{20,30,40}^{\delta}(-)$ and $f(--)$. Bottom (from left to right): $\beta_{20,30,40}^{0, \delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{0, \delta}=W_{n} \beta_{20,30,40}^{0, \delta}(-)$ and $f(--), \beta_{20,30,40}^{\alpha_{o p t}, \delta}(\cdot)$ and $\beta(\circ), f_{20,30,40}^{\alpha_{o p t}, \delta}=W_{n} \beta_{20,30,40}^{\alpha_{o p t}, \delta}(-)$ and $f(-$ -).
$H_{p}^{r+s}\left([0,1]^{2}\right)$. In order to precisely model this operator we first note that each function $f \in L_{2, p}\left([0,1]^{2}\right)$ can be represented by its Fourier series, $f=\sum_{k \in \mathbb{Z}^{2}} f_{k} e^{2 \pi i k}$, where $f_{k}$ denote the $L_{2}$-Fourier coefficients. For $f, g \in H_{p}^{s}\left([0,1]^{2}\right)$ we have $\langle f, g\rangle_{s}=\sum_{k \in \mathbb{Z}^{2}}\left(1+\|k\|^{2}\right)^{s} f_{k} \bar{g}_{k}$. With these inner products the system

$$
\left\{e_{k}^{s}:=\left(1+\|k\|^{2}\right)^{-s / 2} e^{2 \pi i k}, k \in \mathbb{Z}^{2}\right\}
$$

forms an orthonormal basis for $H_{p}^{s}\left([0,1]^{2}\right)$. Now we are ready to define the Sobolev-smoothing operator,

$$
A^{s}: H_{p}^{r}\left([0,1]^{2}\right) \ni f \mapsto A^{s} f=\sum_{k \in \mathbb{Z}^{2}}\left\langle f, e_{k}^{r}\right\rangle_{r} e_{k}^{r+s} \in H_{p}^{r+s}\left([0,1]^{2}\right)
$$

with $\left\langle f, e_{k}^{r}\right\rangle_{r}=\left(1+\|k\|^{2}\right)^{\frac{r}{2}} f_{k}$ and $\left(A^{s} f\right)_{k}=\left(1+\|k\|^{2}\right)^{-\frac{s}{2}} f_{k}$. Since $A^{s}$ is an isometry, i.e. $\left\|A^{s} f\right\|_{r+s}=$ $\|f\|_{r}$, the problem of recovering $f$ is well-posed. Indeed, $f$ can be reconstructed by $f_{k}=\left(A^{s} f\right)_{k}(1+$


Figure 8: Experimental results for $\varepsilon_{r e l}=5 \%: f_{10,40,100}^{\delta}(\mathrm{t} .1),. f_{10,40,100}^{0, \delta}(\mathrm{t} . \mathrm{m}),. f_{10,40,100}^{0.000025, \delta}$ (t.r.), and for $\varepsilon_{\text {rel }}=10 \%: f_{10,40,100}^{\delta}$ (b.l.), $f_{10,40,100}^{0, \delta}$ (b.m.), $f_{10,40,100}^{0.0001, \delta}$ (b.r.). In all subfigures the dashed line (- -) represents the true solution $f$.
$\left.\|k\|^{2}\right)^{\frac{s}{2}}$. However, in practice we are typically confronted with noisy data. Therefore, it is meaningful to consider $A^{s}$ not as an operator between $H_{p}^{r}\left([0,1]^{2}\right)$ and $H_{p}^{r+s}\left([0,1]^{2}\right)$ but rather between $H_{p}^{r}\left([0,1]^{2}\right)$ and $L_{2, p}\left([0,1]^{2}\right)$. This can be adequately modeled by introducing the so-called Sobolev embedding operator, see for a detailed discussion [24],

$$
\iota: H_{p}^{r+s}\left([0,1]^{2}\right) \rightarrow L_{2, p}\left([0,1]^{2}\right)
$$

and by considering the composition $\iota \circ A^{s}$. The embedding operator is compact and has the following singular system,

$$
\sigma_{k}=\left(1+\|k\|^{2}\right)^{-\frac{r+s}{2}}, v_{k}=\left(1+\|k\|^{2}\right)^{-\frac{r+s}{2}} e^{2 \pi i k \cdot}=e_{k}^{r+s}, \text { and } u_{k}=e^{2 \pi i k \cdot}=e_{k}^{0}
$$

Observing that $\left\langle A^{s} f, v_{k}\right\rangle_{r+s}=\left(1+\|k\|^{2}\right)^{r / 2} f_{k}$, we immediately obtain the obvious identities

$$
\begin{align*}
\iota \circ A^{s} f & =\sum_{k \in \mathbb{Z}^{2}} \sigma_{k}\left\langle A^{s} f, v_{k}\right\rangle_{r+s} u_{k}\left(=\sum_{k \in \mathbb{Z}^{2}}\left\langle f, e_{k}^{r}\right\rangle_{r} e_{k}^{r+s}=A^{s} f\right)  \tag{4.19}\\
& =\sum_{k \in \mathbb{Z}^{2}} \sigma_{k}\left\langle f, e_{k}^{r}\right\rangle_{r} u_{k} \tag{4.20}
\end{align*}
$$

Note, (4.19) is considered as an $L_{2}$ function. This implies that a possible further sampling takes places in $L_{2, p}\left([0,1]^{2}\right)$. In principle, we have now different strategies to recover $f$. One straightforward approach is suggested by the singular value representation of $\iota \circ A^{s} f$ in (4.19) allowing to apply one of the methods proposed in Section 4.1 or 4.2 to recover $A^{s} f$ as a function in $H_{p}^{r+s}\left([0,1]^{2}\right)$. Then by inverting $A^{s}$ through $f_{k}=\left(A^{s} f\right)_{k}\left(1+\|k\|^{2}\right)^{\frac{s}{2}}$ we could easily derive $f$. But this, however, implies a recovery of $f$ by its Fourier coefficients and not in a preferred basis or frame. To overcome this possible deficiency we could again apply one of our proposed sampling-recovery methods and sample $A^{s} f$ in $H_{p}^{r+s}\left([0,1]^{2}\right)$ and recover $f$ in $H_{p}^{r}\left([0,1]^{2}\right)$ but now represented in a frame of our choice.

Another (more natural) approach, which we prefer, is given due to (4.20) describing the singular value decomposition of $\iota \circ A^{s}$ as an operator between $H_{p}^{r}\left([0,1]^{2}\right)$ and $L_{2, p}\left([0,1]^{2}\right)$. This equation directly allows us to reconstruct $f$.

In our particular example we choose $r=0$, i.e.

$$
\iota \circ A^{s} f=\sum_{k \in \mathbb{Z}^{2}}\left(1+\|k\|^{2}\right)^{-\frac{s}{2}}\left\langle f, e^{2 \pi i k \cdot}\right\rangle_{0} e^{2 \pi i k}
$$

and specialize $f$ to

$$
f(x, y)=1_{\{0.25 \leq x<0.5,0.25 \leq y<0.75\}}(x, y)
$$



Figure 9: (From top left to bottom right) original function $f$, recovery $f_{n, m, r}^{\alpha, \delta}$, difference $f-f_{n, m, r}^{\alpha, \delta}$, image $\iota \circ A^{s} f$, reconstruction $\left(A^{s} f\right)_{m, r}^{\delta}=g_{m, r}^{\delta}$, difference $\iota \circ A^{s} f-g_{m, r}^{\delta}$, and in the last row, the nonzero entries of $S_{r}^{*} U_{m}$ and $W_{n}^{*} V_{m} V_{m}^{*} W_{n}$.
and $A^{s}$ to $s=1$. Both $f$ and $\iota \circ A^{s} f$ are visualized in Figure 9. As recovery system we choose the Haar basis on $[0,1]^{2}$,

$$
w_{\lambda}(x, y)=\psi_{\lambda}(x, y), \lambda=\lambda\left(q, j, k_{1}, k_{2}\right) \in \Lambda
$$

where $q=0,1,2,3$ prescribes the species of the wavelet ( $q=0$ - generator, $m=1,2,3$ - corresponding wavelets), $j \in \mathbb{Z}$ the scales, and $\left(k_{1}, k_{2}\right) \in I$ the translations. In our numerical experiment, we limit the scales to $-4 \leq j \leq 0$ resulting in a grid $\Lambda$ of size $n=1024$. As sampling system we choose the Fourier system on $[0,1]^{2}$,

$$
s_{k}(x, y)=e^{2 \pi i\left(x k_{1}+y k_{2}\right)},\left(k_{1}, k_{2}\right) \in \mathbb{Z}^{2} .
$$

We restrict the frequencies to $-25 \leq k_{1}, k_{2} \leq 25$ resulting in grid of size $r=(2 \cdot 25+1)^{2}=2601$. Finally, we have to fix the length of the singular value expansion. Similar as for the sampling system we choose here also $m=2601$. Due to the very simple choice for $\left\{w_{\lambda}\right\}_{\lambda \in \Lambda}$ and $\left\{s_{k}\right\}_{k \in \mathbb{Z}^{2}}$, the entries of $S_{r}^{*} U_{m}$ and $V_{m}^{*} W_{n}$ can be derived explicitly (the nonzero entries are visualized in Figure 9). An approximation to $f$ is now obtained through

$$
f_{n, m, r}^{\alpha, \delta}=W_{n}\left(W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Theta_{\alpha, m}^{-1} V_{m}^{*} W_{n}\right)^{\dagger} W_{n}^{*} V_{m} \Theta_{\alpha, m}^{-1} \Sigma_{m}\left(U_{m}^{*} S_{r} S_{r}^{*} U_{m}\right)^{\dagger} U_{m}^{*} S_{r} S_{r}^{*} g^{\delta}
$$

and is visualized 9. In our particular example the relative error is $\varepsilon_{r e l}=5 \%$ and the corresponding Tikhonov stabilization is fine tuned by $\alpha=0.00001$.

## 5 Conclusions and challenges

The purpose of this paper was to extend generalized sampling to sampling and reconstructions in frames, and to inverse and ill-posed problems. The key component is to allow the various parameters (number of samples etc) to vary independently. When done appropriately, it is then possible to prove sampling theorems which give guaranteed error bounds for the reconstruction in terms of constants that can be estimated numerically.

An important theme of this work, as discussed earlier in the paper, is that finite-dimensional tools are not always well suited for formally infinite-dimensional problems. As such, this paper marks only one step in a much larger project on how to reconstruct in infinite dimensions, with many directions for future work. Several of these directions are now described.

All the key theorems in this paper rely on the existence of quantities (e.g. $K_{n, m}$ or $K_{n, m}^{1}, K_{m, r}^{2}, K_{m, r}^{3}$ ) which can be made small by an appropriate choice of parameters. In Section 3.4 an insight was given as to how to numerically estimate the quantity $K_{n, m}$ in the case of generalized sampling with frames, and thus how to compute the so-called stable sampling rate. It is a topic of ongoing investigation to extend this notion to the GS framework for inverse problems of Section 4. Note that the setting here is substantially more complicated - we have to deal with additional parameters, including possibly the regularization $\alpha$. However, the aim of future work is to introduce explicitly computable sampling rates for such problems, which allow one to determine precisely how to balance such parameters. Much of this will be based on the more recent developments of [3] for GS with bases.

An important ill-posed problem is that of X-Ray tomography, which involves inverting the Radon transform. Although the GS framework introduced in Section 4, based on regularization of a compact operator, can be applied to this problem, this is not the only possible approach. Indeed, the Radon transform has a number of key features than one may hope to incorporate into the GS framework in order to get a better reconstruction for this particular problem. This is also a topic of future investigations.

As mentioned at the beginning of this paper, GS can be combined with convex optimization tools to give new techniques, and a completely new theory, for compressed sensing in infinite dimensions $[1,6]$. One may therefore ask, is it possible to incorporate the developments in this paper to allow for sparse recovery for inverse and ill-posed problems? In the near future we intend to investigate this question, with one aim being to extend the previous work of [19] on this topic.

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