# Discrete Hodge Helmholtz DECOMPOSITION 

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#### Abstract

This paper presents a new method using spectral approaches to compute the Discrete Hodge Helmholtz Decomposition (DHHD) of a given vector field. This decomposition consists in extracting the solenoidal (i.e. divergence-free), the non-solenoidal (i.e. rotational-free or, gradient of a scalar field) and the harmonic components (that is divergence-free and rotational-free) of a this vector field. A test case illustrates the proposed method.


Keywords: Discrete Hodge Helmholtz decomposition, spectral element method.
AMS classification: 76D05, 65N35.

## §1. Introduction

The Hodge Helmholtz decomposition of a general vector field $\mathbf{u}=\mathbf{u}_{\psi}+\mathbf{u}_{\phi}+\mathbf{u}_{H}$ is a classical problem in applied and computational physics [11, 14]. Application areas include (among others) electromagnetism, linear elasticity, fluid mechanics, image and video processing. A closed form of this decomposition may be obtained for unbounded domains through BiotSavart type integrals. In finite domains, however such an approach is no longer feasible and computational solutions are the only practical way to perform this decomposition.

With regard to elasticity, work by Brezzi and Fortin [7], and by Arnold and Falk [3] used the Hodge Helmholtz decomposition theorem for the study of the Reissner-Mindlin plate model. With regard to incompressible fluid flows, the scalar potential $\phi$ such that $\mathbf{u}_{\phi}=\nabla \phi$ in the Hodge Helmholtz decomposition is usually related to the pressure field $p$, and the vector potential $\mathbf{u}_{\psi}$ corresponds to the solenoidal velocity field $\mathbf{u}_{S}$ both quantities being involved in the Navier-Stokes equations [6]. Stokes and Navier-Stokes solvers decouple most of the time the computation of the velocity and pressure fields [10]. The family of correction-pressure time splitting methods [12] generates first a tentative velocity field that is not incompressible but contains the right vorticity. The addition of a pressure gradient to this temporary velocity (equivalent to Hodge Helmholtz decomposition) makes it divergence-free [9, 19]. Another approach resorts to pressure penalization [8]. In video processing, the Hodge Helmholtz decomposition allows to detect the fingerprint reference or hurricanes from satellite pictures [16]. In [1] the authors propose a constructive spectral approaches for the Helmholtz decomposition of a Vector Field which consists in projecting the field to be decomposed on the kernel and the ranges of the $-\operatorname{grad}(\mathrm{div})$ operator. Recently, in [17], the authors proposed a meshless approach for the Hodge Helmholtz decomposition while in [13], divergence-free and curl-free wavelets are used.

## §2. Discrete Hodge Helmholtz Decomposition

Consider a given vector field $\mathbf{u}$ defined in some domain $\Omega$ with boundary $\partial \Omega$. The Helmholtz decomposition writes (see [11, Theorem 3.2, p. 40])

$$
\begin{equation*}
\mathbf{u}=\mathbf{u}_{\psi}+\mathbf{u}_{\phi}+\mathbf{u}_{H} \tag{1}
\end{equation*}
$$

The solenoidal component $\mathbf{u}_{\psi}$ satisfies the equations

$$
\left\{\begin{array}{l}
\nabla \cdot \mathbf{u}_{\psi}=0 \text { in } \Omega,  \tag{2}\\
\mathbf{u}_{\psi} \cdot \mathbf{n}=0 \text { on } \partial \Omega,
\end{array}\right.
$$

while the irrotational complement $\mathbf{u}_{\phi}$ is such that

$$
\left\{\begin{array}{l}
\nabla \times \mathbf{u}_{\phi}=\mathbf{0} \text { in } \Omega,  \tag{3}\\
\mathbf{u}_{\phi} \times \mathbf{n}=\mathbf{0} \text { on } \partial \Omega .
\end{array}\right.
$$

Finally, the harmonic component is both solenoidal and irrotational:

$$
\begin{equation*}
\nabla \cdot \mathbf{u}_{H}=0 \quad \text { and } \quad \nabla \times \mathbf{u}_{H}=\mathbf{0} . \tag{4}
\end{equation*}
$$

The main difficulty of the problem (1) consists in satisfying the solenoidal (2) and irrotational constraints (3). The method we present in this paper is based on the construction of a basis satisfying the expected constraints. Its originality lies in the way these bases are built. To illustrate this, we will focus our presentation on how we derive a divergence-free basis. The derivation of the rotational-free basis in 2D will be presented in Section 4. This paper only considers the 2D case and then it is good to distinguish between curl and rot, where

$$
\nabla \times \mathbf{u}=\operatorname{rot} \mathbf{u}=\frac{\partial u_{2}}{\partial x_{1}}-\frac{\partial u_{1}}{\partial x_{2}}, \quad \operatorname{curl} \phi=\binom{\partial \phi / \partial x_{2}}{-\partial \phi / \partial x_{1}}
$$

## §3. Computation of the solenoidal component

In order to state the problem in variational form we introduce the relevant spaces of functions:

$$
\begin{aligned}
H(\operatorname{div}, \Omega) & =\left\{\mathbf{w} \in\left(L^{2}(\Omega)\right)^{2} \mid \nabla \cdot \mathbf{w} \in L^{2}(\Omega)\right\}, \\
L_{0}^{2}(\Omega) & =\left\{q \in L^{2}(\Omega) \mid \int_{\Omega} q d x=0\right\} .
\end{aligned}
$$

Let $\mathbf{v}, \mathbf{w} \in H(\operatorname{div}, \Omega)$. We define the inner product

$$
\begin{equation*}
(\mathbf{v}, \mathbf{w})_{H(\operatorname{div}, \Omega)}=(\mathbf{v}, \mathbf{w})_{\left(L^{2}(\Omega)\right)^{2}}+(\nabla \cdot \mathbf{v}, \nabla \cdot \mathbf{w})_{L^{2}(\Omega)}, \tag{5}
\end{equation*}
$$

and associated norm $\|\mathbf{w}\|_{H(\operatorname{div}, \Omega)}=\left(\|\mathbf{w}\|_{\left(L^{2}(\Omega)\right)^{2}}^{2}+\|\nabla \cdot \mathbf{w}\|_{L^{2}(\Omega)}^{2}\right)^{1 / 2}$.
Consider also the proper subspace $H_{0}(\operatorname{div}, \Omega) \subset H(\operatorname{div}, \Omega)$ :

$$
H_{0}(\operatorname{div}, \Omega)=\{\mathbf{w} \in H(\operatorname{div}, \Omega) \mid \mathbf{w} \cdot \mathbf{n}=0 \text { on } \partial \Omega\}
$$

The admissible space for $\mathbf{u}_{\psi}$ in Problem (1)-(4) is a subspace of $H_{0}(\operatorname{div}, \Omega)$ :

$$
X=\left\{\mathbf{u} \in H_{0}(\operatorname{div}, \Omega) \mid \boldsymbol{\nabla} \cdot \mathbf{u}=0 \text { in } \Omega\right\} .
$$

With the inner product (5) we can define the space $X^{\perp}$ by

$$
X^{\perp}=\left\{\mathbf{u} \in H_{0}(\operatorname{div}, \Omega) \mid(\mathbf{u}, \mathbf{v})_{H(\operatorname{div}, \Omega)}=0, \forall \mathbf{v} \in X\right\}
$$

The space $X$ and $X^{\perp}$ allow us to decompose each element in $\mathbf{u} \in H_{0}(\operatorname{div}, \Omega)$ uniquely as

$$
\mathbf{u}=\mathbf{u}_{X}+\mathbf{u}_{X^{\perp}}, \quad \mathbf{u}_{X} \in X, \quad \mathbf{u}_{X^{\perp}} \in X^{\perp}
$$

The divergence operator is a surjective map from $H_{0}(\operatorname{div}, \Omega)$ onto $L_{0}^{2}(\Omega)$ and it is a bijection from $X^{\perp}$ to $L_{0}^{2}(\Omega)$ (see [11, Corollary 2.4]). This means that, for each $\mathbf{w} \in H_{0}(\operatorname{div}, \Omega)$, there exists a unique $\mathbf{v} \in X^{\perp}$ such that $\operatorname{div} \mathbf{v}=\operatorname{div} \mathbf{w} \in L_{0}^{2}(\Omega)$. In particular, $\operatorname{div}^{-1}$ is a welldefined continuous map from $L_{0}^{2}(\Omega)$ to $X^{\perp}$ and we have the Poincaré inequality [4, p.302]

$$
\|\mathbf{v}\|_{H(\operatorname{div}, \Omega)} \leq c_{P}\|\operatorname{div} \mathbf{v}\|_{L^{2}(\Omega)}
$$

where $c_{P}$ is independent of $\mathbf{v}$. This implies that, if we have a set of linearly independent functions $\left\{q_{1}, \ldots, q_{n}\right\} \in L_{0}^{2}(\Omega)$, then there exist unique vectors $\mathbf{v}_{i} \in X^{\perp}$ with $\operatorname{div} \mathbf{v}_{i}=q_{i}$ and these vectors $\mathbf{v}_{i}$ are linearly independent in $X^{\perp}$ and therefore also in $H_{0}(\operatorname{div}, \Omega)$. The converse only holds for $X^{\perp}$. Indeed, if $\mathbf{u}$ and $\mathbf{v}$ are linearly independent in $X^{\perp}$, then $\operatorname{div} \mathbf{u}$ and $\operatorname{div} \mathbf{v}$ are linearly independent in $L_{0}^{2}(\Omega)$. The generalization for $\mathbf{u}, \mathbf{v} \in H_{0}(\operatorname{div}, \Omega)$ is generally not true: Let $S=\left\{\mathbf{v}_{1}, \ldots, \mathbf{v}_{n}\right\}$ be a set of linear independent vectors in $H_{0}(\operatorname{div}, \Omega)$ with $\operatorname{div} \mathbf{v}_{i} \neq 0$ and let $\operatorname{rank}\left(\operatorname{div} \mathbf{v}_{1}, \ldots, \operatorname{div} \mathbf{v}_{n}\right)=k$, then we can select $k$ vectors from the set $S$ which form a basis for a $k$-dimensional subspace of $X^{\perp}$.

This last result is the keystone to ensure that it is always possible to reduce into a square and invertible sub-part, the algebraic system resulting from any stable discretization of the constraint.

### 3.1. Variational formulation and its discretization

The variational formulation of problem (1) writes: Find $\mathbf{u}_{\psi} \in X$ such that

$$
\begin{equation*}
\int_{\Omega} \mathbf{u}_{\psi} \cdot \mathbf{v} d \mathbf{x}=\int_{\Omega} \mathbf{u} \cdot \mathbf{v} d \mathbf{x}, \quad \forall \mathbf{v} \in X \tag{6}
\end{equation*}
$$

Due to the nature of $X, \mathbf{u}_{\phi}$ and $\mathbf{u}_{H}$ disappear.
We firstly introduce the Raviart-Thomas space [18]

$$
\begin{equation*}
R_{p}=\left(\mathbb{P}_{p}^{0}(\Lambda) \otimes \mathbb{P}_{p-1}(\Lambda)\right) \times\left(\mathbb{P}_{p-1}(\Lambda) \otimes \mathbb{P}_{p}^{0}(\Lambda)\right) \tag{7}
\end{equation*}
$$

where $\mathbb{P}_{N}(\Lambda)$ is the space of polynomials with degree $\leq N$ and $\mathbb{P}_{p}^{0}(\Lambda)$ denotes the space of polynomials of degree $p$ vanishing on $\pm 1$. The dimension of $R_{p}$ is equal to $2(p-1) p$.

The solution is approximated by $\mathbf{u}_{\psi, p}=\left(u_{\psi, p}^{x}, u_{\psi, p}^{y}\right)$ in $R_{p}$ with

$$
\begin{aligned}
& u_{\psi, p}^{x}(x, y)=\sum_{i=1}^{p-1} \sum_{j=1}^{p} u_{\psi, p}^{x}\left(\xi_{i}, \zeta_{j}\right) h_{i}(x) \tilde{h}_{j}(y), \\
& u_{\psi, p}^{y}(x, y)=\sum_{i=1}^{p} \sum_{j=1}^{p-1} u_{\psi, p}^{y}\left(\zeta_{i}, \xi_{j}\right) \tilde{h}_{i}(x) h_{j}(y) .
\end{aligned}
$$

Let $\Sigma_{G L L}=\left\{\left(\xi_{i}, \rho_{i}\right) \mid 0 \leq i \leq N\right\}$ and $\Sigma_{G L}=\left\{\left(\zeta_{i}, \omega_{i}\right) \mid 1 \leq i \leq N\right\}$ denote respectively the sets of Gauss-Lobatto-Legendre and Gauss-Legendre quadrature nodes and weights (see [10]). Likewise, $h_{i}(x) \in \mathbb{P}_{N}(\Lambda)$ and $\widetilde{h}_{j}(x) \in \mathbb{P}_{N-1}(\Lambda)$ are respectively the canonical Lagrange polynomial interpolation basis built on $\Sigma_{G L L}$ and on $\Sigma_{G L}$.

With this choice, the divergence of $\mathbf{u}_{\psi, p}$ is a polynomial of degree $p-1$. Consequently, if the divergence is orthogonal to all polynomial of $\mathbb{P}_{p-1}(\Omega)$, it is necessarily equal to 0 . This point gives a new characterization for $X_{p}=R_{p} \cap X$ :

$$
\begin{aligned}
& X_{p}=\left\{\mathbf{u}_{\psi, p} \in R_{p} \left\lvert\, \frac{\partial u_{\psi, p}^{x}}{\partial x}+\frac{\partial u_{\psi, p}^{y}}{\partial y}=0\right.\right\}, \\
& X_{p}=\left\{\mathbf{u}_{\psi, p} \in R_{p} \left\lvert\, \int_{\Omega}\left(\frac{\partial u_{\psi, p}^{x}}{\partial x}+\frac{\partial u_{\psi, p}^{y}}{\partial y}\right) q d \mathbf{x}=0\right., \forall q \in \mathbb{P}_{p-1}(\Omega)\right\} .
\end{aligned}
$$

### 3.2. A basis for $X_{p}$

The first step consists in determining the dimension of the space $X_{p}$, that we denote by $N_{1}$ :

$$
N_{1}=\operatorname{dim} R_{p}-N_{2}=2(p-1) p-N_{2},
$$

where $N_{2}$ is the dimension of the range of the divergence operator which is also the dimension of $P_{p-1}(\Omega) \cap L_{0}^{2}(\Omega)$ and then equals to $p^{2}-1$. We deduce that $N_{1}=(p-1)^{2}$.
Remark 1. One can also view $N_{2}$ as the number of necessary and sufficient equations to ensure $\nabla \cdot \mathbf{u}_{\psi, p} \equiv 0$ :

$$
\nabla \cdot \mathbf{u}_{\psi, p} \in \mathbb{P}_{p-1}(\Omega) \quad \text { so } \quad N_{2} \leq p^{2}
$$

Due to the boundary conditions (here $\mathbf{u}_{\psi, p} \cdot \mathbf{n}=0$ on $\partial \Omega$ ), there is a dependent equation in the two dimensional case [5], since

$$
\int_{\Omega} \nabla \cdot \mathbf{u}_{\psi, p} L_{0}(x) L_{0}(y) d \mathbf{x}=0, \forall \mathbf{u}_{\psi, p} \in R_{p}
$$

Indeed, the polynomial $L_{0}(x) L_{0}(y)$ is a spurious mode and it reduces the number of independent equations from $p^{2}$ to $p^{2}-1$. Consequently, we require $N_{2}=p^{2}-1$ test functions $q$ to ensure $\int_{\Omega} \nabla \cdot \mathbf{u}_{\psi, p} q d \mathbf{x}=0$.

Once the dimension is known, we describe now how to proceed to derive a divergence free basis from any $N_{1}$ given vectors of $R_{p}$


Figure 1: Algebraic system.


Figure 2: Decomposition of $D$.


Figure 3: Algebraic system.

### 3.2.1. Algebraic characterization of $X_{p}$

Let $\mathbf{u}_{\psi, p}$ be in $X_{p}$. The divergence of $\mathbf{u}_{\psi, p}$ is orthogonal to $p^{2}-1$ polynomials of degree $p-1$. It is equivalent to saying that the divergence of $\mathbf{u}_{\psi, p}$ nullifies into $p^{2}-1$ Gauss points. The algebraic divergence equation writes $D \mathbf{u}_{\psi, p}=\mathbf{0}$, where $D$ is a rectangular matrix with $p^{2}-1$ rows and $2 p(p-1)$ columns (see Figure 1).

One splits $D$ into $D_{1} \oplus D_{2}$ and $\mathbf{u}_{\psi, p}$ into $\mathbf{u}_{\psi, p}^{1} \oplus \mathbf{u}_{\psi, p}^{2}$. The vector $\mathbf{u}_{\psi, p}^{1}$ contains $N_{1}=(p-1)^{2}$ values of $\mathbf{u}_{\psi, p}$, whereas $\mathbf{u}_{\psi, p}^{2}$ contains the $N_{2}=p^{2}-1$ remaining values (see Figure 2). The equation $D \mathbf{u}_{\psi, p}=\mathbf{0}$ becomes $D_{1} \mathbf{u}_{\psi, p}^{1}+D_{2} \mathbf{u}_{\psi, p}^{2}=\mathbf{0}$, as shown in Figure 3 .

Since, the $p^{2}-1$ rows of $D$ are independent, there exists at least one choice of matrix $D_{2}$ invertible and the system leads to a relation between $\mathbf{u}_{\psi, p}^{2}$ and $\mathbf{u}_{\psi, p}^{1}$ :

$$
\begin{equation*}
\mathbf{u}_{\psi, p}^{2}=-D_{2}^{-1} D_{1} \mathbf{u}_{\psi, p}^{1} \tag{8}
\end{equation*}
$$

This equation is very important since it means that, if we have any part $\mathbf{u}_{\psi, p}^{1}$ of $\mathbf{u}_{\psi, p}$, we can build the complementary $\mathbf{u}_{\psi, p}^{2}$ such that divergence of $\mathbf{u}_{\psi, p}$ equals 0 . This argument allows us to build a basis of $X_{p}$.

### 3.2.2. Basis of $X_{p}$

The technique we use to project any vector of $R_{p}$ on $X_{p}$ is in the spirit of that published in [2] and used to solve the Stokes problem.

We consider $\mathbf{v}_{p} \in R_{p}$. Our strategy consists in combining implicitly:

- A reduction from $\mathbf{v}_{p}$ to $\mathbf{v}_{p}^{1}$.
- An extension from $\mathbf{v}_{p}^{1}$ to $\mathbf{w}_{p}=\left(\mathbf{v}_{p}^{1}, \mathbf{v}_{p}^{2}\right)$ such that $\nabla \cdot \mathbf{w}_{p}=0$ ensured by the multiplication
of $\mathbf{v}_{p}^{1}$ by the matrix

$$
M=\left[\begin{array}{c}
I_{N_{1}} \\
-D_{2}^{-1} D_{1}
\end{array}\right]
$$

The first block of $M$ is the identity matrix of order $N_{1}$. The second block contains $N_{2}$ rows and $N_{1}$ columns. It ensures the passage from $\mathbf{v}_{p}^{1}$ to $\mathbf{v}_{p}^{2}$.

- For each $\mathbf{v}_{p} \in R_{p}$ one associates a vector $\mathbf{w}_{p}$ of $X_{p}$.

By consequent, our strategy for the construction of a basis of $X_{p}$ consists in:

- Choosing $N_{1}=(p-1)^{2}$ vectors $\left(\mathbf{v}_{p}^{k}\right)_{k=1 . . N_{1}}$ of the basis of $R_{p}$.
- For each one of these $N_{1}$ vectors, we consider its $N_{1}$-size reduced part denoted by $\mathbf{v}_{p}^{k, 1}$.
- We carry out the divergence-free extension $\left(\mathbf{w}_{p}^{k}\right)_{k=1 . . N_{1}}=\left(M \mathbf{v}_{p}^{k, 1}\right)_{k=1 . . N_{1}}$.

The $\left(\mathbf{w}_{p}^{k}\right)_{k=1 . . N_{1}}$ family is a basis of $X_{p}$. Consequently, $\mathbf{u}_{\psi, p} \in X_{p}$ can be decomposed according to $\mathbf{u}_{\psi, p}=\sum_{k=1}^{N_{1}} \alpha_{k} \mathbf{w}_{p}^{k}$ and the discrete variational formulation similar to (6) writes: Find $\mathbf{u}_{\psi, p} \in X_{p}$ such that

$$
\sum_{k=1}^{N_{1}}\left(\mathbf{w}_{p}^{k}, \mathbf{w}_{p}^{i}\right)_{p} \alpha_{k}=\left(\mathbf{u}_{p}, \mathbf{w}_{p}^{i}\right)_{p} .
$$

This can be written as

$$
\underline{\mathcal{M}} \underline{\alpha}=\underline{\mathcal{F}}
$$

with, for $1 \leq i, k \leq N_{1}$,

$$
\begin{aligned}
\underline{\mathcal{M}}_{i k} & =\left(\mathbf{w}_{p}^{k}, \mathbf{w}_{p}^{i}\right)_{p}, \\
\underline{\mathcal{F}}_{k} & =\left(\mathbf{u}, \mathbf{w}_{p}^{k}\right)_{p} .
\end{aligned}
$$

Finally, this system is equivalent to

$$
M^{T} B M \mathbf{u}_{\psi, p}^{1}=M^{T} B \mathbf{u}_{p}
$$

where B refers to the classical mass matrix computed using the $\left(h_{i} \times \tilde{h}_{j}\right) \otimes\left(\tilde{h}_{i} \times h_{j}\right)$ basis.

## §4. Computation of the irrotational component

A similar strategy to that described previously is used to compute the irrotational component $\mathbf{u}_{\phi}$, so we will limit to the description of the outline and we will not give all the details of its implementation.

Firstly, we introduce three spaces of functions:

$$
\begin{aligned}
H(\mathrm{rot}, \Omega) & =\left\{\mathbf{w} \in\left(L^{2}(\Omega)\right)^{2} \mid \nabla \times \mathbf{w} \in L^{2}(\Omega)\right\}, \\
H_{0}(\mathrm{rot}, \Omega) & =\{\mathbf{w} \in H(\operatorname{rot}, \Omega) \mid \mathbf{w} \times \mathbf{n}=0 \text { on } \partial \Omega\}, \\
Y & =\left\{\mathbf{u} \in H_{0}(\operatorname{rot}, \Omega) \mid \nabla \times \mathbf{u}=0 \text { in } \Omega\right\}
\end{aligned}
$$

The variational formulation of Problem (1) writes: Find $\mathbf{u}_{\phi} \in Y$ such that

$$
\int_{\Omega} \mathbf{u}_{\phi} \cdot \mathbf{v} d \mathbf{x}=\int_{\Omega} \mathbf{u} \cdot \mathbf{v} d \mathbf{x}, \quad \forall \mathbf{v} \in Y
$$

For the discretization, we introduce the Nédélec space [15]

$$
N_{p}=\left(\mathbb{P}_{p-1}(\Lambda) \otimes \mathbb{P}_{p}^{0}(\Lambda)\right) \times\left(\mathbb{P}_{p}^{0}(\Lambda) \otimes \mathbb{P}_{p-1}(\Lambda)\right)
$$

The solution is approximated by $\mathbf{u}_{\phi, p}=\left(u_{\phi, p}^{x}, u_{\phi, p}^{y}\right) \in Y_{p}=N_{p} \cap Y$ with

$$
\begin{aligned}
& u_{\phi, p}^{x}(x, y)=\sum_{i=1}^{p} \sum_{j=1}^{p-1} u_{\phi, p}^{x}\left(\zeta_{i}, \xi_{j}\right) \tilde{h}_{i}(x) h_{j}(y), \\
& u_{\phi, p}^{y}(x, y)=\sum_{i=1}^{p-1} \sum_{j=1}^{p} u_{\phi, p}^{y}\left(\xi_{i}, \zeta_{j}\right) h_{i}(x) \tilde{h}_{j}(y) .
\end{aligned}
$$

As previously, we build a basis of $Y_{p}$. With the same reasoning as for $X_{p}$ and taking into account the same spurious mode $L_{0}(x) L_{0}(y)$, we determine the size of $Y_{p}$ equal to $N_{1}=$ $(p-1)^{2}$.

Then the constraint $\nabla \times \mathbf{u}_{\phi, p}=\mathbf{0}$ is written into $N_{1}=(p-1)^{2}$ Gauss points and gives an algebraic equations $R \mathbf{u}_{\phi, p}=\mathbf{0}$. A splitting strategy for the matrix $R$ into $R_{1} \oplus R_{2}$ and the vector $\mathbf{u}_{\phi, p}$ into $\mathbf{u}_{\phi, p}^{1} \oplus \mathbf{u}_{\phi, p}^{2}$ gives $R_{1} \mathbf{u}_{\phi, p}^{1}+R_{2} \mathbf{u}_{\phi, p}^{2}=\mathbf{0}$ and finally $\mathbf{u}_{\phi, p}^{2}=-R_{2}^{-1} R_{1} \mathbf{u}_{\phi, p}^{1}$.

Thanks to the $\left(N_{1}+N_{2}\right) \times N_{1}$ matrix

$$
N=\left[\begin{array}{c}
I_{N_{1}} \\
-R_{2}^{-1} R_{1}
\end{array}\right],
$$

we can construct a basis of $Y_{p}$.
Finally we obtain the system $N^{T} \tilde{B} N \mathbf{u}_{\phi, p}^{1}=N^{T} \tilde{B} \mathbf{u}_{p}$, where $\tilde{B}$ refers to the classical mass matrix computed using $\left(\tilde{h}_{i} \times h_{j}\right) \otimes\left(h_{i} \times \tilde{h}_{j}\right)$ basis.

## §5. Numerical results

To illustrate the efficiency of our approach for the Hodge Helmholtz decomposition, we have made a numerical experiment in the square $\Omega=(-1,+1)^{2}$ with the case $\mathbf{u}=\mathbf{u}_{\psi}+\mathbf{u}_{\phi}+\mathbf{u}_{H}$ corresponding to the following components:

$$
\begin{aligned}
& \mathbf{u}_{\psi}=(\sin (\pi x) \cos (\pi y),-\sin (\pi y) \cos (\pi x)), \\
& \mathbf{u}_{\phi}=(\sin (\pi y) \cos (\pi x), \sin (\pi x) \cos (\pi y)), \\
& \mathbf{u}_{H}=(0.5,-1)
\end{aligned}
$$

The component $\mathbf{u}_{\psi, p}$ is approximated as outlined in Section 3, while the irrotational part $\mathbf{u}_{\phi, p}$ is computed as outlined in Section 4. Finally, $\mathbf{u}_{H}$ is calculated by the relation $\mathbf{u}_{H}=\mathbf{u}-\mathbf{u}_{\psi}-\mathbf{u}_{\phi}$.

Table 1 gives $\left\|\nabla \cdot \mathbf{u}_{\psi, p}\right\|_{L^{2}(\Omega)}$ and $\left\|\nabla \times \mathbf{u}_{\phi, p}\right\|_{L^{2}(\Omega)}$ as a function of the polynomial degree $p$. As expected, the norms remain close to round-off independently of $p$.

Figure 4 exhibits on a semi-logarithmic scale the $\left(L^{2}(\Omega)\right)^{2}$-norm of the error for the three components as a function of the polynomial degree $p$. We can observe a spectral decrease.

Figure 5 displays the Hodge Helmholtz decomposition of the vector u.

| $p$ | 4 | 8 | 12 | 16 |
| :---: | :---: | :---: | :---: | :---: |
| $\left\\|\nabla \cdot \mathbf{u}_{\psi, p}\right\\|_{L^{2}(\Omega)}$ | $8.56 \times 10^{-16}$ | $1.71 \times 10^{-15}$ | $1.95 \times 10^{-15}$ | $5.23 \times 10^{-15}$ |
| $\left\\|\nabla \times \mathbf{u}_{\phi, p}\right\\|_{L^{2}(\Omega)}$ | $4.90 \times 10^{-16}$ | $2.18 \times 10^{-15}$ | $3.87 \times 10^{-15}$ | $5.09 \times 10^{-15}$ |

Table 1: $L^{2}(\Omega)$-norm of the divergence and the rotational of solenoidal and irrotational components.


Figure 4: $\left(L^{2}(\Omega)\right)^{2}$-norm of the error as a function of the polynomial degree $p$.


Figure 5: Decomposition of the vector and its three components.

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