Recent Development in the CESE Method for the Solution of the Navier-Stokes Equations Using Unstructured Triangular or Tetrahedral Meshes With High Aspect Ratio

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#### Abstract

In the multidimensional CESE development, triangles and tetrahedra turn out to be the most natural building blocks for 2D and 3D spatial meshes. As such the CESE method is compatible with the simplest unstructured meshes and thus can be easily applied to solve problems with complex geometries. However, because the method uses space-time staggered stencils, solution decoupling may become a real nuisance in applications involving unstructured meshes. In this paper we will describe a simple and general remedy which, according to numerical experiments, has removed any possibility of solution decoupling. Moreover, in a real-world viscous flow simulation near a solid wall, one often encounters a case where a boundary with high curvature or sharp corner is surrounded by triangular/tetrahedral meshes of extremely high aspect ratio (up to $10^{6}$ ). For such an extreme case, the spatial projection of a space-time compounded conservation element constructed using the original CESE design may become highly concave and thus its centroid (referred to as a spatial solution point) may lie far outside of the spatial projection. It could even be embedded beyond a solid wall boundary and causes serious numerical difficulties. In this paper we will also present a new procedure for constructing conservation elements and solution elements which effectively overcomes the difficulties associated with the original design. Another difficulty issue which was addressed more recently is the wellknown fact that accuracy of gradient computations involving triangular/tetrahedral grids deteriorates rapidly as the aspect ratio of grid cells increases. The root cause of this difficulty was clearly identified and several remedies to overcome it were found through a rigorous mathematical analysis. However, because of the length of the current paper and the complexity of mathematics involved, this new work will be presented in another paper.


## 1. Introduction

The space-time conservation element and solution element (CESE) method is a high-resolution and genuinely multidimensional method for solving conservation laws [1-70]. Its nontraditional features include: (i) a unified treatment of space and time; (ii) the introduction of conservation elements (CEs) and solution elements (SEs) as the vehicles for enforcing space-time flux conservation; (iii) a time marching strategy that has a space-time staggered stencil at its core and, as such, fluxes at an interface can be evaluated without using
any interpolation or extrapolation procedure (which, in turn, leads to the method's ability to capture shocks without using Riemann solvers, and also partially contributes to the unusually simple CESE non-reflecting boundary conditions); (iv) the requirement that each scheme be built from a nondissipative core scheme and, as a result, the numerical dissipation can be controlled effectively; and (v) the fact that mesh values of the physical dependent variables and their spatial derivatives are considered as independent mesh variables to be solve for simultaneously. Note that CEs are non-overlapping space-time subdomains introduced such that (i) the computational domain can be filled by these subdomains; and (ii) flux conservation can be enforced over each of them and also over the union of any combination of them. On the other hand, SEs are space-time subdomains introduced such that (i) the boundary of each CE can be divided into component parts with each of them belonging to a unique SE; and (ii) within a SE, any physical flux vector is approximated using simple smooth functions constructed from the solution information stored at a space-time solution point. In general, a CE does not coincide with a SE.

Without using flux-splitting, dimensional-splitting, mesh-alignment or other special techniques, since its inception in 1991 [1], the unstructured-mesh compatible 2nd-order CESE method has been used to obtain numerous accurate 1D, 2D and 3D steady and unsteady flow solutions with Mach numbers ranging from 0.0028 to 10 [15]. The physical phenomena modeled include traveling and interacting shocks, acoustic waves, vortex shedding, viscous flows, detonation waves, cavitation, flows in fluid film bearings, heat conduction with melting and/or freezing, MHD vortex, hydraulic jump, crystal growth, chromatographic problems, solar wind and stress waves in elastic solids [3-70]. In particular, its unexpected simple non-reflecting boundary conditions $[12,50]$ and rather unique capability to resolve both strong shocks and small disturbances (e.g., acoustic waves) simultaneously $[19,21,22,66,67]$ make the CESE method an effective tool for attacking computational aeroacoustics (CAA) problems. Note that the fact that the second-order CESE method is capable of solving CAA problems accurately is an exception to the commonly-held belief that a second-order scheme is not adequate for modeling CAA problems. Also note that, while numerical dissipation is needed for shock capturing, it may also result in annihilation of small disturbances. Thus a solver that can handle both strong shocks and small disturbances simultaneously must be able to overcome this difficulty.

In spite of its nontraditional features and robust capabilities, the core ideas of the CESE method are simple. In fact, all of its key features are the results of a pursuit driven by these simple ideas. The first and foremost is the belief that the method must be solid in physics. As such, in the CESE development, conservation laws are enforced locally and globally in their natural space-time unity forms for 1D, 2D and 3D cases. Moreover, because direct physical interaction generally occurs only among the immediate neighbors, use of the simplest stencil also becomes a common CESE feature. Obviously, this requirement has the effect of simplifying boundary-condition implementation.

The second idea emerges from the realization that stability and accuracy are two competing issues in time-accurate computations, i.e., too much numerical dissipation will degrade accuracy while too little of it will cause instability. In other words, to meet both accuracy and stability requirements, computation must be performed away from the edge ("cliff") of instability but not too far from it. This represents a real dilemma in numerical method development. As an example, high order schemes generally have higher accuracy and lower numerical dissipation. However, they are susceptible to computational instabilities. In fact, in complicated real-world applications, not only they seldom live up to their nominal order of accuracy generally they possess only first-order accuracy when shocks are present, stability of these schemes often is highly problem-dependent and difficulty to predict. Moreover, generally speaking, the higher the order of the scheme, the more restritive its stability conditions become. To confront this issue head-on, in the CESE development, generally it is required that a solver be built from a nondissipative (i.e., neutrally stable) core scheme. By definition, computations involving a neutrally stable scheme are performed right on the edge of instability and therefore the numerical results generated are nondissipative. As such numerical dissipation can be controlled effectively if the deviation of a solver from its nondissipative core scheme can be adjusted using some built-in parameters.

Moreover, because an accurate viscous flow simulation requires that the numerical dissipation be much smaller than the physical dissipation which decreases as the Reynolds number increases, in principle, an accurate and robust solver for high-Reynolds-number flows must be able to cut numerical dissipation as the

Reynolds number increases. Obviously this requirement can only be met by a solver built from a nondissipative core solver.

Other CESE ideas are: (i) the flux at an interface be evaluated in a simple and consistent manner; (ii) genuinely multidimensional schemes be built as simple, consistent and straightforward extensions of 1D schemes so that multidimensional shocks can be captured without using any mesh alignment technique; (iii) triangular and tetrahedral meshes be used in 2D and 3D cases, respectively, so that the method is compatible with the simplest unstructured meshes and thus can be easily used to solve problems with complex geometries; and (iv) logical structures and approximation techniques used be as simple as possible, and special techniques that has only limited applicability and may cause undesirable side effects be avoided. Fortunately for the CESE development, realization of the above lesser ideas (i)-(iv) follows naturally from the first two core ideas.

Building on its initial successes, efforts to refine and improve the CESE method have continued in the past few years. As an example, it was shown in [5] that the numerical dissipation of a dissipative extension of a CESE core scheme may increase to an intolerable level as the value of the $C F L$ number decreases from its maximum stability bound. As such, in a case with a large $C F L$ number disparity (e.g., a simulation with a highly non-uniform spatial mesh and a spatially independent time step), the performance sensitivity with respect to the $C F L$ number can lead to a solution that is highly dissipative in a region where the local $C F L$ number $\ll 1$. Even though a remedy was suggested in [5], a simple and robust solution to this problem had not arrived until a family of new Courant number insensitive schemes [25,31, 46-49,51,54,55,58] was developed with fresh insights. Note that these new schemes have one important advantage, i.e., all variables at each mesh point can be evaluated explicitly without resorting to solving a system of linear/nonlinear algebraic equations involving these local mesh variables even in applications where systems of multidimensional nonlinear PDEs are solved.

As another example, even though they are accurate enough to solve CAA problems, second-order CESE solvers are not capable of resolving fine flow structures within a boundary layer without using relatively fine meshes and/or meshes with large aspect ratios. To overcome this limitation, two neutrally stable 4th-order schemes, referred to as the $a(3)$ and $a-4$ schemes [59-61], were developed in 2006 and intended to serve as the core schemes of other high order CESE schemes. Unfortunately, the dissipative extensions of these new schemes turned out to have the same intractable stability problem that afflicts traditional high order schemes. With the aid of a conceptual leap, this difficulty was finally overcome through the development of a new approach [63] by which one can build from a given 2 nd-order CESE scheme its 4th-, 6 th-, 8th-,... order versions which have the same stencil and same stability conditions of the 2nd-order scheme, and also retain all other advantages of the latter scheme. As such, these new high order schemes can avoid the common shortcomings of traditional high order schemes including: (a) susceptibility to computational instabilities; (b) computational inefficiency due to their local implicit nature (i.e., at each mesh points, need to solve a system of linear/nonlinear equations involving all the mesh variables associated with this mesh point); (c) use of large and elaborate stencils which complicates boundary treatments and also makes efficient parallel computing much harder; (d) difficulties in applications involving complex geometries; and (e) use of problem-specific techniques which are needed to overcome stability problems but often cause undesirable side effects.

In the current paper, we will describe other recent developments which successfully address some of the few remaining CESE issues. Some of the new techniques devised have been implemented in the NASA code ez4d (developed by Dr. Chau-Lyan Chang, the second author) and the Jacobs Technology Inc. code JUSTUS (developed by Dr. Joseph C. Yen, the third author). Their effectiveness has been verified by applications involving hypersonic viscous flow over rough solid surfaces $[62,64,65]$ and also by applications involving studies of the WICS (weapons Internal Carriage and separation) test data produced at Arnold Engineering Development Center $[66,67]$. Note that the current paper would represent the first systematic documentation of the new techniques even though the improved numerical results were presented in [62,64$67]$. To explain clearly these new techniques and their significance, we will begin our discussion by providing some basic information about the CESE method.

For simplicity, first we review existing 2nd-order CESE solvers for the simple partial differential equation (PDE)

$$
\begin{equation*}
\frac{\partial u}{\partial t}+a \frac{\partial u}{\partial x}=0 \tag{1.1}
\end{equation*}
$$

where the advection speed $a \neq 0$ is a constant. Let $x_{1}=x$, and $x_{2}=t$ be the coordinates of a two-dimensional Euclidean space $E_{2}$. Then, because Eq. (1.1) can be expressed as $\vec{\nabla} \cdot \vec{h}=0$ with $\vec{h} \stackrel{\text { def }}{=}(a u, u)$ (i.e., au and $u$ are the $x_{1^{-}}$and $x_{2^{-}}$components of $\vec{h}$, respectively), Gauss' divergence theorem in the space-time $E_{2}$ implies that Eq. (1.1) is the differential form of the integral conservation law

$$
\begin{equation*}
\oint_{S(V)} \vec{h} \cdot d \vec{s}=0 \tag{1.2}
\end{equation*}
$$

As depicted in Fig. 1, here (i) $S(V)$ is the boundary of an arbitrary space-time region $V$ in $E_{2}$, and (ii) $d \vec{s}=d \sigma \vec{n}$ with $d \sigma$ and $\vec{n}$, respectively, being the area and the unit outward normal vector of a surface element on $S(V)$. Note that: (i) because $\vec{h} \cdot d \vec{s}$ is the space-time flux of $\vec{h}$ leaving the region $V$ through the surface element $d \vec{s}$, Eq. (1.2) simply states that the total space-time flux of $\vec{h}$ leaving $V$ through $S(V)$ vanishes; (ii) in $E_{2}, d \sigma$ is the length of a differential line segment on the simple closed curve $S(V)$; and (iii) all mathematical operations can be carried out as though $E_{2}$ were an ordinary two-dimensional Euclidean space.

To proceed, let $\Omega_{1}$ denote the set of all space-time staggered mesh points $(j, n)$ in $E_{2}$ (dots in Fig. 2(a)), where $n=0, \pm 1 / 2, \pm 1, \pm 3 / 2, \pm 2, \ldots$, and, for each $n, j=n \pm 1 / 2, n \pm 3 / 2, n \pm 5 / 2, \ldots$ Thus $(j+n)$ is an half integer for each $(j, n) \in \Omega_{1}$. Each $(j, n) \in \Omega_{1}$ is associated with a solution element, i.e., $\mathrm{SE}(j, n)$. By definition, $\mathrm{SE}(j, n)$ is the interior of the space-time region bounded by a dashed curve depicted in Fig. 2(b). It includes a horizontal line segment, a vertical line segment, and their immediate neighborhood. By definition, the end points of the line segments referred to above are excluded from $\mathrm{SE}(j, n)$ so that two SEs will not overlap.

Eq. (1.2) will be simulated numerically assuming that, for any $(x, t) \in \mathrm{SE}(j, n), u(x, t)$ and $\vec{h}(x, t)$, respectively, are approximated by

$$
\begin{equation*}
u^{*}(x, t ; j, n) \stackrel{\text { def }}{=} u_{j}^{n}+\left(u_{x}\right)_{j}^{n}\left(x-x_{j}\right)+\left(u_{t}\right)_{j}^{n}\left(t-t^{n}\right) \quad(j, n) \in \Omega_{1} \tag{1.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\vec{h}^{*}(x, t ; j, n) \stackrel{\text { def }}{=}\left(a u^{*}(x, t ; j, n), u^{*}(x, t ; j, n)\right) \quad(x, t) \in \mathrm{SE}(j, n) \text { and }(j, n) \in \Omega_{1} \tag{1.4}
\end{equation*}
$$

Note that: (i) $u_{j}^{n},\left(u_{x}\right)_{j}^{n}$, and $\left(u_{t}\right)_{j}^{n}$ are constants in $\operatorname{SE}(j, n)$, and (in a rough sense) they can be considered as the numerical analogues of the values of $u, \partial u / \partial x$, and $\partial u / \partial t$ at the mesh point $(j, n)$, respectively, (ii) $\left(x_{j}, t^{n}\right)$ are the coordinates of the mesh point $(j, n)$ with $x_{j}=j \Delta x$ and $t^{n}=n \Delta t$, and (iii) Eq. (1.4) is the numerical analogue of the definition $\vec{h}=(a u, u)$.

Let $u=u^{*}(x, t ; j, n)$ satisfy Eq. (1.1) within $\operatorname{SE}(j, n)$. Then one has

$$
\begin{equation*}
\left(u_{t}\right)_{j}^{n}=-a\left(u_{x}\right)_{j}^{n} \quad(j, n) \in \Omega_{1} \tag{1.5}
\end{equation*}
$$

As a result, Eq. (1.3) reduces to

$$
\begin{equation*}
u^{*}(x, t ; j, n)=u_{j}^{n}+\left(u_{x}\right)_{j}^{n}\left[\left(x-x_{j}\right)-a\left(t-t^{n}\right)\right] \quad(j, n) \in \Omega_{1} \tag{1.6}
\end{equation*}
$$

i.e., $u_{j}^{n}$ and $\left(u_{x}\right)_{j}^{n}$ are the only independent mesh variables associated with $(j, n)$.

Let $E_{2}$ be divided into non-overlapping rectangular regions (see Fig. 2(a)) referred to as basic conservation elements (BCEs). As depicted in Figs. 2(c)-2(e), (i) each $(j, n) \in \Omega_{1}$ is assigned with two BCEs, i.e., $\mathrm{CE}_{-}(j, n)$ and $\mathrm{CE}_{+}(j, n)$; (ii) each BCE has one and only one pair of diagonally opposite vertices which belong to $\Omega_{1}$; (iii) the space-time $E_{2}$ can be filled by $\mathrm{CE}_{-}(j, n)$ and $\mathrm{CE}_{+}(j, n),(j, n) \in \Omega_{1}$; and (iv) $\mathrm{CE}(j, n)$, which is the union of $\mathrm{CE}_{-}(j, n)$ and $\mathrm{CE}_{+}(j, n)$, is referred to as a compounded conservation element (CCE).

Note that, among the line segments forming the boundary of $\mathrm{CE}_{-}(j, n), \overline{A B}$ and $\overline{A D}$ belong to $\mathrm{SE}(j, n)$, while $\overline{C B}$ and $\overline{C D}$ belong to $\mathrm{SE}(j-1 / 2, n-1 / 2)$. Similarly, among the line segments forming the boundary of $\mathrm{CE}_{+}(j, n), \overline{A F}$ and $\overline{A D}$ belong to $\mathrm{SE}(j, n)$, while $\overline{E F}$ and $\overline{E D}$ belong to $\mathrm{SE}(j+1 / 2, n-1 / 2)$. As will be shown, by imposing the two local flux conservation conditions at each $(j, n) \in \Omega_{1}$, i.e.,

$$
\begin{equation*}
\oint_{S\left(\mathrm{CE}_{+}(j, n)\right)} \vec{h}^{*} \cdot d \vec{s}=0 \quad(j, n) \in \Omega_{1} \tag{1.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\oint_{S\left(\mathrm{CE}_{-}(j, n)\right)} \vec{h}^{*} \cdot d \vec{s}=0 \quad(j, n) \in \Omega_{1} \tag{1.8}
\end{equation*}
$$

one can obtain two equations for two unknowns $u_{j}^{n}$ and $\left(u_{x}\right)_{j}^{n}$.
Let

$$
\begin{equation*}
\nu \stackrel{\text { def }}{=} a \frac{\Delta t}{\Delta x} \quad \text { and } \quad\left(u_{\bar{x}}\right)_{j}^{n} \stackrel{\text { def }}{=} \frac{\Delta x}{4}\left(u_{x}\right)_{j}^{n} \quad(j, n) \in \Omega_{1} \tag{1.9}
\end{equation*}
$$

where $\nu$ is the Courant number, and $\left(u_{\bar{x}}\right)_{j}^{n}$ is the normalized version of $\left(u_{x}\right)_{j}^{n}$. Then, with the aid of Eqs. (1.4), (1.6) and (1.9), it can be shown that Eqs. (1.7) and (1.8) are equivalent to

$$
\begin{equation*}
(1-\nu)\left[u+(1+\nu) u_{\bar{x}}\right]_{j}^{n}=(1-\nu)\left[u-(1+\nu) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2} \quad(j, n) \in \Omega_{1} \tag{1.10}
\end{equation*}
$$

and

$$
\begin{equation*}
(1+\nu)\left[u-(1-\nu) u_{\bar{x}}\right]_{j}^{n}=(1+\nu)\left[u+(1-\nu) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2} \quad(j, n) \in \Omega_{1} \tag{1.11}
\end{equation*}
$$

respectively. To simplify notation, in the above and hereafter we adopt a convention that can be explained using an expression on the left side of Eq. (1.10) as an example, i.e.,

$$
\left[u+(1+\nu) u_{\bar{x}}\right]_{j}^{n}=u_{j}^{n}+(1+\nu)\left(u_{\bar{x}}\right)_{j}^{n}
$$

By adding Eqs. (1.10) and (1.11) together, one has

$$
\begin{equation*}
u_{j}^{n}=\frac{1}{2}\left\{(1-\nu)\left[u-(1+\nu) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}+(1+\nu)\left[u+(1-\nu) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega_{1} \tag{1.12}
\end{equation*}
$$

Let $1-\nu^{2} \neq 0$, i.e., $1-\nu \neq 0$ and $1+\nu \neq 0$. Then Eqs. (1.10) and (1.11) can be divided by $(1-\nu)$ and $(1+\nu)$, respectively. By subtracting the resulting equations from each other, one has

$$
\begin{equation*}
\left(u_{\bar{x}}\right)_{j}^{n}=\frac{1}{2}\left\{\left[u-(1+\nu) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}-\left[u+(1-\nu) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega_{1} \tag{1.13}
\end{equation*}
$$

Hereafter, by definition the $a$ scheme is formed by Eq. (1.12) and (1.13) for any $\nu$. According the above derivation, a solution to the $a$ scheme will satisfy Eqs. (1.10) and (1.11) and therefore the two conservation conditions Eqs. (1.7) and (1.8) for any $\nu$. However, because Eqs. (1.7) and (1.8) $\Rightarrow$ (i.e., "imply") Eq. (1.12) for any $\nu$, but they $\Rightarrow$ Eq. (1.13) only if $\nu^{2} \neq 1$, one concludes that a solution to Eqs. (1.7) and (1.8) may not be a solution to the $a$ scheme if $\nu^{2}=1$. As such the $a$ scheme represents a stronger condition than that of Eq. (1.7) and (1.8).

The $a$ scheme is an explicit and neutrally stable (i.e., non-dissipative) solver of Eq. (1.1) in its stability domain $|\nu|<1$ [63]. Also, according to Eqs. (1.12) and (1.13), it has a space-time staggered stencil, i.e., any $(j, n) \in \Omega_{1}$ is associated with a stencil formed by the three mesh points $(j, n),(j-1 / 2, n-1 / 2)$ and $(j+1 / 2, n-1 / 2)$.

Eqs. (1.7) and (1.8), i.e., the flux conservation conditions over the two BCEs $\mathrm{CE}_{-}(j, n)$ and $\mathrm{CE}_{+}(j, n)$ for all $(j, n) \in \Omega_{1}$, are enforced by any solution to the $a$ scheme. Because (i) the vector $\vec{h}^{*}$ at any surface element lying on any interface separating two neighboring BCE is evaluated using the information from a
single SE, and (ii) the unit outward normal vector on the surface element pointing outward from one of these two neighboring BCEs is exactly the negative of that pointing outward from another BCE, one concludes that the flux leaving one of these BCEs through the interface is the negative of that leaving another BCE through the same interface. As a result of this flux cancellation, the local flux conservation relations Eqs. (1.7) and (1.8) lead to a global flux conservation relation, i.e., the total flux of $\vec{h}^{*}$ leaving the boundary of any spacetime region that is the union of any combination of BCEs will also vanish. In particular, because $\mathrm{CE}(j, n)$ is the union of $\mathrm{CE}_{-}(j, n)$ and $\mathrm{CE}_{+}(j, n)$,

$$
\begin{equation*}
\oint_{S(\mathrm{CE}(j, n))} \vec{h}^{*} \cdot d \vec{s}=0 \quad(j, n) \in \Omega_{1} \tag{1.14}
\end{equation*}
$$

must follow from Eqs. (1.7) and(1.8). In fact, it can be shown that Eq. (1.14) is equivalent to Eq. (1.12).
There is a family of the dissipative extensions of the $a$ scheme in which only the less stringent conservation condition Eq. (1.14) is assumed [5,26,31,63]. Because Eq. (1.14) is equivalent to Eq. (1.12), for each of these extensions, $u_{j}^{n}$ is still evaluated using Eq. (1.12) while $\left(u_{\bar{x}}\right)_{j}^{n}$ is evaluated in terms of the mesh variables at $(j-1 / 2, n-1 / 2)$ and $(j+1 / 2, n-1 / 2)$ using an equation different from Eq. (1.13). In fact, (i) the $a-\epsilon$ scheme is formed by Eq. (1.12) and

$$
\begin{equation*}
\left(u_{\bar{x}}\right)_{j}^{n}=\frac{1}{2}\left\{\left[(1-\epsilon) u+(2 \epsilon-1-\nu) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}-\left[(1-\epsilon) u+(1-\nu-2 \epsilon) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega_{1} \tag{1.15}
\end{equation*}
$$

where $\epsilon$ is a real parameter, and (ii) the $c-\tau$ scheme is formed by Eq. (1.12) and

$$
\begin{equation*}
\left(u_{\bar{x}}\right)_{j}^{n}=\frac{1}{2(1+\tau)}\left\{\left[u-(1+2 \nu-\tau) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}-\left[u+(1-2 \nu-\tau) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega_{1} \tag{1.16}
\end{equation*}
$$

where $\tau$ is a real parameter $\neq-1$.
Note that:
(a) By comparing Eqs. (1.13) and (1.15), it is seen that the $a$ scheme is the special case of the $a-\epsilon$ scheme when $\epsilon=0$.
(b) By using Jordan canonical form theorem [71], it can be shown that the $a-\epsilon$ scheme is von Neumann stable $\Leftrightarrow$ (i.e., "if and only if") either

$$
\begin{equation*}
0 \leq \epsilon \leq 1 \quad \text { and } \quad|\nu|<1 \tag{1.17}
\end{equation*}
$$

or

$$
\begin{equation*}
\epsilon=|\nu|=1 \tag{1.18}
\end{equation*}
$$

Also one can show that the $a-\epsilon$ scheme becomes more dissipative as the value of $\epsilon$ increases from 0 to 1 [5].
(c) By using Jordan canonical form theorem [71], it is shown rigorously in [46] that the $c-\tau$ scheme is von Neumann stable $\Leftrightarrow$

$$
\begin{equation*}
\nu^{2} \leq 1, \quad \tau \geq \tau_{o}\left(\nu^{2}\right), \quad \text { and } \quad\left(\nu^{2}, \tau\right) \neq(1,1) \tag{1.19}
\end{equation*}
$$

where

$$
\tau_{o}(s) \stackrel{\text { def }}{=} \begin{cases}0 & \text { if } s=0  \tag{1.20}\\ \frac{4-s-2 \sqrt{2\left(2-s-s^{2}\right)}}{s} & \text { if } 0<s \leq \frac{3}{11} \\ \frac{s-1+\sqrt{1-2 s+5 s^{2}}}{2 s} & \text { if } \frac{3}{11} \leq s \leq 1\end{cases}
$$

It can be shown that: (i) $\tau_{o}(s)$ is continuous at $s=0$, i.e.,

$$
\begin{equation*}
\lim _{s \rightarrow 0^{+}} \tau_{o}(s)=\tau_{o}(0)=0 \tag{1.21}
\end{equation*}
$$

(ii) $\tau_{o}(s)$ is consistently defined at $s=3 / 11$, i.e., both the second and third expressions on the right side of Eq. (1.20) implies that $\tau_{o}(3 / 11)=1 / 3$; (iii) $\tau_{o}(s)$ is continuously differentiable at $s=3 / 11$, i.e.,

$$
\begin{equation*}
\lim _{s \rightarrow \frac{3}{11}-} \tau_{o}^{\prime}(s)=\lim _{s \rightarrow \frac{3}{11}^{+}} \tau_{o}^{\prime}(s)=121 / 90 \tag{1.22}
\end{equation*}
$$

where $\tau_{o}^{\prime}(s) \stackrel{\text { def }}{=} d \tau_{o}(s) / d s$; (iv) $\tau_{o}(s)$ is strictly monotonically increasing in the interval $0<s<1$; (v) $\tau_{o}(1)=1$; and (vi)

$$
\begin{equation*}
s<\tau_{o}(s)<\sqrt{s}, \quad 0<s<1 \tag{1.23}
\end{equation*}
$$

Moreover, (i) the stability conditions Eq. (1.19) agree completely with those generated through numerical experiments [31] (Note: the reader is warned that the condition $\tau \geq \tau_{o}\left(\nu^{2}\right)$ that appears in Eq. (1.19) here is expressed as $\tau \geq \tau_{o}(|\nu|)$ in Eq. (3.12) of [31], i.e., the function $\tau_{o}$ defined here is different from the function $\tau_{o}$ introduced in [31]), and (ii) for any given fixed value of $|\nu|<1$, the $c-\tau$ scheme tends to become more dissipative as the value of $\tau$ increases from its minimum stability bound $\tau_{o}\left(\nu^{2}\right)$.
The $a$ scheme and its extensions were defined over $\Omega_{1}$, the set of space-time staggered mesh points marked by dots in Fig. 2(a). Obviously, the $a$ scheme can also be defined over $\Omega_{2}$, the set of all space-time staggered mesh points $(j, n)$ in $E_{2}$ where $n=0, \pm 1 / 2, \pm 1, \pm 3 / 2, \pm 2, \ldots$, and, for each $n, j=n, n \pm 1, n \pm 2, n \pm 3, \ldots$.. By definition, each $(j, n) \in \Omega_{2}$ is associated with an whole-integer value of $(j+n)$, and represented by a mesh point not marked by dots in Fig. 2(a). Because one can carry out the time marching (using the $a$ scheme or any one of its 1D extensions) exclusively over either $\Omega_{1}$ or $\Omega_{2}$, there is no need to carry out a marching over both $\Omega_{1}$ and $\Omega_{2}$. Note that this conclusion generally remains valid even if an 1D nonuniform spatial mesh is used [31]. The exception may occur for a marching scheme where the time-step size varies from one spatial region to another such as the local time stepping procedure described in [42].

Note that the two disjoint sets $\Omega_{1}$ and $\Omega_{2}$ have the property that, if a space-time mesh point $(j, n)$ belongs to $\Omega_{1}\left(\Omega_{2}\right)$, then its four immediate space-time neighbors $(j \pm 1 / 2, n)$ and $(j, n \pm 1 / 2)$ must belong to $\Omega_{2}\left(\Omega_{1}\right)$. As a result, for a scheme with a space-time staggered stencil such as the $a$ scheme or the classical leapfrog scheme, this implies that a time marching carried out over $\Omega_{1}$ is completely decoupled from that carried out over $\Omega_{2}$. As such one only needs to carry out time marching over one of these two sets. For this reason, the problem of "solution decoupling" (which is famously associated with the leapfrog scheme) usually can be avoided in the 1D case.

However, as will be shown shortly, for a space-time mesh built from a multidimensional unstructured spatial mesh, one generally cannot define two disjoint sets of space-time mesh points which possess the same property of $\Omega_{1}$ and $\Omega_{2}$ discussed above. As such, the problem of solution decoupling may become a real nuisance in multidimensional CESE simulations involving unstructured spatial meshes.

Is was shown earlier that, for an 1D CESE scheme, (i) each mesh point is assigned with two BCEs; and (ii) for each dependent physical variable $u$, each mesh point $(j, n)$ is assigned with two independent mesh variables $u_{j}^{n}$ and $\left(u_{x}\right)_{j}^{n}$. Similarly [12,13], for a 2D (3D) CESE scheme, (i) each space-time solution point (to be defined shortly) is assigned with three (four) BCEs and one CCE, with the CCE being the union of the three (four) BCEs; and (ii) for each dependent physical variable, each space-time solution point is assigned with three (four) independent mesh variables - note that there are two (three) independent spatial derivatives in a problem with two (three) spatial dimension. Because a triangle and a tetrahedron have three and four sides, respectively, triangles and tetrahedra, respectively, become the simplest and most natural building blocks of the spatial meshes for $2 D$ and $3 D$ CESE schemes [10-13,16]. As such, in the following discussion about solution decoupling in a CESE simulation involving an unstructured 2D spatial domain, we consider an unstructured domain built from triangles.

Consider a spatial domain formed by triangles of arbitrary shapes (see Fig. 3). Here (i) any two neighboring triangles share a common side, and (ii) the vertices and centroids of triangles are marked by dots and circles, respectively. As shown in Fig. 3, the triangle $\triangle F B D$ has three neighboring triangles. Let (i) point $G$ be the centroid of $\triangle F B D$, and (ii) points $A, C$ and $E$, respectively, be the centroids of the three neighbors of $\triangle F B D$. Then points $A, B, C, D, E$, and $F$ form a hexagon. The centroid of the hexagon
is denoted by $G^{\prime}$ and marked by a cross. Hereafter point $G^{\prime}$ is referred to as the spatial solution point of $\triangle F B D$. In a similar manner, we define spatial solution points $A^{\prime}, C^{\prime}, E^{\prime}, \ldots$ of other triangles.

A space-time solution point is a mesh point which resides at a time level $n=0,1 / 2,1,3 / 2,2, \ldots(t=$ $t^{n}=n \Delta t$ at the $n$th time level) and has a spatial solution point being its spatial projection. Given any spatial solution point $G^{\prime}$ and any time level $n$, let $\left(G^{\prime}, n\right)$ denote the space-time solution point which resides at the $n$th time level and has point $G^{\prime}$ being its spatial projection. As explained in [13], the three BCEs associated with $\left(G^{\prime}, n\right)$ are constructed such that their spatial projections are the quadrilaterals $G F A B$, $G B C D$, and $G D E F$, respectively. Also, for a reason explained in [13], independent mesh variables are stored at space-time solution points.

As shown in [13], in a 2D CESE scheme using a triangular spatial mesh, the mesh variables at $\left(G^{\prime}, n\right)$ are evaluated in terms of those at $\left(A^{\prime}, n-1 / 2\right),\left(C^{\prime}, n-1 / 2\right)$ and $\left(E^{\prime}, n-1 / 2\right)$, i.e., the scheme has a space-time staggered stencil formed by one point at the $n$th time level and three points at the ( $n-1 / 2$ ) th time level. Also it was shown in [13] that, given an unstructured triangular spatial mesh, the time marching cannot be carried out independently over two disjoint sets of space-time solution points unless the triangles forming the spatial mesh can be divided into two disjoint sets $\Psi_{1}$ and $\Psi_{2}$ such that two triangles sharing a common side always belong to two different sets. As will be shown, because point $F$ is the common vertex of an odd number (5) of triangles, the triangles depicted in Fig. 3 cannot be divided into $\Psi_{1}$ and $\Psi_{2}$.

To prove the last proposition by contradiction, let $\triangle F B D \in \Psi_{1}$. Then because $\triangle F H B$ and $\triangle F B D$ share a common side, by definition, $\triangle F H B \in \Psi_{2}$. Similarly, because (i) $\triangle F N H$ and $\triangle F H B$ share a common side, (ii) $\triangle F M N$ and $\triangle F N H$ share a common side, and (iii) $\triangle F D M$ and $\triangle F M N$ share a common side, one concludes that $\triangle F N H \in \Psi_{1}, \triangle F M N \in \Psi_{2}$, and $\triangle F D M \in \Psi_{1}$. As such, we have arrive at the contradiction that both $\triangle F B D$ and $\triangle F D M$ belong to $\Psi_{1}$ in spite of the fact that they share a common side. Because we will reach a similar contradiction by assuming $\triangle F B D \in \Psi_{2}$, the proposition has been proved.

The triangles depicted in Fig. 3 cannot be divided into two disjoint sets $\Psi_{1}$ and $\Psi_{2}$. As such, according to another proposition stated earlier, the CESE time marching cannot be carried out independently over two disjoint sets of space-time solution points constructed from the triangular spatial mesh depicted in Fig. 3. Moreover, because the incidence that an odd number of triangles sharing a common vertex occurs regularly in an unstructured triangular mesh, the CESE time marching generally cannot be carried out independently over two disjoint sets of space-time solution points if they are constructed from an unstructured triangular spatial mesh.

Similarly, for a 3D CESE scheme using an unstructured spatial mesh built from tetrahedra, it can be shown that generally the time marching cannot be carried out independently over two disjoint sets of space-time staggered solution points.

According to the above discussions, a CESE multidimensional time marching using an unstructured triangular or tetrahedral spatial mesh generally will involve all space-time solution points. This coupled with the fact that a CESE scheme has a space-time staggered stencil implies that solution decoupling may become a real nuisance. To deal with this problem once and for all, recently new CESE schemes have been developed such that possibility of solution decoupling can be removed completely. The description of these new schemes is one of the topics covered in the current paper. As will be shown, (i) these new schemes are constructed without compromising the basic CESE ideas, and (ii) their effectiveness has been verified by real-world applications. [62, 64-67].

To describe another topic of the current paper, note that the spatial solution point $G^{\prime}$ depicted in Fig. 3 is the centroid of the hexagon $A B C D E F$ where points $A, C$ and $E$, respectively, are rather arbitrarily designated as the centroids of the three neighbors of $\triangle F B D$. In fact, to insure that CEs will not overlap in space and thus no loss of local and global flux conservation will occur, it only requires that points such as $G, A, C$ and $E$ lie in the interiors of $\triangle F B D, \triangle F H B, \triangle B L D$, and $\triangle F D M$, respectively.

For the rather benign triangular domain depicted in Fig. 3, the hexagon $A B C D E F$ happens to be convex, i.e., a line segment joining any two points in this hexagon always lies entirely in it. Because the centroid of a convex polygon always lies within its interior, the centroid (i.e., point $G^{\prime}$ ) of the hexagon $A B C D E F$ lies within its interior. However, this may not be true for a more pathological triangular domain
formed from triangles with very large aspect ratios and sharply different orientations. As an example, consider Fig. 4 where $\triangle F B D$ is again surrounded by three neighboring triangles $\triangle F H B, \triangle B L D$ and $\triangle F D M$. However, in Fig. 4, the aspect ratio of each of the four triangles depicted is much larger than that of its counterpart depicted in Fig. 3. Moreover, in Fig. 4, the orientations of the sides $\overline{B D}$ and $\overline{F D}$ of $\triangle F B D$ are sharply different from those of the sides $\overline{B H}$ and $\overline{F H}$ of $\triangle F H B$, respectively. Let points $G, A$, $C$ and $E$ depicted in Fig. 4 again be the centroids of $\triangle F B D, \triangle F H B \triangle B L D$, and $\triangle F D M$, respectively. Then the resulting hexagon $A B C D E F$ may become highly concave and thus its centroid $G^{\prime}$ (not shown in Fig. 4) may no longer lie in its interior.

Similarly, a CESE tetrahedral mesh is built from tetrahedra where any two neighboring tetrahedra share a common triangular face. Given any two neighboring tetrahedra, in the original CESE tetrahedralmesh construction, the 5-vertex polyhedron with its vertices being (i) the centroids of these two neighbors and (ii) the three vertices of the triangular face shared by them, becomes the spatial projection of a BCE. Because each tetrahedron is associated with four neighbors, it is associated with four 5 -vertex polyhedra and their union-a 8-vertex polyhedron. The 8 -vertex polyhedron and its centroid, respectively, are the spatial projections of a CCE and a spatial solution point. For some pathological case where the 8 -vertex polyhedron becomes highly concave, its centroid (i.e., the spatial solution point associated with this polyhedron) may lie outside of its interior.

In a real-world CFD simulation near a solid wall, one often encounters a case where a boundary with high curvature or sharp corner is surrounded by triangular/tetrahedral meshes of extremely high aspect ratio (up to $10^{6}$ ). For such an extreme case, a spatial solution point may lie far outside of the associated hexagon/8-vertex polyhedron. It could even be embedded beyond a solid-wall boundary and causes serious numerical difficulties. To overcome this problem, a new way of building spatial meshes from triangles and tetrahedra has been developed. Specifically, the centroids of triangles depicted in Fig. 3 have been replaced by the corresponding in-centers or other specially-defined interior points in the new construction. However, the solution point is still the centroid of the new hexagon formed by vertices and in-centers (or other specially-defined interior points). Similarly, centroids of tetrahedra also are replaced by corresponding incenters or other specially-defined interior points in the new construction. According to the numerical results obtained $[62,64,65]$, the new construction has the following advantages: (i) it can properly model curved boundaries-without it, computed surface gradient such as heat flux will be much less accurate (and noisier in many cases); (ii) it can reduce aspect ratios of triangular or tetrahedral cells used in gradient calculations in a highly stretched mesh and, as a result, can improve the overall accuracy for triangular and tetrahedral meshes; (iii) for vortex-shedding problems, better resolution is observed by using the in-center formulation. In this paper we will describe this new development and also provide the theoretical justification for the new construction.

As mentioned earlier, mesh values of the physical dependent variables and their spatial derivatives are treated as independent mesh variables in a CESE scheme. As it turns out, for a case where a triangular or tetrahedral mesh is used, numerical evaluation of the spatial derivatives using the established CESE approach tends to become less accurate as the mesh cell aspect ratio becomes larger. In the latest development, the root cause behind this tendency was clearly identified and several remedies to overcome it were found through a rigorous mathematical analysis. However, because of the length of the current paper and the complexity of the mathematics involved, this latest work will be presented in another paper.

Because it is rather simple and straightforward, in the following, first we will describe the new way of building spatial meshes from triangles and tetrahedra. The new flux-conserving solution-coupling procedures will be presented in Sec. 3 .

## 2. New construction of triangular and tetrahedral meshes

### 2.1. Triangular meshes

Again consider the triangle $\triangle F B D$ and the hexagon $A B C D E F$ depicted in Fig. 4. As noted earlier, (i) because it is the centroid of the hexagon $A B C D E F$, the spatial solution point $G^{\prime}$ (not shown in Fig. 4) always lies in the interior of the hexagon if it is convex but may lie outside of it if it is concave; and (ii) serious numerical difficulties may arise if point $G^{\prime}$ lies outside of the hexagon $A B C D E F$. As such, to avoid
numerical difficulties, it is essential to ensure that $A B C D E F$ is a convex hexagon.
Because any triangle such as $\triangle F B D$ is always convex, the hexagon $A B C D E F$ will also be convex if its boundary more or less coincides with that of $\triangle F B D$. In other words, the hexagon is convex if points $A, C$, and $E$ are located close enough to the sides $\overline{B F}, \overline{B D}$ and $\overline{F D}$, respectively. Unfortunately, because (i) $A$ is the centroid of $\triangle F H B$, and (ii)

$$
\begin{equation*}
|\overline{B F}| \ll|\overline{B H}| \quad \text { and } \quad|\overline{B F}| \ll|\overline{F H}| \tag{2.1}
\end{equation*}
$$

(hereafter, as an example, the length of $\overline{B F}$ is denoted by $|\overline{B F}|$ ), point $A$ generally is far away from the side $\overline{B F}$. As such the hexagon $A B C D E F$ may become concave. For this reason and the previously stated requirement that point $A, C$, and $E$ must lie in the interiors of $\triangle F H B, \triangle B L D$, and $\triangle F D M$, respectively, in the following, we will describe a new approach in which the locations of $A, C$, and $E$ are chosen so that: (i) each of them again is an interior point of the associated triangle, and (ii) a point such as $A$ will remain close to the side $\overline{B F}$ even for the pathological case Eq. (2.1).

Consider Fig. 5. Here (i) points $P_{1}, P_{2}, P_{3}$ are the vertices of the triangle $\triangle P_{1} P_{2} P_{3}$ which lies on the $x-y$ plane and has its area

$$
\begin{equation*}
A\left(\triangle P_{1} P_{2} P_{3}\right)>0 \tag{2.2}
\end{equation*}
$$

(ii) $P_{0}$ is a point on the $x$ - $y$ plane to be defined shortly; (iii) for each $k=0,1,2,3, x_{k}$ and $y_{k}$ are the $x$ - and $y$ - coordinates of point $P_{k}$, respectively; (iv) point $Q_{1}$ is the point on $\overline{P_{2} P_{3}}$ (or its extension) which meets the condition

$$
\begin{equation*}
\overline{P_{0} Q_{1}} \perp \overline{P_{2} P_{3}} \tag{2.3}
\end{equation*}
$$

i.e., $Q_{1}$ is the perpendicular projection of $P_{0}$ on the straight line which contains $\overline{P_{2} P_{3}} ;(\mathrm{v})$ point $Q_{2}$ is the point on $\overline{P_{3} P_{1}}$ (or its extension) which meets the condition

$$
\begin{equation*}
\overline{P_{0} Q_{2}} \perp \overline{P_{3} P_{1}} \tag{2.4}
\end{equation*}
$$

i.e., $Q_{2}$ is the perpendicular projection of $P_{0}$ on the straight line which contains $\overline{P_{3} P_{1}}$; (vi) point $Q_{3}$ is the point on $\overline{P_{1} P_{2}}$ (or its extension) which meets the condition

$$
\begin{equation*}
\overline{P_{0} Q_{3}} \perp \overline{P_{1} P_{2}} \tag{2.5}
\end{equation*}
$$

i.e., $Q_{3}$ is the perpendicular projection of $P_{0}$ on the straight line which contains $\overline{P_{1} P_{2}}$; (vii) $\alpha_{1}, \alpha_{2}$, and $\alpha_{3}$ are the internal angles of $\triangle P_{1} P_{2} P_{3}$ facing the sides $\overline{P_{2} P_{3}}, \overline{P_{3} P_{1}}$, and $\overline{P_{1} P_{2}}$, respectively; (viii)

$$
\begin{equation*}
\ell_{1} \stackrel{\text { def }}{=}\left|\overline{P_{2} P_{3}}\right|, \quad \ell_{2} \stackrel{\text { def }}{=}\left|\overline{P_{3} P_{1}}\right| \quad \text { and } \quad \ell_{3} \stackrel{\text { def }}{=}\left|\overline{P_{1} P_{2}}\right| \tag{2.6}
\end{equation*}
$$

(ix)

$$
\begin{equation*}
h_{k} \stackrel{\text { def }}{=}\left|\overline{P_{0} Q_{k}}\right| \quad k=1,2,3 \tag{2.7}
\end{equation*}
$$

and $(\mathrm{x})$ the $x-y-z$ coordinate system depicted in Fig. 5 is a right-handed system, i.e.,

$$
\begin{equation*}
\vec{e}_{x} \times \vec{e}_{y}=\vec{e}_{z}, \quad \vec{e}_{y} \times \vec{e}_{z}=\vec{e}_{x} \quad \text { and } \quad \vec{e}_{z} \times \vec{e}_{x}=\vec{e}_{y} \tag{2.8}
\end{equation*}
$$

where $\vec{e}_{x}, \vec{e}_{y}$, and $\vec{e}_{z}$ are the unit vectors in the $x$-, $y$-, and $z$ - directions, respectively.
Note that, according to Eq. (2.6) and Fig. 5, for each $k=1,2,3, \ell_{k}$ denotes the length of the side of $\triangle P_{1} P_{2} P_{3}$ facing the vertex $P_{k}$. Also from the definition of points $Q_{1}, Q_{2}$, and $Q_{3}$, and that of $h_{k}, k=1,2,3$, one concludes that $h_{1}, h_{2}$, and $h_{3}$, respectively, are the (shortest) distances from $P_{0}$ to the straight lines that contain the sides $\overline{P_{2} P_{3}}, \overline{P_{3} P_{1}}$, and $\overline{P_{1} P_{2}}$, respectively.

For each $k=0,1,2,3$, let $\vec{P}_{k}$ be the position vector of point $P_{k}$, i.e.,

$$
\begin{equation*}
\vec{P}_{k} \stackrel{\text { def }}{=} x_{k} \vec{e}_{x}+y_{k} \vec{e}_{y} \quad k=0,1,2,3 \tag{2.9}
\end{equation*}
$$

Also let $\vec{P}_{0}$ be an weighted average of $\vec{P}_{k}, k=1,2,3$, i.e.,

$$
\begin{equation*}
\vec{P}_{0}=w_{1} \vec{P}_{1}+w_{2} \vec{P}_{2}+w_{3} \vec{P}_{3} \tag{2.10}
\end{equation*}
$$

where $w_{1}, w_{2}$, and $w_{3}$ are real weight factors with

$$
\begin{equation*}
w_{1}+w_{2}+w_{3}=1 \tag{2.11}
\end{equation*}
$$

Moreover, given a pair of integers $k_{1}$ and $k_{2}$ with $k_{1} \neq k_{2}$ and $k_{1}, k_{2}=0,1,2,3$, let $\overrightarrow{P_{k_{1}} P_{k_{2}}}$ be the displacement vector joining points $P_{k_{1}}$ and $P_{k_{2}}$ and pointing in the direction from $P_{k_{1}}$ to $P_{k_{2}}$, i.e.,

$$
\begin{equation*}
\overrightarrow{P_{k_{1}} P_{k_{2}}} \stackrel{\text { def }}{=} \vec{P}_{k_{2}}-\vec{P}_{k_{1}}, \quad k_{1} \neq k_{2} \quad \text { and } \quad k_{1}, k_{2}=0,1,2,3 \tag{2.12}
\end{equation*}
$$

Then, by using Eq. (2.12) and other well-known properties of the vector product such as

$$
\begin{equation*}
\vec{A} \times \vec{A}=0 \quad \text { and } \quad \vec{A} \times \vec{B}=-\vec{B} \times \vec{A} \tag{2.13}
\end{equation*}
$$

for any vectors $\vec{A}$ and $\vec{B}$, one can show that

$$
\begin{equation*}
\overrightarrow{P_{1} P_{2}} \times \overrightarrow{P_{1} P_{3}}=\left(\vec{P}_{2}-\vec{P}_{1}\right) \times\left(\vec{P}_{3}-\vec{P}_{1}\right)=\vec{P}_{1} \times \vec{P}_{2}+\vec{P}_{2} \times \vec{P}_{3}+\vec{P}_{3} \times \vec{P}_{1} \tag{2.14}
\end{equation*}
$$

Similarly, with the aid of Eqs. (2.10)-(2.14), one can show that

$$
\begin{align*}
& \overrightarrow{P_{0} P_{2}} \times \overrightarrow{P_{0} P_{3}}=w_{1}\left(\vec{P}_{1} \times \vec{P}_{2}+\vec{P}_{2} \times \vec{P}_{3}+\vec{P}_{3} \times \vec{P}_{1}\right)=w_{1}\left(\overrightarrow{P_{1} P_{2}} \times \overrightarrow{P_{1} P_{3}}\right)  \tag{2.15}\\
& \overrightarrow{P_{0} P_{3}} \times \overrightarrow{P_{0} P_{1}}=w_{2}\left(\vec{P}_{1} \times \vec{P}_{2}+\vec{P}_{2} \times \vec{P}_{3}+\vec{P}_{3} \times \vec{P}_{1}\right)=w_{2}\left(\overrightarrow{P_{1} P_{2}} \times \overrightarrow{P_{1} P_{3}}\right)  \tag{2.16}\\
& \overrightarrow{P_{0} P_{1}} \times \overrightarrow{P_{0} P_{2}}=w_{3}\left(\vec{P}_{1} \times \vec{P}_{2}+\vec{P}_{2} \times \vec{P}_{3}+\vec{P}_{3} \times \vec{P}_{1}\right)=w_{3}\left(\overrightarrow{P_{1} P_{2}} \times \overrightarrow{P_{1} P_{3}}\right) \tag{2.17}
\end{align*}
$$

Next, by using (i) a well-known property of the vector product, (ii) the fact that $A\left(\triangle P_{1} P_{2} P_{3}\right)$ denotes the area of $\triangle P_{1} P_{2} P_{3}$, and (iii) Fig. 5, one has

$$
\begin{equation*}
A\left(\triangle P_{1} P_{2} P_{3}\right)=\left(\left|\overrightarrow{P_{1} P_{2}}\right|\left|\overrightarrow{P_{1} P_{3}}\right| \sin \alpha_{1}\right) / 2=\left|\overrightarrow{P_{1} P_{2}} \times \overrightarrow{P_{1} P_{3}}\right| / 2 \tag{2.18}
\end{equation*}
$$

Similarly, with the aid of Eqs. (2.15)-(2.18), one can show that

$$
\begin{align*}
& A\left(\triangle P_{0} P_{2} P_{3}\right)=\left|\overrightarrow{P_{0} P_{2}} \times \overrightarrow{P_{0} P_{3}}\right| / 2=\left|w_{1}\right| \times A\left(\triangle P_{1} P_{2} P_{3}\right)  \tag{2.19}\\
& A\left(\triangle P_{0} P_{3} P_{1}\right)=\left|\overrightarrow{P_{0} P_{3}} \times \overrightarrow{P_{0} P_{1}}\right| / 2=\left|w_{2}\right| \times A\left(\triangle P_{1} P_{2} P_{3}\right) \tag{2.20}
\end{align*}
$$

and

$$
\begin{equation*}
A\left(\triangle P_{0} P_{1} P_{2}\right)=\left|\overrightarrow{P_{0} P_{1}} \times \overrightarrow{P_{0} P_{2}}\right| / 2=\left|w_{3}\right| \times A\left(\triangle P_{1} P_{2} P_{3}\right) \tag{2.21}
\end{equation*}
$$

where $A\left(\triangle P_{0} P_{2} P_{3}\right), A\left(\triangle P_{0} P_{3} P_{1}\right)$, and $A\left(\triangle P_{0} P_{1} P_{2}\right)$ denote the areas of $\triangle P_{0} P_{2} P_{3}, \triangle P_{0} P_{3} P_{1}$, and $\triangle P_{0} P_{1} P_{2}$, respectively. In turn, Eqs. (2.19)-(2.21) $\Rightarrow$

$$
\begin{equation*}
A\left(\triangle P_{0} P_{2} P_{3}\right)+A\left(\triangle P_{0} P_{3} P_{1}\right)+A\left(\triangle P_{0} P_{1} P_{2}\right)=\left(\left|w_{1}\right|+\left|w_{2}\right|+\left|w_{3}\right|\right) \times A\left(\triangle P_{1} P_{2} P_{3}\right) \tag{2.22}
\end{equation*}
$$

Because a triangle such as $\triangle P_{1} P_{2} P_{3}$ is intrinsically convex, according to a theorem on convex sets, point $P_{0}$ defined by Eqs. (2.10) and (2.11) lies on the boundary or in the interior of $\triangle P_{1} P_{2} P_{3} \Leftrightarrow$

$$
\begin{equation*}
0 \leq w_{k} \leq 1 \quad k=1,2,3 \tag{2.23}
\end{equation*}
$$

Moreover, because Eq. (2.2) implies that the interior of $\triangle P_{1} P_{2} P_{3}$ is not empty, $P_{0}$ lies in the interior of $\triangle P_{1} P_{2} P_{3} \Leftrightarrow$

$$
\begin{equation*}
0<w_{k}<1 \quad k=1,2,3 \tag{2.24}
\end{equation*}
$$

For the case where $P_{0}$ lies on the boundary or in the interior of $\triangle P_{1} P_{2} P_{3}$, Eqs. (2.11), (2.22), and (2.23) now $\Rightarrow$ the following obvious result:

$$
\begin{equation*}
A\left(\triangle P_{0} P_{2} P_{3}\right)+A\left(\triangle P_{0} P_{3} P_{1}\right)+A\left(\triangle P_{0} P_{1} P_{2}\right)=A\left(\triangle P_{1} P_{2} P_{3}\right) \tag{2.25}
\end{equation*}
$$

On the other hand, for the case where $P_{0}$ lies outside of $\triangle P_{1} P_{2} P_{3}$, Eq. (2.23) is no longer true, i.e., some of the weight factors must be negative. Because Eq. (2.11) requires that some of the weight factors must be positive, we arrive at the conclusion that

$$
\begin{equation*}
\left|w_{1}\right|+\left|w_{2}\right|+\left|w_{3}\right|>\left|w_{1}+w_{2}+w_{3}\right|=1 \tag{2.26}
\end{equation*}
$$

if $P_{0}$ lies outside of $\triangle P_{1} P_{2} P_{3}$. In turn, Eqs. (2.22) and $(2.26) \Rightarrow$

$$
\begin{equation*}
A\left(\triangle P_{0} P_{2} P_{3}\right)+A\left(\triangle P_{0} P_{3} P_{1}\right)+A\left(\triangle P_{0} P_{1} P_{2}\right)>A\left(\triangle P_{1} P_{2} P_{3}\right) \tag{2.27}
\end{equation*}
$$

if $P_{0}$ lies outside of $\triangle P_{1} P_{2} P_{3}$.
Note that, for any real parameter $\zeta$, Eqs. (2.2) and (2.6) $\Rightarrow$

$$
\begin{equation*}
\ell_{k}>0 \quad \text { and } \quad\left(\ell_{k}\right)^{\zeta}>0 \quad k=1,2,3 ; \quad-\infty<\zeta<+\infty \tag{2.28}
\end{equation*}
$$

As such the weight factors

$$
\begin{equation*}
w_{k}(\zeta) \stackrel{\text { def }}{=} \frac{\left(\ell_{k}\right)^{\zeta}}{\left(\ell_{1}\right)^{\zeta}+\left(\ell_{2}\right)^{\zeta}+\left(\ell_{3}\right)^{\zeta}} \quad k=1,2,3 ; \quad-\infty<\zeta<+\infty \tag{2.29}
\end{equation*}
$$

are well defined and satisfy the conditions

$$
\begin{equation*}
0<w_{k}(\zeta)<1 \quad k=1,2,3, \quad-\infty<\zeta<+\infty \tag{2.30}
\end{equation*}
$$

and

$$
\begin{equation*}
w_{1}(\zeta)+w_{2}(\zeta)+w_{3}(\zeta)=1 \quad-\infty<\zeta<+\infty \tag{2.31}
\end{equation*}
$$

In this subsection hereafter we consider only the special case where

$$
\begin{equation*}
w_{k}=w_{k}(\zeta) \quad k=1,2,3 ; \quad-\infty<\zeta<+\infty \tag{2.32}
\end{equation*}
$$

For this special case, Eqs. (2.30)-(2.32) imply that both Eqs. (2.11) and (2.24) are satisfied for any real number $\zeta$. Thus, for the case Eq. (2.32), (i) point $P_{0}$ defined by Eq. (2.10) must always lie in the interior of $\triangle P_{1} P_{2} P_{3}$, and (ii) Eqs. (2.19)-(2.22) reduce to

$$
\begin{align*}
& A\left(\triangle P_{0} P_{2} P_{3}\right)=w_{1} \times A\left(\triangle P_{1} P_{2} P_{3}\right)  \tag{2.33}\\
& A\left(\triangle P_{0} P_{3} P_{1}\right)=w_{2} \times A\left(\triangle P_{1} P_{2} P_{3}\right)  \tag{2.34}\\
& A\left(\triangle P_{0} P_{1} P_{2}\right)=w_{3} \times A\left(\triangle P_{1} P_{2} P_{3}\right) \tag{2.35}
\end{align*}
$$

and Eq. (2.25), respectively. Moreover, for the special case $\zeta=0$, we have

$$
\begin{equation*}
w_{1}=w_{2}=w_{3}=1 / 3 \quad \text { and } \quad \vec{P}_{0}=\left(\vec{P}_{1}+\vec{P}_{2}+\vec{P}_{3}\right) / 3 \quad(\zeta=0) \tag{2.36}
\end{equation*}
$$

Because the position vector of the centroid of a triangle is the simple average of those of its three vertices, Eq. $(2.36) \Rightarrow$ point $P_{0}$ is the centroid of $\triangle P_{1} P_{2} P_{3}$ if $\zeta=0$.

Note that, for a polygon other than a triangle (which is the simplest polygon), generally it is not true that the position vector of its centroid is the simple average of those of its vertices. As an example, the position vector of the centroid of a quadrilateral may not be the simple average of those of its four vertices unless the quadrilateral is a parallelogram.

Next, with the aid of Eqs. (2.6) and (2.7), and Fig. 5, one can see that

$$
\begin{equation*}
A\left(\triangle P_{0} P_{2} P_{3}\right)=\frac{h_{1} \ell_{1}}{2}, \quad A\left(\triangle P_{0} P_{3} P_{1}\right)=\frac{h_{2} \ell_{2}}{2} \quad \text { and } \quad A\left(\triangle P_{0} P_{1} P_{2}\right)=\frac{h_{3} \ell_{3}}{2} \tag{2.37}
\end{equation*}
$$

In turn, by substituting Eq. (2.37) into Eqs. (2.33)-(2.35) and using Eqs. (2.29) and (2.32), one has

$$
\begin{equation*}
\frac{h_{k} \ell_{k}}{2}=w_{k} A\left(\triangle P_{1} P_{2} P_{3}\right)=\frac{\left(\ell_{k}\right)^{\zeta} \times A\left(\triangle P_{1} P_{2} P_{3}\right)}{\left(\ell_{1}\right)^{\zeta}+\left(\ell_{2}\right)^{\zeta}+\left(\ell_{3}\right)^{\zeta}} \quad k=1,2,3 ; \quad-\infty<\zeta<+\infty \tag{2.38}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\frac{h_{1}}{\left(\ell_{1}\right)^{\zeta-1}}=\frac{h_{2}}{\left(\ell_{2}\right)^{\zeta-1}}=\frac{h_{3}}{\left(\ell_{3}\right)^{\zeta-1}}=\frac{2 \times A\left(\triangle P_{1} P_{2} P_{3}\right)}{\left(\ell_{1}\right)^{\zeta}+\left(\ell_{2}\right)^{\zeta}+\left(\ell_{3}\right)^{\zeta}} \quad-\infty<\zeta<+\infty \tag{2.39}
\end{equation*}
$$

Because (i) $h_{1}, h_{2}$ and $h_{3}$ are the distances from point $P_{0}$ to the sides $\overline{P_{2} P_{3}}, \overline{P_{3} P_{1}}$, and $\overline{P_{1} P_{2}}$, respectively, and (ii) $\ell_{1}, \ell_{2}$, and $\ell_{3}$ are the lengths of $\overline{P_{2} P_{3}}, \overline{P_{3} P_{1}}$, and $\overline{P_{1} P_{2}}$, respectively, Eq. (2.39) states that, for a given $\triangle P_{1} P_{2} P_{3}$ and a given $\zeta$, the distance from point $P_{0}$ to any side of $\triangle P_{1} P_{2} P_{3}$ is proportional to the $(\zeta-1)$ th power of the length of this side.

In fact, for the case $\zeta=0$ (i.e., when $P_{0}$ is the centroid of $\triangle P_{1} P_{2} P_{3}$ ), Eq. (2.39) reduces to

$$
\begin{equation*}
h_{k}=\left[\frac{2 \times A\left(\triangle P_{1} P_{2} P_{3}\right)}{3}\right] \frac{1}{\ell_{k}} \quad k=1,2,3 \quad \text { if } \quad \zeta=0 \tag{2.40}
\end{equation*}
$$

i.e., the distance between the centroid $P_{0}$ and a side of $\triangle P_{1} P_{2} P_{3}$ is inversely proportional to the side length. Thus the centroid is farthest away from the shortest side and closest to the longest side - a result consistent with what we observe from the triangle $\triangle F H B$ depicted in Fig. 4.

On the other hand, for the case $\zeta=1$, Eq. (2.39) reduces to

$$
\begin{equation*}
h_{1}=h_{2}=h_{3}=\frac{2 \times A\left(\triangle P_{1} P_{2} P_{3}\right)}{\ell_{1}+\ell_{2}+\ell_{3}} \quad \text { if } \quad \zeta=1 \tag{2.41}
\end{equation*}
$$

i.e., point $P_{0}$ has the same distance to the three sides of $\triangle P_{1} P_{2} P_{3}$ if $\zeta=1$. By definition, this implies that $P_{0}$ is the in-center of $\triangle P_{1} P_{2} P_{3}$ for the case $\zeta=1$.

In general one can infer from Eq. (2.39) that:
(a) For the case $\zeta>1$, point $P_{0}$ is farthest away from the longest side of $\triangle P_{1} P_{2} P_{3}$ and closest to its shortest side. The trend becomes stronger as the value of $\zeta$ increases from 1.
(b) For the case $\zeta<1$, point $P_{0}$ is farthest away from the shortest side of $\triangle P_{1} P_{2} P_{3}$ and closest to its longest side. The trend becomes stronger as the value of $\zeta$ decreases from 1.
Let the interior point $P_{0}$ of any triangle $\triangle P_{1} P_{2} P_{3}$ defined by Eqs. (2.10) and (2.32) be referred to as the interior point of this triangle associated with the parameter $\zeta$. Then, from the above observations and the discussions given in Sec. 1, it becomes clear that a spatial solution point such as the point $G^{\prime}$ depicted in Fig. 3 will be much less likely to lie outside of the associated hexagon $A B C D E F$ and causes numerical difficulties if each of the centroidal mesh points $A, C, G$ and $E$ depicted in Figs. 3 and 4 is replaced by a an interior point associated with a parameter $\zeta \geq 1$. In fact, it has been shown through numerical experiments that the numerical difficulties associated with spatial solution points being embedded beyond a solid-wall boundary can be overcome by replacing centroids of triangles with in-centers (i.e., interior points associated with $\zeta=1$ ) in the construction of a spatial triangular mesh.

### 2.2. Tetrahedral meshes

Consider a 3D spatial domain formed by tetrahedra with the understanding that any two neighboring tetrahedra share a common triangular face. As an example, consider the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ depicted in Fig. 6(a) where $P_{1}, P_{2}, P_{3}$ and $P_{4}$ are the vertices of the tetrahedron. The tetrahedron has four triangular faces, i.e., $\triangle P_{1} P_{2} P_{3}, \triangle P_{2} P_{3} P_{4}, \triangle P_{3} P_{4} P_{1}$, and $\triangle P_{4} P_{1} P_{2}$. As shown in Fig. 6(b) , the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ shares these faces with four neighboring tetrahedra, i.e., $P_{1} P_{2} P_{3} D, P_{2} P_{3} P_{4} A, P_{3} P_{4} P_{1} B$, and $P_{4} P_{1} P_{2} C$, respectively. Note that, to simplify notation, the same symbols used in Sec. 2.1 may also be used in this subsection albeit that those used here generally denote objects in the 3D $x-y-z$ space.

Because a tetrahedron is intrinsically convex in a 3D $x-y-z$ space, its centroid must lie in its interior if its volume $>0$. Hereafter, only tetrahedra of volume $>0$ will be considered. In particular, we assume that

$$
\begin{equation*}
V\left(P_{1} P_{2} P_{3} P_{4}\right) \stackrel{\text { def }}{=} \text { the volume of the tetrahedron } P_{1} P_{2} P_{3} P_{4}>0 \tag{2.42}
\end{equation*}
$$

Let the points $P_{0}, A_{0}, B_{0}, C_{0}$, and $D_{0}$ depicted in Fig. 6(b) be the centroids of the tetrahedra, $P_{1} P_{2} P_{3} P_{4}$, $P_{2} P_{3} P_{4} A, P_{3} P_{4} P_{1} B, P_{4} P_{1} P_{2} C$, and $P_{1} P_{2} P_{3} D$, respectively. Then each of these five centroids lies in the interior of the associated tetrahedron. As such the four polyhedra $P_{0} P_{2} P_{3} P_{4} A_{0}, P_{0} P_{3} P_{4} P_{1} B_{0}, P_{0} P_{4} P_{1} P_{2} C_{0}$, and $P_{0} P_{1} P_{2} P_{3} D_{0}$ (each of them has five vertices, six triangular faces, and nine edges) do not overlap. In the original CESE 3D tetrahedral mesh construction, (i) each of these 5 -vertex polyhedra is the spatial projection of a BCE, (ii) the union of the above four 5 -vertex polyhedra is a polyhedron with 8 vertices, 12 triangular faces, and 18 edges, and also is the spatial projection of a CCE, and (iii) the centroid of the 8 -vertex polyhedron referred to above is a spatial solution point, i.e., the spatial projection of a space-time solution point.

As in the 2D case, for some pathological case where an 8 -vertex polyhedron become highly concave, its centroid (i.e., the associated spatial solution point) may lie outside of the polyhedron and causes numerical difficulties. In the following, we explain how these difficulties can be overcome by a new construction of 3D tetrahedral meshes.

Consider Fig. 6(a) again. Here (i) $P_{0}$ is a point in the $x-y-z$ space to be defined shortly and not necessary the centroid of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$; (ii) for each $k=0,1,2,3,4, x_{k}, y_{k}$, and $z_{k}$ are the $x$-, $y$ - and $z$ coordinates of point $P_{k}$, respectively; (iii) point $Q_{1}$ (not shown in Fig. 6(a)) is the point that lies on the plane containing the face $\triangle P_{2} P_{3} P_{4}$, and also meets the condition

$$
\begin{equation*}
\overline{P_{0} Q_{1}} \perp \triangle P_{2} P_{3} P_{4} \tag{2.43}
\end{equation*}
$$

i.e., $Q_{1}$ is the perpendicular projection of $P_{0}$ on the plane containing $\triangle P_{2} P_{3} P_{4}$; (iv) point $Q_{2}$ (not shown) is the point that lies on the plane containing the face $\triangle P_{3} P_{4} P_{1}$, and also meets the condition

$$
\begin{equation*}
\overline{P_{0} Q_{2}} \perp \triangle P_{3} P_{4} P_{1} \tag{2.44}
\end{equation*}
$$

i.e., $Q_{2}$ is the perpendicular projection of $P_{0}$ on the plane containing $\triangle P_{3} P_{4} P_{1} ;(\mathrm{v})$ point $Q_{3}$ (not shown) is the point that lies on the plane containing the face $\triangle P_{4} P_{1} P_{2}$, and also meets the condition

$$
\begin{equation*}
\overline{P_{0} Q_{3}} \perp \triangle P_{4} P_{1} P_{2} \tag{2.45}
\end{equation*}
$$

i.e., $Q_{3}$ is the perpendicular projection of $P_{0}$ on the plane containing $\triangle P_{4} P_{1} P_{2}$; (vi) point $Q_{4}$ (not shown) is the point that lies on the plane containing the face $\triangle P_{1} P_{2} P_{3}$, and also meets the condition

$$
\begin{equation*}
\overline{P_{0} Q_{4}} \perp \triangle P_{1} P_{2} P_{3} \tag{2.46}
\end{equation*}
$$

i.e., $Q_{4}$ is the perpendicular projection of $P_{0}$ on the plane containing $\triangle P_{1} P_{2} P_{3}$; (vii)

$$
\begin{equation*}
A_{1} \stackrel{\text { def }}{=} A\left(\triangle P_{2} P_{3} P_{4}\right), \quad A_{2} \stackrel{\text { def }}{=} A\left(\triangle P_{3} P_{4} P_{1}\right), \quad A_{3} \stackrel{\text { def }}{=} A\left(\triangle P_{4} P_{1} P_{2}\right), \quad \text { and } \quad A_{4} \stackrel{\text { def }}{=} A\left(\triangle P_{1} P_{2} P_{3}\right) \tag{2.47}
\end{equation*}
$$

where $A\left(\triangle P_{2} P_{3} P_{4}\right), A\left(\triangle P_{3} P_{4} P_{1}\right), A\left(\triangle P_{4} P_{1} P_{2}\right)$, and $A\left(\triangle P_{1} P_{2} P_{3}\right)$ are the areas of the four triangular faces $\triangle P_{2} P_{3} P_{4}, \triangle P_{3} P_{4} P_{1}, \triangle P_{4} P_{1} P_{2}$, and $\triangle P_{1} P_{2} P_{3}$, respectively; (viii)

$$
\begin{equation*}
h_{k} \stackrel{\text { def }}{=}\left|\overline{P_{0} Q_{k}}\right| \quad k=1,2,3,4 \tag{2.48}
\end{equation*}
$$

and (ix) the $x-y-z$ coordinate system depicted in Fig. 6(a) is also a right-handed system which satisfies Eq. (2.8).

Note that, according to Eq. (2.47) and Fig. 6(a), for each $k=1,2,3,4, A_{k}$ denotes the area of the face of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ which is opposite to the vertex $P_{k}$. Also from the definition of points $Q_{1}, Q_{2}, Q_{3}$, and $Q_{4}$ and that of $h_{k}, k=1,2,3,4$, one concludes that $h_{1}, h_{2}, h_{3}$, and $h_{4}$, respectively, are the (shortest) distances from $P_{0}$ to the planes that contain the faces $\triangle P_{2} P_{3} P_{4}, \triangle P_{3} P_{4} P_{1}, \triangle P_{4} P_{1} P_{2}$, and $\triangle P_{1} P_{2} P_{3}$, respectively.

For each $k=0,1,2,3,4$, let $\vec{P}_{k}$ be the position vector of point $P_{k}$, i.e.,

$$
\begin{equation*}
\vec{P}_{k} \stackrel{\text { def }}{=} x_{k} \vec{e}_{x}+y_{k} \vec{e}_{y}+z_{k} \vec{e}_{z} \quad k=0,1,2,3,4 \tag{2.49}
\end{equation*}
$$

Also let $\vec{P}_{0}$ be an weighted average of $\vec{P}_{k}, k=1,2,3,4$, i.e.,

$$
\begin{equation*}
\vec{P}_{0}=w_{1} \vec{P}_{1}+w_{2} \vec{P}_{2}+w_{3} \vec{P}_{3}+w_{4} \vec{P}_{4} \tag{2.50}
\end{equation*}
$$

where $w_{1}, w_{2}, w_{3}$, and $w_{4}$ are real weight factors with

$$
\begin{equation*}
w_{1}+w_{2}+w_{3}+w_{4}=1 \tag{2.51}
\end{equation*}
$$

Moreover, given a pair of integers $k_{1}$ and $k_{2}$ with $k_{1} \neq k_{2}$ and $k_{1}, k_{2}=0,1,2,3,4$, let $\overrightarrow{P_{k_{1}} P_{k_{2}}}$ be the displacement vector joining points $P_{k_{1}}$ and $P_{k_{2}}$ and pointing in the direction from $P_{k_{1}}$ to $P_{k_{2}}$, i.e.,

$$
\begin{equation*}
\overline{P_{k_{1}} P_{k_{2}}} \stackrel{\text { def }}{=} \vec{P}_{k_{2}}-\vec{P}_{k_{1}}, \quad k_{1} \neq k_{2} \quad \text { and } \quad k_{1}, k_{2}=0,1,2,3,4 \tag{2.52}
\end{equation*}
$$

Then, by using Eq. (2.52) and other well-known properties of the vector product such as Eq. (2.13), one can show that Eq. (2.14) is still valid in the current 3D case. In turn, with the aid of Eqs. (2.14) and (2.52), and other well-known properties of the scalar triple product such as

$$
\begin{equation*}
(\vec{A} \times \vec{B}) \cdot \vec{C}=(\vec{B} \times \vec{C}) \cdot \vec{A}=(\vec{C} \times \vec{A}) \cdot \vec{B}=-(\vec{A} \times \vec{C}) \cdot \vec{B} \tag{2.53}
\end{equation*}
$$

and

$$
\begin{equation*}
(\vec{A} \times \vec{A}) \cdot \vec{B}=(\vec{A} \times \vec{B}) \cdot \vec{A}=(\vec{B} \times \vec{A}) \cdot \vec{A}=0 \tag{2.54}
\end{equation*}
$$

for any vectors $\vec{A}, \vec{B}$, and $\vec{C}$, one can show that

$$
\begin{equation*}
\Delta \stackrel{\text { def }}{=}\left(\overrightarrow{P_{1} P_{2}} \times \overrightarrow{P_{1} P_{3}}\right) \cdot \overrightarrow{P_{1} P_{4}}=\left(\vec{P}_{1} \times \vec{P}_{2}\right) \cdot \vec{P}_{4}+\left(\vec{P}_{2} \times \vec{P}_{3}\right) \cdot \vec{P}_{4}+\left(\vec{P}_{3} \times \vec{P}_{1}\right) \cdot \vec{P}_{4}+\left(\vec{P}_{3} \times \vec{P}_{2}\right) \cdot \vec{P}_{1} \tag{2.55}
\end{equation*}
$$

Moreover, by using Eqs. $(2.13)$, (2.50) and (2.52), one has

$$
\begin{equation*}
\overrightarrow{P_{0} P_{1}} \times \overrightarrow{P_{0} P_{2}}=\left(1-w_{1}-w_{2}\right) \vec{P}_{1} \times \vec{P}_{2}+\left(w_{3} \vec{P}_{3}+w_{4} \vec{P}_{4}\right) \times\left(\vec{P}_{1}-\vec{P}_{2}\right) \tag{2.56}
\end{equation*}
$$

Next, by using Eqs. (2.50)-(2.56), one can show that

$$
\begin{align*}
& \left(\overrightarrow{P_{0} P_{1}} \times \overrightarrow{P_{0} P_{2}}\right) \cdot \overrightarrow{P_{0} P_{3}} \\
& =w_{4}\left[\left(\vec{P}_{1} \times \vec{P}_{2}\right) \cdot \vec{P}_{3}+\left(\vec{P}_{2} \times \vec{P}_{1}\right) \cdot \vec{P}_{4}+\left(\vec{P}_{3} \times \vec{P}_{4}\right) \cdot \vec{P}_{1}+\left(\vec{P}_{4} \times \vec{P}_{3}\right) \cdot \vec{P}_{2}\right]=-w_{4} \Delta \tag{2.57}
\end{align*}
$$

In turn, by carrying out the index mapping

$$
(1,2,3,4) \rightarrow(2,3,4,1)
$$

of Eq. (2.57) repeatedly and using Eqs. (2.53)-(2.55), one obtains

$$
\begin{align*}
& \left(\overrightarrow{P_{0} P_{2}} \times \overrightarrow{P_{0} P_{3}}\right) \cdot \overrightarrow{P_{0} P_{4}} \\
& =w_{1}\left[\left(\vec{P}_{2} \times \vec{P}_{3}\right) \cdot \vec{P}_{4}+\left(\vec{P}_{3} \times \vec{P}_{2}\right) \cdot \vec{P}_{1}+\left(\vec{P}_{4} \times \vec{P}_{1}\right) \cdot \vec{P}_{2}+\left(\vec{P}_{1} \times \vec{P}_{4}\right) \cdot \vec{P}_{3}\right]=w_{1} \Delta  \tag{2.58}\\
& \left(\overrightarrow{P_{0} P_{3}} \times \overrightarrow{P_{0} P_{4}}\right) \cdot \overrightarrow{P_{0} P_{1}} \\
& =w_{2}\left[\left(\vec{P}_{3} \times \vec{P}_{4}\right) \cdot \vec{P}_{1}+\left(\vec{P}_{4} \times \vec{P}_{3}\right) \cdot \vec{P}_{2}+\left(\vec{P}_{1} \times \vec{P}_{2}\right) \cdot \vec{P}_{3}+\left(\vec{P}_{2} \times \vec{P}_{1}\right) \cdot \vec{P}_{4}\right]=-w_{2} \Delta \tag{2.59}
\end{align*}
$$

and

$$
\begin{align*}
& \left(\overrightarrow{P_{0} P_{4}} \times \overrightarrow{P_{0} P_{1}}\right) \cdot \overrightarrow{P_{0} P_{2}} \\
& =w_{3}\left[\left(\vec{P}_{4} \times \vec{P}_{1}\right) \cdot \vec{P}_{2}+\left(\vec{P}_{1} \times \vec{P}_{4}\right) \cdot \vec{P}_{3}+\left(\vec{P}_{2} \times \vec{P}_{3}\right) \cdot \vec{P}_{4}+\left(\vec{P}_{3} \times \vec{P}_{2}\right) \cdot \vec{P}_{1}\right]=w_{3} \Delta \tag{2.60}
\end{align*}
$$

For any four different integers $k_{1}, k_{2}, k_{3}, k_{4}=0,1,2,3,4$, let

$$
\begin{equation*}
V\left(P_{k_{1}} P_{k_{2}} P_{k_{3}} P_{k_{4}}\right) \stackrel{\text { def }}{=} \text { the volume of the tetrahedron with the vertices } P_{k_{1}}, P_{k_{2}}, P_{k_{3}} \text {, and } P_{k_{4}} \tag{2.61}
\end{equation*}
$$

Then by using the definition of the scalar triple product along with Eqs. (2.55) and (2.57)-(2.60), one has

$$
\begin{gather*}
V\left(P_{1} P_{2} P_{3} P_{4}\right)=\left|\left(\overrightarrow{P_{1} P_{2}} \times \overrightarrow{P_{1} P_{3}}\right) \cdot \overrightarrow{P_{1} P_{4}}\right| / 6=|\Delta| / 6  \tag{2.62}\\
V\left(P_{0} P_{2} P_{3} P_{4}\right)=\left|\left(\overrightarrow{P_{0} P_{2}} \times \overrightarrow{P_{0} P_{3}}\right) \cdot \overrightarrow{P_{0} P_{4}}\right| / 6=\left|w_{1} \Delta\right| / 6=\left|w_{1}\right| \times V\left(P_{1} P_{2} P_{3} P_{4}\right)  \tag{2.63}\\
V\left(P_{0} P_{3} P_{4} P_{1}\right)=\left|\left(\overrightarrow{\overrightarrow{P_{0} P_{3}}} \times \overrightarrow{P_{0} P_{4}}\right) \cdot \overrightarrow{P_{0} P_{1}}\right| / 6=\left|w_{2} \Delta\right| / 6=\left|w_{2}\right| \times V\left(P_{1} P_{2} P_{3} P_{4}\right)  \tag{2.64}\\
V\left(P_{0} P_{4} P_{1} P_{2}\right)=\left|\left(\overrightarrow{\overrightarrow{P_{0} P_{4}}} \times \overrightarrow{P_{0} P_{1}}\right) \cdot \overrightarrow{P_{0} P_{2}}\right| / 6=\left|w_{3} \Delta\right| / 6=\left|w_{3}\right| \times V\left(P_{1} P_{2} P_{3} P_{4}\right) \tag{2.65}
\end{gather*}
$$

and

$$
\begin{equation*}
V\left(P_{0} P_{1} P_{2} P_{3}\right)=\left|\left(\overrightarrow{P_{0} P_{1}} \times \overrightarrow{P_{0} P_{2}}\right) \cdot \overrightarrow{P_{0} P_{3}}\right| / 6=\left|w_{4} \Delta\right| / 6=\left|w_{4}\right| \times V\left(P_{1} P_{2} P_{3} P_{4}\right) \tag{2.66}
\end{equation*}
$$

Eqs. $(2.63)-(2.66) \Rightarrow$

$$
\begin{align*}
& V\left(P_{0} P_{2} P_{3} P_{4}\right)+V\left(P_{0} P_{3} P_{4} P_{1}\right)+V\left(P_{0} P_{4} P_{1} P_{2}\right)+V\left(P_{0} P_{1} P_{2} P_{3}\right) \\
& =\left(\left|w_{1}\right|+\left|w_{2}\right|+\left|w_{3}\right|+\left|w_{4}\right|\right) \times V\left(P_{1} P_{2} P_{3} P_{4}\right) \tag{2.67}
\end{align*}
$$

Because the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ is intrinsically convex in the $x-y-z$ space, point $P_{0}$ defined by Eqs. (2.50) and (2.51) lies on the boundary or in the interior of the tetrahedron $P_{1} P_{2} P_{3} P_{4} \Leftrightarrow$

$$
\begin{equation*}
0 \leq w_{k} \leq 1 \quad k=1,2,3,4 \tag{2.68}
\end{equation*}
$$

Moreover, because Eq. (2.42) implies that the interior of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ is not empty, $P_{0}$ lies in the interior of the tetrahedron $P_{1} P_{2} P_{3} P_{4} \Leftrightarrow$

$$
\begin{equation*}
0<w_{k}<1 \quad k=1,2,3,4 \tag{2.69}
\end{equation*}
$$

For the case where $P_{0}$ lies on the boundary or in the interior of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$, Eqs. (2.51), (2.67), and (2.68) now $\Rightarrow$ the following obvious result:

$$
\begin{equation*}
V\left(P_{0} P_{2} P_{3} P_{4}\right)+V\left(P_{0} P_{3} P_{4} P_{1}\right)+V\left(P_{0} P_{4} P_{1} P_{2}\right)+V\left(P_{0} P_{1} P_{2} P_{3}\right)=V\left(P_{1} P_{2} P_{3} P_{4}\right) \tag{2.70}
\end{equation*}
$$

On the other hand, for the case where $P_{0}$ lies outside of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$, Eq. (2.68) is no longer true, i.e., some of the weight factors must be negative. Because Eq. (2.51) requires that some of the weight factors must be positive, we arrive at the conclusion that

$$
\begin{equation*}
\left|w_{1}\right|+\left|w_{2}\right|+\left|w_{3}\right|+\left|w_{4}\right|>\left|w_{1}+w_{2}+w_{3}+w_{4}\right|=1 \tag{2.71}
\end{equation*}
$$

if $P_{0}$ lies outside of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$. In turn, Eqs. (2.67) and (2.71) $\Rightarrow$

$$
\begin{equation*}
V\left(P_{0} P_{2} P_{3} P_{4}\right)+V\left(P_{0} P_{3} P_{4} P_{1}\right)+V\left(P_{0} P_{4} P_{1} P_{2}\right)+V\left(P_{0} P_{1} P_{2} P_{3}\right)>V\left(P_{1} P_{2} P_{3} P_{4}\right) \tag{2.72}
\end{equation*}
$$

if $P_{0}$ lies outside of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$.
Note that, for any real parameter $\zeta$, Eqs. (2.42) and (2.47) $\Rightarrow$

$$
\begin{equation*}
A_{k}>0 \quad \text { and } \quad\left(A_{k}\right)^{\zeta}>0 \quad k=1,2,3,4 ; \quad-\infty<\zeta<+\infty \tag{2.73}
\end{equation*}
$$

As such the weight factors

$$
\begin{equation*}
w_{k}(\zeta) \stackrel{\text { def }}{=} \frac{\left(A_{k}\right)^{\zeta}}{\left(A_{1}\right)^{\zeta}+\left(A_{2}\right)^{\zeta}+\left(A_{3}\right)^{\zeta}+\left(A_{4}\right)^{\zeta}} \quad k=1,2,3,4 ; \quad-\infty<\zeta<+\infty \tag{2.74}
\end{equation*}
$$

are well defined and satisfy the conditions

$$
\begin{equation*}
0<w_{k}(\zeta)<1 \quad k=1,2,3,4 \quad-\infty<\zeta<+\infty \tag{2.75}
\end{equation*}
$$

and

$$
\begin{equation*}
w_{1}(\zeta)+w_{2}(\zeta)+w_{3}(\zeta)+w_{4}(\zeta)=1 \quad-\infty<\zeta<+\infty \tag{2.76}
\end{equation*}
$$

Note that, even though the same notation is used, the weight factors $w_{k}(\zeta), k=1,2,3,4$, defined in Eq. (2.74) are different from those defined in Eq. (2.29). The definition Eq. (2.74) will be used in this subsection.

Moreover, hereafter we consider only the special case where

$$
\begin{equation*}
w_{k}=w_{k}(\zeta) \quad k=1,2,3,4 ; \quad-\infty<\zeta<+\infty \tag{2.77}
\end{equation*}
$$

For this special case, Eqs. (2.75)-(2.77) imply that both Eqs. (2.51) and (2.69) are satisfied for any real number $\zeta$. Thus, for the case Eq. (2.77), (i) point $P_{0}$ defined by Eq. (2.50) must always lie in the interior of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$, and (ii) Eqs. (2.63)-(2.67) reduce to

$$
\begin{align*}
& V\left(P_{0} P_{2} P_{3} P_{4}\right)=w_{1} \times V\left(P_{1} P_{2} P_{3} P_{4}\right)  \tag{2.78}\\
& V\left(P_{0} P_{3} P_{4} P_{1}\right)=w_{2} \times V\left(P_{1} P_{2} P_{3} P_{4}\right)  \tag{2.79}\\
& V\left(P_{0} P_{4} P_{1} P_{2}\right)=w_{3} \times V\left(P_{1} P_{2} P_{3} P_{4}\right)  \tag{2.80}\\
& V\left(P_{0} P_{1} P_{2} P_{3}\right)=w_{4} \times V\left(P_{1} P_{2} P_{3} P_{4}\right) \tag{2.81}
\end{align*}
$$

and Eq. (2.70), respectively. Moreover, for the special case $\zeta=0$, we have

$$
\begin{equation*}
w_{1}=w_{2}=w_{3}=w_{4}=1 / 4 \quad \text { and } \quad \vec{P}_{0}=\left(\vec{P}_{1}+\vec{P}_{2}+\vec{P}_{3}+\vec{P}_{4}\right) / 4 \quad(\zeta=0) \tag{2.82}
\end{equation*}
$$

Because the position vector of the centroid of a tetrahedron is the simple average of those of its four vertices, Eq. $(2.82) \Rightarrow$ point $P_{0}$ is the centroid of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ if $\zeta=0$.

Note that, for a polyhedron other than a tetrahedron (which is the simplest polyhedron), generally it is not true that the position vector of its centroid is the simple average of those of its vertices. As an example,
the position vector of the centroid of a hexahedron may not be the simple average of those of its eight vertices unless the hexahedron is a parallelepiped.

Next, with the aid of Eq. (2.47) and (2.48), and Fig. 6(a), one can see that

$$
\begin{equation*}
V\left(P_{0} P_{2} P_{3} P_{4}\right)=\frac{h_{1} A_{1}}{3}, V\left(P_{0} P_{3} P_{4} P_{1}\right)=\frac{h_{2} A_{2}}{3}, V\left(P_{0} P_{4} P_{1} P_{2}\right)=\frac{h_{3} A_{3}}{3} \text { and } V\left(P_{0} P_{1} P_{2} P_{3}\right)=\frac{h_{4} A_{4}}{3} \tag{2.83}
\end{equation*}
$$

In turn, by substituting Eq. (2.83) into Eqs. (2.78)-(2.81) and using Eqs. (2.74) and (2.77), one has

$$
\begin{equation*}
\frac{h_{k} A_{k}}{3}=w_{k} V\left(P_{1} P_{2} P_{3} P_{4}\right)=\frac{\left(A_{k}\right)^{\zeta} \times V\left(P_{1} P_{2} P_{3} P_{4}\right)}{\left(A_{1}\right)^{\zeta}+\left(A_{2}\right)^{\zeta}+\left(A_{3}\right)^{\zeta}+\left(A_{4}\right)^{\zeta}} \quad k=1,2,3,4 ; \quad-\infty<\zeta<+\infty \tag{2.84}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\frac{h_{1}}{\left(A_{1}\right)^{\zeta-1}}=\frac{h_{2}}{\left(A_{2}\right)^{\zeta-1}}=\frac{h_{3}}{\left(A_{3}\right)^{\zeta-1}}=\frac{h_{4}}{\left(A_{4}\right)^{\zeta-1}}=\frac{3 \times V\left(P_{1} P_{2} P_{3} P_{4}\right)}{\left(A_{1}\right)^{\zeta}+\left(A_{2}\right)^{\zeta}+\left(A_{3}\right)^{\zeta}+\left(A_{4}\right)^{\zeta}} \quad-\infty<\zeta<+\infty \tag{2.85}
\end{equation*}
$$

Because (i) $h_{1}, h_{2}, h_{3}$, and $h_{4}$ are the distances from point $P_{0}$ to the triangular faces $\triangle P_{2} P_{3} P_{4}, \triangle P_{3} P_{4} P_{1}$, $\triangle P_{4} P_{1} P_{2}$, and $\triangle P_{1} P_{2} P_{3}$ respectively, and (ii) $A_{1}, A_{2}, A_{3}$, and $A_{4}$ are the areas of $\triangle P_{2} P_{3} P_{4}, \triangle P_{3} P_{4} P_{1}$, $\triangle P_{4} P_{1} P_{2}$, and $\triangle P_{1} P_{2} P_{3}$, respectively, Eq. (2.85) states that, for a given tetrahedron $P_{1} P_{2} P_{3} P_{4}$ and a given $\zeta$, the distance from point $P_{0}$ to any face of this tetrahedron is proportional to the $(\zeta-1)$ th power of the area of the face.

In fact, for the case $\zeta=0$ (i.e., when $P_{0}$ is the centroid of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ ), Eq. (2.85) reduces to

$$
\begin{equation*}
h_{k}=\left[\frac{3 \times V\left(P_{1} P_{2} P_{3} P_{4}\right)}{4}\right] \frac{1}{A_{k}} \quad k=1,2,3,4, \quad \text { if } \quad \zeta=0 \tag{2.86}
\end{equation*}
$$

i.e., the distance between the centroid $P_{0}$ and a face of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ is inversely proportional to the area of this face. Thus the centroid is farthest away from the smallest face and closest to the largest face.

On the other hand, for the case $\zeta=1$, Eq. (2.85) reduces to

$$
\begin{equation*}
h_{1}=h_{2}=h_{3}=h_{4}=\frac{3 \times V\left(P_{1} P_{2} P_{3} P_{4}\right)}{A_{1}+A_{2}+A_{3}+A_{4}} \quad \quad \text { if } \quad \zeta=1 \tag{2.87}
\end{equation*}
$$

i.e., point $P_{0}$ has the same distance to the four faces of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ if $\zeta=1$. By definition, this implies that $P_{0}$ is the in-center of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ for the case $\zeta=1$.

In general one can infer from Eq. (2.85) that:
(a) For the case $\zeta>1$, point $P_{0}$ is farthest away from the largest face of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ and closest to its smallest face. The trend becomes stronger as the value of $\zeta$ increases from 1 .
(b) For the case $\zeta<1$, point $P_{0}$ is farthest away from the smallest face of the tetrahedron $P_{1} P_{2} P_{3} P_{4}$ and closest to its largest face. The trend becomes stronger as the value of $\zeta$ decreases from 1.
Let the interior point $P_{0}$ of any tetrahedron $P_{1} P_{2} P_{3} P_{4}$ defined by Eqs. (2.50) and (2.77) be referred to as the interior point of this tetrahedron associated with the parameter $\zeta$. Then, from the above observations and the discussions given in Sec. 1, it becomes clear that a spatial solution point which is the centroid of the 8-vertex polyhedron referred to earlier will be much less likely to lie outside of the associated 8-vertex polyhedron and causes numerical difficulties if each of the centroidal mesh points $P_{0}, A_{0}, B_{0}, C_{0}$ and $D_{0}$ depicted in Fig. 6(b) is replaced by an interior point associated with a parameter $\zeta \geq 1$. In fact, it has been shown through numerical experiments that the numerical difficulties associated with spatial solution points being embedded beyond a solid-wall boundary can be overcome by replacing centroids of tetrahedra with in-centers (i.e., interior points associated with $\zeta=1$ ) in the construction of a spatial tetrahedral mesh.

## 3. Flux-conserving solution-coupling procedures

In this section, it will be described how an 1D CESE scheme can be modified so that the two decoupled solutions of the original scheme can become coupled in the new scheme while still preserving a flux conservation condition to be defined shortly. This description is then followed by a sketch of the 2D and 3D extensions.

### 3.1. 1D CESE solution-coupling schemes

As a preliminary, note that the CEs, SEs and the 1D schemes, which were defined over $\Omega_{1}$ in Sec. 1, hereafter will be defined over

$$
\begin{equation*}
\Omega \stackrel{\text { def }}{=}\{(j, n) \mid j, n=0, \pm 1 / 2, \pm 1, \pm 3 / 2, \ldots\}=\Omega_{1} \cup \Omega_{2} \tag{3.1}
\end{equation*}
$$

In particular, it should be understood that the symbol $\Omega_{1}$ from now on is replaced by $\Omega$ in each of Eqs. (1.3)(1.16). Also, as an example, the extension (defined over $\Omega$ ) of the $a-\epsilon$ scheme will be referred to as the dual $a-\epsilon$ scheme.

To clarify the significance of the extension, we offer the following observations:
(a) The solution of a dual scheme is formed by two decoupled solutions which are defined over $\Omega_{1}$ and $\Omega_{2}$, respectively.
(b) The same space-time region can be designated as a BCE of a mesh point $\in \Omega_{1}$ and also as a BCE of another mesh point $\in \Omega_{2}$. For examples, the space-time region $A B C D$ depicted in Fig. 2(c) can be designated as $\mathrm{CE}_{-}(j, n)$ and $\mathrm{CE}_{+}(j-1 / 2, n)$, respectively, while the region $A F E D$ depicted in Fig. 2(d) as $\mathrm{CE}_{+}(j, n)$ and $\mathrm{CE}_{-}(j+1 / 2, n)$, respectively. Hereafter, as an example, points $(j, n)$ and $(j-1 / 2, n)$ (i.e., points $A$ and $B$ ) will be referred to as the hosts of the region $A B C D$.
(c) Let the region $A B C D$ be designated as $\mathrm{CE}_{-}(j, n)$. Then, on its boundary, the line segments $\overline{A B}$ and $\overline{A D}$ belong to $\mathrm{SE}(j, n)$ while $\overline{C B}$ and $\overline{C D}$ belong to $\mathrm{SE}(j-1 / 2, n-1 / 2)$, i.e., over $\overline{A B}$ and $\overline{A D}$, the vector $\vec{h}^{*}$ should be evaluated using the solution data stored at point $A$ while, over $\overline{C B}$ and $\overline{C D}$, it should be evaluated using those stored at point $C$.
Alternately, let the region $A B C D$ be designated as $\mathrm{CE}_{+}(j-1 / 2, n)$. Then, on its boundary, the line segments $\overline{B A}$ and $\overline{B C}$ belong to $\mathrm{SE}(j-1 / 2, n)$ while $\overline{D A}$ and $\overline{D C}$ belong to $\mathrm{SE}(j, n-1 / 2)$, i.e., over $\overline{B A}$ and $\overline{B C}$, the vector $\vec{h}^{*}$ should be evaluated using the solution data stored at point $B$ while, over $\overline{D A}$ and $\overline{D C}$, it should be evaluated using those stored at point $D$. Note that, in the above and hereafter, as an example, a line segment joining points $A$ and $B$ is denoted by $\overline{A B}$ if it belongs to $S E(A)$ (i.e., the $S E$ centered at point $A$ ) while the same line segment is denoted by $\overline{B A}$ if it belongs $S E(B)$ (i.e., the $S E$ centered at point $B$ ).
(d) From the discussions given in the above item (c), it becomes clear that

$$
\begin{equation*}
\oint_{S\left(\mathrm{CE}_{+}(j, n)\right)} \vec{h}^{*} \cdot d \vec{s}=0 \quad \text { and } \quad \oint_{S\left(\mathrm{CE}_{-}(j+1 / 2, n)\right)} \vec{h}^{*} \cdot d \vec{s}=0 \quad(j, n) \in \Omega \tag{3.2}
\end{equation*}
$$

represent two different flux conservation conditions even though $\mathrm{CE}_{+}(j, n)$ and $\mathrm{CE}_{-}(j+1 / 2, n)$ occupy the same space-time region for any $(j, n) \in \Omega$.
As will be shown, each marching step of the 1D solution-coupling scheme is formed by two sub-steps. In the first sub-step, referred to as the dual-scheme marching step, the intermediate mesh values at the $n$th time level are evaluated from the given mesh values at the $(n-1 / 2)$ th time level using a dual-scheme (see Fig. 7). In the second sub-step, referred to as the solution-coupling step, the final mesh values at the $n$th time level are updated from the intermediate mesh values through a flux-conserving weighted-averaging procedure. Let the intermediate mesh values of $u$ and its normalized spatial derivative $u_{\bar{x}}$ at any $(j, n) \in \Omega$ be denoted by $\hat{u}_{j}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j}^{n}$, respectively. Then, as an example, $\hat{u}_{j}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j}^{n}$ can be evaluated in terms of the given mesh values at the $(n-1 / 2)$ th time level by using the dual $a-\epsilon$ scheme defined by Eqs. (1.12) and (1.15), i.e.,

$$
\begin{equation*}
\hat{u}_{j}^{n}=\frac{1}{2}\left\{(1-\nu)\left[u-(1+\nu) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}+(1+\nu)\left[u+(1-\nu) u_{\bar{x}}^{n}\right]_{j-1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega \tag{3.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(\hat{u}_{\bar{x}}\right)_{j}^{n}=\frac{1}{2}\left\{\left[(1-\epsilon) u+(2 \epsilon-1-\nu) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}-\left[(1-\epsilon) u+(1-\nu-2 \epsilon) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega \tag{3.4}
\end{equation*}
$$

Note that, by using a relation similar to the second equation in Eq. (1.9), one can define the "unnormalized" version of $\left(\hat{u}_{\bar{x}}\right)_{j}^{n}$, i.e.,

$$
\begin{equation*}
\left(\hat{u}_{x}\right)_{j}^{n} \stackrel{\text { def }}{=} \frac{4}{\Delta x}\left(\hat{u}_{\bar{x}}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.5}
\end{equation*}
$$

To avoid introducing unnecessary new notations, in the following description of the solution-coupling step, at each $(j, n) \in \Omega$, the final mesh values after updating from the intermediate values are simply denoted as $u_{j}^{n}$ and $\left(u_{\bar{x}}\right)_{j}^{n}$, respectively.

Note that the line segment $\overline{A F}$ depicted in Fig. 7 belongs to $\mathrm{SE}(j, n)$. Thus, by using (i) Eqs. (1.4), (1.6) and (1.9), and (ii) Fig. 2(d), one concludes that the flux of $\vec{h}^{*}$ leaving $\mathrm{CE}_{+}(j, n)$ through $\overline{A F}$ evaluated using the final solution data stored at point $A$ is given by

$$
\begin{equation*}
F(\overline{A F}) \stackrel{\text { def }}{=} \int_{x_{j}}^{x_{j}+(\Delta x / 2)} u^{*}\left(x, t^{n} ; j, n\right) d x=\frac{\Delta x}{2}\left(u_{+}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.6}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(u_{+}\right)_{j}^{n} \stackrel{\text { def }}{=}\left(u+u_{\bar{x}}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.7}
\end{equation*}
$$

Let $\hat{F}(\overline{A F})$ be the intermediate value of $F(\overline{A F})$, i.e., the value of $F(\overline{A F})$ when $u_{j}^{n}$ and $\left(u_{\bar{x}}\right)_{j}^{n}$ are replaced by $\hat{u}_{j}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j}^{n}$, respectively. Then Eqs. (3.5)-(3.7) $\Rightarrow$

$$
\begin{equation*}
\hat{F}(\overline{A F})=\frac{\Delta x}{2}\left(\hat{u}_{+}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.8}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(\hat{u}_{+}\right)_{j}^{n} \stackrel{\text { def }}{=}\left(\hat{u}+\hat{u}_{\bar{x}}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.9}
\end{equation*}
$$

At this juncture, note that: (i) at the midpoint $M_{+}$of $\overline{A F}, x=x_{j}+(\Delta x / 4)$ and $t=t^{n}$; and (ii) Eqs. (1.6), (1.9), and (3.7) $\Rightarrow$

$$
\begin{equation*}
u^{*}\left(x_{j}+(\Delta x / 4), t^{n} ; j, n\right)=\left(u_{+}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.10}
\end{equation*}
$$

Thus, $\left(u_{+}\right)_{j}^{n}$ is the value of $u$ at the midpoint $M_{+}$of $\overline{A F}$ evaluated using the values $u_{j}^{n}$ and $\left(u_{\bar{x}}\right)_{j}^{n}$ with the aid of Eqs. (1.6) and (1.9). Similarly $\left(\hat{u}_{+}\right)_{j}^{n}$ is the value of $u$ at the midpoint $M_{+}$evaluated using the intermediate values $\hat{u}_{j}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j}^{n}$ with the aid of Eqs. (1.6), (1.9), and (3.5).

Because $\left(u_{+}\right)_{j}^{n}$ and $\left(\hat{u}_{+}\right)_{j}^{n}$ are evaluated using the solution data stored at point $(j, n)$ (i.e., point $A$ ), hereafter, respectively, they will be referred to as the final and intermediate values of $u$ at the midpoint $M_{+}$associated with the host point $A$ of the region $A F E D$. Similarly, $F(\overline{A F})$ and $\hat{F}(\overline{A F})$, respectively, are referred to as the final and intermediate values of the flux leaving the top face of the same region associated with the same host.

Moreover, with the aid of Eqs. (3.6) and (3.8), and Fig. 7, one concludes that $F(\overline{A F})(\hat{F}(\overline{A F}))$ is equal to the length of $\overline{A F}$ multiplied by the approximated value $\left(u_{+}\right)_{j}^{n}\left(\left(\hat{u}_{+}\right)_{j}^{n}\right)$ of $u$ at the midpoint $M_{+}$. In other words, the value of $u$ at the midpoint $M_{+}$is the average value of $u$ over $\overline{A F}$-a result stemmed from the fact that $u^{*}(x, t ; j, n)$ is a first-order Taylor's expansion in $x$ when $t=t^{n}$ (see Eqs. (1.3) and (1.6)).

Similarly, with the aid of Fig. 7, one concludes that:
(a) The flux of $\vec{h}^{*}$ leaving $\mathrm{CE}_{-}(j, n)$ through the line segment $\overline{A B}$ evaluated using the final solution data stored at point $A$ is given by

$$
\begin{equation*}
F(\overline{A B}) \stackrel{\text { def }}{=} \int_{x_{j}}^{x_{j}-(\Delta x / 2)} u^{*}\left(x, t^{n} ; j, n\right) d x=\frac{\Delta x}{2}\left(u_{-}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.11}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(u_{-}\right)_{j}^{n} \stackrel{\text { def }}{=}\left(u-u_{\bar{x}}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.12}
\end{equation*}
$$

is the value of $u$ at the midpoint $M_{-}$of $\overline{A B}$ evaluated using the values $u_{j}^{n}$ and $\left(u_{\bar{x}}\right)_{j}^{n}$ with the aid of Eqs. (1.6) and (1.9).
Also $\hat{F}(\overline{A B})$, the intermediate value of $F(\overline{A B})$, is given by

$$
\begin{equation*}
\hat{F}(\overline{A B})=\frac{\Delta x}{2}\left(\hat{u}_{-}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.13}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(\hat{u}_{-}\right)_{j}^{n} \stackrel{\text { def }}{=}\left(\hat{u}-\hat{u}_{\bar{x}}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.14}
\end{equation*}
$$

is the value of $u$ at the midpoint $M_{-}$evaluated using the intermediate values $\hat{u}_{j}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j}^{n}$ with the aid of Eqs. (1.6), (1.9), and (3.5).
Because $\left(u_{-}\right)_{j}^{n}$ and $\left(\hat{u}_{-}\right)_{j}^{n}$ are evaluated using the solution data stored at point $(j, n)$ (i.e., point $A$ ), hereafter, respectively, they will be referred to as the final and intermediate values of $u$ at the midpoint $M_{-}$associated with the host point $A$ of the region $A B C D$. Similarly, $F(\overline{A B})$ and $\hat{F}(\overline{A B})$, respectively, are referred to as the final and intermediate values of the flux leaving the top face of the same region associated with the same host.
Moreover, with the aid of Eqs. (3.11) and (3.13), one concludes that $F(\overline{A B})(\hat{F}(\overline{A B}))$ is equal to the length of $\overline{A B}$ multiplied by the approximated value $\left(u_{-}\right)_{j}^{n}\left(\left(\hat{u}_{-}\right)_{j}^{n}\right)$ of $u$ at the midpoint $M_{-}$.
(b) The flux of $\vec{h}^{*}$ leaving $\mathrm{CE}_{-}(j+1 / 2, n)$ through the line segment $\overline{F A}$ evaluated using the final solution data stored at point $F$ is given by

$$
\begin{equation*}
F(\overline{F A}) \stackrel{\text { def }}{=} \int_{x_{j}}^{x_{j}+(\Delta x / 2)} u^{*}\left(x, t^{n} ; j+1 / 2, n\right) d x=\frac{\Delta x}{2}\left(u_{-}\right)_{j+1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.15}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(u_{-}\right)_{j+1 / 2}^{n} \stackrel{\text { def }}{=}\left(u-u_{\bar{x}}\right)_{j+1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.16}
\end{equation*}
$$

is the value of $u$ at the midpoint $M_{+}$of $\overline{F A}$ evaluated using the values $u_{j+1 / 2}^{n}$ and $\left(u_{\bar{x}}\right)_{j+1 / 2}^{n}$ with the aid of Eqs. (1.6) and (1.9). Note that: (i) because $(j, n) \in \Omega \Leftrightarrow(j+1 / 2, n) \in \Omega$, Eq. (3.12) $\Leftrightarrow$ Eq. (3.16); and (ii) according to Eqs. (3.6) and (3.15), $F(\overline{A F}) \neq F(\overline{F A})$ even though $\overline{A F}$ and $\overline{F A}$ denote the same line segment. In fact, $F(\overline{A F})$ is a flux evaluated using the solution data stored at point $A$ while $F(\overline{F A})$ is another flux evaluated using those stored at point $F$. The reader is warned that other similar pairs of notations will appear throughout the paper.
Also $\hat{F}(\overline{F A})$, the intermediate value of $F(\overline{F A})$, is given by

$$
\begin{equation*}
\hat{F}(\overline{F A})=\frac{\Delta x}{2}\left(\hat{u}_{-}\right)_{j+1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.17}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(\hat{u}_{-}\right)_{j+1 / 2}^{n} \stackrel{\text { def }}{=}\left(\hat{u}-\hat{u}_{\bar{x}}\right)_{j+1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.18}
\end{equation*}
$$

is the value of $u$ at the midpoint $M_{+}$evaluated using the intermediate values $\hat{u}_{j+1 / 2}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j+1 / 2}^{n}$ with the aid of Eqs. (1.6), (1.9), and (3.5). Obviously Eq. (3.14) $\Leftrightarrow$ Eq. (3.18).
Because $\left(u_{-}\right)_{j+1 / 2}^{n}$ and $\left(\hat{u}_{-}\right)_{j+1 / 2}^{n}$ are evaluated using the solution data stored at point $(j+1 / 2, n)$ (i.e., point $F$ ), hereafter, respectively, they will be referred to as the final and intermediate values of $u$ at the midpoint $M_{+}$associated with the host point $F$ of the region $A F E D$. Similarly, $F(\overline{F A})$ and $\hat{F}(\overline{F A})$, respectively, are referred to as the final and intermediate values of the flux leaving the top face of the same region associated with the same host.

Moreover, with the aid of Eqs. (3.15) and (3.17), one concludes that $F(\overline{F A})(\hat{F}(\overline{F A}))$ is equal to the length of $\overline{F A}$ multiplied by the approximated value $\left(u_{-}\right)_{j+1 / 2}^{n}\left(\left(\hat{u}_{-}\right)_{j+1 / 2}^{n}\right)$ of $u$ at the midpoint $M_{+}$.
(c) The flux of $\vec{h}^{*}$ leaving $\mathrm{CE}_{+}(j-1 / 2, n)$ through the line segment $\overline{B A}$ evaluated using the final solution data stored at point $B$ is given by

$$
\begin{equation*}
F(\overline{B A}) \stackrel{\text { def }}{=} \int_{x_{j}-(\Delta x / 2)}^{x_{j}} u^{*}\left(x, t^{n} ; j-1 / 2, n\right) d x=\frac{\Delta x}{2}\left(u_{+}\right)_{j-1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.19}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(u_{+}\right)_{j-1 / 2}^{n} \stackrel{\text { def }}{=}\left(u+u_{\bar{x}}\right)_{j-1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.20}
\end{equation*}
$$

is the value of $u$ at the midpoint $M_{-}$of $\overline{B A}$ evaluated using the values $u_{j-1 / 2}^{n}$ and $\left(u_{\bar{x}}\right)_{j-1 / 2}^{n}$ with the aid of Eqs. (1.6) and (1.9). Obviously, Eq. (3.7) $\Leftrightarrow$ Eq. (3.20).
Also $\hat{F}(\overline{B A})$, the intermediate value of $F(\overline{B A})$, is given by

$$
\begin{equation*}
\hat{F}(\overline{B A})=\frac{\Delta x}{2}\left(\hat{u}_{+}\right)_{j-1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.21}
\end{equation*}
$$

where

$$
\begin{equation*}
\left(\hat{u}_{+}\right)_{j-1 / 2}^{n} \stackrel{\text { def }}{=}\left(\hat{u}+\hat{u}_{\bar{x}}\right)_{j-1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.22}
\end{equation*}
$$

is the value of $u$ at the midpoint $M_{-}$evaluated using the intermediate values $\hat{u}_{j-1 / 2}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j-1 / 2}^{n}$ with the aid of Eqs. (1.6), (1.9), and (3.5). Obviously Eq. (3.9) $\Leftrightarrow$ Eq. (3.22).
Because $\left(u_{+}\right)_{j-1 / 2}^{n}$ and $\left(\hat{u}_{+}\right)_{j-1 / 2}^{n}$ are evaluated using the solution data stored at point $(j-1 / 2, n)$ (i.e., point $B$ ), hereafter, respectively, they will be referred to as the final and intermediate values of $u$ at the midpoint $M_{-}$associated with the host point $B$ of the region $A B C D$. Similarly, $F(\overline{B A})$ and $\hat{F}(\overline{B A})$, respectively, are referred to as the final and intermediate values of the flux leaving the top face of the same region associated with the same host.
Moreover, with the aid of Eqs. (3.19) and (3.21), one concludes that $F(\overline{B A})(\hat{F}(\overline{B A}))$ is equal to the length of $\overline{B A}$ multiplied by the approximated value $\left(u_{+}\right)_{j-1 / 2}^{n}\left(\left(\hat{u}_{+}\right)_{j-1 / 2}^{n}\right)$ of $u$ at the midpoint $M_{-}$.
Next recall that two different final fluxes $F(\overline{A F})$ and $F(\overline{F A})$, respectively associated with the hosts points $A$ and $F$ of the region $A F E D$, are defined over the top face of this region. For a reason explained in [42], the dual flux (denoted by $F(A ; F)$ or $F(F ; A)$ ) associated with these two fluxes is defined as the simple average of them, i.e.,

$$
\begin{equation*}
F(A ; F) \equiv F(F ; A) \stackrel{\text { def }}{=}[F(\overline{A F})+F(\overline{F A})] / 2 \tag{3.23}
\end{equation*}
$$

Similarly, we have

$$
\begin{equation*}
\hat{F}(A ; F) \equiv \hat{F}(F ; A) \stackrel{\text { def }}{=}[\hat{F}(\overline{A F})+\hat{F}(\overline{F A})] / 2 \tag{3.24}
\end{equation*}
$$

The solution-coupling step will be constructed so that, at the top face of any BCE (such as the space-time region $A F E D$ ) sandwiched between the $n$th and $(n-1 / 2)$ th time levels, the dual flux evaluated using $u_{j}^{n}$ and $\left(u_{\bar{x}}\right)_{j}^{n},(j, n) \in \Omega$, is equal to that evaluated using $\hat{u}_{j}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j}^{n},(j, n) \in \Omega$, i.e.,

$$
\begin{equation*}
F(A ; F)=\hat{F}(A ; F) \quad(j, n) \in \Omega \tag{3.25}
\end{equation*}
$$

By substituting Eqs. (3.6), (3.8), (3.15), and (3.17) into Eqs. (3.23) and (3.24), and then imposing the condition Eq. (3.25), one has

$$
\begin{equation*}
\left(u_{+}\right)_{j}^{n}+\left(u_{-}\right)_{j+1 / 2}^{n}=\left(\hat{u}_{+}\right)_{j}^{n}+\left(\hat{u}_{-}\right)_{j+1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.26}
\end{equation*}
$$

Recall that: (i) $\left(u_{+}\right)_{j}^{n}$ and $\left(u_{-}\right)_{j+1 / 2}^{n}$, respectively, are the final values of $u$ at the midpoint $M_{+}$evaluated using the solution data stored at the two hosts (points $A$ and $F$ ) of the region $A F E D$; and (ii) $\left(\hat{u}_{+}\right)_{j}^{n}$ and
$\left(\hat{u}_{-}\right)_{j+1 / 2}^{n}$, respectively, are the intermediate-value counterparts of $\left(u_{+}\right)_{j}^{n}$ and $\left(u_{-}\right)_{j+1 / 2}^{n}$. As such, Eq. (3.26) states that the average of the two final values of $u$ at the midpoint $M_{+}$evaluated, respectively, using the solution data stored at the two hosts of the region $A F E D$, is identical to its intermediate-value counterpart. Let

$$
\begin{equation*}
\left(u_{+}\right)_{j}^{n}=(1-\eta)\left(\hat{u}_{+}\right)_{j}^{n}+\eta\left(\hat{u}_{-}\right)_{j+1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.27}
\end{equation*}
$$

where $\eta$ is a real adjustable parameter. Then $\left(u_{+}\right)_{j}^{n}$ is an weighted average of $\left(\hat{u}_{+}\right)_{j}^{n}$ and $\left(\hat{u}_{-}\right)_{j+1 / 2}^{n}$. In turn, by substituting Eq. (3.27) into Eq. (3.26), one has

$$
\begin{equation*}
\left(u_{-}\right)_{j+1 / 2}^{n}=\eta\left(\hat{u}_{+}\right)_{j}^{n}+(1-\eta)\left(\hat{u}_{-}\right)_{j+1 / 2}^{n} \quad(j, n) \in \Omega \tag{3.28}
\end{equation*}
$$

i.e., $\left(u_{-}\right)_{j+1 / 2}^{n}$ is also an weighted average of $\left(\hat{u}_{+}\right)_{j}^{n}$ and $\left(\hat{u}_{-}\right)_{j+1 / 2}^{n}$. In fact, Eqs. (3.27) and (3.28) $\Rightarrow$ (i)

$$
\begin{equation*}
\left(u_{+}\right)_{j}^{n}=\left(\hat{u}_{+}\right)_{j}^{n} \quad \text { and } \quad\left(u_{-}\right)_{j+1 / 2}^{n}=\left(\hat{u}_{-}\right)_{j+1 / 2}^{n} \quad \text { if } \quad \eta=0 \quad(j, n) \in \Omega \tag{3.29}
\end{equation*}
$$

(ii)

$$
\begin{equation*}
\left(u_{+}\right)_{j}^{n}=\left(u_{-}\right)_{j+1 / 2}^{n}=\frac{1}{2}\left[\left(\hat{u}_{+}\right)_{j}^{n}+\left(\hat{u}_{-}\right)_{j+1 / 2}^{n}\right] \quad \text { if } \quad \eta=1 / 2 \quad(j, n) \in \Omega \tag{3.30}
\end{equation*}
$$

and (iii)

$$
\begin{equation*}
\left(u_{+}\right)_{j}^{n}=\left(\hat{u}_{-}\right)_{j+1 / 2}^{n} \quad \text { and } \quad\left(u_{-}\right)_{j+1 / 2}^{n}=\left(\hat{u}_{+}\right)_{j}^{n} \quad \text { if } \quad \eta=1 \quad(j, n) \in \Omega \tag{3.31}
\end{equation*}
$$

Because $(j, n) \in \Omega \Leftrightarrow(j-1 / 2, n) \in \Omega$, Eq. (3.27)-(3.31) remain valid if each index $j$ in these equations is replaced by $j-1 / 2$. In particular, Eq. (3.28) $\Leftrightarrow$

$$
\begin{equation*}
\left(u_{-}\right)_{j}^{n}=\eta\left(\hat{u}_{+}\right)_{j-1 / 2}^{n}+(1-\eta)\left(\hat{u}_{-}\right)_{j}^{n} \quad(j, n) \in \Omega \tag{3.32}
\end{equation*}
$$

In turn, by substituting Eqs. (3.7), (3.9), (3.12), (3.14), (3.18), and (3.22) into Eqs.(3.27) and (3.32), one has

$$
\begin{array}{llrl}
\left(u+u_{\bar{x}}\right)_{j}^{n}=(1-\eta)\left(\hat{u}+\hat{u}_{\bar{x}}\right)_{j}^{n}+\eta\left(\hat{u}-\hat{u}_{\bar{x}}\right)_{j+1 / 2}^{n} & & (j, n) \in \Omega \\
\left(u-u_{\bar{x}}\right)_{j}^{n}=(1-\eta)\left(\hat{u}-\hat{u}_{\bar{x}}\right)_{j}^{n}+\eta\left(\hat{u}+\hat{u}_{\bar{x}}\right)_{j-1 / 2}^{n} & & (j, n) \in \Omega \tag{3.34}
\end{array}
$$

By adding Eqs. (3.33) and (3.34) together, one has

$$
\begin{equation*}
u_{j}^{n}=(1-\eta) \hat{u}_{j}^{n}+(\eta / 2)\left[\left(\hat{u}-\hat{u}_{\bar{x}}\right)_{j+1 / 2}^{n}+\left(\hat{u}+\hat{u}_{\bar{x}}\right)_{j-1 / 2}^{n}\right] \quad(j, n) \in \Omega \tag{3.35}
\end{equation*}
$$

On the other hand, by subtracting Eq. (3.34) from Eq. (3.33), one has

$$
\begin{equation*}
\left(u_{\bar{x}}\right)_{j}^{n}=(1-\eta)\left(\hat{u}_{\bar{x}}\right)_{j}^{n}+(\eta / 2)\left[\left(\hat{u}-\hat{u}_{\bar{x}}\right)_{j+1 / 2}^{n}-\left(\hat{u}+\hat{u}_{\bar{x}}\right)_{j-1 / 2}^{n}\right] \quad(j, n) \in \Omega \tag{3.36}
\end{equation*}
$$

Eqs. (3.35) and (3.36) form the solution coupling step, i.e., with the aid of these two equations, at each mesh point $(j, n) \in \Omega, u_{j}^{n}$ and $\left(u_{\bar{x}}\right)_{j}^{n}$ can be determined in terms of the intermediate mesh values at the mesh points $(j-1 / 2, n),(j, n)$, and $(j+1 / 2, n)$.

As expected from Eq. (3.29), we have

$$
\begin{equation*}
u_{j}^{n}=\hat{u}_{j}^{n} \quad \text { and } \quad\left(u_{\bar{x}}\right)_{j}^{n}=\left(\hat{u}_{\bar{x}}\right)_{j}^{n} \quad \text { if } \quad \eta=0 \quad(j, n) \in \Omega \tag{3.37}
\end{equation*}
$$

i.e., solution coupling does not occur when $\eta=0$. On the other hand, because of Eq. (3.30), the case with $\eta=1 / 2$ is referred as the full coupling case.

Let the dual-scheme marching step be executed using the $a-\epsilon$ scheme, i.e., Eqs. (3.3) and (3.4) are assumed. Then, by substituting Eqs. (3.3) and (3.4) into Eqs. (3.35) and (3.36), one obtains

$$
\begin{align*}
& u_{j}^{n}=\frac{\eta}{2}\left[(2-\epsilon) u-\nu u_{\bar{x}}\right]_{j}^{n-1 / 2} \\
& +\frac{(1-\eta)}{2}\left\{\left[(1+\nu) u+\left(1-\nu^{2}\right) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}+\left[(1-\nu) u-\left(1-\nu^{2}\right) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega  \tag{3.38}\\
& +\frac{\eta}{4}\left\{\left[(\epsilon+\nu) u+\left(2 \epsilon+\nu-\nu^{2}\right) u_{\bar{x}}\right]_{j-1}^{n-1 / 2}+\left[(\epsilon-\nu) u+\left(\nu+\nu^{2}-2 \epsilon\right) u_{\bar{x}}\right]_{j+1}^{n-1 / 2}\right\}
\end{align*}
$$

and

$$
\begin{align*}
& \left(u_{\bar{x}}\right)_{j}^{n}=\frac{\eta}{2}\left[\nu u+\left(2-2 \epsilon-\nu^{2}\right) u_{\bar{x}}\right]_{j}^{n-1 / 2} \\
& +\frac{(1-\eta)}{2}\left\{\left[(\epsilon-1) u+(2 \epsilon-1+\nu) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}+\left[(1-\epsilon) u+(2 \epsilon-1-\nu) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega  \tag{3.39}\\
& -\frac{\eta}{4}\left\{\left[(\epsilon+\nu) u+\left(2 \epsilon+\nu-\nu^{2}\right) u_{\bar{x}}\right]_{j-1}^{n-1 / 2}-\left[(\epsilon-\nu) u+\left(\nu+\nu^{2}-2 \epsilon\right) u_{\bar{x}}\right]_{j+1}^{n-1 / 2}\right\}
\end{align*}
$$

Because Eqs. (3.38) and (3.39) are derived assuming that (i) the dual-mesh marching step is carried out using the dual $a-\epsilon$ scheme, and (ii) the solution-coupling step is formed using an weighted-averaging procedure involving the parameter $\eta$ (see Eq. (3.27)), hereafter the full solution-coupling scheme formed by Eqs. (3.38) and (3.39) will be referred to as the $a-\epsilon-\eta$ scheme.

Alternatively, the dual marching step can be executed using the dual $c-\tau$ scheme which are defined by Eqs. (1.12) and (1.16). For this case, the intermediate values $\hat{u}_{j}^{n}$ and $\left(\hat{u}_{\bar{x}}\right)_{j}^{n}$ should be evaluated using Eq. (3.3) and

$$
\begin{equation*}
\left(\hat{u}_{\bar{x}}\right)_{j}^{n}=\frac{1}{2(1+\tau)}\left\{\left[u-(1+2 \nu-\tau) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}-\left[u+(1-2 \nu-\tau) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}\right\} \quad(j, n) \in \Omega \tag{3.40}
\end{equation*}
$$

By substituting Eqs. (3.3) and (3.40) into Eqs. (3.35) and (3.36), one obtains the $c-\tau-\eta$ scheme, i.e.,

$$
\begin{align*}
& u_{j}^{n}=\frac{\eta}{2(1+\tau)}\left[(2+\tau) u-2 \nu u_{\bar{x}}\right]_{j}^{n-1 / 2} \\
& +\frac{(1-\eta)}{2}\left\{\left[(1+\nu) u+\left(1-\nu^{2}\right) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}+\left[(1-\nu) u-\left(1-\nu^{2}\right) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}\right\}  \tag{3.41}\\
& +\frac{\eta}{4(1+\tau)}\left\{\left[(\tau+\nu+\tau \nu) u+\left[2(\nu+\tau)-(1+\tau) \nu^{2}\right] u_{\bar{x}}\right]_{j-1}^{n-1 / 2}\right. \\
& \left.+\left[(\tau-\nu-\tau \nu) u+\left[2(\nu-\tau)+(1+\tau) \nu^{2}\right] u_{\bar{x}}\right]_{j+1}^{n-1 / 2}\right\}
\end{align*}
$$

and

$$
\begin{align*}
& \left(u_{\bar{x}}\right)_{j}^{n}=\frac{\eta}{2(1+\tau)}\left[(1+\tau) \nu u+\left[2-(1+\tau) \nu^{2}\right] u_{\bar{x}}\right]_{j}^{n-1 / 2} \\
& -\frac{(1-\eta)}{2(1+\tau)}\left\{\left[u+(1-2 \nu-\tau) u_{\bar{x}}\right]_{j-1 / 2}^{n-1 / 2}-\left[u-(1+2 \nu-\tau) u_{\bar{x}}\right]_{j+1 / 2}^{n-1 / 2}\right\}  \tag{3.42}\\
& -\frac{\eta}{4(1+\tau)}\left\{\left[(\tau+\nu+\tau \nu) u+\left[2(\nu+\tau)-(1+\tau) \nu^{2}\right] u_{\bar{x}}\right]_{j-1}^{n-1 / 2}\right. \\
& \left.-\left[(\tau-\nu-\tau \nu) u+\left[2(\nu-\tau)+(1+\tau) \nu^{2}\right] u_{\bar{x}}\right]_{j+1}^{n-1 / 2}\right\}
\end{align*}
$$

The von Neumann stability analysis was carried out for both $a-\epsilon-\eta$ and $c-\tau-\eta$ scheme. In particular, the amplification matrices and the associated amplification factors were obtained for both schemes involving arbitrary values of $\epsilon, \tau, \eta$, and the phase angle $\theta[72]$. Because the $a-\epsilon-\eta$ and $c-\tau-\eta$ schemes reduce to the dual $a-\epsilon$ and $c-\tau$ schemes, respectively when $\eta=0$, as expected, the stability conditions of these two schemes also reduce to those of the $a-\epsilon$ and $c-\tau$ schemes, respectively when $\eta=0$. For the general case involving an arbitrary value of $\eta$, it turns out that the stability conditions of these schemes generally become so complicated that they can only be expressed numerically. As such, the discussion of their stability conditions for the general case is beyond the scope of the current paper.

However, for the full coupling case, i.e., when $\eta=1 / 2$, it can be shown that [72] the $a-\epsilon-\eta$ scheme is von Neumann stable $\Leftrightarrow$

$$
\begin{equation*}
\nu^{4} \leq \epsilon \leq 1 \quad(\text { the } a-\epsilon-\eta \text { scheme with } \eta=1 / 2) \tag{3.43}
\end{equation*}
$$

Let

$$
\begin{equation*}
\epsilon=\beta|\nu| \quad \beta>0 \tag{3.44}
\end{equation*}
$$

where $\beta>0$ is an adjustable parameter. Then Eq. (3.43) becomes

$$
\begin{equation*}
\nu^{4} \leq \beta|\nu| \leq 1 \tag{3.45}
\end{equation*}
$$

Excluding the meaningless case $\nu=0$, Eq. (3.45) $\Leftrightarrow\left(|\nu|^{3}-\beta\right) \leq 0$ and $|\nu| \leq 1 / \beta$, i.e.,

$$
|\nu| \leq \min \left\{\beta^{1 / 3}, 1 / \beta\right\}=\left\{\begin{array}{ll}
1 & \text { if } \beta=1  \tag{3.46}\\
\beta^{1 / 3}<1 & \text { if } 0<\beta<1 \\
1 / \beta<1 & \text { if } \beta>1
\end{array} \quad(\nu \neq 0)\right.
$$

In turn, Eqs. (3.44) and (3.46) $\Rightarrow$

$$
\epsilon=\beta|\nu|=\left\{\begin{array}{ll}
1 & \text { if } \beta \geq 1  \tag{3.47}\\
\beta^{4 / 3}<1 & \text { if } 0<\beta<1
\end{array} \quad(\nu \neq 0)\right.
$$

Also, for the case $\eta=1 / 2$, it has been shown that the stability conditions of the $c-\tau-\eta$ scheme is

$$
\begin{equation*}
f\left(\nu^{2}\right) \leq \tau \quad\left(\nu^{2} \leq 2\right) \tag{3.48}
\end{equation*}
$$

if $\nu^{2} \leq 2$, where

$$
\begin{equation*}
f(s) \stackrel{\text { def }}{=} \frac{s(s+2)}{1+2 s-s^{2}} \quad(0 \leq s \leq 2) \tag{3.49}
\end{equation*}
$$

Note that (i)

$$
\begin{equation*}
1+2 s-s^{2}=-(s+\sqrt{2}-1)(s-\sqrt{2}-1)>0 \quad \text { if } \quad 0 \leq s \leq 2 \tag{3.50}
\end{equation*}
$$

i.e., the denominator of the expression on the right side of Eq. (3.49) is not zero throughout the domain of $f(s)$ and therefore the function is well defined over its domain; and (ii)

$$
\begin{equation*}
\frac{d f(s)}{d s}=\frac{2\left(2 s^{2}+s+1\right)}{\left(1+2 s-s^{2}\right)^{2}}>0 \quad \text { if } \quad 0 \leq s \leq 2 \tag{3.51}
\end{equation*}
$$

i.e., $f(s)$ is monotonically increasing in its domain and therefore

$$
\begin{equation*}
0=f(0) \leq f(s) \leq f(2)=8 \quad(0 \leq s \leq 2) \tag{3.52}
\end{equation*}
$$

The stability conditions of the $c-\tau-\eta$ scheme for the case $\eta=1 / 2$ and $\nu^{2}>2$ can also be cast into an analytical form. However, the form is much more complicated and will not be presented here.

### 3.2. 2D and 3D extensions

For a 2D case using triangular meshes, or a 3D case using tetrahedral meshes, there is a natural extension of each of the 1D $a-\epsilon-\eta$ and $c-\tau-\eta$ schemes. As an example, consider Fig. 8 where the triangle $\triangle B D F$ and its three neighbors $\triangle B F H, \triangle B I D$, and $\triangle D J F$ (which all lie on the $x-y$ plane) are shown. As explained in Sec. 2, each of the centroid mesh points $G, A, C$, and $E$ shown in Fig. 3 is replaced in Fig. 8 by another associated interior point which is (i) defined by Eqs. (2.10) and (2.29) with $\zeta \geq 1$; and (ii) denoted by the same symbol and also marked by an open circle. The spatial solution point $G^{\prime}$ (also marked by a cross) associated with $\triangle B D F$ again is defined as the centroid of the hexagon $A B C D E F$. Other spatial solution points $A^{\prime}, C^{\prime}$, and $E^{\prime}$ are defined similarly. In addition, points $K, L$, and $M$ (marked by triangles) depicted in Fig. 8, by definition, are the centroids of the quadrilaterals $G F A B, G B C D$, and $G D E F$, respectively.

Let each spatial solution point be assigned an unique identification index (an integer). As an example, let points $G^{\prime}, A^{\prime}, C^{\prime}$, and $E^{\prime}$ be identified by the indices $j, j_{1}, j_{2}$, and $j_{3}$, respectively. As such, as an example, the x - and y - coordinates of point $G^{\prime}$ will be denoted by $x_{j}$ and $y_{j}$, respectively. Moreover, the space-time solution points $\left(G^{\prime}, n\right),\left(A^{\prime}, n\right),\left(C^{\prime}, n\right)$, and $\left(E^{\prime}, n\right)$, which lie on the $n$th time level and have the points $G^{\prime}, A^{\prime}, C^{\prime}$, and $E^{\prime}$ being their spatial projections, respectively, will be denoted by $(j, n),\left(j_{1}, n\right)$, $\left(j_{2}, n\right)$, and $\left(j_{3}, n\right)$, respectively. Also each space-time solution point $(j, n)$ is associated with [13] (i) one SE denoted by $\mathrm{SE}(j, n)$; and (ii) three BCEs denoted by $\mathrm{CE}_{(k)}(j, n), k=1,2,3$, respectively, and a CCE which is the union of $\mathrm{CE}_{(k)}(j, n), k=1,2,3$, and denoted by $\mathrm{CE}(j, n)$. By definition, $\mathrm{CE}_{(1)}(j, n)$ is the space-time cylinder sandwiched between the $n$th and $(n-1 / 2)$ th time levels with the quadrilateral $G F A B$ being its spatial projection. $\mathrm{CE}_{(2)}(j, n)$ and $\mathrm{CE}_{(3)}(j, n)$ are defined identically except that the quadrilateral $G F A B$ are replaced by the quadrilaterals $G B C D$ and $G D E F$, respectively. It is seen that (i) the solution points $(j, n)$ and $\left(j_{1}, n\right)$ are the hosts of the space-time region occupied by $\mathrm{CE}_{(1)}(j, n)$, (ii) the solution points $(j, n)$ and $\left(j_{2}, n\right)$ are the hosts of the space-time region occupied by $\mathrm{CE}_{(2)}(j, n)$, and (iii) the solution points $(j, n)$ and $\left(j_{3}, n\right)$ are the hosts of the space-time region occupied by $\mathrm{CE}_{(3)}(j, n)$,

Let the 2 D counterpart of Eq. (1.1) be

$$
\begin{equation*}
\frac{\partial u}{\partial t}+a_{x} \frac{\partial u}{\partial x}+a_{y} \frac{\partial u}{\partial y}=0 \tag{3.53}
\end{equation*}
$$

where $a_{x}$, and $a_{y}$ are constants. Then the 2D counterparts of Eqs. (1.4) and (1.6) are

$$
\begin{equation*}
\vec{h}^{*}(x, y, t ; j, n) \stackrel{\text { def }}{=}\left[a_{x} u^{*}(x, y, t ; j, n), a_{y} u^{*}(x, y, t ; j, n), u^{*}(x, y, t ; j, n)\right] \tag{3.54}
\end{equation*}
$$

and

$$
\begin{equation*}
u^{*}(x, y, t ; j, n)=u_{j}^{n}+\left(u_{x}\right)_{j}^{n}\left[\left(x-x_{j}\right)-a_{x}\left(t-t^{n}\right)\right]+\left(u_{y}\right)_{j}^{n}\left[\left(y-y_{j}\right)-a_{y}\left(t-t^{n}\right)\right] \tag{3.55}
\end{equation*}
$$

respectively. Thus there are three independent marching variables, i.e., $u_{j}^{n},\left(u_{x}\right)_{j}^{n}$, and $\left(u_{y}\right)_{j}^{n}$ associated with a space-time solution point $(j, n)$ (i.e., point $\left.\left(G^{\prime}, n\right)\right)$.

Next let (i) $(K, n),(L, n)$, and $(M, n)$ denote the points on the $n$th time level with their spatial projections being points $K, L$, and $M$, respectively; (ii) the intermediate solution data stored at $(j, n),\left(j_{1}, n\right)$, $\left(j_{2}, n\right)$, and $\left(j_{3}, n\right)$ be evaluated from the given solution data at the ( $n-1 / 2$ )th time level using the 2D versions of the $a-\epsilon$ or $c-\tau$ scheme $[7,13,31]$; (iii) $[\hat{u}(K, n)]_{j}^{n}$ and $[\hat{u}(K, n)]_{j_{1}}^{n}$ denote the intermediate values of $u$ at $(K, n)$ evaluated using the intermediate solution data stored at the hosts $(j, n)$ and $\left(j_{1}, n\right)$ of $\mathrm{CE}_{(1)}(j, n)$, respectively; (iv) $[\hat{u}(L, n)]_{j}^{n}$ and $[\hat{u}(L, n)]_{j_{2}}^{n}$ denote the intermediate values of $u$ at $(L, n)$ evaluated using the intermediate solution data stored at the hosts $(j, n)$ and $\left(j_{2}, n\right)$ of $\mathrm{CE}_{(2)}(j, n)$, respectively; (v) $[\hat{u}(M, n)]_{j}^{n}$ and $[\hat{u}(M, n)]_{j_{3}}^{n}$ denote the intermediate values of $u$ at $(L, n)$ evaluated using the intermediate solution data stored at the hosts $(j, n)$ and $\left(j_{3}, n\right)$ of $\mathrm{CE}_{(3)}(j, n)$, respectively; (vi) $[u(K, n)]_{j}^{n}$ and $[u(K, n)]_{j_{1}}^{n}$ denote the final values of $u$ at $(K, n)$ evaluated using the final solution data stored at the hosts $(j, n)$ and ( $j_{1}, n$ ) of $\mathrm{CE}_{(1)}(j, n)$, respectively; (vii) $[u(L, n)]_{j}^{n}$ and $[u(L, n)]_{j_{2}}^{n}$ denote the final values of $u$ at $(L, n)$ evaluated using
the final solution data stored at the hosts $(j, n)$ and $\left(j_{2}, n\right)$ of $\mathrm{CE}_{(2)}(j, n)$, respectively; and (viii) $[u(M, n)]_{j}^{n}$ and $[u(M, n)]_{j_{3}}^{n}$ denote the final values of $u$ at $(M, n)$ evaluated using the final solution data stored at the hosts $(j, n)$ and $\left(j_{3}, n\right)$ of $\mathrm{CE}_{(3)}(j, n)$, respectively.

Moreover, for each $k=1,2,3$, let (i) $A_{k}(j, n)$ denote the area of the spatial projection of $\mathrm{CE}_{(k)}(j, n)$; (ii) $\hat{F}_{k}(j, n)$ denote the intermediate value of the flux of $\vec{h}^{*}$ leaving $\mathrm{CE}_{(k)}(j, n)$ through its top face evaluated using the intermediate solution data stored at the space-time solution point $(j, n)$; (iii) $\hat{F}_{k}\left(j_{k}, n\right)$ denote the intermediate value of the flux of $\vec{h}^{*}$ leaving $\mathrm{CE}_{(k)}(j, n)$ through its top face evaluated using the intermediate solution data stored at the space-time solution point $\left(j_{k}, n\right)$; (iv) $F_{k}(j, n)$ denote the final value of the flux of $\vec{h}^{*}$ leaving $\mathrm{CE}_{(k)}(j, n)$ through its top face evaluated using the final solution data stored at the space-time solution point $(j, n)$; and (v) $F_{k}\left(j_{k}, n\right)$ denote the final value of the flux of $\vec{h}^{*}$ leaving $\mathrm{CE}_{(k)}(j, n)$ through its top face evaluated using the final solution data stored at the space-time solution point $\left(j_{k}, n\right)$. According to the above definitions, $A_{1}(j, n), A_{2}(j, n)$, and $A_{3}(j, n)$ are the areas of the quadrilaterals $G F A B, G B C D$, and $G D E F$, respectively. Because these quadrilaterals are the spatial projections of $\mathrm{CE}_{(k)}(j, n), k=1,2,3$, respectively, for each $k=1,2,3, A_{k}(j, n)$ is the area of the top face of $\mathrm{CE}_{(k)}(j, n)$.

With the aid of (i) the above definitions, (ii) the fact that point $K$ depicted in Fig. 8 is the centroid of the quadrilateral $G F A B$ (i.e., the spatial projection of $\mathrm{CE}_{(1)}(j, n)$ ), and (iii) Eq. (3.55) is a first-order Taylor's expansion, by using Eqs. (3.54) and (3.55), one can establish the following relations for several fluxes of $\vec{h}^{*}$ leaving $\mathrm{CE}_{(1)}(j, n)$ through its top face:

$$
\begin{align*}
\hat{F}_{1}(j, n) & =[\hat{u}(K, n)]_{j}^{n} A_{1}(j, n)  \tag{3.56}\\
\hat{F}_{1}\left(j_{1}, n\right) & =[\hat{u}(K, n)]_{j_{1}}^{n} A_{1}(j, n)  \tag{3.57}\\
F_{1}(j, n) & =[u(K, n)]_{j}^{n} A_{1}(j, n) \tag{3.58}
\end{align*}
$$

and

$$
\begin{equation*}
F_{1}\left(j_{1}, n\right)=[u(K, n)]_{j_{1}}^{n} A_{1}(j, n) \tag{3.59}
\end{equation*}
$$

The current solution-coupling step will be constructed assuming the following dual-flux conservation relation:

$$
\begin{equation*}
F_{1}(j, n)+F_{1}\left(j_{1}, n\right)=\hat{F}_{1}(j, n)+\hat{F}_{1}\left(j_{1}, n\right) \tag{3.60}
\end{equation*}
$$

By substituting Eqs. (3.56)-(3.59) into Eq. (3.60), one obtains

$$
\begin{equation*}
[u(K, n)]_{j}^{n}+[u(K, n)]_{j_{1}}^{n}=[\hat{u}(K, n)]_{j}^{n}+[\hat{u}(K, n)]_{j_{1}}^{n} \tag{3.61}
\end{equation*}
$$

Similarly, by considering the fluxes at the top faces of $\mathrm{CE}_{(2)}(j, n)$ and $\mathrm{CE}_{(3)}(j, n)$ and imposing the dual-flux conservation relations similar to Eq. (3.60), one has

$$
\begin{align*}
{[u(L, n)]_{j}^{n}+[u(L, n)]_{j_{2}}^{n} } & =[\hat{u}(L, n)]_{j}^{n}+[\hat{u}(L, n)]_{j_{2}}^{n}  \tag{3.62}\\
{[u(M, n)]_{j}^{n}+[u(M, n)]_{j_{3}}^{n} } & =[\hat{u}(M, n)]_{j}^{n}+[\hat{u}(M, n)]_{j_{3}}^{n} \tag{3.63}
\end{align*}
$$

We assume the following weighted-averaging relations:

$$
\begin{align*}
{[u(K, n)]_{j}^{n} } & =(1-\eta)[\hat{u}(K, n)]_{j}^{n}+\eta[\hat{u}(K, n)]_{j_{1}}^{n}  \tag{3.64}\\
{[u(L, n)]_{j}^{n} } & =(1-\eta)[\hat{u}(L, n)]_{j}^{n}+\eta[\hat{u}(L, n)]_{j_{2}}^{n} \tag{3.65}
\end{align*}
$$

and

$$
\begin{equation*}
[u(M, n)]_{j}^{n}=(1-\eta)[\hat{u}(M, n)]_{j}^{n}+\eta[\hat{u}(M, n)]_{j_{3}}^{n} \tag{3.66}
\end{equation*}
$$

where $\eta$ is a real adjustable parameter. In turn, by substituting Eqs. (3.64)-(3.66) into Eqs. (3.61)-(3.63), one has

$$
\begin{align*}
{[u(K, n)]_{j_{1}}^{n} } & =\eta[\hat{u}(K, n)]_{j}^{n}+(1-\eta)[\hat{u}(K, n)]_{j_{1}}^{n}  \tag{3.67}\\
{[u(L, n)]_{j_{2}}^{n} } & =\eta[\hat{u}(L, n)]_{j}^{n}+(1-\eta)[\hat{u}(L, n)]_{j_{2}}^{n} \tag{3.68}
\end{align*}
$$

and

$$
\begin{equation*}
[u(M, n)]_{j_{3}}^{n}=\eta[\hat{u}(M, n)]_{j}^{n}+(1-\eta)[\hat{u}(M, n)]_{j_{3}}^{n} \tag{3.69}
\end{equation*}
$$

respectively.
According to Eqs. (3.64)-(3.66), $[u(K, n)]_{j}^{n},[u(L, n)]_{j}^{n}$ and $[u(M, n)]_{j}^{n}$ can be determined in terms of $[\hat{u}(K, n)]_{j}^{n},[\hat{u}(K, n)]_{j_{1}}^{n},[\hat{u}(L, n)]_{j}^{n},[\hat{u}(L, n)]_{j_{2}}^{n},[\hat{u}(M, n)]_{j}^{n},[\hat{u}(M, n)]_{j_{3}}^{n}$, and a given value of $\eta$. On the other hand, because $t=t^{n}$ at points $(j, n),(K, n),(L, n)$, and $(M, n)$, Eq. (3.55) along with the definitions of $[u(K, n)]_{j}^{n},[u(L, n)]_{j}^{n}$, and $[u(M, n)]_{j}^{n}, \Rightarrow$

$$
\begin{align*}
u_{j}^{n}+\left(u_{x}\right)_{j}^{n}\left[x(K)-x_{j}\right]+\left(u_{y}\right)_{j}^{n}\left[y(K)-y_{j}\right] & =[u(K, n)]_{j}^{n}  \tag{3.70}\\
u_{j}^{n}+\left(u_{x}\right)_{j}^{n}\left[x(L)-x_{j}\right]+\left(u_{y}\right)_{j}^{n}\left[y(L)-y_{j}\right] & =[u(L, n)]_{j}^{n} \tag{3.71}
\end{align*}
$$

and

$$
\begin{equation*}
u_{j}^{n}+\left(u_{x}\right)_{j}^{n}\left[x(M)-x_{j}\right]+\left(u_{y}\right)_{j}^{n}\left[y(M)-y_{j}\right]=[u(M, n)]_{j}^{n} \tag{3.72}
\end{equation*}
$$

Here (i) $x(K)$ and $y(K)$ are the $x$ - and $y$ - coordinates of point $K$, (ii) $x(L)$ and $y(L)$ are the $x$ - and $y$ coordinates of point $L$, and (iii) $x(M)$ and $y(M)$ are the $x$ - and $y$-coordinates of point $M$.

Let

$$
\operatorname{det}\left(\begin{array}{ccc}
1 & x(K)-x_{j} & y(K)-y_{j}  \tag{3.73}\\
1 & x(L)-x_{j} & y(L)-y_{j} \\
1 & x(M)-x_{j} & y(M)-y_{j}
\end{array}\right) \equiv \operatorname{det}\left(\begin{array}{ll}
x(L)-x(K) & y(L)-y(K) \\
x(M)-x(K) & y(M)-y(K)
\end{array}\right) \neq 0
$$

i.e., points $K, L$, and $M$ do not lie on a straight line. Then Eqs. (3.70)-(3.72) can be inverted and, as a result, $u_{j}^{n},\left(u_{x}\right)_{j}^{n}$, and $\left(u_{y}\right)_{j}^{n}$ can be explicitly determined in terms of $[u(K, n)]_{j}^{n},[u(L, n)]_{j}^{n}$, and $[u(M, n)]_{j}^{n}$ if the spatial coordinates of points $(j, n), K, L$, and $M$ are given. As such, given all the spatial geometric data, the final solution values at any $(j, n)$ can be determined in terms of the given solution values at the ( $n-1 / 2$ ) time level by using the 2D solution-coupling procedure describe above.

In a CESE 3D scheme using a spatial mesh constructed from tetrahedra, generally a solution point is associated with four independent mesh variables, four BCEs and one CCE. Also each BCE has two hosts. As such, construction of the 3D solution-coupling procedure is a straightforward extension of its 2D counterpart.

Finally, it should be emphasized that the use of a solution-coupling procedure generally will introduce numerical dissipation and, as a result, degrade accuracy. As such it is strongly recommended that, as long as it serves the purpose of preventing solution decoupling, this procedure should be applied as sparingly as possible, i.e., it should be used with a small value of $\eta$ and/or applied once after a preset number of marching steps.

## 4. Conclusions and discussions

As a preliminary for the later development, a review of basic CESE ideas and schemes along with a discussion of several issues related to the application of the CESE method involving unstructured triangular/tetrahedral meshes was presented in Sec. 1. Among the issues discussed are: (A) geometric difficulties associated with a boundary which has high curvature or sharp corner and is also surrounded by triangular/tetrahedral meshes of extremely high aspect ratio, and (B) possible occurrences of solution decoupling in CESE simulations involving unstructured triangular/tetrahedral meshes.

To tackle the issue (A), Sec. 2 began by identifying the root cause of the geometry difficulty for the triangular-mesh case, i.e., the hexagon $A B C D E F$ depicted in Fig. 4 may become highly concave if (i) the
mesh points $A, C, E$, and $G$ are designated, respectively, as the centroids of the associated triangles, and (ii) two neighboring triangles with high aspect ratio also have sharply different orientations. A rigorous mathematical analysis was then presented to validate the observation. According to this analysis, the centroid and in-center of a triangle are but two special cases with $\zeta=0$ and $\zeta=1$, respectively, of the family of the interior points defined by Eqs. (2.10), (2.29), and (2.32). By using the results presented in Eqs. (2.39)-(2.41), it becomes clear that, by replacing the centroid with another interior point with $\zeta \geq 1$, the hexagon $A B C D E F$ referred to above will become convex or much less concave. This conclusion has been verified by numerous numerical experiments. In fact, the geometric difficulty can be overcome by simply designating points $A, C, E$, and $G$ (see Fig. 4), respectively, as the in-centers $(\zeta=1)$ of the associated triangles. Moreover, it was shown in Sec. 2 that the tetrahedral-mesh case can be dealt with in a completely parallel manner.

The new solution-coupling procedures described in Sec. 3 have become an integral part of the CESE code development. In fact, the nuisance of solution decoupling is avoided completely with the use of these new procedures.

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Figure 1.-A surface element on the boundary $S(V)$ of an arbitrary space-time volume $V$.




2(e) CE(j,n)
Figure 2.-The SEs and CEs.


Figure 3.-A spatial domain formed by triangles.


Figure 4.-Neighboring triangles with high aspect ratios and sharply different orientations.


Figure 5.-The perpendicular projections of paint $P_{0}$ on the three sides (or their extensions) of a triangle $\Delta \mathrm{P}_{1}, \mathrm{P}_{2}, \mathrm{P}_{3}$.


Figure 6.-Five neighboring tetrahedra in the $x-y-z$ space.


Figure 7.-Construction of the $a-\varepsilon-\eta$ and $c-\tau-\eta$ schemes.


Figure 8.-Construction of the 2D $a-\varepsilon-\eta$ and $c-\tau-\eta$ schemes.


