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Citation for published version (APA):

Hou, Y., & Li, K. (1999). Post Galerkin method for the Navier-Stokes equations. (RANA: reports on applied and numerical analysis; Vol. 9917). Technische Universiteit Eindhoven.

#### Document status and date:

Published: 01/01/1999

#### Document Version:

Publisher's PDF, also known as Version of Record (includes final page, issue and volume numbers)

#### Please check the document version of this publication:

- A submitted manuscript is the version of the article upon submission and before peer-review. There can be important differences between the submitted version and the official published version of record. People interested in the research are advised to contact the author for the final version of the publication, or visit the DOI to the publisher's website.
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Download date: 16. Nov. 2023

# Post Galerkin Method for the Navier-Stokes Equations \*

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Abstract: A kind of post Galerkin method based on the virtue of inertial manifold and approximate inertial manifold for the two dimensional Navier-Stokes equations is constructed in this paper. This kind of post Galerkin method also leads to a kind of new construction of approximate inertial manifold. We investigate the property of this manifold and derive the error estimation of our scheme. According to our method, one can get a much more accuracy approximate solution at any time once the standard Galerkin approximate solution is at hand. Obviously, this method will yield a significant gain in computing time.

Key Words: Approximate inertial manifold, Galerkin method, error estimation, Navier-Stokes equations

AMS Subject Classification 65M15, 65M70, 76D05, 35Q30

## 1 Motivation

Although the computing facilities improved in the last decades, directly simulating the Navier-Stokes Equations (NSE) still remains an open problem because of its large computing scale and long time integrations. Therefore, how to construct high effective and high accuracy numerical scheme is still an important and practical problem attracting people. Many authors derived new techniques and methods. For example, Lin Qun[1], W. Layton[5] and J. Xu[6] used extrapolation and two level meshes respectively. Especially, it is worth mentioning the applications of Inertial Manifolds (IMs) and Approximate Inertial Manifolds (AIMs) theory which were firstly introduced in 1988 by C. Foias, G. R. Sell, R. Temam[7] and C. Foias, O. Manley, R. Temam[8]. Based upon the finite dimensional behavior of the solutions, they show that there must be at least some approximate interactive rules between large eddy components and small eddy components of the solutions of many dissipative partial differential equations. From then on, many papers were contributed to this subject on constructing related new algorithms, that is all sort of nonlinear Galerkin methods, and their numerical analysis. For example, we refer readers to [9], [10], [11], [12], [13] and references therein.

Suppose H be a Hilbert space and  $H = H_m \oplus \hat{H}$  with  $\dim(H_m) = m < +\infty$ ,  $u(t) \in H$  be the solution of two dimensional NSE. Decomposing u(t) as

$$u(t) = p(t) + q(t)$$
, with  $p(t) \in H_m$ ,  $q(t) \in \hat{H}$ ,

AIMs believes that there must be some approximate interactive rule  $\Phi: H_m \to \hat{H}$  such that  $q(t) \approx \Phi(p(t))$ . Then its related nonlinear Galerkin method aims to search the approximate solution of u in form of  $\tilde{u}_m = \tilde{p} + \tilde{q}$  with  $\tilde{q} = \Phi(\tilde{p})$  such that it can generate a more accuracy approximation of u than that of Galerkin approximation  $u_m$ . In fact, for some positive sequence

<sup>\*</sup>Project supported by NSF of CHINA and State Key Major Basic Research Project

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 $\{\lambda_m\}_{m\in\mathcal{N}}$  which tends to  $\infty$  as  $m\to\infty$ , general nonlinear Galerkin solution  $\tilde{u}_m$  admits

$$(1.1) |u(t) - \tilde{u}_m(t)|_{H_0^1} \le C(t)\lambda_{m+1}^{-1}.$$

And the Galerkin solution  $u_m$  satisfies

$$|u(t) - u_m(t)|_{H_0^1} \le C(t)\lambda_{m+1}^{-\frac{1}{2}}.$$

Here C(t) is some positive constant depending on various data. Obviously, nonlinear Galerkin method can greatly improve the convergence rate of Galerkin method. That is to say we could get more accuracy approximation of u(t) with lower computing price compared with Galerkin method. But its defects is also obvious:

- 1 At each time step, nonlinear Galerkin method must solve  $\tilde{p}$  and  $\tilde{q}$  simultaneously, that is, nonlinear Galerkin method can not obtain  $\tilde{p}$  without  $\tilde{q}$ . This leads to solving a coupled equations and increasing computing price.
- 2 Nonlinear Galerkin methods use large eddy component  $\tilde{p}$  to correct small eddy component  $\tilde{q}$ . Noticing  $\tilde{q} \in \hat{H}$  and  $\tilde{p} \in H_m$ , the final accuracy depends on both  $p \tilde{p}$  and  $q \tilde{q}$ . Therefore, the final accuracy can not exceed  $p \tilde{p}$ . This may restrict the high performance of nonlinear Galerkin method.
- **3** Nonlinear Galerkin solution  $\tilde{u}_m$  takes no information from Galerkin solution  $u_m$ .

Being aimed at the above shortages, our paper intends to find a new  $\Phi$  such that the related algorithms can overcome those defects.

## 2 NSE and Its Galerkin Approximation

Let us consider the following two dimensional NSE confined on a bounded domain  $\Omega \subset \mathcal{R}^2$ 

(2.1) 
$$\begin{cases} \frac{\partial u}{\partial t} - \nu \Delta u + (u \cdot \nabla)u + \nabla p = F, & (x, t) \in \Omega \times \mathcal{R}^+, \\ \nabla \cdot u = 0, & (x, t) \in \Omega \times \mathcal{R}^+, \\ u(x, 0) = a(x), & x \in \Omega, \\ u(x, t) = 0, & (x, t) \in \partial \Omega \times \mathcal{R}^+ \end{cases}$$

Here u is fluid velocity and a the initial velocity satisfying  $\nabla \cdot a = 0$ , p the pressure.  $\nu > 0$  stands for kinetic viscosity and F is external force which is assumed to be time independent. For the sake of simplicity, we also assume that  $\partial \Omega$  is of class  $C^2$ .

Now we introduce a Hilbert space

$$H = \{ u \in L^2(\Omega)^2, \nabla \cdot u = 0 \text{ in weak sense }, u \cdot n |_{\partial \Omega} = 0 \},$$

where n stands for unit out normal vector of  $\Omega$ . If we denote P the Leray orthogonal projection from  $L^2(\Omega)^2$  onto H, by projecting (2.1) onto H, we can obtain the abstract NSE

(2.2) 
$$\begin{cases} \frac{du}{dt} + \nu Au + B(u, u) = f, \\ u(0) = a, \end{cases}$$

where  $A = -P\Delta$  is Stokes operator,  $B(u, u) = P[(u \cdot \nabla)u]$  and f = PF. It is well known that A is an unbounded, self-adjoint and positive definite operator with compact inverse. Thus we have that there exists two sequences

$$0 < \lambda_1 \le \lambda_2 \le \cdots \le \lambda_n \le \cdots \to \infty$$
, and  $\phi_1, \phi_2, \cdots, \phi_n \cdots \in H$ ,

such that

$$A\phi_i = \lambda_i \phi_i, \quad \forall i \in \mathcal{N}.$$

At the same time we can define its powers  $A^{\alpha}$  for  $\alpha \in \mathcal{R}$ . In fact,

$$D(A^{\alpha}) = \{ v \in H : v = \sum_{j=1}^{+\infty} v_j \phi_j, \sum_{j=1}^{+\infty} \lambda_j^{2\alpha} |v_j|^2 < +\infty, v_j \in \mathcal{R} \}$$

is a closed subspace of  $H^{2\alpha}(\Omega)^2$  and  $|A^{\alpha}|$  is an equivalent norm of it at least for  $\alpha < \frac{5}{4}$ , where  $|\cdot|$  stands for the  $L^2-$  norm. In the rest, sometimes we denote  $V=D(A^{\frac{1}{2}})$ . In addition,  $-\nu A$  generates an analytic semigroup on H, denoted by  $\{e^{-\nu At}\}_{t\geq 0}$ , with following estimation

(2.3) 
$$||A^{\alpha}e^{-\nu At}||_{\mathcal{L}(H,H)} \le c_0(\nu t)^{-\alpha}e^{-\nu \delta t}, \quad t > 0, \ \alpha > 0,$$

where  $\delta > 0$  is a constant related only on A,  $c_0 > 0$  is a constant. For the sake of convenience, we always use  $c_j > 0, j \in \mathcal{N}$ , to denote constants which have different meaning in different places appearing in the analysis. We also introduce the notation  $||u||_t = \sup_{0 \le s \le t} |u(s)|$ , especially,

$$|||u||| = \sup_{s \in \mathcal{R}^+} |u(s)|.$$

Now for any  $m \in \mathcal{N}$ , we introduce an orthogonal projection  $P_m$ 

$$\forall v = \sum_{j=1}^{+\infty} v_j \phi_j \in H, \quad P_m v = \sum_{j=1}^m v_j \phi_j.$$

Meanwhile, we use  $Q_m$  to denote  $I - P_m$ . The following inequalities are classical and we just state them out.

$$|Q_m \phi| \le \lambda_{m+1}^{-\frac{1}{2}} |A^{\frac{1}{2}} \phi|, \quad \forall \phi \in V,$$

$$|\phi|_{L^{\infty}} \le L_m |A^{\frac{1}{2}}\phi|, \quad \forall \phi \in H_m,$$

where  $L_m \sim (1 + \ln \lambda_m)^{\frac{1}{2}}$ .

By using these two orthogonal projections, we could decompose H as

$$H = H_m \oplus \hat{H}, \quad H_m = P_m H, \hat{H} = Q_m H.$$

To numerically solve (2.2), we use following Galerkin scheme,

(2.4) 
$$\begin{cases} \frac{du_m}{dt} + \nu A u_m + P_m B(u_m, u_m) = P_m f, \\ u_m(0) = P_m a. \end{cases}$$

Here we use  $u_m(t)$  to represent the Galerkin approximation of u(t). As well known, its stability and convergence results are classical. For example, we have the following convergence results

$$(2.5) |u(t) - u_m(t)| \le C_1(t)\lambda_{m+1}^{-1}, |A^{\frac{1}{2}}(u(t) - u_m(t))| \le C_2(t)\lambda_{m+1}^{-\frac{1}{2}},$$

where  $C_1(t)$  and  $C_2(t)$  are positive constants depending on  $\nu$ , f and a.

How to construct more accuracy approximate solution of u(t) by taking advantages of known information of it, that is  $u_m(t)$ , is the main problem to be solved in the rest. Of course, nonlinear Galerkin method can partly solve the problem although it has some defects as we just said in section 1. Here, we aim to get a new method which can overcome those shortages as well as improve convergence rate by using the virtue of AIMs to construct a new finite dimensional mapping. We call it post Galerkin method. It should be able to correct  $u_m$  at any time without introducing any extra computation except for computing  $u_m$ .

## 3 Further Properties of $u_m$

In this section, we will give some properties of  $u_m$  which will be very important for our further discussion. First of all, let us give a kind of new decomposition of true solution u. In fact, after we get its Galerkin approximation  $u_m$ , the very nature decomposition is to decompose u(t) as  $u_m(t)$  and its residue  $\hat{u}(t) = u(t) - u_m(t)$ , that is

(3.1) 
$$u(t) = u_m(t) + \hat{u}(t).$$

For the convenience of stating, we identify  $u_m(t) \in H_m$  and  $\hat{u}(t) \in H$  as large and small eddy components respectively. Obviously, the Galerkin approximation  $u_m(t)$  reaches the large eddy component of true solution exactly. Then, how to approximate small eddy components is the only problem. According to the ideal of AIMs, we suppose there also exist some kind of approximate interactive rule between large and small eddies. That is, there should exist a finite dimensional mapping  $\Phi$  from  $H_m$  into H. Before we begin to construct it, we need some further property of  $u_m$ .

Subtracting (2.4) from (2.2), we can easily get

(3.2) 
$$\begin{cases} \frac{d\hat{u}}{dt} + \nu A \hat{u} + B(\hat{u}, u_m) + B(u_m, \hat{u}) + B(\hat{u}, \hat{u}) = Q_m [f - B(u_m, u_m)], \\ \hat{u}(0) = Q_m a. \end{cases}$$

This is a nonlinear evolutionary equation of  $\hat{u}(t)$ .

Now let us consider some properties of  $\hat{u}$ . To do so, we decompose  $\hat{u}$  as

$$\hat{u} = P_m \hat{u} + Q_m \hat{u} \stackrel{\triangle}{=} p + q.$$

Then p satisfies

(3.4) 
$$\begin{cases} \frac{dp}{dt} + \nu Ap + P_m B(u_m, p+q) + P_m B(p+q, u_m) + P_m B(p+q, p+q) = 0, \\ p(0) = 0. \end{cases}$$

If we set  $\phi$ , v, w be any vectors in V which have the forms of

$$\phi = (\phi_1, \phi_2)^T, \quad v = (v_1, v_2)^T, \quad w = (w_1, w_2)^T,$$

we denote by b the trilinear form [2]

$$b(\phi, v, w) = (B(\phi, v), w) = \int_{\Omega} (\phi \cdot \nabla) v \cdot w dx,$$

which has the following estimations

$$(3.5) b(\phi, v, w) \le c_b |A^{\frac{s_1}{2}}\phi| |A^{\frac{s_2+1}{2}}v| |A^{\frac{s_3}{2}}w|, \forall \phi \in D(A^{\frac{s_1}{2}}), v \in D(A^{\frac{s_2+1}{2}}), w \in D(A^{\frac{s_3}{2}})$$

Here,  $s_1, s_2, s_3 \ge 0$  satisfies  $s_1 + s_2 + s_3 \ge 1$  with  $(s_1, s_2, s_3) \ne (1, 0, 0), (0, 1, 0)$  and (0, 0, 1). On the other hand, we can alter the form of b. In fact

$$(\phi \cdot \nabla)v \cdot w = \begin{bmatrix} \begin{pmatrix} \phi_1 \\ \phi_2 \end{pmatrix} \cdot \begin{pmatrix} \partial_x \\ \partial_y \end{bmatrix} \end{bmatrix} \begin{pmatrix} v_1 \\ v_2 \end{pmatrix} \cdot \begin{pmatrix} w_1 \\ w_2 \end{pmatrix}$$
$$= \phi_1 \partial_x v_1 w_1 + \phi_2 \partial_y v_1 w_1 + \phi_1 \partial_x v_2 w_2 + \phi_2 \partial_y v_2 w_2$$
$$= (w \cdot \nabla v) \cdot \phi,$$

where  $w \cdot \nabla v$  means  $w_1 \nabla v_1 + w_2 \nabla v_2$ . Denoting  $\mathcal{B}(w,v) = P(w \cdot \nabla v)$ , then we have

(3.6) 
$$b(\phi, v, w) = (B(\phi, v), w) = (\mathcal{B}(w, v), \phi).$$

We are familiar with the general properties of B, about which we refer readers to [2]. But those are not enough for our discussion, we still need some further properties of B and B which we state as following two lemmas.

**Lemma 3.1** For any  $w \in D(A^{\frac{1}{2}})$ ,  $v \in D(A)$  and  $0 < r < \frac{1}{2}$ , it holds

$$|A^r B(w,v)| \le c_1 |A^{\frac{1}{2}} w| |Av|.$$

And for  $r = \frac{1}{2}$  and  $w \in D(A^{\frac{3}{4}})$ , we have

$$|A^{\frac{1}{2}}B(w,v)| \le c_1|A^{\frac{3}{4}}w||Av|.$$

Here  $c_1$  is some positive constant depending only on  $\Omega$  and r.

This property of B was proven in [13]. Noticing the form of B and B are quite alike, we can easily get the similar property of B by the same method used in [13]. So we only state the property in the following without proving.

**Lemma 3.2** For any  $w \in D(A^{\frac{1}{2}})$ ,  $v \in D(A)$  and  $0 < r < \frac{1}{2}$ , it holds

$$|A^r \mathcal{B}(w,v)| \le c_1 |A^{\frac{1}{2}} w| |Av|.$$

And for  $r = \frac{1}{2}$  and  $w \in D(A^{\frac{3}{4}})$ , we have

$$|A^{\frac{1}{2}}\mathcal{B}(w,v)| \le c_1 |A^{\frac{3}{4}}w| |Av|.$$

Here  $c_1$  has the same meaning as in lemma 3.1.

Now we are ready to study the important property of  $\hat{u}$ . By using the semi-group presentation, we can rewrite (3.4) as

$$\begin{split} p(t) &= -\int_0^t e^{-\nu A(t-s)} P_m \{B(u_m, p) + B(p, u_m) + B(p, p) + B(p, q) + B(q, p) \\ &+ B(u_m, q) + B(q, u_m) + B(q, q) \} ds \\ &= -\int_0^t e^{-\nu A(t-s)} P_m B_1(p) ds - \int_0^t e^{-\nu A(t-s)} P_m B_2(q) ds \,, \end{split}$$

where

$$B_1(p) = B(u_m, p) + B(p, u_m) + B(p, p) + B(p, q) + B(q, p),$$
  

$$B_2(q) = B(u_m, q) + B(q, u_m) + B(q, q).$$

Then by using (2.3), we have

$$(3.7) \quad |A^{-\frac{1}{2}}p(t)| \leq \int_{0}^{t} |A^{-\frac{1}{2}}e^{-\nu A(t-s)}P_{m}B_{1}(p)|ds + \int_{0}^{t} |A^{-\frac{1}{2}}e^{-\nu A(t-s)}P_{m}B_{2}(q)|ds$$

$$= \int_{0}^{t} |A^{\frac{1}{2}}e^{-\nu A(t-s)}A^{-1}P_{m}B_{1}(p)|ds + \int_{0}^{t} |A^{\frac{1}{2}}e^{-\nu A(t-s)}A^{-1}P_{m}B_{2}(q)|ds$$

$$\leq c_{0}\nu^{-\frac{1}{2}}\int_{0}^{t} (t-s)^{-\frac{1}{2}}e^{-\delta(t-s)}|A^{-1}P_{m}B_{1}(p)|ds$$

$$+ c_{0}\nu^{-\frac{1}{2}}\int_{0}^{t} (t-s)^{-\frac{1}{2}}e^{-\delta(t-s)}|A^{-1}P_{m}B_{2}(q)|ds$$

Let us estimate each term of  $|A^{-1}B_1(p)|$  and  $|A^{-1}B_2(q)|$ . Firstly, we consider each term of  $|A^{-1}B_1(p)|$ . For example, consider  $|A^{-1}P_mB(u_m,p)|$ . Noticing lemma 3.1, for any  $v \in H_m$ 

$$|b(u_m, p, A^{-1}v)| = |b(u_m, A^{-1}v, p)| = |(A^{\frac{1}{2}}B(u_m, A^{-1}v), A^{-\frac{1}{2}}p)|$$

$$\leq c_1 |A^{\frac{3}{4}}u_m| |A^{-\frac{1}{2}}p| |v|.$$

Thus we have

$$(3.8) |A^{-1}P_mB(u_m,p)| < c_1|A^{\frac{3}{4}}u_m||A^{-\frac{1}{2}}p|.$$

Similarly, we can derive

$$(3.9) |A^{-1}P_mB(p,p)| \le c_1|A^{\frac{3}{4}}u||A^{-\frac{1}{2}}p|,$$

$$(3.10) |A^{-1}P_mB(q,p)| \le c_1|A^{\frac{3}{4}}u||A^{-\frac{1}{2}}p|.$$

For the other two terms, we will use lemma 3.2 to cope with them. For any  $v \in H_m$ 

$$|b(p, u_m, A^{-1}v)| = |b(p, A^{-1}v, u_m)| = |(A^{\frac{1}{2}}\mathcal{B}(u_m, A^{-1}v), A^{-\frac{1}{2}}p)|$$

$$< c_1 |A^{\frac{3}{4}}u_m| |A^{-\frac{1}{2}}p| |v|.$$

So we get

$$(3.11) |A^{-1}P_mB(p,u_m)| \le c_1|A^{\frac{3}{4}}u_m||A^{-\frac{1}{2}}p|.$$

Do the same thing to the last term, we have

$$(3.12) |A^{-1}P_mB(p,q)| < c_1|A^{\frac{3}{4}}u||A^{-\frac{1}{2}}p|.$$

Combining  $(3.8) \sim (3.12)$ , we derive the first estimation

$$(3.13) |A^{-1}P_mB_1(p)| \le 3c_1(||A^{\frac{3}{4}}u|| + ||A^{\frac{3}{4}}u_m|||)|A^{-\frac{1}{2}}p|.$$

For the estimation of  $|A^{-1}P_mB_2(q)|$ , the method is completely the same as the above one. We also use lemma 3.1 to deal with  $B(u_m,q) + B(q,q)$  and lemma 3.2 to deal with  $B(q,u_m)$ . So we just give the result in the following

$$(3.14) |A^{-1}P_mB_2(q)| \le 2c_1(||A^{\frac{3}{4}}u|| + ||A^{\frac{3}{4}}u_m|||)|A^{-\frac{1}{2}}q|.$$

Obviously,

$$\sup_{t>0} \int_0^t (t-s)^{-\frac{1}{2}} e^{-\delta(t-s)} ds < \delta^{-1/2} \gamma_{\frac{1}{2}} < +\infty$$

where

$$\gamma_{\alpha} = \int_{0}^{\infty} s^{-\alpha} e^{-s} ds = \Gamma(-1 - \alpha).$$

By introducing the following constants

$$c_2 = 2c_0c_1\delta^{-\frac{1}{2}}\gamma_{\frac{1}{2}}(|||A^{\frac{3}{4}}u||| + ||||A^{\frac{3}{4}}u_m|||),$$

$$c_3 = 3c_0c_1(||||A^{\frac{3}{4}}u||| + ||||A^{\frac{3}{4}}u_m|||),$$

we can get a new integration inequality from (3.7). That is

$$|A^{-\frac{1}{2}}p(t)| \le c_3 \nu^{-\frac{1}{2}} \int_0^t (t-s)^{-\frac{1}{2}} e^{-\delta(t-s)} |A^{-\frac{1}{2}}p| ds + c_2 \nu^{-\frac{1}{2}} |||A^{-\frac{1}{2}}q|||_t.$$

Set

$$g(s) = |A^{-\frac{1}{2}}p(s)|e^{\delta s},$$

we have

$$g(t) \le c_2 \nu^{-\frac{1}{2}} e^{\delta t} |||A^{-\frac{1}{2}}q|||_t + c_3 \nu^{-\frac{1}{2}} \int_0^t (t-s)^{-\frac{1}{2}} g(s) ds.$$

To give the estimation of g, we must introducing an inequality. Many inequalities of this type can be found in Henry[4]. The following special version, lemma 3.3, was proven in [15].

**Lemma 3.3** Let T,  $\alpha$ ,  $\beta$  and  $\theta$  be positive constants,  $0 < \theta < 1$ . Then for any continuous function  $f : [0, T] \to [0, +\infty)$  that satisfies

$$f(t) \le \alpha + \beta \int_0^t (t-s)^{-\theta} f(s) ds, \quad 0 \le t \le T,$$

we have

$$f(t) \le c_4 \alpha \exp\{c_4 \beta^{1/(1-\theta)}t\}, \quad 0 \le t \le T,$$

with a positive constant  $c_4$  that depends only on  $\theta$ .

Now by using lemma 3.3, we can immediately obtain

$$g(t) \le c_2 c_4 \nu^{-\frac{1}{2}} e^{\delta t} \exp\{c_4 c_3^2 \nu^{-1} t\} |||A^{-\frac{1}{2}} q|||_{t}.$$

Denoting by  $T_1(t) > 0$  the constant  $c_2 c_4 \exp\{c_4 c_3^2 \nu^{-1} t\} + \nu^{\frac{1}{2}}$ , we have

$$(3.15) |A^{-\frac{1}{2}}p(t)| \le (\nu^{-\frac{1}{2}}T_1(t) - 1)|||A^{-\frac{1}{2}}q||_{t}.$$

Now we summarize the above deducing into the following

**Theorem 3.1** For any given data  $a \in D(A)$  and  $f \in H$ , we know the Navier-Stokes equations (2.6) and its Galerkin approximate equations (2.4) have unique solutions

$$u(t) \in L^{\infty}(\mathcal{R}^+, D(A)), \quad u_m(t) \in L^{\infty}(\mathcal{R}^+, D(A)).$$

And there also exist some positive constants  $M_0$  and  $M_1$  related on  $\nu$ , a and f such that

$$|||A^{\frac{1}{2}}u|||, |||A^{\frac{1}{2}}u_m||| \le \frac{M_0}{2}, \quad ||||Au|||, |||Au_m||| \le M_1.$$

Then for any t > 0, we have

$$|A^{-\frac{1}{2}}(u-u_m)| \le \nu^{-\frac{1}{2}} \lambda_{m+1}^{-\frac{1}{2}} T_1(t) |||u-u_m|||_t.$$

**Proof** From [14] and (3.15), we can immediately get the result.

# 4 Finite Dimensional Mapping $\Phi$

As we said before, the main task in our paper is to construct some kind of approximate interactive rule between  $\hat{u}$  and  $u_m$ . That is to find some kind of finite dimensional mapping  $\Phi: H_m \to V$  such that  $\hat{u} \approx \Phi(u_m)$ .

To do so, we introduce an smooth function

$$g(s) \in C^{\infty}(\mathcal{R}^+)$$

with the following properties

$$0 < g(s) < 1, |g'(s)| < 2,$$
 and  $\forall s \in [0, 1], g(s) = 1, \forall s \in [2, \infty), g(s) = 0.$ 

Now let us recall (3.2). Of course, (3.2) is a kind of rule. We may exactly get  $\hat{u}$  from  $u_m$ . This is, in fact, to solve the Navier-Stokes equations. It is not suitable for our purpose because it is as complex as (2.2). But, enlightened by this equations and noticing  $\frac{d\hat{u}}{dt}$  and  $B(\hat{u}, \hat{u})$  are smaller quantities compared with other terms in (3.2), we introduce the following finite dimensional mapping

(4.1) 
$$\begin{cases} \forall \phi \in H_m, & \text{find } w \in V \text{ such that} \\ h^{-1}w + \nu Aw + g(\frac{2|A^{\frac{1}{2}}\phi|}{M_0})[B(w,\phi) + B(\phi,w)] = g(\frac{2|A^{\frac{1}{2}}\phi|}{M_0})Q_m[f - B(\phi,\phi)]. \end{cases}$$

Here, h > 0 is a small constant which will be given soon. Let us introduce bilinear form

$$\mathcal{L}(w,v) := h^{-1}(w,v) + \nu (A^{\frac{1}{2}}w, A^{\frac{1}{2}}v) + g(\frac{2|A^{\frac{1}{2}}\phi|}{M_0})[b(w,\phi,v) + b(\phi,w,v)].$$

It is clear that  $\mathcal{L}(\cdot, \cdot)$  is a continuous bilinear form from  $V \times V$  to R. Furthermore, we have **Lemma 4.1**  $\mathcal{L}(\cdot, \cdot)$  is continuous and coercive if h is small enough such that

$$(4.2) h \le \frac{2\nu}{c_b^2 M_0^2}.$$

Proof Indeed

$$(h^{-1}w + \nu Aw + g(\frac{2|A^{\frac{1}{2}}\phi|}{M_0})[B(w,\phi) + B(\phi,w)], w) = \frac{|w|^2}{h} + \nu |A^{\frac{1}{2}}w|^2 + g(\frac{2|A^{\frac{1}{2}}\phi|}{M_0})b(w,\phi,w)$$

$$\geq \frac{|w|^2}{h} + \nu |A^{\frac{1}{2}}w|^2 - c_b|A^{\frac{1}{2}}\phi||w||A^{\frac{1}{2}}w| \geq \frac{|w|^2}{h} + \nu |A^{\frac{1}{2}}w|^2 - c_bM_0|w||A^{\frac{1}{2}}w|$$

$$\geq (\frac{1}{h} - \frac{c_b^2 M_0^2}{2\nu})|w|^2 + \frac{\nu}{2}|A^{\frac{1}{2}}w|^2 \geq \frac{\nu}{2}|A^{\frac{1}{2}}w|^2.$$

Notice we used an implicit condition  $|A^{\frac{1}{2}}\phi| < M_0$  in the above inequality. For  $|A^{\frac{1}{2}}\phi| \ge M_0$ , g will be equal to zero and the above result is obvious. Then we can get the result.

Consequently, by using Lax-Milgram theorem, we know, for any give  $\phi \in H_m$ , there exists an unique solution  $w = \Phi(\phi)$  of the following variational problem corresponding to (4.1):

$$\left\{ \begin{array}{l} \forall \phi \in H_m, \ \ \mathrm{find} \ w \in V \ \mathrm{such \ that} \\ \mathcal{L}(w,v) = g(\frac{2|A^{\frac{1}{2}}\phi|}{M_0})[(Q_mf,v) - b(\phi,\phi,Q_mv)] \quad \forall v \in V \end{array} \right.$$

**Theorem 4.1** Assume h satisfies (4.2). Then (4.1) can determine a finite dimensional mapping  $\Phi$  from  $H_m$  to V which has the following properties

- i)  $\Phi(\phi) = 0 \ for \ |A^{\frac{1}{2}}\phi| \ge M_0$ .
- ii) For any  $\phi \in H_m$ ,

$$|A^{\frac{1}{2}}\Phi(\phi)| \le \rho_m = \frac{2}{\nu}(|f| + c_b M_0^2 L_m)\lambda_{m+1}^{-\frac{1}{2}},$$

of course  $\rho_m \to 0$  when  $m \to \infty$ .

iii)  $\Phi$  is a Lipschitz smooth mapping. That is, there exists some constant  $l_m > 0$  such that

$$|A^{\frac{1}{2}}(\Phi(\phi_1) - \Phi(\phi_2))| \le l_m |A^{\frac{1}{2}}(\phi_1 - \phi_2)|.$$

And  $l_m \to 0$  when  $m \to \infty$ .

**Proof** By virtue of lemma 4.1, it asserts that (4.1) can determine a finite dimensional mapping.

i) Let us consider  $|A^{\frac{1}{2}}\phi| \geq M_0$ . Notice the definition of g, we know  $g(\frac{2|A^{\frac{1}{2}}\phi|}{M_0}) = 0$  at this time. So (4.1) becomes

$$h^{-1}w + \nu Aw = 0.$$

Of course, it only has zero solution. That is, under this circumstance,

$$w = \Phi(\phi) = 0.$$

ii) We only need to consider  $|A^{\frac{1}{2}}\phi| < M_0$ . Just as being done to prove the uniqueness of the solution, we can get

$$\frac{\nu}{2} |A^{\frac{1}{2}} \Phi(\phi)|^{2} \leq |(f, Q_{m} \Phi(\phi)) + b(\phi, \phi, Q_{m} \Phi(\phi))|$$

$$\leq (|f| + c_{b} M_{0}^{2} L_{m}) \lambda_{m+1}^{-\frac{1}{2}} |A^{\frac{1}{2}} \Phi(\phi)|.$$

Then we can get the result.

iii) At last, we will show this mapping is also Lipschitz continuous. For any given  $\phi_1, \phi_2 \in H_m$ , we can get  $w_1 = \Phi(\phi_1)$  and  $w_2 = \Phi(\phi_2)$  from (4.1). That is

$$h^{-1}w_1 + \nu Aw_1 + g(\frac{2|A^{\frac{1}{2}}\phi_1|}{M_0})[B(w_1,\phi_1) + B(\phi_1,w_1)] = g(\frac{2|A^{\frac{1}{2}}\phi_1|}{M_0})Q_m[f - B(\phi_1,\phi_1)],$$
  
$$h^{-1}w_2 + \nu Aw_2 + g(\frac{2|A^{\frac{1}{2}}\phi_2|}{M_0})[B(w_2,\phi_2) + B(\phi_2,w_2)] = g(\frac{2|A^{\frac{1}{2}}\phi_2|}{M_0})Q_m[f - B(\phi_2,\phi_2)].$$

If we denote  $\phi_e = \phi_1 - \phi_2$ ,  $w_e = w_1 - w_2$  and  $\Delta g = g(\frac{2|A^{\frac{1}{2}}\phi_1|}{M_0}) - g(\frac{2|A^{\frac{1}{2}}\phi_2|}{M_0})$ , we can derive from the above two equations that

$$(4.4) h^{-1}w_e + \nu Aw_e + g(\frac{2|A^{\frac{1}{2}}\phi_1|}{M_0})[B(w_e, \phi_1) + B(w_2, \phi_e)]$$

$$+ g(\frac{2|A^{\frac{1}{2}}\phi_1|}{M_0})[B(\phi_e, w_1) + B(\phi_2, w_e)] + \Delta g[B(w_2, \phi_2) + B(\phi_2, w_2)]$$

$$= \Delta g Q_m f + g(\frac{2|A^{\frac{1}{2}}\phi_1|}{M_0})Q_m[B(\phi_e, \phi_1) + B(\phi_2, \phi_e)] + \Delta g Q_m B(\phi_2, \phi_2).$$

For different values of  $|A^{\frac{1}{2}}\phi_1|$  and  $|A^{\frac{1}{2}}\phi_2|$ , we divided our proof into several cases.

Case 1)  $|A^{\frac{1}{2}}\phi_1|$ ,  $|A^{\frac{1}{2}}\phi_2| \geq M_0$ . Noticing the definition of g, (4.4) becomes

$$h^{-1}w_e + \nu Aw_e = 0.$$

We can get the Lipschitz continuous result for any  $l_m \in \mathcal{R}^+$ .

Case 2) One of them exceeds  $M_0$ .

Without loss of generality, we suppose  $|A^{\frac{1}{2}}\phi_1| \geq M_0$  and  $|A^{\frac{1}{2}}\phi_2| < M_0$ . Then (4.4) reads

$$(4.5) h^{-1}w_e + \nu Aw_e + \Delta gB(w_2, \phi_2) + \Delta gB(\phi_2, w_2) = \Delta gQ_m f + \Delta gQ_m B(\phi_2, \phi_2).$$

Notice that  $\Delta g = -g(\frac{2|A^{\frac{1}{2}}\phi_2|}{M_0})$  at this time. But we pretend that  $g(\frac{2|A^{\frac{1}{2}}\phi_1|}{M_0})$  is still there. Then by using the property of g, we have

$$|\Delta g| \le \frac{4}{M_0} (|A^{\frac{1}{2}}\phi_1| - |A^{\frac{1}{2}}\phi_2|) \le \frac{4}{M_0} |A^{\frac{1}{2}}\phi_e|.$$

Multiply (4.5) by  $w_e$  and integrate it on  $\Omega$ , we get

$$(4.7) h^{-1}|w_e|^2 + \nu|A^{\frac{1}{2}}w_e|^2 \le |\Delta gb(w_2, \phi_2, w_e)| + |\Delta gb(\phi_2, w_2, w_e)| + |\Delta g(f, Q_m w_e)| + |\Delta gb(\phi_2, \phi_2, Q_m w_e)|.$$

Noticing (4.6) and the result of ii), we majorize each term on the right hand side of (4.7) as:

$$\begin{split} |\Delta g b(w_2,\phi_2,w_e)| &\leq \frac{4c_b}{M_0} |A^{\frac{1}{2}}w_2| \, |A^{\frac{1}{2}}\phi_2| \, |A^{\frac{1}{2}}w_e| \, |A^{\frac{1}{2}}\phi_e| \leq 4c_b\rho_m |A^{\frac{1}{2}}w_e| \, |A^{\frac{1}{2}}\phi_e|, \\ |\Delta g b(\phi_2,w_2,w_e)| &\leq 4c_b\rho_m |A^{\frac{1}{2}}w_e| \, |A^{\frac{1}{2}}\phi_e|, \\ |\Delta g(f,Q_mw_e)| &\leq \frac{4|f|}{M_0\lambda_{m+1}^{\frac{1}{2}}} \, |A^{\frac{1}{2}}w_e| \, |A^{\frac{1}{2}}\phi_e|, \\ |\Delta g b(\phi_2,\phi_2,Q_mw_e)| &\leq \frac{4c_b}{M_0} |\phi_2|_{L^{\infty}} |A^{\frac{1}{2}}\phi_2| \, |Q_mw_e| \, |A^{\frac{1}{2}}\phi_e| \end{split}$$

$$\Delta g v(\phi_2, \phi_2, Q_m w_e) | \leq \frac{1}{M_0} |\phi_2|_{L^{\infty}} |A^2 \phi_2| |Q_m w_e| |A^2 \phi_e|_{L^{\infty}}$$

$$\leq 4c_b M_0 L_m \lambda_{m+1}^{-\frac{1}{2}} |A^{\frac{1}{2}} w_e| |A^{\frac{1}{2}} \phi_e|_{L^{\infty}}.$$

Combining the above inequalities and omitting  $h^{-1}|w_e|^2$  on the left hand side of (4.7), it yields

$$\nu |A^{\frac{1}{2}}w_e|^2 \le (8c_b\rho_m + 4M_0|f|\lambda_{m+1}^{-\frac{1}{2}} + 4c_bM_0L_m\lambda_{m+1}^{-\frac{1}{2}})|A^{\frac{1}{2}}w_e||A^{\frac{1}{2}}\phi_e|.$$

Denoting  $l_m = \nu^{-1} (8c_b \rho_m + 4M_0 | f | \lambda_{m+1}^{-\frac{1}{2}} + 4c_b M_0 L_m \lambda_{m+1}^{-\frac{1}{2}})$ , we can get the result.

Case 3)  $|A^{\frac{1}{2}}\phi_1|$ ,  $|A^{\frac{1}{2}}\phi_2| < M_0$ . Multiply (4.4) with  $w_e$  and integrate it on  $\Omega$ , we have

$$(4.8) h^{-1}|w_{e}|^{2} + \nu |A^{\frac{1}{2}}w_{e}|^{2} \leq |b(w_{e}, \phi_{1}, w_{e})| + |b(\phi_{e}, w_{1}, w_{e})| + |\Delta gb(w_{2}, \phi_{2}, w_{e})| + |\Delta gb(\phi_{2}, w_{2}, w_{e})| + |\Delta g(f, Q_{m}w_{e})| + |b(\phi_{e}, \phi_{1}, Q_{m}w_{e})| + |b(\phi_{2}, \phi_{e}, Q_{m}w_{e})| + |\Delta gb(\phi_{2}, \phi_{2}, Q_{m}w_{e})|.$$

Just like the previous case, we majorize each term on the right hand side of (4.8) as:

$$|b(w_e, \phi_1, w_e)| \le c_b M_0 |A^{\frac{1}{2}} w_e| |w_e| \le \frac{\nu}{2} |A^{\frac{1}{2}} w_e|^2 + \frac{c_b^2 M_0^2}{2\nu} |w_e|^2,$$

$$\begin{split} |b\left(\phi_{e},w_{1},w_{e}\right)| &\leq c_{b}\rho_{m}|A^{\frac{1}{2}}w_{e}|\,|A^{\frac{1}{2}}\phi_{e}|,\\ |\Delta gb\left(w_{2},\phi_{2},w_{e}\right)| &\leq 4c_{b}\rho_{m}|A^{\frac{1}{2}}w_{e}|\,|A^{\frac{1}{2}}\phi_{e}|,\\ |\Delta gb\left(\phi_{2},w_{2},w_{e}\right)| &\leq 4c_{b}\rho_{m}|A^{\frac{1}{2}}w_{e}|\,|A^{\frac{1}{2}}\phi_{e}|,\\ |\Delta g(f,Q_{m}w_{e})| &\leq \frac{4|f|}{M_{0}}\lambda_{m+1}^{-\frac{1}{2}}|A^{\frac{1}{2}}w_{e}|\,|A^{\frac{1}{2}}\phi_{e}|,\\ |b\left(\phi_{e},\phi_{1},Q_{m}w_{e}\right)| &\leq c_{b}|\phi_{e}|_{L^{\infty}}|A^{\frac{1}{2}}\phi_{1}||Q_{m}w_{e}| &\leq c_{b}M_{0}L_{m}\lambda_{m+1}^{-\frac{1}{2}}|A^{\frac{1}{2}}w_{e}||A^{\frac{1}{2}}\phi_{e}|,\\ |b\left(\phi_{2},\phi_{e},Q_{m}w_{e}\right)| &\leq c_{b}M_{0}L_{m}\lambda_{m+1}^{-\frac{1}{2}}|A^{\frac{1}{2}}w_{e}||A^{\frac{1}{2}}\phi_{e}|,\\ |\Delta gb\left(\phi_{2},\phi_{2},Q_{m}w_{e}\right)| &\leq 4c_{b}M_{0}L_{m}\lambda_{m+1}^{-\frac{1}{2}}|A^{\frac{1}{2}}w_{e}||A^{\frac{1}{2}}\phi_{e}|. \end{split}$$

Then, from (4.8), we have

$$\frac{\nu}{2}|A^{\frac{1}{2}}w_e| \le (9c_b\rho_m + 6c_bM_0L_m\lambda_{m+1}^{-\frac{1}{2}} + 4|f|M_0^{-1}\lambda_{m+1}^{-\frac{1}{2}})|A^{\frac{1}{2}}\phi_e|.$$

Once again, if we denote

$$l_m = \frac{2}{\nu} (9c_b \rho_m + 6c_b M_0 L_m \lambda_{m+1}^{-\frac{1}{2}} + 4|f|M_0^{-1} \lambda_{m+1}^{-\frac{1}{2}}),$$

we can derive the result again.

## 5 Post Galerkin Method

In previous section, we construct a finite dimensional mapping  $\Phi: H_m \to V$ . Now, we will use it to give our post Galerkin scheme. In fact, once we get  $\Phi$ , the construction of post Galerkin scheme is obvious. We state it as following three steps.

(Step 1) Galerkin approximation

$$\begin{cases} & \text{find } u_m(t) \text{ such that} \\ & \frac{du_m}{dt} + \nu A u_m + P_m B(u_m, u_m) = P_m f, \\ & u_m(0) = P_m a. \end{cases}$$

(Step 2) Postprocess  $u_m(t)$  at any time  $t \in \mathcal{R}^+$ 

$$\left\{ \begin{array}{l} \mbox{find } \bar{u}(t) = \Phi(u_m(t)) \mbox{ such that} \\ h^{-1}\bar{u} + \nu A\bar{u} + g(\frac{2|A^{\frac{1}{2}}u_m|}{M_0})[B(\bar{u},u_m) + B(u_m,\bar{u})] = g(\frac{2|A^{\frac{1}{2}}u_m|}{M_0})Q_m[f - B(u_m,u_m)]. \end{array} \right.$$

(Step 3) Post Galerkin approximation

$$u^*(t) = u_m(t) + \bar{u}(t) = u_m(t) + \Phi(u_m(t))$$

In the rest, we will investigate the accuracy presented by this scheme. First of all, we give some classical result as

**Lemma 5.1**[3] Under the conditions of theorem 3.1, the solution u(t) of (2.2) is analytic in time, in a neighborhood of the positive real axis, as a D(A) – valued function.

Denoting by d the distance between the boundary of analytic region and the positive real axis, we could derive the following estimation of  $\frac{d\hat{u}}{dt}$  at any time by Cauchy theorem.

For any given t > 0, we know from the above lemma and Cauchy theorem

$$\frac{d\hat{u}(t)}{dt} = \frac{1}{2\pi i} \int_{|z-t|=d} \frac{\hat{u}(z)}{(t-z)^2} dz.$$

Thus,

$$|A^{-\frac{1}{2}} \frac{d\hat{u}(t)}{dt}| \leq \frac{1}{2\pi} \left| \int_{|z-t|=d} \frac{A^{-\frac{1}{2}} \hat{u}(z)}{(t-z)^2} dz \right|$$

$$\leq d^{-1} |||A^{-\frac{1}{2}} \hat{u}|||_{t} \leq d^{-1} \nu^{-\frac{1}{2}} \lambda_{m+1}^{-\frac{1}{2}} T_{1}(t) |||\hat{u}|||_{t}.$$

Generally, d is a small constant related to  $\nu$  and f. For its concrete express, we refer readers to [3]. Now, we give our main result as

**Theorem 5.1** Under the conditions of theorem 3.1 and (4.2), we have

$$|A^{\frac{1}{2}}(u(t) - u^*(t))| \le T_2(t)\lambda_{m+1}^{-\frac{3}{2}},$$

where

$$T_2(t) = 2\nu^{-1}C_1(t)((h^{-1} + d^{-1})\nu^{-\frac{1}{2}}T_1(t) + c_1C_2(t)).$$

**Proof** Notice that

$$u^*(t) = u_m(t) + \Phi(u_m(t)), \quad u(t) = u_m(t) + \hat{u}(t).$$

Thus, to get the estimation, we only need to concern about

$$\hat{u}(t) - \Phi(u_m(t)) = \hat{u}(t) - \bar{u}(t).$$

Because of  $||A^{\frac{1}{2}}u_m|| \le \frac{1}{2}M_0$ , we have

$$g(\frac{2|A^{\frac{1}{2}}u_m|}{M_0}) = 1.$$

Thus  $\bar{u}$  satisfies the following equations at any give time

(5.2) 
$$h^{-1}\bar{u} + \nu A\bar{u} + B(\bar{u}, u_m) + B(u_m, \bar{u}) = Q_m[f - B(u_m, u_m)].$$

Subtracting (5.2) from (3.2) and denoting  $e = \hat{u} - \bar{u}$ , we have

$$h^{-1}e + \nu Ae + B(e, u_m) + B(u_m, e) = h^{-1}\hat{u} - \frac{d\hat{u}}{dt} - B(\hat{u}, \hat{u}).$$

From lemma 4.1, we know

$$||h^{-1}e + \nu Ae + B(e, u_m) + B(u_m, e)||_{\mathcal{L}(V, V^*)} \ge \frac{\nu}{2} |A^{\frac{1}{2}}e|.$$

And from (2.5), theorem 3.1 and (5.1), we have

$$||h^{-1}\hat{u} - \frac{d\hat{u}}{dt} - B(\hat{u}, \hat{u})||_{\mathcal{L}(V, V^*)} \le (h^{-1} + d^{-1})|||A^{-\frac{1}{2}}\hat{u}|||_t + c_1|\hat{u}||A^{\frac{1}{2}}\hat{u}||$$

$$\le ((h^{-1} + d^{-1})\nu^{-\frac{1}{2}}T_1(t)\lambda_{m+1}^{-\frac{1}{2}} + c_1|A^{\frac{1}{2}}\hat{u}|)|\hat{u}|.$$
(5.3)

By using (2.5) we conclude

$$|A^{\frac{1}{2}}e| \le 2\nu^{-1}C_1(t)((h^{-1}+d^{-1})\nu^{-\frac{1}{2}}T_1(t) + c_1C_2(t))\lambda_{m+1}^{-\frac{3}{2}}.$$

End of the proof.  $\Box$ 

**Remark.** Notice that (Step 2) of the post Galerkin scheme is solved in whole space V. When we consider the realistic implementation of this scheme, we should restrict this step in a larger finite dimensional subspace of V, namely,  $H_M$  with  $M \gg m$  and to get a finite dimensional approximation  $u_M^*(t)$ . That is we should modify the (Step 2) and (Step 3) as (Step 2') Postprocess  $u_m(t)$  at any time  $t \in \mathcal{R}^+$ 

$$\begin{cases} & \text{find } u^M(t) = \Phi_M(u_m(t)) \in H_M \text{ such that} \\ & h^{-1}u^M + \nu A u^M + g(\frac{2|A^{\frac{1}{2}}u_m|}{M_0})P_M[B(u^M, u_m) + B(u_m, u^M)] \\ & = g(\frac{2|A^{\frac{1}{2}}u_m|}{M_0})P_MQ_m[f - B(u_m, u_m)]. \end{cases}$$

(Step 3') Post Galerkin approximation

$$u_M^*(t) = u_m(t) + u^M(t) = u_m(t) - \Phi_M(u_m(t)).$$

It is easy to show that this finite dimensional scheme has the following error estimation

$$|A^{\frac{1}{2}}(u(t) - u_M^*(t))| \le T_2(t)\lambda_{m+1}^{-\frac{3}{2}} + C\lambda_{M+1}^{-\frac{1}{2}}.$$

Of course, to balance the two terms on the right hand side of the above inequality, we should choose  $M \sim m^3$ . That is the performance of our proposed scheme is just like that of a standard Galerkin scheme with very large computing scale.

On the other hand, the results here are also valid when we consider the periodic boundary conditions case.

## 6 Numerical experiment

In this section, we will present a simple numerical experiment for our scheme. For the sake of simplicity, we will consider problem (2.2) in a square domain  $\Omega = (-\pi, \pi)^2$  with periodical boundary conditions. Under this circumstance, H is

$$H = \{u = \sum_{k \in Z, k \neq 0} u_k e^{ix \cdot k}, \ u_k = \overline{u_{-k}}, \ \mathrm{div} u = 0 \ \text{ under weak sense }, \ \sum_{k \in Z, k \neq 0} |u_k|^2 < +\infty \}.$$

We assume that we have a true solution of (2.2). In fact, we give some function  $u(t) \in H$  with

$$|u_k(t)| \sim |k|^{-2}$$
.

And for given  $m \in \mathcal{N}$ , we define  $P_m$  is the orthogonal projection from H onto

$$H_m = \{ u = \sum_{|k_1|, |k_2| \le m} u_k e^{ix \cdot k}, \ u_k = \overline{u_{-k}}, \ \mathrm{div} u = 0 \},$$

where  $k = (k_1, k_2)^T$ . Meanwhile, for the periodical case, it is very easy to get the divergence free projection P. Then we use this u(t) to compute  $P_m f(t)$  in (2.4) and solve it to get the standard Galerkin approximation.

Concerning about the computing scale, we only give a small scale simulation here. For example, in our numerical implementation, we take m=9, M=2m and  $h=\nu=1$ . Following table indicates the relative error of standard Galerkin method and post Galerkin method defined as

$$SGM = \frac{\|u - u_m\|}{\|u\|},$$

and

$$PGM = \frac{\|u - u_M^*\|}{\|u\|},$$

where SGM and PGM mean the relative error of standard Galerkin method and post Galerkin method respectively.

Time	$\operatorname{SGM}$	PGM
2.0	2.08%	0.400%
4.0	2.86%	1.33%
6.0	4.87%	2.71%
8.0	10.6%	5.98%
10.0	9.06%	4.90%
12.0	4.87%	2.38%
14.0	3.43%	1.47%
16.0	2.83%	1.06%
18.0	2.56%	0.845%
20.0	2.50%	0.757%
22.0	2.56%	0.747%

It seems that  $\frac{PGM}{SGM} \sim \frac{1}{2}$ . From the remark at the end of last section, we know this ratio is restricted by the truncation error because we just take M=2m instead of  $M\sim m^3$  when concerning about the large computing scale. So this ratio is reasonable because  $\frac{\lambda_{M^2+1}^{-\frac{1}{2}}}{\lambda_{m^2+1}^{-\frac{1}{2}}}$  should be close to  $\frac{M^{-1}}{m^{-1}}=\frac{1}{2}$ .

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